## MECHANICS OF SOLIDS

M.A./M.Sc. Mathematics (Final)

## MM-504 \& MM 505 <br> (Option-A ${ }_{1}$ )

Directorate of Distance Education
Maharshi Dayanand University ROHTAK - 124001

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# M.A./M.Sc. Mathematics (Final) MECHANICS OF SOLIDS <br> MM-504 \& 505 ( $\mathrm{A}_{1}$ ) 


#### Abstract

Max. Marks : 100 Time: 3 Hours Note: Question paper will consist of three sections. Section I consisting of one question with ten parts covering whole of the syllabus of 2 marks each shall be compulsory. From Section II, 10 questions to be set selecting two questions from each unit. The candidate will be required to attempt any seven questions each of five marks. Section III, five questions to be set, one from each unit. The candidate will be required to attempt any three questions each of fifteen marks.


## Unit I

Analysis of Strain: Affine transformation. Infinite simal affine deformation. Geometrical interpretation of the components of strain. Strain quadric of Cauchy. Principal strains and invariants. General infinitesimal deformation. Saint-Venant's equations of Compatibility. Finite deformations.

## Unit II

Equations of Elasticity: Generalized Hooke's law. Homogeneous isotropic media. Elasticity moduli for isotropic media. Equilibrium and dynamic equations for an isotropic elastic solid. Strain energy function and its connection with Hooke's law. Uniquness of solution. Beltrami-Michell compatibility equations. Saint-Venant's principle.

## Unit III

Two - dimensional Problems: Plane stress. Generalized plane stress. Airy stress function. General solution of Biharmonic equation. Stresses and displacements in terms of complex potentials. Simple problems. Stress function appropriate to problems of plane stress. Problems of semi-infinite solids with displacements or stresses prescribed on the plane boundary.

## Unit IV

Torsional Problem: Torsion of cylindrical bars. Tortional rigidity. Torsion and stress functions. Lines of shearing stress. Simple problems related to circle, elipse and equilateral triangle.

Variational Methods: Theorems of minimum potential energy. Theorems of minimum complementary energy. Reciprocal theorem of Betti and Rayleigh. Deflection of elastic string, central line of a beam and elastic membrane. Torsion of cylinders. Variational problem related to biharmonic equation. Solution of Euler's equation by Ritz, Galerkin and Kantorovich methods.

## Unit V

Elastic Waves: Propagation of waves in an isotropic elastic solid medium. Waves of dilatation and distortion Plane waves. Elastic surface waves such as Rayleigh and Love waves.

## Chapter-1 Cartesian Tensors

### 1.1 INTRODUCTION

There are physical quantities which are independent or invariant of any particular coordinate system that may be used to describe them. Mathematically, such quantities are represented by tensors. That is, a tensor is a quantity which describes a physical state or a physical phenomenon.

As a mathematical entity, a tensor has an existence independent of any coordinate system. Yet it may be specified in a particular coordinate system by a certain set of quantities, known as its components. Specifying the components of a tensor in one coordinate system determines the components in any other system according to some definite law of transformation.

In dealing with general coordinate transformations between arbitrary curvilinear coordinate systems, the tensors defined are known as general tensors. When one is dealing with cartesian rectangular frames of reference only, the tensor involved are referred to as cartesian tensors. From now onwards, the word "tensor" means "cartesian tensors" unless specifically stated otherwise.

### 1.2 COORDINATE TRANSFORMATIONS

Let us consider a right handed system of rectangular cartesian axes $x_{i}$ with a fixed origin O . Let P be a general point whose coordinates with respect to this system $\mathrm{Ox}_{1} \mathrm{x}_{2} \mathrm{x}_{3}$ are $\left(\mathrm{x}_{1}, \mathrm{x}_{2}, \mathrm{x}_{3}\right)$.


Fig (1.1)
Let $\overline{\mathrm{r}}$ be the position vector of P w.r.t. O .

Then

$$
\begin{equation*}
\overline{\mathrm{r}}=\mathrm{x}_{1} \hat{\mathrm{e}}_{1}+\mathrm{x}_{2} \hat{\mathrm{e}}_{2}+\mathrm{x}_{3} \hat{\mathrm{e}}_{3} \tag{1}
\end{equation*}
$$

and $\left(\mathrm{x}_{1}, \mathrm{x}_{2}, \mathrm{x}_{3}\right)$ are the components of the vector $\overline{\mathrm{OP}}$. Here, $\hat{\mathrm{e}}_{1}, \hat{\mathrm{e}}_{2}, \hat{\mathrm{e}}_{3}$ are unit vectors along axes.
Let a new system $\mathrm{Ox}_{1}{ }^{\prime} \mathrm{x}_{2}{ }^{\prime} \mathrm{x}_{3}$ ' of axes be obtained by rotating the "old system" of axes about some line in space through $O$. The position vector $\overline{\mathbf{r}}$ of P has the following representation in the new system

$$
\begin{equation*}
\overline{\mathrm{r}}=\mathrm{x}_{1}{ }^{\prime} \hat{\mathrm{e}}_{1}{ }^{\prime}+\mathrm{x}_{2}{ }^{\prime} \hat{\mathrm{e}}_{2}^{\prime}+\mathrm{x}_{3}{ }^{\prime} \hat{\mathrm{e}}_{3}^{\prime} \tag{2}
\end{equation*}
$$

where $\hat{\mathrm{e}}_{\mathrm{i}}$ is the unit vector directed along the positive $\mathrm{x}_{\mathrm{i}}{ }^{\prime}$-axis, and

$$
\hat{\mathrm{e}}_{\mathrm{i}} \cdot{ }^{\prime} \hat{\mathrm{e}}_{\mathrm{j}}^{\prime}=\left\{\begin{array}{l}
1 \text { for } \mathrm{i}=\mathrm{j} \\
0 \text { for } \mathrm{i} \neq \mathrm{j}
\end{array}\right.
$$

and

$$
\hat{\mathrm{e}}_{1}^{\prime} \times \hat{\mathrm{e}}_{2}{ }^{\prime}=\hat{\mathrm{e}}_{3}^{\prime} \text {, etc. }
$$

and $\left(\mathrm{x}_{1}{ }^{\prime}, \mathrm{x}_{2}{ }^{\prime}, \mathrm{x}_{3}{ }^{\prime}\right)$ are the new components of $\overline{\mathrm{OP}}=\overline{\mathrm{r}}$ relative to the new axes $\mathrm{Ox}_{1}{ }^{\prime} \mathrm{x}_{2}{ }^{\prime} \mathrm{x}_{3}{ }^{\prime}$. Let $\mathrm{a}_{\mathrm{pi}}$ be the direction cosines of new $\mathrm{x}_{\mathrm{p}}{ }^{\prime}$-axis w.r.t. the old $\mathrm{x}_{\mathrm{i}}-$ axis.
That is,

$$
\mathrm{a}_{\mathrm{pi}}=\cos \left(\mathrm{x}_{\mathrm{p}}^{\prime}, \mathrm{x}_{\mathrm{i}}\right)=\text { cosine of the angle between the positive } \mathrm{x}_{\mathrm{p}}{ }^{\prime} \text {-axis }
$$ (new axis) and the positive $\mathrm{x}_{\mathrm{i}}$-axis (old axis)

$$
\begin{equation*}
=\overline{\mathrm{e}}_{\mathrm{p}}^{\prime} . \overline{\mathrm{e}}_{\mathrm{i}} \tag{3}
\end{equation*}
$$

Form (2), we write

$$
\begin{align*}
& \overline{\mathrm{r}} \cdot \hat{\mathrm{e}}_{\mathrm{p}}^{\prime}=\mathrm{x}_{\mathrm{p}}^{\prime} \\
\Rightarrow & \mathrm{x}_{\mathrm{p}}^{\prime}={\overline{\mathrm{r}} \cdot \hat{e}_{\mathrm{p}}^{\prime}=\left(\mathrm{x}_{1} \hat{\mathrm{e}}_{1}+\mathrm{x}_{2} \hat{\mathrm{e}}_{2}+\mathrm{x}_{3} \hat{\mathrm{e}}_{3}\right) \cdot \hat{\mathrm{e}}_{\mathrm{p}}^{\prime}} \quad \\
\Rightarrow \mathrm{x}_{\mathrm{p}}^{\prime}= & \mathrm{a}_{\mathrm{p} 1} \mathrm{x}_{1}+\mathrm{a}_{\mathrm{p} 2} \mathrm{x}_{2}+\mathrm{a}_{\mathrm{p}_{3}} \mathrm{x}_{3}=\mathrm{a}_{\mathrm{pi}} \mathrm{x}_{\mathrm{i}} . \tag{4}
\end{align*}
$$

Here $p$ is the free suffix and $i$ is dummy.
In the above, the following Einstein summation convection is used.
"Unless otherwise stated specifically, whenever a suffix is repeated, it is to be given all possible values $(1,2,3)$ and that the terms are to be added for all".

Similarly

$$
\begin{aligned}
\mathrm{x}_{\mathrm{i}} & =\overline{\mathrm{r}} \cdot \hat{\mathrm{e}}_{\mathrm{i}} \\
& =\left(\mathrm{x}_{1}{ }^{\prime} \hat{\mathrm{e}}_{1}^{\prime}+\mathrm{x}_{2} \hat{\mathrm{e}}^{\prime}{ }_{2}+\mathrm{x}_{3} \hat{\mathrm{e}}^{\prime}{ }_{3}\right) \cdot \hat{\mathrm{e}}_{\mathrm{i}}, \\
& =\mathrm{a}_{1 \mathrm{i}} \mathrm{x}^{\prime}{ }_{1}+\mathrm{a}_{2 \mathrm{i}} \mathrm{x}^{\prime}{ }_{2}+\mathrm{a}_{3 \mathrm{i}} \mathrm{x}^{\prime}{ }_{3}
\end{aligned}
$$

$$
\begin{equation*}
=\mathrm{a}_{\mathrm{pi}} \mathrm{X}_{\mathrm{p}}^{\prime} . \tag{5}
\end{equation*}
$$

Here i is a free suffix and p is dummy. When the orientations of the new axes w.r.t. the old axes are known, the coefficients $\mathrm{a}_{\mathrm{pi}}$ are known. Relation (4) represent the law that transforms the old triplet $\mathrm{x}_{\mathrm{i}}$ to the new triplet $\mathrm{x}_{\mathrm{p}}{ }^{\prime}$ and (5) represent the inverse law, giving old system in terms of new system.
Remark 1. The transformation rules (4) and (5) may be displaced in the following table

|  | $\mathrm{X}_{1}$ | $\mathrm{X}_{2}$ | $\mathrm{X}_{3}$ |
| :---: | :---: | :---: | :---: |
| $\mathrm{X}^{\prime}{ }_{1}$ | $\mathrm{a}_{11}$ | $\mathrm{a}_{12}$ | $\mathrm{a}_{13}$ |
| $\mathrm{X}^{\prime}{ }_{2}$ | $\mathrm{a}_{21}$ | $\mathrm{a}_{22}$ | $\mathrm{a}_{23}$ |
| $\mathrm{X}^{\prime}{ }_{3}$ | $\mathrm{a}_{31}$ | $\mathrm{a}_{32}$ | $\mathrm{a}_{33}$ |

Remark 2. The transformation (4) is a linear transformation given by

$$
\left[\begin{array}{l}
x_{1}^{\prime}  \tag{7}\\
x_{2}^{\prime} \\
x_{3}^{\prime}
\end{array}\right]=\left[\begin{array}{lll}
a_{11} & a_{12} & a_{13} \\
a_{21} & a_{22} & a_{23} \\
a_{31} & a_{32} & a_{33}
\end{array}\right]\left[\begin{array}{l}
x_{1} \\
x_{2} \\
x_{3}
\end{array}\right]
$$

The matrix

$$
\begin{equation*}
[\mathrm{L}]=\left(\mathrm{a}_{\mathrm{ij}}\right)_{3 \times 3} \tag{8}
\end{equation*}
$$

may be thought as an operator operating on the vector $\left[\begin{array}{l}x_{1} \\ x_{2} \\ x_{3}\end{array}\right]$ and giving the vector $\left[\begin{array}{l}\mathrm{x}_{1} \\ \mathrm{x}_{2} \\ \mathrm{x}_{3}\end{array}\right]$.
Remark 3. Since this transformation is rotational only, so the matrix $L$ of the transformation is non-symmetric.

Remark 4. Relations (4) and (5) yield

$$
\begin{equation*}
\frac{\partial \mathrm{x}_{\mathrm{p}}^{\prime}}{\partial \mathrm{x}_{\mathrm{i}}}=\mathrm{a}_{\mathrm{pi}} \tag{9}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{\partial \mathrm{x}_{\mathrm{i}}}{\partial \mathrm{x}_{\mathrm{p}}^{\prime}}=\mathrm{a}_{\mathrm{pi}} . \tag{10}
\end{equation*}
$$

### 1.3. THE SYMBOL $\delta_{\mathrm{ij}}$

It is defined as

$$
\delta_{\mathrm{ij}}=\begin{align*}
& 1 \text { if } \mathrm{i}=\mathrm{j}  \tag{1}\\
& 0 \text { if } \mathrm{i} \neq \mathrm{j}
\end{align*} .
$$

That is,

$$
\begin{aligned}
& \delta_{11}=\delta_{22}=\delta_{33}=1, \\
& \delta_{12}=\delta_{21}=\delta_{31}=\delta_{23}=\delta_{32}=0 .
\end{aligned}
$$

The symbol $\delta_{\mathrm{ij}}$ is known as the Kronecker $\delta$ symbol, named after the German mathematician Leopold Kronecker (1827-1891). The following property is inherent in the definition of $\delta_{\mathrm{ij}}$

$$
\delta_{\mathrm{ij}}=\delta_{\mathrm{ji}} .
$$

By summation convention

$$
\begin{equation*}
\delta_{\mathrm{ii}}=\delta_{11}+\delta_{22}+\delta_{33}=3 . \tag{2}
\end{equation*}
$$

The unit matrix of order 3 is

$$
\mathrm{I}_{3}=\left(\delta_{\mathrm{ij}}\right) \text { and } \operatorname{det}\left(\delta_{\mathrm{ij}}\right)=1 .
$$

The orthonormality of the base unit vectors $\hat{\mathrm{e}}_{\mathrm{i}}$ can be written as

$$
\begin{equation*}
\hat{\mathrm{e}}_{\mathrm{i}} \cdot \hat{\mathrm{e}}_{\mathrm{j}}=\delta_{\mathrm{ij}} . \tag{3}
\end{equation*}
$$

We know that

$$
\frac{\partial x_{i}}{\partial x_{j}}=\left\{\begin{array}{l}
1 \text { if } i=j \\
0 \text { if } i \neq j
\end{array}\right.
$$

Therefore,

$$
\begin{equation*}
\frac{\partial \mathrm{x}_{\mathrm{i}}}{\partial \mathrm{x}_{\mathrm{j}}}=\delta_{\mathrm{ij}} . \tag{4}
\end{equation*}
$$

Theorem 1.1. Prove the following (known as substitution properties of $\delta_{i j}$ )

$$
\begin{equation*}
\mathrm{u}_{\mathrm{j}}=\delta_{\mathrm{ij}} \mathrm{u}_{\mathrm{i}} \tag{i}
\end{equation*}
$$

(ii) $\delta_{\mathrm{ij}} \mathrm{u}_{\mathrm{jk}}=\mathrm{u}_{\mathrm{ik}}, \delta_{\mathrm{ij}} \mathrm{u}_{\mathrm{ik}}=\mathrm{u}_{\mathrm{jk}}$
(iii) $\delta_{\mathrm{ij}} \mathrm{u}_{\mathrm{ij}}=\mathrm{u}_{11}+\mathrm{u}_{22}+\mathrm{u}_{33}=\mathrm{u}_{\mathrm{ii}}$

Proof. (i) Now $\delta_{i j} \mathbf{u}_{\mathrm{i}}=\delta_{1 \mathrm{j}} \mathrm{u}_{1}+\delta_{2 \mathrm{j}} \mathrm{u}_{2}+\delta_{3 \mathrm{j}} \cdot \mathrm{u}_{3}$
(ii)

$$
\begin{aligned}
&=\mathrm{u}_{\mathrm{j}}+\sum_{\substack{\mathrm{i}=1 \\
\mathrm{i} \neq \mathrm{j}}}^{3} \delta_{\mathrm{ij}} \mathrm{u}_{\mathrm{i}} \\
&=\mathrm{u}_{\mathrm{j}} \\
& \delta_{\mathrm{ij}} \mathrm{u}_{\mathrm{jk}}=\sum_{\mathrm{j}=1}^{3} \delta_{\mathrm{ij}} \mathrm{u}_{\mathrm{jk}} \\
&=\delta_{\mathrm{ii}} \mathrm{u}_{\mathrm{ik}} \quad\left(\text { for } \mathrm{j} \neq \mathrm{i}, \delta_{\mathrm{ij}}=0\right), \text { here summation over } \mathrm{i} \text { is } \\
& \text { not taken }
\end{aligned}
$$

(iii)

$$
\begin{aligned}
& =\mathrm{u}_{\mathrm{ik}} \\
\delta_{\mathrm{ij}} \mathrm{u}_{\mathrm{ij}} & =\sum_{\mathrm{i}}\left[\sum_{\mathrm{j}} \delta_{\mathrm{ij}} \mathrm{u}_{\mathrm{ij}}\right] \\
& =\sum_{\mathrm{i}}\left(1 \cdot \mathrm{u}_{\mathrm{ii}}\right), \text { in } \mathrm{u}_{\mathrm{ii}} \text { summation is not being taken } \\
& =\sum_{\mathrm{i}} \mathrm{u}_{\mathrm{ii}} \\
& =\mathrm{u}_{\mathrm{ii}}=\mathrm{u}_{11}+\mathrm{u}_{22}+\mathrm{u}_{33} .
\end{aligned}
$$

Question. Given that

$$
\mathrm{a}_{\mathrm{ij}}=\alpha \delta_{\mathrm{ij}} \mathrm{~b}_{\mathrm{kk}}+\beta \mathrm{b}_{\mathrm{ij}}
$$

where $\beta \neq 0,3 \alpha+\beta \neq 0$, find $\mathrm{b}_{\mathrm{ij}}$ in terms of $\mathrm{a}_{\mathrm{ij}}$.
Solution. Setting $\mathrm{j}=\mathrm{i}$ and summing accordingly, we obtain

$$
\begin{aligned}
& \mathrm{a}_{\mathrm{ii}}=\alpha \cdot 3 \cdot \mathrm{~b}_{\mathrm{kk}}+\beta \mathrm{b}_{\mathrm{ii}}=(3 \alpha+\beta) \mathrm{b}_{\mathrm{kk}} \\
\Rightarrow \quad & \mathrm{~b}_{\mathrm{kk}}=\frac{1}{3 \alpha+\beta} \mathrm{a}_{\mathrm{kk}}
\end{aligned}
$$

Hence

$$
\begin{aligned}
\mathrm{b}_{\mathrm{ij}} & =\frac{1}{\beta} \mathrm{a}_{\mathrm{ij}}-\alpha \delta_{\mathrm{ij}} \mathrm{~b}_{\mathrm{kk}} . \\
\Rightarrow \quad \mathrm{b}_{\mathrm{ij}} & =\frac{1}{\beta}\left[\mathrm{a}_{\mathrm{ij}}-\frac{\alpha}{3 \alpha+\beta} \delta_{\mathrm{ij}} \mathrm{a}_{\mathrm{kk}}\right] .
\end{aligned}
$$

Theorem 1.2. Prove that

$$
\begin{equation*}
\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qi}}=\delta_{\mathrm{pq}} \tag{i}
\end{equation*}
$$

(ii) $\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{pj}}=\delta_{\mathrm{ij}}$

$$
\begin{equation*}
\left|\mathrm{a}_{\mathrm{ij}}\right|=1, \quad\left(\mathrm{a}_{\mathrm{ij}}\right)^{-1}=\left(\mathrm{a}_{\mathrm{ij}}\right)^{\prime} . \tag{iii}
\end{equation*}
$$

Proof. From the transformation rules of coordinate axes, we have

$$
\begin{align*}
& \mathrm{x}_{\mathrm{p}}^{\prime}=\mathrm{a}_{\mathrm{pi}} \mathrm{x}_{\mathrm{i}}  \tag{1}\\
& \mathrm{x}_{\mathrm{i}}=\mathrm{a}_{\mathrm{pi}} \mathrm{x}_{\mathrm{p}}^{\prime} \tag{2}
\end{align*}
$$

where

$$
\begin{equation*}
\mathrm{a}_{\mathrm{pi}}=\cos \left(\mathrm{x}_{\mathrm{p}}^{\prime}, \mathrm{x}_{\mathrm{i}}\right) \tag{3}
\end{equation*}
$$

(i) Now

$$
\mathrm{x}_{\mathrm{p}}^{\prime}=\mathrm{a}_{\mathrm{pi}} \mathrm{x}_{\mathrm{i}}
$$

$$
\begin{align*}
& =\mathrm{a}_{\mathrm{pi}}\left(\mathrm{a}_{\mathrm{qi}} \mathrm{x}_{\mathrm{q}}^{\prime}\right) \\
& =\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qi}} \mathrm{X}_{\mathrm{q}}^{\prime} \tag{4}
\end{align*}
$$

Also

$$
\begin{equation*}
\mathrm{x}_{\mathrm{p}}^{\prime}=\delta_{\mathrm{pq}} \mathrm{x}^{\prime}{ }_{\mathrm{q}} \tag{5}
\end{equation*}
$$

Therefore,

$$
\left(\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qi}}-\delta_{\mathrm{pq}}\right) \mathrm{x}_{\mathrm{q}}^{\prime}=0
$$

or

$$
\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qi}} \delta_{\mathrm{pq}}=0
$$

or

$$
\begin{equation*}
\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qi}}=\delta_{\mathrm{pq}} \tag{6}
\end{equation*}
$$

This proves that (i).
(ii) Similarly, $\mathrm{x}_{\mathrm{i}}=\mathrm{a}_{\mathrm{pi}} \mathrm{x}_{\mathrm{p}}^{\prime}$

$$
\begin{equation*}
=\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{pj}} \mathrm{x}_{\mathrm{j}} \tag{7}
\end{equation*}
$$

Also

$$
\begin{equation*}
\mathrm{x}_{\mathrm{i}}=\delta_{\mathrm{ij}} \mathrm{x}_{\mathrm{j}} \tag{8}
\end{equation*}
$$

Hence,

$$
\begin{equation*}
\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{pj}}=\delta_{\mathrm{ij}} . \tag{9}
\end{equation*}
$$

(iii) Relation (6) gives, in the expanded form,

$$
\begin{align*}
& a_{11}^{2}+a_{12}^{2}+a_{13}^{2}=1, \\
& a_{21}^{2}+a_{22}^{2}+a_{23}^{2}=1, \\
& a_{31}^{2}+a_{32}^{2}+a_{33}^{2}=1 \\
& a_{11} a_{21}+a_{12} a_{22}+a_{13} a_{23}=0, \\
& a_{21} a_{31}+a_{22} a_{32}+a_{23} a_{33}=0, \\
& a_{31} a_{11}+a_{32} a_{12}+a_{33} a_{13}=0 \tag{10}
\end{align*}
$$

The relations (6) and (9) are referred to as the orthonormal relations for $\mathrm{a}_{\mathrm{ij}}$. In matrix notation, relations (6) and (9) may be represented respectively, as follows

$$
\begin{align*}
& {\left[\begin{array}{lll}
a_{11} & a_{12} & a_{13} \\
a_{21} & a_{22} & a_{23} \\
a_{31} & a_{32} & a_{33}
\end{array}\right]\left[\begin{array}{lll}
a_{11} & a_{21} & a_{31} \\
a_{12} & a_{22} & a_{32} \\
a_{13} & a_{23} & a_{33}
\end{array}\right]=\left[\begin{array}{lll}
1 & 0 & 0 \\
0 & 1 & 0 \\
0 & 0 & 1
\end{array}\right]}  \tag{11}\\
& {\left[\begin{array}{lll}
a_{11} & a_{21} & a_{31} \\
a_{12} & a_{22} & a_{32} \\
a_{13} & a_{23} & a_{33}
\end{array}\right]\left[\begin{array}{lll}
a_{11} & a_{12} & a_{13} \\
a_{21} & a_{22} & a_{23} \\
a_{31} & a_{32} & a_{33}
\end{array}\right]=\left[\begin{array}{lll}
1 & 0 & 0 \\
0 & 1 & 0 \\
0 & 0 & 1
\end{array}\right]} \tag{12}
\end{align*}
$$

or

$$
\begin{equation*}
\mathrm{LL}^{\prime}=\mathrm{L}^{\prime} \mathrm{L}=1 \tag{13}
\end{equation*}
$$

These expressions show that the matrix $\mathrm{L}=\left(\mathrm{a}_{\mathrm{ij}}\right)$ is non-singular and that

$$
\begin{equation*}
\left(\mathrm{a}_{\mathrm{ij}}\right)^{-1}=\left(\mathrm{a}_{\mathrm{ij}}\right)^{\prime} \text { and }\left|\mathrm{a}_{\mathrm{ij}}\right|=1 . \tag{14}
\end{equation*}
$$

The transformation matrix $\mathrm{L}=\left(\mathrm{a}_{\mathrm{ij}}\right)$ is called the proper orthogonal matrix. For this reason, the transformation laws (3) and (4), determined by the matrix L $=\left(\mathrm{a}_{\mathrm{ij}}\right)$, are called orthogonal transformations.

Example. The $\mathrm{x}^{\prime}{ }_{i}$-system is obtained by rotating the $\mathrm{x}_{\mathrm{i}}$-system about the $\mathrm{x}_{3}$ axis through an angle $\theta$ in the sense of right handed screw. Find the transformation matrix. If a point P has coordinates $(1,1,1)$ in the $\mathrm{x}_{\mathrm{i}}$-system, find its coordinate in the $\mathrm{x}^{\prime}{ }_{\mathrm{i}}$-system. If a point Q has coordinate $(1,1,1)$ in the $\mathrm{x}_{\mathrm{i}}^{\prime}$-system, find its coordinates in the $\mathrm{x}_{\mathrm{i}}$-system.

Solution. The figure (1.2) shows how the $\mathrm{x}^{\prime}{ }_{i}$-system is related to the $\mathrm{x}_{\mathrm{i}}{ }^{-}$ system. The table of direction cosines for the given transform is


Fig. (1.2)

|  | $\hat{\mathrm{e}}_{1}$ | $\hat{\mathrm{e}}_{2}$ | $\hat{\mathrm{e}}_{3}$ |
| :---: | :---: | :---: | :---: |
| $\hat{\mathrm{e}}^{\prime}{ }_{1}$ | $\cos \theta$ | $\sin \theta$ | 0 |
| $\hat{\mathrm{e}}^{\prime}{ }_{2}$ | $-\sin \theta$ | $\cos \theta$ | 0 |
| $\hat{\mathrm{e}}^{\prime}{ }_{3}$ | 0 | 0 | 1 |

Hence, the matrix of this transformation is

$$
\left(\mathrm{a}_{\mathrm{ij}}\right)=\left(\begin{array}{ccc}
\cos \theta & \sin \theta & 0  \tag{1}\\
-\sin \theta & \cos \theta & 0 \\
0 & 0 & 1
\end{array}\right)
$$

The transformation rules for coordinates are

$$
\begin{align*}
& \mathrm{x}_{\mathrm{p}}^{\prime}=\mathrm{a}_{\mathrm{pi}} \mathrm{x}_{\mathrm{i}}  \tag{2}\\
& \mathrm{x}_{\mathrm{i}}=\mathrm{a}_{\mathrm{pi}} \mathrm{x}_{\mathrm{p}}^{\prime} . \tag{3}
\end{align*}
$$

The coordinates $\mathrm{P}\left(\mathrm{x}^{\prime}{ }_{1}, \mathrm{x}^{\prime}{ }_{2}, \mathrm{x}^{\prime}{ }_{3}\right)$ of the point $\mathrm{P}(1,1,1)$ in the new system are given by

$$
\begin{align*}
& \mathrm{x}^{\prime}{ }_{1}=\mathrm{a}_{1 \mathrm{i}} \mathrm{X}_{\mathrm{i}}=\mathrm{a}_{11} \mathrm{x}_{1}+\mathrm{a}_{12} \mathrm{X}_{2}+\mathrm{a}_{13} \mathrm{x}_{3}=\cos \theta+\sin \theta \\
& \mathrm{x}^{\prime}{ }_{2}=\mathrm{a}_{2 \mathrm{i}} \mathrm{X}_{\mathrm{i}}=\mathrm{a}_{21} \mathrm{x}_{1}+\mathrm{a}_{22} \mathrm{x}_{2}+\mathrm{a}_{23} \mathrm{x}_{3}=\cos \theta-\sin \theta \\
& \mathrm{x}^{\prime}{ }_{3}=\mathrm{a}_{3 \mathrm{i}} \mathrm{X}_{\mathrm{i}}=\mathrm{a}_{31} \mathrm{x}_{1}+\mathrm{a}_{32} \mathrm{X}_{2}+\mathrm{a}_{33} \mathrm{x}_{3}=1 . \tag{4}
\end{align*}
$$

Therefore, coordinates of P in the $\mathrm{x}_{\mathrm{i}}^{\prime}$-system are $(\cos \theta+\sin \theta, \cos \theta-\sin \theta, 1)$. The coordinates ( $\mathrm{x}_{1}, \mathrm{x}_{2}, \mathrm{x}_{3}$ ) of a w.r.t. old system are given as

$$
\begin{align*}
& \mathrm{x}_{1}=\mathrm{a}_{\mathrm{p} 1} \mathrm{x}_{\mathrm{p}}^{\prime}=\mathrm{a}_{11} \mathrm{x}^{\prime}{ }_{1}+\mathrm{a}_{21} \mathrm{x}^{\prime}{ }_{2}+\mathrm{a}_{31} \mathrm{x}^{\prime}{ }_{3}=\cos \theta-\sin \theta \\
& \mathrm{x}_{2}=\mathrm{a}_{\mathrm{p} 2} \mathrm{x}^{\prime}{ }_{\mathrm{p}}=\mathrm{a}_{12} \mathrm{x}^{\prime}{ }_{1}+\mathrm{a}_{22} \mathrm{x}^{\prime}{ }_{2}+\mathrm{a}_{32} \mathrm{x}_{3}^{\prime}=\cos \theta+\sin \theta \\
& \mathrm{x}_{3}=\mathrm{a}_{\mathrm{p} 3} \mathrm{x}_{\mathrm{p}}=\mathrm{a}_{13} \mathrm{x}^{\prime}{ }_{1}+\mathrm{a}_{23} \mathrm{x}^{\prime}{ }_{2}+\mathrm{a}_{33} \mathrm{x}^{\prime}{ }_{3}=1 . \tag{5}
\end{align*}
$$

Hence, the coordinates of the point Q in the old $\mathrm{x}_{\mathrm{i}}$-system are $(\cos \theta-\sin \theta, \cos \theta$ $+\sin \theta, 1$ ).

### 1.4 SCALARS AND VECTORS

Under a transformation of certesian coordinate axes, a scalar quantity, such as the density or the temperature, remains unchanged. This means that a scalar is an invariant under a coordinate transformation. Scalaras are called tensors of zero rank.

We know that a scalar is represented by a single quantity in any coordinate system. Accordingly, a tensor of zero rank (or order) is specified in any coordinate system in three-dimensional space by one component or a single number.

All physical quantities having magnitude only are tensors of zero order.

## Transformation of a Vector

Let $\bar{u}$ be any vector having components $\left(u_{1}, u_{2}, u_{3}\right)$ along the $x_{i}$-axes and components ( $\mathrm{u}_{1}^{\prime}, \mathrm{u}^{\prime}{ }_{2}, \mathrm{u}^{\prime}{ }_{3}$ ) along the $\mathrm{x}^{\prime}{ }_{i}$-axes so that vector u is represented by three components/quantities. Then we have

$$
\begin{align*}
\overline{\mathrm{u}} & =\mathrm{u}_{\mathrm{i}} \hat{\mathrm{e}}_{\mathrm{i}}  \tag{1}\\
\text { and } \quad \mathrm{u} & =\mathrm{u}_{\mathrm{i}}^{\prime} \hat{\mathrm{e}}_{\mathrm{i}}^{\prime} \tag{2}
\end{align*}
$$

where $\hat{\mathrm{e}}_{\mathrm{i}}$ is the unit vector along $\mathrm{x}_{\mathrm{i}}$-direction and $\hat{\mathrm{e}}_{\mathrm{i}}^{\prime}$ is the unit vector along $\mathrm{x}_{\mathrm{i}^{\prime}}{ }^{-}$ direction.
Now

$$
\mathrm{u}_{\mathrm{p}}^{\prime}=\overline{\mathrm{u}} \cdot \hat{\mathrm{e}}_{\mathrm{p}}^{\prime}
$$

$$
\begin{align*}
& =\left(u_{i} \hat{e}_{\mathrm{i}}\right) \cdot \mathrm{e}_{\mathrm{p}}^{\prime} \\
& =\left(\hat{\mathrm{e}}_{\mathrm{p}}^{\prime} \cdot \hat{\mathrm{e}}_{\mathrm{i}}\right) \mathrm{u}_{\mathrm{i}} \\
\Rightarrow \quad \mathrm{u}_{\mathrm{p}}^{\prime} & =\mathrm{a}_{\mathrm{pi}} \mathrm{u}_{\mathrm{i}} . \tag{3}
\end{align*}
$$

where

$$
\begin{equation*}
\mathrm{a}_{\mathrm{pi}}=\hat{\mathrm{e}}_{\mathrm{p}}^{\prime} \cdot \hat{\mathrm{e}}_{\mathrm{i}}=\cos \left(\mathrm{x}_{\mathrm{p}}^{\prime}, \mathrm{x}_{\mathrm{i}}\right) . \tag{4}
\end{equation*}
$$

Also

$$
\begin{align*}
\mathrm{u}_{\mathrm{i}} & =\overline{\mathrm{u}} \cdot \hat{\mathrm{e}}_{\mathrm{i}} \\
& =\left(\mathrm{u}_{\mathrm{p}}^{\prime} \hat{\mathrm{e}}_{\mathrm{p}}^{\prime}\right) \cdot \hat{\mathrm{e}}_{1} \\
& =\left(\hat{\mathrm{e}}_{\mathrm{p}}^{\prime} \cdot \hat{e}_{\mathrm{i}}^{\prime}\right) \mathrm{u}_{\mathrm{p}}^{\prime} \\
& =\mathrm{a}_{\mathrm{pi}} \mathrm{u}_{\mathrm{p}}^{\prime}, \tag{5}
\end{align*}
$$

where ( $\mathrm{a}_{\mathrm{pi}}$ ) is the proper orthogonal transformation matrix.
Relations (3) and (5) are the rules that determine $u_{p}^{\prime}$ in terms of $u_{i}$ and viceversa. Evidently, these relations are analogous to rules of transformation of coordinates.

## Definition (Tensor of rank one)

A cartesian tensor of rank one is an entity that may be represented by a set of three quantities in every cartesian coordinate system with the property that its components $u_{i}$ relative to the system $\mathrm{ox}_{1} \mathrm{x}_{2} \mathrm{x}_{3}$ are related/connected with its components $\mathrm{u}_{\mathrm{p}}$ relative to the system $\mathrm{ox}^{\prime}{ }_{1} \mathrm{x}^{\prime}{ }_{2} \mathrm{x}^{\prime}{ }_{3}$ by the relation

$$
\mathrm{u}_{\mathrm{p}}^{\prime}=\mathrm{a}_{\mathrm{p} i} \mathbf{u}_{\mathrm{i}}
$$

where the law of transformation of coordinates of points is

$$
\mathrm{x}_{\mathrm{p}}^{\prime}=\mathrm{a}_{\mathrm{p} i} \mathrm{x}_{\mathrm{i}} \quad \text { and } \mathrm{a}_{\mathrm{pi}}=\cos \left(\mathrm{x}_{\mathrm{p}}^{\prime}, \mathrm{x}_{\mathrm{i}}\right)=\hat{\mathrm{e}}_{\mathrm{p}}^{\prime} . \hat{\mathrm{e}}_{\mathrm{i}} .
$$

Note: We note that every vector in space is a tensor of rank one. Thus, physical quantities possessing both magnitude and direction such as force, displacement, velocity, etc. are all tensors of rank one. In three-dimensional space, 3 real numbers are needed to represent a tensor of order 1.

## Definition (Tensor of order 2)

Any entity representable by a set of 9 (real) quantities relative to a system of rectangular axes is called a tensor of rank two if its components $\mathrm{w}_{\mathrm{ij}}$ relative to system $\mathrm{ox}_{1} \mathrm{x}_{2} \mathrm{x}_{3}$ are connected with its components $\mathrm{w}_{\mathrm{pq}}^{\prime}$ relative to the system $\mathrm{ox}^{\prime}{ }_{1} \mathrm{X}^{\prime}{ }_{2} \mathrm{x}^{\prime}{ }_{3}$ by the transformation rule

$$
\mathrm{w}_{\mathrm{pq}}^{\prime}=\mathrm{a}_{\mathrm{pi}} \mathrm{q}_{\mathrm{qj}} \mathrm{w}_{\mathrm{ij}}
$$

when the law of transformation of coordinates is

$$
\mathrm{x}_{\mathrm{p}}^{\prime}=\mathrm{a}_{\mathrm{pi} i} \mathrm{x}_{\mathrm{i}},
$$

$$
\mathrm{a}_{\mathrm{pi}}=\cos \left(\mathrm{x}_{\mathrm{p}}^{\prime}, \mathrm{x}_{\mathrm{i}}\right)=\hat{\mathrm{e}}_{\mathrm{p}}^{\prime} . \hat{\mathrm{e}}_{\mathrm{i}} .
$$

Note : Tensors of order 2 are also called dyadics. For example, strain and stress tensors of elasticity are, each of rank 2. In the theory of elasticity, we shall use tensors of rank 4 also.

Example. In the $\mathrm{x}_{\mathrm{i}}$-system, a vector $\overline{\mathrm{u}}$ has components $(-1,0,1)$ and a second order tensor has the representation

$$
\left(\mathrm{w}_{\mathrm{ij}}\right)=\left(\begin{array}{ccc}
0 & 1 & 0 \\
-1 & 0 & 2 \\
0 & -2 & 0
\end{array}\right)
$$

The $\mathrm{x}^{\prime}{ }_{\mathrm{i}}$-system is obtained by rotating the $\mathrm{x}_{\mathrm{i}}$-system about the $\mathrm{x}_{3}$-axis through an angle of $45^{\circ}$ in the sense of the right handed screw. Find the components of the vector $\overline{\mathrm{u}}$ and the second ordered tensor in the $\mathrm{x}_{\mathrm{i}}{ }_{\mathrm{i}}$-system.

Solution. The table of transformation of coordinates is

|  | $\mathrm{x}_{1}$ | $\mathrm{x}_{2}$ | $\mathrm{x}_{3}$ |
| :--- | :--- | :--- | :--- |
| $\mathrm{x}^{\prime}{ }_{1}$ | $\frac{1}{\sqrt{2}}$ | $\frac{1}{\sqrt{2}}$ | 0 |
| $\mathrm{x}^{\prime}{ }_{2}$ | $-\frac{1}{\sqrt{2}}$ | $\frac{1}{\sqrt{2}}$ | 0 |
| $\mathrm{x}^{\prime}{ }_{3}$ | 0 | 0 | 1 |

If $u_{p}^{\prime}$ are the components of vector in the new system, then

$$
\begin{equation*}
\mathrm{u}_{\mathrm{p}}^{\prime}=\mathrm{a}_{\mathrm{pi}} \mathrm{u}_{\mathrm{i}}, \mathrm{a}_{\mathrm{pi}}=\hat{\mathrm{e}}_{\mathrm{p}}^{\prime} . \hat{e}_{\mathrm{i}} . \tag{1}
\end{equation*}
$$

This gives $\quad \mathrm{u}^{\prime}{ }_{1}=-\frac{1}{\sqrt{2}}$

$$
\begin{aligned}
& \mathrm{u}_{2}^{\prime}=\frac{1}{\sqrt{2}}, \\
& \mathrm{u}_{3}^{\prime}=1 .
\end{aligned}
$$

Let $\mathrm{w}^{\prime}{ }_{\mathrm{pq}}$ be the components of the given second order tensor in the $\mathrm{x}^{\prime}{ }_{\mathrm{i}}$-system. Then the transformation law for second order tensor yields

$$
\begin{equation*}
\mathrm{w}_{\mathrm{pq}}^{\prime}=\mathrm{a}_{\mathrm{pi} \mathrm{i}} \mathrm{a}_{\mathrm{ij}} \mathrm{w}_{\mathrm{ij}} \tag{2}
\end{equation*}
$$

We find (left an exercise to readers)

$$
\mathrm{w}_{\mathrm{pq}}^{\prime}=\left(\begin{array}{ccc}
0 & 1 & \sqrt{2}  \tag{3}\\
-1 & 0 & \sqrt{2} \\
-\sqrt{2} & -\sqrt{2} & 0
\end{array}\right)
$$

Definition 1. A second order tensor $\mathrm{u}_{\mathrm{ij}}$ is said to be symmetric if

$$
\mathrm{u}_{\mathrm{ij}}=\mathrm{u}_{\mathrm{ij}} \text { for all } \mathrm{i} \text { and } \mathrm{j} .
$$

Definition 2. A second order tensor $\mathrm{u}_{\mathrm{ij}}$ is said to be skew-symmetric if

$$
\mathrm{u}_{\mathrm{ij}}=-\mathrm{u}_{\mathrm{ji}} \text { for all } \mathrm{i} \text { and } \mathrm{j} \text {. }
$$

## Tensors of order $n$

A tensor of order $n$ has $3^{n}$ components. If $u_{i j k} \ldots$...are components of a tensor of order n , then, the transformation law is

$$
\mathrm{u}_{\mathrm{pqr}}^{\prime} \ldots \ldots=\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qj}} \mathrm{a}_{\mathrm{rk}} \ldots . \mathrm{u}_{\mathrm{ijk}} \ldots \ldots
$$

where the law of transformation of coordinates is

$$
\mathrm{x}_{\mathrm{p}}^{\prime}=\mathrm{a}_{\mathrm{pi}} \mathrm{x}_{\mathrm{i}}
$$

and

$$
\mathrm{a}_{\mathrm{pi}}=\cos \left(\mathrm{x}_{\mathrm{p}}^{\prime}, \mathrm{x}_{\mathrm{i}}\right) .
$$

## Importance of the Concept of Tensors

(a) Tensors are quantities describing the same phenomenon regardless of the coordinate system used. Therefore, tensors provide an important guide in the formulation of the correct forms of physical laws.
(b) The tensor concept gives us a convenient means of transforming an equation from one system of coordinates to another.
(c) An advantage of the use of cartesian tensors is that once the properties of a tensor of order n have been established, they hold for all such tensors regardless of the physical phenomena they present.

Note : For example, in the study of strain, stress, inertia properties of rigid bodies, the common bond is that they are all symmetric tensors of rank two.
(d) With the use of tensors, equations are condensed, such as

$$
\tau_{\mathrm{ij}, \mathrm{j}}+\mathrm{f}_{\mathrm{i}}=0
$$

is the equation of equilibrium in tensor form. It consists of 3 equations and each equation has 4 terms.
(e) Equations describing physical laws must be tensorially homogeneous, which means that every term of the equation must be a tensor of the same rank.

### 1.5 PROPERTIES OF TENSORS

Property 1 : If all components of a tensor are 0 in one coordinate system, then they are 0 in all coordinate systems.

Proof : Let $\mathrm{u}_{\mathrm{ijkl}} \ldots$. and $\mathrm{u}_{\mathrm{pqrs}}^{\prime} \ldots$. .be the components of a nth order tensor in two systems $0 \mathrm{x}_{1} \mathrm{X}_{2} \mathrm{x}_{3}$ and $0 \mathrm{x}_{1}{ }^{\prime} \mathrm{x}_{2}{ }^{\prime} \mathrm{x}_{3}{ }^{\prime}$, respectively.

Suppose that

$$
\mathrm{u}_{\mathrm{ij} k l} / \ldots . .=0 .
$$

Then, the transformation rule yields

$$
\mathrm{u}_{\mathrm{pqrs}}^{\prime} \ldots \ldots=\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qj}} \mathrm{a}_{\mathrm{rk}} \mathrm{a}_{\mathrm{sm}} \ldots \mathrm{u}_{\mathrm{ijk} /} \ldots
$$

giving

$$
\mathrm{u}_{\mathrm{pqr}}^{\prime} \ldots . .=0
$$

This proves the result.

## Zero Tensor

A tensor whose all components in one Cartesian coordinates system are 0 is called a zero tensor.

A zero tensor may have any order $n$.
Property 2 : If the corresponding components of two tensors of the same order are equal in one coordinate system, then they are equal in all coordinate systems.

Corollary : A tensor equation which holds in one Cartesian coordinate system also holds in every other Cartesian coordinate system.

## Equality of Tensors

Two tensors of the same order whose corresponding components are equal in a coordinate system (and hence in all coordinates) are called equal tensors.

Note : Thus, in order to show that two tensors are equal, it is sufficient to show that their corresponding components are equal in any one of the coordinate systems.

## Property 3 (Scalar multiplication of a tensor)

If the components of a tensor of order $n$ are multiplied by a scalar $\alpha$, then the resulting components form a tensor of the same order $n$.

Proof : Let $\mathrm{u}_{\mathrm{ijk}} \ldots$. be a tensor of order $n$. Let $\mathrm{u}^{\prime}{ }_{\mathrm{pqr}} \ldots$. be the corresponding components in the dashed system $0 \mathrm{x}_{1}{ }^{\prime} \mathrm{x}_{2}{ }^{\prime} \mathrm{x}_{3}{ }^{\prime}$. The transformation rule for a tensor of order n yields

$$
\begin{equation*}
\mathrm{u}_{\mathrm{pqr}}^{\prime} \ldots . .=\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qj}} \ldots . . \mathrm{u}_{\mathrm{ijkl}} \ldots . \tag{1}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathrm{a}_{\mathrm{pi}}=\cos \left(\mathrm{x}_{\mathrm{p}}{ }^{\prime}, \mathrm{x}_{\mathrm{i}}\right) . \tag{2}
\end{equation*}
$$

Now

$$
\begin{equation*}
\left(\alpha \mathrm{u}_{\mathrm{pqr}}^{\prime} \cdots . .\right)=\left(\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qj}} \ldots \ldots\right)\left(\alpha \mathrm{u}_{\mathrm{ijk}} \ldots .\right) \tag{3}
\end{equation*}
$$

This shows that components $\alpha u_{i j k} \ldots$ form a tensor of rank $n$.

## Tensor Equations

An equation of the form

$$
\alpha_{\mathrm{ijk}}-\beta_{\mathrm{ij}} \mathrm{u}_{\mathrm{k}}=0
$$

is called a tensor equation.
For checking the correctness of a tensor equation, there are following two rules :

Rule (i) In a correctly written tensor equation, no suffix shall appear more than twice in a term, otherwise, the operation will not be defined. For example, an equation

$$
\mathrm{u}_{\mathrm{j}}^{\prime}=\alpha_{\mathrm{ij}} \mathrm{u}_{\mathrm{j}} \mathrm{v}_{\mathrm{j}}
$$

is not a tensor equation.
Rule (ii) If a suffix appears only once in a term, then it must appear only once in the remaining terms also. For example, an equation

$$
\mathrm{u}_{\mathrm{j}}^{\prime}-l_{\mathrm{ij}} \mathrm{u}_{\mathrm{j}}=0
$$

is not a tensor equation.
Here j appears once in the first term while it appears twice in the second term.

## Property 4 (Sum and Difference of tensors)

If $\mathrm{u}_{\mathrm{ijk}} \ldots$ and $\mathrm{v}_{\mathrm{ijk}} \ldots \ldots$ are two tensors of the same rank n then the sums

$$
\left(\mathrm{u}_{\mathrm{ijk}} \ldots \ldots+\mathrm{v}_{\mathrm{ijk}} \ldots . .\right)
$$

of their components are components of a tensor of the same order $n$.

Proof : Let

$$
\begin{equation*}
\mathrm{w}_{\mathrm{ijk}} \ldots . .=\mathrm{u}_{\mathrm{ijk}} \ldots .+\mathrm{v}_{\mathrm{ijk}} \ldots \ldots \tag{1}
\end{equation*}
$$

Let $\mathrm{u}^{\prime}{ }_{\mathrm{pqr}} \ldots .$. and $\mathrm{v}^{\prime}{ }_{\mathrm{pqr}} \ldots$. be the components of the given tensors of order n relative to the new dashed system $0 \mathrm{x}_{1}{ }^{\prime} \mathrm{x}_{2}{ }^{\prime} \mathrm{x}_{3}{ }^{\prime}$. Then, transformation rules for these tensors are

$$
\begin{equation*}
\mathrm{u}_{\mathrm{pqr}}^{\prime} \ldots . .=\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qj}} \ldots . \mathrm{u}_{\mathrm{ijk}} \ldots \tag{2}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathrm{v}_{\mathrm{pqr}}^{\prime} \ldots \ldots=\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qj}} \ldots . \mathrm{v}_{\mathrm{ijk} \mathrm{l}} \ldots \tag{3}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathrm{a}_{\mathrm{pi}}=\cos \left(\mathrm{x}_{\mathrm{p}}^{\prime}, \mathrm{x}_{\mathrm{i}}\right) \tag{4}
\end{equation*}
$$

Let

$$
\begin{equation*}
\mathrm{w}_{\mathrm{pqr}}^{\prime} \ldots \ldots=\mathrm{u}_{\mathrm{pqr}}^{\prime} \ldots \ldots \mathrm{v}_{\mathrm{pqr}}^{\prime} \ldots \ldots \tag{5}
\end{equation*}
$$

Then equations (2) - (5) give

$$
\begin{equation*}
\mathrm{w}_{\mathrm{pqr}}^{\prime} \ldots \ldots=\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{ij}} \ldots . \mathrm{w}_{\mathrm{ijk}} \ldots \ldots . \tag{6}
\end{equation*}
$$

Thus quantities $w_{i j k} \ldots$. obey the transformation rule of a tensor of order $n$. Therefore, they are components of a tensor of rank $n$.

Corollary : Similarly, $\mathrm{u}_{\mathrm{ijk}} \ldots-\mathrm{v}_{\mathrm{ikl}} \ldots$ are components of a tensor of rank n .

## Property 5 (Tensor Multiplication)

The product of two tensors is also a tensor whose order is the sum of orders of the given tensors.

Proof : Let $\mathrm{u}_{\mathrm{ijk}} \ldots \ldots$ and $\mathrm{v}_{\mathrm{pqr}} \ldots \ldots$ be two tensors of order m and n respectively. We shall show that the product

$$
\begin{equation*}
\mathrm{w}_{\mathrm{ijk}} \cdots \cdots \mathrm{pqr} \cdots=\mathrm{u}_{\mathrm{ijk}} \cdots . \mathrm{v}_{\mathrm{pqr}} \cdots \tag{1}
\end{equation*}
$$

is a tensor of order $m+n$.
Let $u_{i_{1} j_{1}}^{\prime} \ldots$ and $v_{p_{1} q_{1}}^{\prime} \ldots$. be the components of the given tensors of orders $m$ and n relative to the new system $0 \mathrm{x}_{1}{ }^{\prime} \mathrm{x}_{2}{ }^{\prime} \mathrm{x}_{3}{ }^{\prime}$. Then

$$
\begin{equation*}
\mathrm{u}_{\mathrm{i}_{1} \mathrm{j}_{1}}^{\prime} \ldots=\mathrm{a}_{\mathrm{i}_{1} \mathrm{i}} \mathrm{a}_{\mathrm{j}_{1} \mathrm{j}} \ldots \ldots \mathrm{u}_{\mathrm{ijk}} \ldots \ldots \ldots \tag{2}
\end{equation*}
$$

$$
\begin{equation*}
\mathrm{v}_{\mathrm{p}_{1} \mathrm{q}_{1}}^{\prime} \ldots \ldots=\mathrm{a}_{\mathrm{p}_{1} \mathrm{p}} \mathrm{a}_{\mathrm{q}_{1} \mathrm{q}} \ldots \ldots . \mathrm{v}_{\mathrm{pqr}} \ldots \ldots \tag{3}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathrm{a}_{\mathrm{pi}}=\cos \left(\mathrm{x}_{\mathrm{p}}^{\prime}, \mathrm{x}_{\mathrm{i}}\right) . \tag{4}
\end{equation*}
$$

Let

$$
\begin{equation*}
w_{i_{1} j_{1}}^{\prime} \cdots \cdot{ }_{p_{1} q_{1}} \ldots .=u_{i_{1} j_{1}}^{\prime} \ldots \cdots v_{p_{1} q_{1}}^{\prime} \cdots \cdots \tag{5}
\end{equation*}
$$

Multiplying (2) an (3), we get

This shows that components $\mathrm{W}_{\mathrm{ijk}} \ldots$ pqr $\ldots$. obey the transformation rule of a tensor of order $(\mathrm{m}+\mathrm{n})$. Hence $\mathrm{u}_{\mathrm{ijk}} \ldots . \mathrm{v}_{\mathrm{pqr}} \ldots$. are components of a $(\mathrm{m}+\mathrm{n})$ order tensor.

Exercise 1 : If $u_{i}$ and $v_{i}$ are components of vectors, then show that $u_{i} v_{j}$ are components of a second - order tensor.

Exercise 2: If $a_{i j}$ are components of a second - order tensor and $b_{i}$ are components of a vector, show that $\mathrm{a}_{\mathrm{ij}} \mathrm{b}_{\mathrm{k}}$ are components of a third order tensor.

Exercise 3 : If $a_{i j}$ and $b_{i j}$ are components of two second - order tensors show that $\mathrm{a}_{\mathrm{ij}} \mathrm{b}_{\mathrm{km}}$ are components of fourth - order tensor.

Exercise 4 : Let $u_{i}$ and $v_{i}$ be two vectors. Let $w_{i j}=u_{i} v_{j}+u_{j} v_{i}$ and $\alpha_{i j}=u_{i} v_{j}-$ $u_{j} v_{i}$. Show that each of $w_{i j}$ and $\alpha_{i j}$ is a second order tensor.

## Property 6 (Contraction of a tensor)

The operation or process of setting two suffices equal in a tensor and then summing over the dummy suffix is called a contraction operation or simply a contraction.

The tensor resulting from a contraction operation is called a contraction of the original tensor.

Contraction operations are applicable to tensor of all orders (higher than 1) and each such operation reduces the order of a tensor by 2.

Theorem : Prove that the result of applying a contraction to a tensor of order $n$ is a tensor of order $\mathrm{n}-2$.

Proof : Let $u_{i j k} \ldots \ldots$ and $u^{\prime}{ }_{p q r} \ldots$. be the components of the given tensor of order n relative to two cartesian coordinate systems $0 \mathrm{x}_{1} \mathrm{x}_{2} \mathrm{x}_{3}$ and $0 \mathrm{x}_{1}{ }^{\prime} \mathrm{x}_{2}{ }^{\prime} \mathrm{x}_{3}{ }^{\prime}$. The rule of transformation of tensors is

$$
\begin{equation*}
\mathrm{u}_{\mathrm{pqr}}^{\prime} \ldots \ldots=\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qj}} \mathrm{a}_{\mathrm{rk}} \ldots \ldots . \mathrm{u}_{\mathrm{ijk}} \ldots \ldots . \tag{1}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathrm{a}_{\mathrm{pi}}=\cos \left(\mathrm{x}_{\mathrm{p}}^{\prime}, \mathrm{x}_{\mathrm{i}}\right) . \tag{2}
\end{equation*}
$$

Without loss of generality, we contract the given tensor by setting $j=i$ and using summation convention . Let

Now

$$
\begin{align*}
\mathrm{v}_{\mathrm{k} l} \ldots & =\mathrm{u}_{\mathrm{iik} l} \ldots \ldots  \tag{3}\\
\mathrm{u}_{\mathrm{pqr}}^{\prime} \ldots \ldots & =\left(\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qij}}\right) \mathrm{a}_{\mathrm{rk}} \ldots \ldots . \mathrm{u}_{\mathrm{iik} l} \ldots \ldots \\
& =\delta_{\mathrm{pq}} \mathrm{a}_{\mathrm{rk}} \ldots . \mathrm{u}_{\mathrm{k} l} \ldots
\end{align*}
$$

This gives

$$
\mathrm{u}_{\mathrm{pprs}}^{\prime} \ldots .=\mathrm{a}_{\mathrm{rk}} \mathrm{a}_{\mathrm{s} l} \ldots \ldots \mathrm{v}_{\mathrm{k} l} \ldots \ldots
$$

or

$$
\begin{equation*}
\mathrm{v}_{\mathrm{rs}}^{\prime} \ldots \ldots . .=\mathrm{a}_{\mathrm{rk}} \ldots \mathrm{v}_{\mathrm{k} l} \ldots \tag{4}
\end{equation*}
$$

## Property 7 (Quotient laws)

## Quotient law is the partial converse of the contraction law.

Theorem : If there is an entity representable by the set of 9 quantities $u_{i j}$ relative to any given system of cartesian axes and if $u_{i j} \mathrm{v}_{\mathrm{j}}$ is a vector for an arbitrary vector $\mathrm{v}_{\mathrm{i}}$, then show that $\mathrm{u}_{\mathrm{ij}}$ is second order tensor.

Proof: Let

$$
\begin{equation*}
\mathrm{w}_{\mathrm{i}}=\mathrm{u}_{\mathrm{ij}} \mathrm{v}_{\mathrm{j}} \tag{1}
\end{equation*}
$$

Suppose that $\mathrm{u}_{\mathrm{pq}}^{\prime}, \mathrm{u}_{\mathrm{p}}^{\prime}, \mathrm{w}_{\mathrm{p}}^{\prime}$ be the corresponding components in the dashed system $0 \mathrm{x}_{1}{ }^{\prime} \mathrm{x}_{2}{ }^{\prime} \mathrm{x}_{3}{ }^{\prime}$. Then

$$
\begin{align*}
& \mathrm{v}_{\mathrm{q}}^{\prime}=\mathrm{a}_{\mathrm{qj}} \mathrm{v}_{\mathrm{j}},  \tag{2}\\
& \mathrm{w}_{\mathrm{p}}^{\prime}=\mathrm{a}_{\mathrm{pi}} \mathrm{w}_{\mathrm{i}}, \tag{3}
\end{align*}
$$

where

$$
\begin{equation*}
\mathrm{a}_{\mathrm{pi}}=\cos \left(\mathrm{x}_{\mathrm{p}}^{\prime}, \mathrm{x}_{\mathrm{i}}\right) . \tag{4}
\end{equation*}
$$

Equation (1) in the dashed system is

$$
\begin{equation*}
\mathrm{w}_{\mathrm{p}}^{\prime}=\mathrm{u}_{\mathrm{pq}}^{\prime} \mathrm{v}_{\mathrm{q}}^{\prime} \tag{5}
\end{equation*}
$$

Inverse laws of (2) and (3) are

$$
\begin{align*}
\mathrm{v}_{\mathrm{j}} & =\mathrm{a}_{\mathrm{qj}} \mathrm{v}_{\mathrm{q}}^{\prime},  \tag{6}\\
\mathrm{w}_{\mathrm{i}} & =\mathrm{a}_{\mathrm{pi}} \mathrm{w}_{\mathrm{p}}^{\prime} . \tag{7}
\end{align*}
$$

Now

$$
\begin{aligned}
\mathrm{u}_{\mathrm{pq}}^{\prime} \mathrm{v}_{\mathrm{q}}{ }^{\prime} & =\mathrm{w}_{\mathrm{p}}{ }^{\prime} \\
& =\mathrm{a}_{\mathrm{pi}} \mathrm{w}_{\mathrm{i}} \\
& =\mathrm{a}_{\mathrm{pi}}\left(\mathrm{u}_{\mathrm{ij}} \mathrm{v}_{\mathrm{j}}\right) \\
& =\mathrm{a}_{\mathrm{pi}}\left(\mathrm{a}_{\mathrm{qj}} \mathrm{v}_{\mathrm{q}}\right) \mathrm{u}_{\mathrm{ij}} \\
& =\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qj}} \mathrm{u}_{\mathrm{ij}} \mathrm{v}_{\mathrm{q}}{ }^{\prime} .
\end{aligned}
$$

This gives

$$
\begin{equation*}
\left(\mathrm{u}_{\mathrm{pq}}^{\prime}-\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qj}} \mathrm{u}_{\mathrm{ij}}\right) \mathrm{v}_{\mathrm{q}}^{\prime}=0, \tag{8}
\end{equation*}
$$

for an arbitrary vector $\mathrm{v}_{\mathrm{q}}{ }^{\prime}$. Therefore, we must have

$$
\begin{equation*}
\mathrm{u}_{\mathrm{pq}}^{\prime}=\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qj}} \mathrm{u}_{\mathrm{ij}} . \tag{9}
\end{equation*}
$$

This rule shows that components $\mathrm{u}_{\mathrm{ij}}$ obey the tensor law of transformation of a second order.

Hence, $\mathrm{u}_{\mathrm{ij}}$ is a tensor of order two.
Question : Show that $\delta_{\mathrm{ij}}$ and $\mathrm{a}_{\mathrm{ij}}$ are tensors, each of order two.
Solution : Let $u_{i}$ be any tensor of order one.
(a) By the substitution property of the Kronecker delta tensor $\delta_{\mathrm{ij}}$, we have

$$
\begin{equation*}
\mathrm{u}_{\mathrm{i}}=\delta_{\mathrm{ij}} \mathrm{u}_{\mathrm{j}} . \tag{1}
\end{equation*}
$$

Now $u_{i}$ and $v_{j}$ are, each of tensor order 1. Therefore, by quotient law, we conclude that $\delta_{\mathrm{ij}}$ is a tensor of rank two.
(b) The transformation law for the first order tensor $u_{i}$ is

$$
\begin{equation*}
\mathrm{u}_{\mathrm{p}}{ }^{\prime}=\mathrm{a}_{\mathrm{pi}} \mathrm{u}_{\mathrm{i}}, \tag{2}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathrm{a}_{\mathrm{pi}}=\cos \left(\mathrm{x}_{\mathrm{p}}^{\prime}, \mathrm{x}_{\mathrm{i}}\right) \tag{3}
\end{equation*}
$$

Now $u_{i}$ is a vector and $a_{p i} u_{i}$ is a vector by contraction property. Therefore, by quotient law, the quantities $\mathrm{a}_{\mathrm{pi}}$ are components of a second order tensor.

Hence the result.
Note (1) The tensor $\delta_{\mathrm{ij}}$ is called a unit tensor or an identity tensor of order two.

Note (2) We may call the tensor $a_{i j}$ as the transformation tensor of rank two.

Exercise 1 : Let $a_{i}$ be an ordered triplet and $b_{i}$ be a vector, referred to the $x_{i}-$ axis. If $a_{i} b_{i}$ is a scalar, show that $a_{i}$ are components of a vector.

Exercise 2 : If there is an entity representable by a set of 27 quantities $\mathrm{u}_{\mathrm{ijk}}$ relative to $o x_{1} x_{2} x_{3}$ system and if $u_{i j k} v_{j k}$ is a tensor of order one for an arbitrary tensor $\mathrm{v}_{\mathrm{jk}}$ of order 2 , show that $\mathrm{u}_{\mathrm{ijk}}$ is a tensor of order 3 .

Exercise 3 : If $u_{i j k} v_{k}$ is a tensor of order 2 for an arbitrary tensor $v_{k}$ of order one, show that $\mathrm{u}_{\mathrm{ijk}}$ is tensor of order 3 .

### 1.6 THE SYMBOL $\epsilon_{i j k}$

The symbol $\epsilon_{\mathrm{ijk}}$ is known as the Levi - civita $\in$ - symbol, named after the Italian mathematician Tullio Levi - civita (1873-1941).

The $\in$ - symbol is also referred to as the permutation symbol / alternating symbol or alternator.

In terms of mutually orthogonal unit vectors $\overline{\mathrm{e}}_{1}, \overline{\mathrm{e}}_{2}, \overline{\mathrm{e}}_{3}$ along the cartesian axes , it is defined as

$$
\overline{\mathrm{e}}_{\mathrm{i}} \cdot\left(\overline{\mathrm{e}}_{\mathrm{j}} \times \overline{\mathrm{e}}_{\mathrm{k}}\right)=\epsilon_{\mathrm{ijk}},
$$

for $\mathrm{i}, \mathrm{j}, \mathrm{k}=1,2,3$. Thus, the symbol $\epsilon_{\mathrm{ijk}}$ gives

$$
\epsilon_{\mathrm{ijk}}=\left\{\begin{aligned}
1 & \text { if } \mathrm{i}, \mathrm{j}, \mathrm{k} \text { takevalues in the cyclic order } \\
-1 & \text { if } \mathrm{i}, \mathrm{j}, \mathrm{k} \text { takesvalue in the acyclic order } \\
0 & \text { if twoorallof } \mathrm{i}, \mathrm{j}, \mathrm{k} \text { take the same value }
\end{aligned}\right.
$$

These relations are 27 in number.
The $\epsilon-$ symbol is useful in expressing the vector product of two vectors and scalar triple product.
(i) We have

$$
\overline{\mathrm{e}}_{\mathrm{i}} \times \overline{\mathrm{e}}_{\mathrm{j}}=\epsilon_{\mathrm{ijk}} \overline{\mathrm{e}}_{\mathrm{k}} .
$$

(ii) For two vectors $a_{i}$ and $b_{i}$, we write

$$
\bar{a} \times \bar{b}=\left(a_{i} e_{i}\right) \times\left(b_{j} e_{k}\right)=a_{i} b_{j}\left(e_{i} \times e_{j}\right)=\epsilon_{i j k} a_{i} b_{j} e_{k} .
$$

(iii) For vectors

$$
a=a_{i} e_{i}, b=b_{j} e_{j}, c=c_{k} e_{k},
$$

we have

$$
\begin{aligned}
{[\overline{\mathrm{a}}, \overline{\mathrm{~b}}, \overline{\mathrm{c}}] } & =(\overline{\mathrm{a}} \times \overline{\mathrm{b}}) \cdot \overline{\mathrm{c}} \\
& =\left(\epsilon_{\mathrm{ijk}} \mathrm{a}_{\mathrm{i}} \mathrm{~b}_{\mathrm{j}} \mathrm{e}_{\mathrm{k}}\right) \cdot\left(\mathrm{c}_{\mathrm{k}} \mathrm{e}_{\mathrm{k}}\right) \\
& =\in_{\mathrm{ijk}} \mathrm{a}_{\mathrm{i}} \mathrm{~b}_{\mathrm{j}} \mathrm{c}_{\mathrm{k}} \\
& =\left|\begin{array}{ccc}
a_{1} & a_{2} & a_{3} \\
b_{1} & b_{2} & b_{3} \\
c_{1} & c_{2} & c_{3}
\end{array}\right| .
\end{aligned}
$$

Question: Show that $\epsilon_{\mathrm{ijk}}$ is a tensor of order 3 .
Solution: Let $\bar{a}=a_{i}$ and $\bar{b}=b_{j}$ be any two vectors. Let

$$
\overline{\mathrm{c}}=\mathrm{c}_{\mathrm{i}}=\overline{\mathrm{a}} \times \overline{\mathrm{b}} .
$$

Then

$$
\begin{equation*}
c_{i}=\epsilon_{i j k} a_{j} b_{k} . \tag{i}
\end{equation*}
$$

Now $a_{j} b_{k}$ is a tensor of order 2 and $\epsilon_{i j k} a_{j} b_{k}$ is a tensor of order one. Therefore , by quotient law, $\epsilon_{\mathrm{ijk}}$ is a tensor of order 3 .

Note (1) Due to tensorial character of the $\in$ - symbol, it is called an alternating tensor or permutation tensor.

Note (2) The symbols $\delta_{\mathrm{ij}}$ and $\epsilon_{\mathrm{ijk}}$ were introduced earlier to simplifying the writing of some equations.

## Vector of a Second Order Tensor

Let $\mathrm{u}_{\mathrm{ij}}$ be a second order tensor. The vector

$$
\in_{\mathrm{ijk}} \mathrm{u}_{\mathrm{jk}}
$$

is called the vector of the tensor $\mathrm{u}_{\mathrm{jk}}$.
Example 1: Show that $\mathrm{w}_{\mathrm{ij}}=\epsilon_{\mathrm{ijk}} u_{\mathrm{k}}$ is a skew - symmetric tensor, where $u_{\mathrm{k}}$ is a vector and $\epsilon_{\mathrm{ijk}}$ is an alternating tensor.

Solution: Since $\epsilon_{\mathrm{ijk}}$ is a tensor of order 3 and $\mathrm{u}_{\mathrm{k}}$ is a tensor of order one, so by contraction, the product $\in_{\mathrm{ijk}} \mathrm{u}_{\mathrm{k}}$ is a tensor of order 2 . Further

$$
\begin{aligned}
\mathrm{w}_{\mathrm{ji}} & =\epsilon_{\mathrm{jik}} \mathrm{u}_{\mathrm{k}} \\
& =-\epsilon_{\mathrm{ijk}} \mathbf{u}_{\mathrm{k}} \\
& =-\mathrm{w}_{\mathrm{ji}} .
\end{aligned}
$$

This shows that $\mathrm{w}_{\mathrm{ij}}$ is a tensor which is skew - symmetric.
Example 2: Show that $\mathrm{u}_{\mathrm{ij}}$ is symmetric iff $\epsilon_{\mathrm{ijk}} \mathrm{u}_{\mathrm{ij}}=0$.
Solution: We find

$$
\begin{aligned}
& \epsilon_{\mathrm{ij} 1} \mathrm{u}_{\mathrm{ij}}=\epsilon_{231} \mathrm{u}_{23}+\epsilon_{321} \mathrm{u}_{32}=\mathrm{u}_{23}-\mathrm{u}_{32} \\
& \epsilon_{\mathrm{ij} 2} \mathrm{u}_{\mathrm{ij}}=\mathrm{u}_{31}-\mathrm{u}_{13}, \epsilon_{\mathrm{ij} 3} \mathrm{u}_{\mathrm{ij}}=\mathrm{u}_{12}-\mathrm{u}_{21}
\end{aligned}
$$

Thus, $\mathrm{u}_{\mathrm{ij}}$ is symmetric iff

$$
\mathrm{u}_{\mathrm{ij}}=\mathrm{u}_{\mathrm{ji}}
$$

or

$$
u_{23}=u_{32}, u_{13}=u_{31}, u_{12}=u_{21} .
$$

### 1.7 ISOTROPIC TENSORS

Definition: A tensor is said to be an isotropic tensor if its components remain unchanged / invariant however the axes are rotated.

Note (1) An isotropic tensor possesses no directional properties. Therefore a non - zero vector (or a non - zero tensor of rank 1) can never be an isotropic tensor.

Tensors of higher orders, other than one, can be isotropic tensors.
Note (2) Zero tensors of all orders are isotropic tensors.
Note (3) By definition, a scalar (or a tensor of rank zero) is an isotropic tensor.
Note (4) A scalar multiple of an isotropic tensor is an isotropic tensor.
Note (5) The sum and the differences of two isotropic tensors is an isotropic tensor.

Theorem: Prove that substitution tensor $\delta_{\mathrm{ij}}$ and alternating tensor $\epsilon_{\mathrm{ijk}}$ are isotropic tensors.

Proof: Let the components $\delta_{\mathrm{ij}}$ relative to $\mathrm{x}_{\mathrm{i}}$ system are transformed to quantities $\delta^{\prime}{ }_{p q}$ relative to $\mathrm{x}_{\mathrm{i}}{ }^{\prime}$ - system. Then, the tensorial transformation rule is

$$
\begin{equation*}
\delta_{\mathrm{pq}}^{\prime}=\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qj}} \delta_{\mathrm{ij}} \tag{1}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathrm{a}_{\mathrm{pi}}=\cos \left(\mathrm{x}_{\mathrm{p}}^{\prime}, \mathrm{x}_{\mathrm{i}}\right) . \tag{2}
\end{equation*}
$$

Now

$$
\begin{align*}
\text { RHS of }(1) & =\mathrm{a}_{\mathrm{pi}}\left[\mathrm{a}_{\mathrm{qj}} \delta_{\mathrm{ij}}\right] \\
& =\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qi}} \\
& =\delta_{\mathrm{pq}} \\
& = \begin{cases}0 & \text { if } p \neq q \\
1 & \text { if } p=q\end{cases} \tag{3}
\end{align*}
$$

Relations (1) and (3) show that the components $\delta_{\mathrm{ij}}$ are transformed into itself under all co-ordinate transformations. Hence, by definition, $\delta_{\mathrm{ij}}$ is an isotropic tensor.

We know that $\epsilon_{\mathrm{ijk}}$ is a system of 27 numbers. Let

$$
\begin{equation*}
\epsilon_{\mathrm{ijk}}=\left[\overline{\mathrm{e}}_{\mathrm{i}}, \overline{\mathrm{e}}_{\mathrm{j}}, \overline{\mathrm{e}}_{\mathrm{k}}\right]=\overline{\mathrm{e}}_{\mathrm{i}} .\left(\overline{\mathrm{e}}_{\mathrm{j}} \times \overline{\mathrm{e}}_{\mathrm{k}}\right), \tag{4}
\end{equation*}
$$

be related to the $x_{i}$ axes. Suppose that these components are transformed to $\epsilon_{\text {pqr }}^{\prime}$ relative to $\mathrm{x}_{\mathrm{i}}{ }^{\prime}$ - axis. Then, the third order tensorial law of transformation gives

$$
\begin{equation*}
\epsilon_{\mathrm{pqr}}^{\prime}=\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qj}} \mathrm{a}_{\mathrm{rk}} \in_{\mathrm{ijk}} \tag{5}
\end{equation*}
$$

where $l_{\mathrm{pi}}$ is defined in (2)
we have already checked that (exercise)

$$
\in_{\mathrm{ijk}} \mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qj}} \mathrm{a}_{\mathrm{rk}}=\left|\begin{array}{ccc}
\mathrm{a}_{\mathrm{p} 1} & \mathrm{a}_{\mathrm{p} 2} & \mathrm{a}_{\mathrm{p} 3}  \tag{6}\\
\mathrm{a}_{\mathrm{q} 1} & \mathrm{a}_{\mathrm{q} 2} & \mathrm{a}_{\mathrm{q} 3} \\
\mathrm{a}_{\mathrm{r} 1} & \mathrm{a}_{\mathrm{r} 2} & \mathrm{a}_{\mathrm{r} 3}
\end{array}\right|
$$

and

$$
\left[\overline{\mathrm{e}}_{\mathrm{p}}^{\prime}, \overline{\mathrm{e}}_{\mathrm{q}}^{\prime}, \overline{\mathrm{e}}_{\mathrm{r}}^{\prime}\right]=\left|\begin{array}{ccc}
\mathrm{a}_{\mathrm{p} 1} & \mathrm{a}_{\mathrm{p} 2} & \mathrm{a}_{\mathrm{p} 3}  \tag{7}\\
\mathrm{a}_{\mathrm{q} 1} & \mathrm{a}_{\mathrm{q} 2} & a_{\mathrm{q} 3} \\
\mathrm{a}_{\mathrm{r} 1} & \mathrm{a}_{\mathrm{r} 2} & a_{\mathrm{r} 3}
\end{array}\right|
$$

From (5) - (7), we get

$$
\begin{aligned}
\in_{\mathrm{pqr}}^{\prime} & =\left[\overline{\mathrm{e}}_{\mathrm{p}}^{\prime}, \overline{\mathrm{e}}_{\mathrm{q}}^{\prime}, \overline{\mathrm{e}}_{\mathrm{r}}^{\prime}\right] \\
& =\overline{\mathrm{e}}_{\mathrm{p}}^{\prime} \cdot\left(\overline{\mathrm{e}}_{\mathrm{q}}^{\prime} \times \overline{\mathrm{e}}_{\mathrm{r}}^{\prime}\right)
\end{aligned}
$$

$$
=\left\{\begin{align*}
1 & \text { if } p, q, \text { rareincyclic order }  \tag{8}\\
-1 & \text { if } p, q, \text { rareincyclicorder } \\
0 & \text { if anytwoorall sufficesare equal }
\end{align*}\right.
$$

This shows that components $\epsilon_{\mathrm{ijk}}$ are transformed into itself under all coordinate transformations. Thus, the third order tensor $\epsilon_{\mathrm{ijk}}$ is an isotropic tensor.

Theorem: If $\mathrm{u}_{\mathrm{ij}}$ is an isotropic tensor of second order, then show that

$$
\mathrm{u}_{\mathrm{ij}}=\alpha \delta_{\mathrm{ij}}
$$

for some scalar $\alpha$.
Proof: As the given tensor is isotropic, we have

$$
\begin{equation*}
\mathrm{u}_{\mathrm{ij}}{ }^{\prime}=\mathrm{u}_{\mathrm{ij}}, \tag{1}
\end{equation*}
$$

for all choices of the $\mathrm{x}_{\mathrm{i}}{ }^{\prime}$ - system. In particular, we choose

$$
\begin{equation*}
\mathrm{x}_{1}{ }^{\prime}=\mathrm{x}_{2}, \mathrm{x}_{2}{ }^{\prime}=\mathrm{x}_{3}, \mathrm{x}_{3}{ }^{\prime}=\mathrm{x}_{1} \tag{2}
\end{equation*}
$$

Then

$$
\mathrm{a}_{\mathrm{ij}}=\left|\begin{array}{lll}
0 & 1 & 0  \tag{3}\\
0 & 0 & 1 \\
1 & 0 & 0
\end{array}\right|
$$

and

$$
\begin{equation*}
\mathrm{u}_{\mathrm{pq}}^{\prime}=\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qj}} \mathrm{u}_{\mathrm{ij}} . \tag{4}
\end{equation*}
$$

Now

$$
\begin{aligned}
u_{11}=u_{11^{\prime}} & =a_{1 \mathrm{i}} a_{1 \mathrm{j}} u_{\mathrm{ij}} \\
& =\mathrm{a}_{12} \mathrm{a}_{12} u_{23}=u_{22}, \\
\mathrm{u}_{22}=\mathrm{u}_{22^{\prime}} & =\mathrm{a}_{2 \mathrm{i}} \mathrm{a}_{2 \mathrm{j}} \mathrm{u}_{\mathrm{ij}} \\
& =\mathrm{a}_{23} \mathrm{a}_{23} u_{33}=\mathrm{u}_{33}, \\
\mathrm{u}_{12}=\mathrm{u}_{12^{\prime}} & =\mathrm{a}_{1 \mathrm{i}} \mathrm{a}_{2 \mathrm{j}} \mathrm{u}_{\mathrm{ij}} \\
& =\mathrm{a}_{12} \mathrm{a}_{23} \mathrm{u}_{23}=\mathrm{u}_{23}, \\
\mathrm{u}_{23}=\mathrm{u}_{23^{\prime}} & =\mathrm{a}_{2 \mathrm{i}} \mathrm{a}_{3 \mathrm{j}} \mathrm{u}_{\mathrm{ij}} \\
& =\mathrm{a}_{23} \mathrm{a}_{31} \mathrm{u}_{31}=\mathrm{u}_{31}, \\
\mathrm{u}_{13}=\mathrm{u}_{13^{\prime}} & =\mathrm{a}_{1 \mathrm{i}} \mathrm{a}_{3 \mathrm{j}} \mathrm{u}_{\mathrm{ij}}
\end{aligned}
$$

$$
\begin{gathered}
=\mathrm{a}_{12} \mathrm{a}_{31} \mathrm{u}_{21}=\mathrm{u}_{21} \\
\mathrm{u}_{21}=\mathrm{u}_{21}^{\prime}=\mathrm{a}_{2 \mathrm{i}} \mathrm{a}_{1 \mathrm{j}} \mathrm{u}_{\mathrm{ij}} \\
=\mathrm{a}_{23} \mathrm{a}_{12} \mathrm{u}_{32}=\mathrm{u}_{32}
\end{gathered}
$$

Thus

$$
\begin{align*}
& \mathrm{u}_{11}=\mathrm{u}_{22}=\mathrm{u}_{33} \\
& \mathrm{u}_{12}=\mathrm{u}_{23}=\mathrm{u}_{31} \\
& \mathrm{u}_{21}=\mathrm{u}_{32}=\mathrm{u}_{13} \tag{5}
\end{align*}
$$

Now, we consider the transformation

$$
\begin{equation*}
\mathrm{x}_{1}^{\prime}=\mathrm{x}_{2}, \mathrm{x}_{2}^{\prime}=-\mathrm{x}_{1}, \mathrm{x}_{3}^{\prime}=\mathrm{x}_{3} . \tag{6}
\end{equation*}
$$

Then

$$
\begin{align*}
& \left(\mathrm{a}_{\mathrm{ij}}\right)=\left(\begin{array}{rrr}
0 & 1 & 0 \\
-1 & 0 & 0 \\
0 & 0 & 1
\end{array}\right),  \tag{7}\\
& \mathrm{u}_{\mathrm{pq}}^{\prime}=\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qj}} \mathrm{u}_{\mathrm{ij}} \tag{8}
\end{align*}
$$

This gives

$$
\begin{aligned}
\mathrm{u}_{13}=\mathrm{u}^{\prime}{ }_{13} & =\mathrm{a}_{1 \mathrm{i}} \mathrm{a}_{3 \mathrm{j}} \mathrm{u}_{\mathrm{ij}} \\
& =\mathrm{a}_{12} \mathrm{a}_{33} \mathrm{u}_{23}=\mathrm{u}_{23}, \\
\mathrm{u}_{23}=\mathrm{u}^{\prime}{ }_{23} & =\mathrm{a}_{2 \mathrm{i}} \mathrm{a}_{3 \mathrm{j}} \mathrm{u}_{\mathrm{ij}} \\
& =\mathrm{a}_{21} \mathrm{a}_{33} \mathrm{u}_{13}=-\mathrm{u}_{13} .
\end{aligned}
$$

Thus

$$
\begin{equation*}
\mathrm{u}_{13}=\mathrm{u}_{23}=0 \tag{9}
\end{equation*}
$$

From (5) and (9), we obtain

$$
\begin{equation*}
\mathrm{a}_{\mathrm{ij}}=\alpha \delta_{\mathrm{ij}} \tag{10}
\end{equation*}
$$

where

$$
\begin{equation*}
\alpha=\mathrm{a}_{11}=\mathrm{a}_{22}=\mathrm{a}_{33} \tag{11}
\end{equation*}
$$

Note 1: If $\mathrm{a}_{\mathrm{ijk}}$ are components of an isotropic tensor of third order, then

$$
\mathrm{a}_{\mathrm{ijk}}=\alpha \cdot \epsilon_{\mathrm{ijk}}
$$

for some scalar $\alpha$.
Note 2: If $\mathrm{a}_{\mathrm{ijkm}}$ are components of a fourth - order isotropic tensor , then

$$
\mathrm{a}_{\mathrm{ijkm}}=\alpha \delta_{\mathrm{ij}} \delta_{\mathrm{km}}+\beta \delta_{\mathrm{ik}} \delta_{\mathrm{jm}}+\gamma \delta_{\mathrm{im}} \delta_{\mathrm{jk}}
$$

for some scalars $\alpha, \beta, \gamma$.

## Definition: (Gradient)

If $\mathrm{u}_{\text {pqr... }}\left(\mathrm{x}_{1}, \mathrm{x}_{2}, \mathrm{x}_{3}\right)$ is a tensor of order n , then

$$
\begin{aligned}
\mathrm{v}_{\mathrm{spqr} \ldots \ldots .} & =\frac{\partial}{\partial x_{s}} u_{p q r \ldots \ldots} \\
& =\mathrm{u}_{\mathrm{pqr} \ldots \ldots, \mathrm{~s}}
\end{aligned}
$$

is defined as the gradient of the tensor field $\mathrm{u}_{\mathrm{pqr} . \ldots .}\left(\mathrm{x}_{1}, \mathrm{x}_{2}, \mathrm{x}_{3}\right)$.
Theorem: Show that the gradient of a scalar point function is a tensor of order 1.

Proof: Suppose that $\mathrm{U}=\mathrm{U}\left(\mathrm{x}_{1}, \mathrm{x}_{2}, \mathrm{x}_{3}\right)$ be a scalar point function and

$$
\begin{equation*}
\mathrm{v}_{\mathrm{i}}=\frac{\partial U}{\partial x_{i}}=\mathrm{U}_{\mathrm{i}}=\text { gradient of } \mathrm{U} . \tag{1}
\end{equation*}
$$

Let the components of the gradient of $U$ in the dashed system $\mathrm{ox}_{1}{ }^{\prime} \mathrm{x}_{2}{ }^{\prime} \mathrm{x}_{3}{ }^{\prime}$ be $\mathrm{v}_{\mathrm{p}}$, so that,

$$
\begin{equation*}
\mathrm{v}_{\mathrm{p}}{ }^{\prime}=\frac{\partial U}{\partial x_{p^{\prime}}} \tag{2}
\end{equation*}
$$

where the transformation rule of coordinates is

$$
\begin{align*}
& \mathrm{x}_{\mathrm{p}}^{\prime}=\mathrm{a}_{\mathrm{pi}} \mathrm{x}_{\mathrm{i}}  \tag{3}\\
& \mathrm{x}_{\mathrm{i}}=\mathrm{a}_{\mathrm{pi}} \mathrm{x}_{\mathrm{p}}^{\prime}  \tag{4}\\
& \mathrm{a}_{\mathrm{pi}}=\cos \left(\mathrm{x}_{\mathrm{p}}^{\prime}, \mathrm{x}_{\mathrm{i}}\right) \tag{5}
\end{align*}
$$

By chain rule

$$
\mathrm{v}_{\mathrm{p}}{ }^{\prime}=\frac{\partial U}{\partial x_{p^{\prime}}}
$$

$$
\begin{aligned}
& =\frac{\partial U}{\partial x_{i}} \frac{\partial x_{i}}{\partial x_{p^{\prime}}} \\
& =\mathrm{a}_{\mathrm{pi}} \frac{\partial U}{\partial x_{i}} \\
& =\mathrm{a}_{\mathrm{pi}} \mathrm{v}_{\mathrm{i}}
\end{aligned}
$$

which is a transformation rule for tensors of order 1.
Hence gradient of the scalar point function U is a tensor of order one.
Theorem: Show that the gradient of a vector $u_{i}$ is a tensor of order 2 . Deduce that $\delta_{\mathrm{ij}}$ is a tensor of order 2.

Proof: The gradient of the tensor $u_{i}$ is defined as

$$
\begin{equation*}
\mathrm{w}_{\mathrm{ij}}=\frac{\partial \mathrm{u}_{\mathrm{i}}}{\partial \mathrm{x}_{\mathrm{j}}}=\mathrm{u}_{\mathrm{i}, \mathrm{j}} . \tag{1}
\end{equation*}
$$

Let the vector $u_{i}$ be transformed to the vector $u_{p}{ }^{\prime}$ relative to the new system o $\mathrm{x}_{1}{ }^{\prime} \mathrm{x}_{2}{ }^{\prime} \mathrm{X}_{3}{ }^{\prime}$. Then the transformation law for tensors of order 1 yields

$$
\begin{equation*}
\mathrm{u}_{\mathrm{p}}^{\prime}=\mathrm{a}_{\mathrm{pi}} \mathrm{u}_{\mathrm{i}}, \tag{2}
\end{equation*}
$$

where the law of transformation of coordinates is

$$
\begin{align*}
& \mathrm{x}_{\mathrm{q}}^{\prime}=\mathrm{a}_{\mathrm{qj}} \mathrm{x}_{\mathrm{j}},  \tag{3}\\
& \mathrm{x}_{\mathrm{j}}=\mathrm{a}_{\mathrm{qj}} \mathrm{x}_{\mathrm{q}}^{\prime},  \tag{4}\\
& \mathrm{a}_{\mathrm{pi}}=\cos \left(\mathrm{x}_{\mathrm{p}}^{\prime}, \mathrm{x}_{\mathrm{i}}\right) . \tag{5}
\end{align*}
$$

Suppose that the 9 quantities $\mathrm{w}_{\mathrm{ij}}$ relative to new system are transformed to $\mathrm{w}_{\mathrm{pq}}^{\prime}$. Then

$$
\begin{aligned}
\mathrm{w}_{\mathrm{pq}}^{\prime} & =\frac{\partial \mathrm{u}_{\mathrm{p}}}{\partial \mathrm{x}_{\mathrm{q}}^{\prime}} \\
& =\frac{\partial}{\partial \mathrm{x}_{\mathrm{q}}{ }^{\prime}}\left(\mathrm{a}_{\mathrm{pi}} \mathrm{u}_{\mathrm{i}}\right) \\
& =\mathrm{a}_{\mathrm{pi}} \frac{\partial u_{i}}{\partial x_{q}{ }^{\prime}}
\end{aligned}
$$

$$
\begin{align*}
& =\mathrm{a}_{\mathrm{pi}} \frac{\partial u_{i}}{\partial x_{j}} \frac{\partial x_{j}}{\partial x_{q}{ }^{\prime}} \\
& =\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qj}} \frac{\partial u_{i}}{\partial x_{j}} \\
& =\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qj}} \mathrm{w}_{\mathrm{ij}}, \tag{6}
\end{align*}
$$

which is the transformation rule for tensors of order 2.
Hence, $w_{i j}$ is a tensor of order 2 . Consequently, the gradient of the vector $u_{i}$ is a tensor of order 2.

Deduction : We know that

$$
\delta_{\mathrm{ij}}=\frac{\partial x_{i}}{\partial x_{j}},
$$

and that $\mathrm{x}_{\mathrm{i}}$ is a vector. So, $\delta_{\mathrm{ij}}$ is a gradient of the vector $\mathrm{x}_{\mathrm{i}}$. It follows that 9 quantities $\delta_{\mathrm{ij}}$ are components of a tensor of order 2.

### 1.8 EIGENVALUES AND EIGEN VECTORS OF A SECOND ORDER SYMMETRIC TENSOR.

Definition: Let $\mathrm{u}_{\mathrm{ij}}$ be a second order symmetric tensor. A scalar $\lambda$ is called an eigenvalue of the tensor $u_{i j}$ if there exists a non - zero vector $v_{i}$ such that

$$
\mathrm{u}_{\mathrm{ij}} \mathrm{v}_{\mathrm{j}}=\lambda \mathrm{v}_{\mathrm{i}}, \quad \text { for } \mathrm{i}=1,2,3 .
$$

The non - zero vector $\mathbf{v}_{\mathbf{i}}$ is then called an eigenvector of tensor $u_{i j}$ corresponding to the eigen vector $\lambda$.

We observe that every (non - zero) scalar multiple of an eigenvector is also an eigen vector.

Article: Show that it is always possible to find three mutually orthogonal eigenvectors of a second order symmetric tensor.

Proof: Let $u_{i j}$ be a second order symmetric tensor and $\lambda$ be an eigen value of $u_{i j}$ . Let $\mathrm{v}_{\mathrm{i}}$ be an eigenvector corresponding to $\lambda$.

Then

$$
\mathrm{u}_{\mathrm{ij}} \mathrm{v}_{\mathrm{j}}=\lambda \mathrm{v}_{\mathrm{i}}
$$

or

$$
\begin{equation*}
\left(\mathrm{u}_{\mathrm{ij}}-\lambda \delta_{\mathrm{ij}}\right) \mathrm{v}_{\mathrm{j}}=0 \tag{1}
\end{equation*}
$$

This is a set of three homogeneous simultaneous linear equations in three unknown $\mathrm{v}_{1}, \mathrm{v}_{2}, \mathrm{v}_{3}$. These three equations are

$$
\left.\begin{array}{l}
\left(u_{11}-\lambda\right) v_{1}+u_{12} v_{2}+u_{13} v_{3}=0  \tag{2}\\
u_{21} v_{1}+\left(u_{22}-\lambda\right) v_{2}+u_{23} v_{3}=0 \\
u_{31} v_{1}+u_{32} v_{2}+\left(u_{33}-\lambda\right) v_{3}=0
\end{array}\right\}
$$

This set of equations possesses a non - zero solution when

$$
\left|\begin{array}{ccc}
u_{11}-\lambda & u_{12} & u_{13} \\
u_{21} & u_{22}-\lambda & u_{23} \\
u_{31} & u_{32} & u_{33}-\lambda
\end{array}\right|=0
$$

or

$$
\begin{equation*}
\left|\mathrm{u}_{\mathrm{ij}}-\lambda \delta_{\mathrm{ij}}\right|=0 \tag{3}
\end{equation*}
$$

Expanding the determinant in (3), we find

$$
\begin{aligned}
& \left(u_{11}-\lambda\right)\left[\left(u_{22}-\lambda\right)\left(u_{33}-\lambda\right)-u_{32} u_{23}\right] \\
& -u_{12}\left[u_{21}\left(u_{33}-\lambda\right)-u_{31} u_{23}\right]+u_{13}\left[u_{21} u_{32}-u_{31}\left(u_{22}-\right.\right. \\
& \quad \lambda)]=0
\end{aligned}
$$

or

$$
\begin{align*}
& -\lambda^{3}+\left(u_{11}+u_{22}+u_{33}\right) \lambda^{2}-\left(u_{11} u_{22}+u_{22} u_{33}+u_{33} u_{11}-\right. \\
& \left.u_{23} u_{32}-u_{31} u_{13}-u_{12} u_{21}\right) \lambda+\left[u _ { 1 1 } \left(u_{22} u_{33}-u_{23}\right.\right. \\
& \left.u_{32}\right)-u_{12}\left(u_{21} u_{33}-u_{31} u_{23}\right) \\
& \left.+u_{13}\left(u_{21} u_{32}-u_{31} u_{22}\right)\right]=0 \tag{4}
\end{align*}
$$

We write (4) as

$$
\begin{equation*}
-\lambda^{3}+\lambda^{2} \mathrm{I}_{1}-\lambda \mathrm{I}_{2}+\mathrm{I}_{3}=0 \tag{5}
\end{equation*}
$$

where

$$
\begin{aligned}
\mathrm{I}_{1} & =\mathrm{u}_{11}+\mathrm{u}_{22}+\mathrm{u}_{33} \\
& =\mathrm{u}_{\mathrm{ii}} \\
\mathrm{I}_{2} & =\mathrm{u}_{11} \mathrm{u}_{22}+\mathrm{u}_{22} \mathrm{u}_{33}+\mathrm{u}_{33} \mathrm{u}_{11}-\mathrm{u}_{12} \mathrm{u}_{21}-\mathrm{u}_{23} \mathrm{u}_{32}-\mathrm{u}_{31} \mathrm{u}_{13}
\end{aligned}
$$

$$
\begin{align*}
& =\frac{1}{2}\left[u_{i \mathrm{i}} \mathrm{u}_{\mathrm{jj}}-\mathrm{u}_{\mathrm{ij}} \mathrm{u}_{\mathrm{j} 1}\right],  \tag{7}\\
\mathrm{I}_{3} & =\left|\mathrm{u}_{\mathrm{ij}}\right| \\
& =\epsilon_{\mathrm{ijk}} \mathrm{u}_{\mathrm{i} 1} \mathrm{u}_{\mathrm{j} 2} \mathrm{u}_{\mathrm{k} 3} . \tag{8}
\end{align*}
$$

Equation (5) is a cubic equation in $\lambda$. Therefore, it has three roots, say, $\lambda_{1}, \lambda_{2}$ , $\lambda_{3}$ which may not be distinct (real or imaginary). These roots (which are scalar) are the three eigenvalues of the symmetric tensor $\mathrm{u}_{\mathrm{ij}}$.

Further

$$
\begin{align*}
& \lambda_{1}+\lambda_{2}+\lambda_{3}=\mathrm{I}_{1}  \tag{9}\\
& \lambda_{1} \lambda_{2}+\lambda_{2} \lambda_{3}+\lambda_{3} \lambda_{1}=\mathrm{I}_{2}  \tag{10}\\
& \lambda_{1} \lambda_{2} \lambda_{3}=\mathrm{I}_{3} \tag{11}
\end{align*}
$$

Each root $\lambda_{i}$, when substituted in equation (2), gives a set of three linear equations (homogeneous) which are not all independent. By discarding one of equations and using the condition

$$
v_{1}^{2}+v_{2}^{2}+v_{3}^{2}=1
$$

for unit vectors, the eigenvector $\overline{\mathrm{v}}_{\mathrm{i}}$ is determined.
Before proceeding further, we state and prove two important lemmas.

## Lemma 1: Eigenvalues of a real symmetric tensor $\mathbf{u}_{\mathrm{ij}}$ are real.

Proof: Let $\lambda$ be an eigenvalue with corresponding eigenvector $u_{i}$.
Then

$$
\begin{equation*}
\mathrm{u}_{\mathrm{ij}} \mathrm{v}_{\mathrm{j}}=\lambda \mathrm{v}_{\mathrm{i}} . \tag{I}
\end{equation*}
$$

Taking the complex conjugate on both sides of (I), we find

$$
\begin{array}{ll}
\overline{\mathrm{u}_{\mathrm{ij}} \mathrm{v}_{\mathrm{j}}}=\overline{\lambda \mathrm{v}_{\mathrm{i}}} & \overline{\mathrm{u}_{\mathrm{ij}}} \overline{\mathrm{v}_{\mathrm{j}}}=\bar{\lambda} \overline{\mathrm{v}_{\mathrm{i}}} \\
\mathrm{u}_{\mathrm{ij}} \overline{\mathrm{v}}_{\mathrm{j}}=\bar{\lambda} \overline{\mathrm{v}}_{\mathrm{i}} & \tag{II}
\end{array}
$$

since $\mathrm{u}_{\mathrm{ij}}$ is a real tensor. Now

$$
\begin{align*}
\mathrm{u}_{\mathrm{ij}} \overline{\mathrm{v}}_{\mathrm{j}} \mathrm{v}_{\mathrm{i}} & =\left(\mathrm{u}_{\mathrm{ij}} \overline{\mathrm{v}}_{\mathrm{j}}\right) \mathrm{v}_{\mathrm{i}} \\
& =\left(\begin{array}{ll}
\bar{\lambda} & \left.\overline{\mathrm{v}}_{\mathrm{i}}\right) \mathrm{v}_{\mathrm{i}} \\
& =\bar{\lambda} \quad \overline{\mathrm{v}}_{\mathrm{i}} \mathrm{v}_{\mathrm{i}}
\end{array} .\right.
\end{align*}
$$

Also

$$
\begin{aligned}
\overline{u_{i j} \overline{v_{j}}} v_{i} & =\overline{u_{i j}} v_{j} \overline{v_{i}} \\
= & u_{i j} v_{j} \overline{v_{i}} \\
& =u_{j i} v_{i} \bar{v}_{j} \\
& =u_{i j} \bar{v}_{j} v_{i} .
\end{aligned}
$$

This shows that quantity $u_{i j} \bar{v}_{j} v_{i}$ is real. Hence $\bar{\lambda} \bar{v}_{i} v_{i}$ is real.
Since $\bar{v}_{i} \mathrm{v}_{\mathrm{i}}$ is always real, it follows that $\bar{\lambda}$ is real.
Therefore $\lambda$ is real.
Lemma 2: Eigen vector corresponding to two distinct eigen values of the symmetric tensor $\mathbf{u}_{\mathbf{i j}}$ are orthogonal.

Proof: Let $\lambda_{1} \neq \lambda_{2}$ be two distinct eigenvalues of $u_{i j}$. Let $A_{i}$ and $B_{i}$ be the corresponding non - zero eigenvectors. Then

$$
\begin{align*}
& \mathrm{u}_{\mathrm{ij}} \mathrm{~A}_{\mathrm{j}}=\lambda_{1} \mathrm{~A}_{\mathrm{i}}, \\
& \mathrm{u}_{\mathrm{ij}} \mathrm{~B}_{\mathrm{j}}=\lambda_{2} \mathrm{~B}_{\mathrm{i}} . \tag{I}
\end{align*}
$$

We obtain

$$
\begin{align*}
& \mathrm{u}_{\mathrm{ij}} \mathrm{~A}_{\mathrm{j}} \mathrm{~B}_{\mathrm{i}}=\lambda_{1} \mathrm{~A}_{\mathrm{i}} \mathrm{~B}_{\mathrm{i}}, \\
& \mathrm{u}_{\mathrm{ij}} \mathrm{~B}_{\mathrm{j}} \mathrm{~A}_{\mathrm{i}}=\lambda_{2} \mathrm{~A}_{\mathrm{i}} \mathrm{~B}_{\mathrm{i}} . \tag{II}
\end{align*}
$$

Now

$$
\begin{align*}
\mathrm{u}_{\mathrm{ij}} \mathrm{~A}_{\mathrm{j}} \mathrm{~B}_{\mathrm{i}} & =\mathrm{u}_{\mathrm{ji}} \mathrm{~A}_{\mathrm{i}} \mathrm{~B}_{\mathrm{j}} \\
& =\mathrm{u}_{\mathrm{ij}} \mathrm{~B}_{\mathrm{j}} \mathrm{~A}_{\mathrm{i}} . \tag{III}
\end{align*}
$$

From (II) \& (III), we get

$$
\begin{aligned}
& \lambda_{1} \mathrm{~A}_{\mathrm{i}} \mathrm{~B}_{\mathrm{i}}=\lambda_{2} \mathrm{~A}_{\mathrm{i}} \mathrm{~B}_{\mathrm{i}} \\
& \left(\lambda_{1}-\lambda_{2}\right) \mathrm{A}_{\mathrm{i}} \mathrm{~B}_{\mathrm{i}}=0
\end{aligned}
$$

$\mathrm{A}_{\mathrm{i}} \mathrm{B}_{\mathrm{i}}=0$.
$\left(\because \lambda_{1} \neq \lambda_{2}\right)$

Hence, eigenvectors $A_{i}$ and $B_{i}$ are mutually orthogonal.
This completes the proof of lemma 2.
Now we consider various possibilities about eigenvalues $\lambda_{1}, \lambda_{2}, \lambda_{3}$ of the main theorem.

Case 1: If $\lambda_{1} \neq \lambda_{2} \neq \lambda_{3}$, i. e., when all eigenvalues are different and real.
Then, by lemma 2 , three eigenvectors corresponding to $\lambda_{\mathrm{i}}$ are mutually orthogonal. Hence the result holds.

Case 2: If $\lambda_{1} \neq \lambda_{2}=\lambda_{3}$. Let $\stackrel{1}{v}_{i}$ be the eigenvector of the tensor $\mathrm{u}_{\mathrm{ij}}$ corresponding to the eigenvalue $\lambda_{1}$ and $v_{i}^{2}$ be the eigenvector corresponding to $\lambda_{2}$. Then

$$
\begin{array}{cc}
1 & 2 \\
v_{i} & v_{i}
\end{array}=0
$$



Let $\mathrm{p}_{\mathrm{i}}$ be a vector orthogonal to both $\stackrel{1}{v}_{i}$ and ${\stackrel{2}{v_{i}}}_{i}$. Then

$$
\begin{equation*}
\mathrm{p}_{\mathrm{i}} \stackrel{1}{v_{i}}=\mathrm{p}_{\mathrm{i}} \stackrel{2}{v_{i}}=0 \tag{12}
\end{equation*}
$$

and

$$
\begin{gather*}
\mathrm{u}_{\mathrm{ij}} \stackrel{1}{v_{j}}=\lambda_{1} \stackrel{1}{v_{i}}, \\
\mathrm{u}_{\mathrm{ij}} \stackrel{2}{v}_{j}=\lambda_{2} \stackrel{2}{v_{i}} . \tag{13}
\end{gather*}
$$

Let

$$
\begin{equation*}
\mathrm{u}_{\mathrm{ij}} \mathrm{p}_{\mathrm{j}}=\mathrm{q}_{\mathrm{i}}=\mathrm{a} \text { tensor of order } 1 \tag{14}
\end{equation*}
$$

we shall show that $q_{i}$ and $p_{i}$ are parallel.
Now

$$
\begin{align*}
& \mathrm{q}_{\mathrm{i}} \stackrel{1}{v_{i}}=\mathrm{u}_{\mathrm{ij}} \mathrm{p}_{\mathrm{j}} \stackrel{1}{v_{i}} \\
& =\mathrm{u}_{\mathrm{ji}} \stackrel{1}{v_{j}} \mathrm{p}_{\mathrm{i}} \\
& =\mathrm{u}_{\mathrm{ij}} \stackrel{1}{v_{j}} \mathrm{p}_{\mathrm{i}} \\
& =\lambda_{1} \stackrel{1}{v_{i}} \mathrm{p}_{\mathrm{i}} \\
& =0 \text {. } \tag{15}
\end{align*}
$$

## Similarly

$$
\begin{equation*}
\mathrm{q}_{\mathrm{i}} \stackrel{2}{2}_{i}=0 \tag{16}
\end{equation*}
$$

Thus $\mathrm{q}_{\mathrm{i}}$ is orthogonal to both orthogonal eigenvectors $\stackrel{1}{v}_{i}$ and $\stackrel{2}{v_{i}}$.
Thus $q_{i}$ must be parallel to $\mathrm{p}_{\mathrm{i}}$. So , we may write

$$
\begin{equation*}
\mathrm{u}_{\mathrm{ij}} \mathrm{p}_{\mathrm{j}}=\mathrm{q}_{\mathrm{i}}=\alpha \mathrm{p}_{\mathrm{i}} \tag{17}
\end{equation*}
$$

for some scalar $\alpha$.
Equation (10) shows tha $\alpha$ must be an eigenvalue and $p_{i}$ must be the corresponding eigenvector of $u_{i j}$.

Let

$$
\begin{equation*}
\stackrel{3}{v_{i}}=\frac{\mathrm{p}_{\mathrm{i}}}{\left|\mathrm{p}_{\mathrm{i}}\right|} \tag{18}
\end{equation*}
$$

Since $\mathrm{u}_{\mathrm{ij}}$ has only three eigenvalues $\lambda_{1}, \lambda_{2}=\lambda_{3}$, so $\alpha$ must be equal to $\lambda_{2}=\lambda_{3}$.
Thus $\stackrel{3}{3}_{i}$ is an eigenvector which is orthogonal to both $v_{i}^{1}$ and $v_{i}^{2}$ where $v_{i}^{1} \perp v_{i}^{2}$.
Thus, there exists three mutually orthogonal eigenvectors.
Further, let $\mathrm{w}_{\mathrm{i}}$ be any vector which lies in the plane containing the two eigenvectors $\stackrel{2}{v_{i}}$ and $\stackrel{3}{v}_{\mathrm{i}}$ corresponding to the repeated eigenvalues. Then

$$
\mathrm{w}_{\mathrm{i}}=\mathrm{k}_{1} \stackrel{2}{v_{i}}+\mathrm{k}_{2} \stackrel{3}{v_{i}}
$$

for some scalars $k_{1}$ and $k_{2}$ and

$$
\mathrm{w}_{\mathrm{i}} \stackrel{1}{v_{i}}=0,
$$

and

$$
\begin{aligned}
\mathrm{u}_{\mathrm{ij}} \mathrm{w}_{\mathrm{j}} & =\mathrm{u}_{\mathrm{ij}}\left(\mathrm{k}_{1} v_{j}^{2}+\mathrm{k}_{2} v_{j}^{3}\right) \\
& =\mathrm{k}_{1} \mathrm{u}_{\mathrm{ij}} v_{j}^{2}+\mathrm{k}_{2} \mathrm{u}_{\mathrm{ij}} v_{j}^{3} \\
& =\mathrm{k}_{1} \lambda_{2} v_{i}^{2}+\mathrm{k}_{2} \lambda_{3} v_{i}^{3}
\end{aligned}
$$

$$
\begin{gather*}
=\lambda_{2}\left(\mathrm{k}_{1} \stackrel{2}{v}_{i}^{2}+\mathrm{k}_{2} \stackrel{3}{v}_{i}\right) \quad\left(\because \lambda_{2}=\lambda_{3}\right) \\
=\lambda_{2} \mathrm{w}_{\mathrm{i}} . \tag{19}
\end{gather*}
$$

Thus $\mathrm{w}_{\mathrm{i}}$ is orthogonal to $v_{i}$ and $\mathrm{w}_{\mathrm{i}}$ is an eigenvector corresponding to $\lambda_{2}$.

Hence, any two orthogonal vectors that lie on the plane normal to $v_{i}^{1}$ can be chosen as the other two eigenvectors of $u_{i j}$.

Case 3: If $\lambda_{1}=\lambda_{2}=\lambda_{3}$.
In this case , the cubic equation in $\lambda$ becomes

$$
\left(\lambda-\lambda_{1}\right)^{3}=0,
$$

or

$$
\left|\begin{array}{ccc}
\lambda_{1}-\lambda & 0 & 0  \tag{20}\\
0 & \lambda_{1}-\lambda & 0 \\
0 & 0 & \lambda_{1}-\lambda
\end{array}\right|=0 .
$$

Comparing it with equation (3), we find

$$
\mathrm{u}_{\mathrm{ij}}=0 \text { for } \mathrm{i} \neq \mathrm{j}
$$

and

$$
u_{11}=u_{22}=u_{33}=\lambda_{1} .
$$

Thus

$$
\begin{equation*}
\mathbf{u}_{\mathrm{ij}}=\lambda_{1} \delta_{\mathrm{ij}} \tag{21}
\end{equation*}
$$

Let $\mathrm{v}_{\mathrm{i}}$ be any non - zero vector. Then

$$
\begin{align*}
\mathrm{u}_{\mathrm{ij}} \mathrm{v}_{\mathrm{j}} & =\lambda_{1} \delta_{\mathrm{ij}} \mathrm{v}_{\mathrm{j}}, \\
& =\lambda_{1} \mathrm{v}_{\mathrm{i}} . \tag{22}
\end{align*}
$$

This shows that $v_{i}$ is an eigenvector corresponding to $\lambda_{1}$. Thus, every non zero vector in space is an eigenvector which corresponds to the same eigenvalue $\lambda_{1}$. Of these vectors, we can certainly chose (at least) there vectors $v_{i}^{1}, v_{i}^{2}, v_{i}^{3}$ that are mutually orthogonal.

Thus, in every case, there exists (at least) three mutually orthogonal eigenvectors of $\mathrm{u}_{\mathrm{ij}}$.

Example: Consider a second order tensor $\mathrm{u}_{\mathrm{ij}}$ whose matrix representation is

$$
\left[\begin{array}{rrr}
1 & 0 & -1 \\
1 & 2 & 1 \\
2 & 2 & 3
\end{array}\right] .
$$

It is clear, the tensor $\mathrm{u}_{\mathrm{ij}}$ is not symmetric. We shall find eigenvalues and eigenvectors of $\mathrm{u}_{\mathrm{ij}}$.

The characteristic equation is

$$
\left|\begin{array}{ccc}
1-\lambda & 0 & -1 \\
1 & 2-\lambda & 1 \\
2 & 2 & 3-\lambda
\end{array}\right|=0
$$

or

$$
(1-\lambda)[(2-\lambda)(3-\lambda)-2]-1[2-2(2-\lambda)]=0
$$

or

$$
(1-\lambda)(2-\lambda)(3-\lambda)=0 .
$$

Hence, eigenvalues are

$$
\lambda_{1}=1, \lambda_{2}=2, \lambda_{3}=3
$$

which are all different.
We find that an unit eigenvector corresponding to $\lambda=1$ is

$$
\stackrel{1}{v_{i}}=\left(\frac{1}{\sqrt{2}}, \frac{1}{\sqrt{2}}, 0\right)
$$

an unit eigenvector corresponding to $\lambda=2$ is

$$
\stackrel{2}{v}_{i}=\left(\frac{2}{3},-\frac{1}{3},-\frac{2}{3}\right),
$$

and an unit eigenvector corresponding to $\lambda=3$ is

$$
\stackrel{3}{v_{i}}=\left(\frac{1}{\sqrt{6}},-\frac{1}{\sqrt{6}},-\frac{2}{\sqrt{6}}\right) .
$$

We note that

This happens due to non - symmetry of the tensor $\mathrm{u}_{\mathrm{ij}}$.
Example 2: Let the matrix of the components of the second order tensor $\mathrm{u}_{\mathrm{ij}}$ be

$$
\left[\begin{array}{lll}
2 & 2 & 0 \\
2 & 2 & 0 \\
0 & 0 & 1
\end{array}\right] .
$$

Find eigenvalues and eigenvectors of it.
We note that the tensor is symmetric. The characteristic equation is

$$
\left|\begin{array}{ccc}
2-\lambda & 2 & 0 \\
2 & 2-\lambda & 0 \\
0 & 0 & 1-\lambda
\end{array}\right|=0
$$

or

$$
\lambda(\lambda-1)(\lambda-4)=0 .
$$

Thus eigenvalues are

$$
\lambda_{1}=0, \lambda_{2}=1, \lambda_{3}=4,
$$

which are all different.
Let $\stackrel{1}{v}_{i}$ be the unit eigenvector corresponding to eigenvalue $\lambda_{1}=0$. Then, the system of homogegeous equations is

$$
\left[\begin{array}{lll}
2 & 2 & 0 \\
2 & 2 & 0 \\
0 & 0 & 1
\end{array}\right]\left[\begin{array}{c}
1 \\
v_{1} \\
1 \\
v_{2} \\
1 \\
v_{3}
\end{array}\right]=0
$$

This gives $\stackrel{1}{v_{1}}+\stackrel{1}{v_{2}}=0, \stackrel{1}{v_{1}}+\stackrel{2}{v_{2}}=0, \stackrel{1}{v_{3}}=0$.
we find

$$
\stackrel{1}{v_{i}}=\left(\frac{1}{\sqrt{2}}, \frac{-1}{\sqrt{2}}, 0\right) .
$$

## Similarly

$$
\stackrel{2}{v_{i}}=(0,0,1),
$$

and

$$
\stackrel{3}{v_{i}}=\left(\frac{1}{\sqrt{2}}, \frac{1}{\sqrt{2}}, 0\right),
$$

are eigen vectors corresponding to $\lambda_{2}=1$ and $\lambda_{3}=4$, respectively.
Moreover, these vector are mutually orthogonal.

## Books Recommended

1. Y.C. Fung Foundations of Solid Mechanics, Prentice Hall, Inc., New Jersey, 1965
2. T.M. Atanackovic, Theory of Elasticity for Scientists and Engineers, A. Guran Birkhauser, Boston, 2000
3. Saada, A.S. Elasticity - Theory and Applications, Pergamon Press, Inc., NY, 1974.
4. Sokolnikoff, I.S. Mathematical Theory of Elasticity, Tata McGraw Hill Publishing Company, Ltd., New Delhi, 1977.
5. Garg, N.R. and Generation of Displacements and Stresses in a Sharma, R.K. Multilayered Half-Space due to a Strip-Loading, Journal ISET, Vol 28, 1991, pp 1-26.

## Chapter-2 Analysis of Stress

### 2.1 INTRODUCTION

Deformation and motion of an elastic body are generally caused by external forces such as surface loads or internal forces such as earthquakes, nuclear explosions, etc. When an elastic body is subjected to such forces, its behaviour depends upon the magnitude of the forces, upon their direction , and upon the inherent strength of the material of which the body is made. Such forces give rise to interactions between neighbouring portions in the interior parts of the elastic solid. Such interactions are studied through the concept of stress. The concepts of stress vector on a surface and state of stress at a point of the medium shall be discussed.

An approach to the solution of problems in elastic solid mechanics is to examine deformations initially, and then consider stresses and applied loads. Another approach is to establish relationships between applied loads and internal stresses first and then to consider deformations. Regardless of the approach selected, it is necessary to derive the component relations individually.

### 2.2. BODY FORCES AND SURFACE FORCES

Consider a continuous medium. We refer the points of this medium to a rectangular cartesian coordinate system. Let $\tau$ represents the region occupied by the body in the deformed state. A deformable body may be acted upon by two different types of external forces.
(i) Body forces : These forces are those forces which act on every volume element of the body and hence on the entire volume of the body. For example , gravitational force is a body force (magnetic forces are also body forces). Let $\rho$ denotes the density of a volume element $\Delta \tau$ of the body $\tau$. Let $g$ be the gravitational force / acceleration. Then , the force acting on the mass $\rho \Delta \tau$ contained in volume $\Delta \tau$ is $\mathbf{g} . \rho \Delta \tau$.
(ii) Surface forces : These forces are those which act upon every surface element of the body. Such forces are also called contact forces. Loads applied over the exterior surface or bounding surface are examples of surface forces. A hydrostatic pressure acting on the surface of a body submerged in a liquid / water is a surface force.

Internal forces: In addition to the external forces, there are internal forces (such as earthquakes, nuclear explosions) which arise from the mutual interaction between various parts of the elastic body.

Now, we consider an elastic body in its undeformed state with no forces acting on it. Let a system of forces be applied to it. Due to these forces, the body is deformed and a system of internal forces is set up to oppose this deformation. These internal forces give rise to stress within the body. It is therefore necessary to consider how external forces are transmitted through the medium.

### 2.3 STRESS VECTOR ON A PLANE AT A POINT

Let us consider an elastic body in equilibrium under the action of a system of external forces. Let us pass a fictitious plane $\pi$ through a point $\mathbf{P}\left(\mathbf{x}_{1}, \mathbf{x}_{2}\right.$, $x_{3}$ ) in the interior of this body. The body can be considered as consisting of two parts, say , A
and $B$ and these parts are in welded contact at the interface $\pi$. Part $A$ of the body is in equilibrium under forces (external) and the effect of part $B$ on the plane $\pi$. We assume that this effect is continuously distributed over the surface of intersection.


Fig. (2.1)
Around the point $P$, let us consider a small surface $\delta S$ (on the place $\pi$ ) and let $\hat{v}$ be an outward unit normal vector (for the part $A$ of the body). The effect of part $B$ on this small surface element can be reduces to a force $\overline{\mathbf{Q}}$ and a vector couple $\overline{\mathbf{C}}$. Now, let $\delta S$ shrink in size towards zero in a manner such that the point $P$ always remain inside $\delta S$ and $\hat{v}$ remains the normal vector.
We assume that $\frac{\bar{Q}}{\delta S}$ tends to a definite limit $\overline{\mathbf{T}}\left(\mathbf{x}_{1}, \mathbf{x}_{2}, \mathbf{x}_{\mathbf{3}}\right)$ and that $\frac{\bar{C}}{\delta S}$ tends to zero as $\delta S$ tends to zero. Thus

$$
\begin{aligned}
& \lim _{\delta S \rightarrow 0} \frac{\bar{Q}}{\delta S}=\overline{\mathbf{T}}\left(\mathbf{x}_{\mathbf{1}}, \mathbf{x}_{\mathbf{2}}, \mathbf{x}_{\mathbf{3}}\right), \\
& \lim _{\delta S \rightarrow 0} \frac{\bar{C}}{\delta S}=\mathbf{0}
\end{aligned}
$$

Now $\overline{\mathbf{T}}$ is a surface force per unit area.
This force,$\overline{\mathbf{T}}$, is called the stress vector or traction on the plane $\pi$ at P .

Note 1: Forces acting over the surface of a body are never idealized point forces, they are, in reality, forces per unit area applied over some finite area. These external forces per unit area are called also tractions.

Note 2: Cauchy's stress postulate
If we consider another oriented plane $\pi^{\prime}$ containing the same point $P\left(\mathbf{x}_{\mathbf{i}}\right)$, then the stress vector is likely to have a different direction. For this purpose, Cauchy made the following postulated - known as Cauchy's stress postulate.
"The stress vector T depends on the orientation of the plane upon which it acts".

Let $\hat{v}$ be the unit normal to plane $\pi$ through the point $P$. This normal characterize the orientation of the plane upon which the stress vector acts. For this reason, we write the stress vector as $\stackrel{\hat{v}}{\hat{\sim}}$, indicating its dependence on the orientation $\hat{v}$.

## Cauchy's Reciprocal relation

When the plane $\pi$ is in the interior of the elastic body, the normal $\hat{v}$ has two possible directions that are opposite to each other and we choose one of these directions.


Fig. (2.2)
For a chosen $\hat{v}$, the stress vector $\hat{v}^{\hat{v}}$ is interpreted as the internal surface force per unit area acting on plane $\pi$ due to the action of part $B$ of the material / body which $\hat{v}$ is directed upon the part $\mathbf{A}$ across the plane $\pi$.

Consequently, $\stackrel{-\hat{\mathrm{v}}}{\mathrm{T}}$ is the internal surface force per unit area acting on $\pi$ due to the action of part $\mathbf{A}$ for which $\hat{v}$ is the outward drawn unit normal.

By Newton's third law of motion, vector $\underset{\sim}{T}$ and $\underset{\sim}{T} \underset{\sim}{\hat{v}}$ balance each other as the body is in equilibrium. Thus

$$
\underset{\sim}{T}=-\underset{\sim}{\hat{v}},
$$

which is known as Cauchy's reciprocal relation.
Homogeneous State of Stress
If $\pi$ and $\pi^{\prime}$ are any two parallel planes through any two points $P$ and $P^{\prime}$ of a continuous elastic body, and if the stress vector on $\pi$ at $P$ is equal to the stress on $\pi^{\prime}$ at $\mathbf{P}^{\prime}$, then the state of stress in the body is said to be a homogeneous state of stress.

### 2.4 NORMAL AND TANGENTIAL STRESSES

In general , the stress vector $\stackrel{v}{T}$ is inclined to the plane on which it acts and need not be in the direction of unit normal $\hat{v}$. The projection of ${ }^{v}$ on the normal $\hat{v}$ is called the normal stress. It is denoted by $\sigma$ or $\sigma_{n}$. The projection of $\underset{\sim}{v}$ on the plane $\pi$, in the plane of $\underset{\sim}{\hat{v}}$ and $\hat{v}$, is called the tangential or shearing stress. It is denoted by $\tau$ or $\sigma_{t}$.


Fig. (2.3)
Thus,

$$
\begin{align*}
& \sigma=\sigma_{n}=\underset{\sim}{v} \cdot \hat{v} \quad, \quad \tau=\sigma_{t}=\underset{\sim}{\underset{\sim}{v}} \cdot \hat{t},  \tag{1}\\
& |\underset{\sim}{T}|^{2}=\sigma_{n}^{2}+\sigma_{t}^{2} \tag{2}
\end{align*}
$$

where $\hat{t}$ is a unit vector normal to $\hat{v}$ and lies in the place $\pi$.
A stress in the direction of the outward normal is considered positive (i.e. $\sigma>0)$ and is called a tensile stress. A stress in the opposite direction is considered negative ( $\sigma<0$ ) and is called a compressible stress.

If $\sigma=0, \stackrel{\hat{v}}{\mathrm{~T}}$ is perpendicular to $\hat{v}$. Then , the stress vector $\underset{\sim}{\hat{v}}$ is called a pure shear stress or a pure tangential stress.

If $\boldsymbol{\tau}=\mathbf{0}$, then $\frac{\hat{v}}{\mathrm{v}}$ is parallel to $\hat{v}$. The stress vector $\mathrm{T}^{\hat{\mathrm{v}}}$ is then called pure normal stress.

When $\underset{\sim}{\text { T }}$ acts opposite to the normal $\hat{v}$, then the pure normal stress is called pressure ( $\sigma<\mathbf{0}, \tau=\mathbf{0}$ ).

From (1), we can write

$$
\begin{equation*}
\hat{\mathrm{v}}=\sigma \hat{v}+\tau \hat{t} \tag{3}
\end{equation*}
$$

and

$$
\begin{equation*}
\tau=\sqrt{|\underset{\sim}{\hat{\mathrm{T}}}|^{2}-\sigma^{2}} \tag{4}
\end{equation*}
$$

Note : $\sigma_{\mathbf{t}}=\tau=|\underset{\sim}{T}| \sin \alpha \Rightarrow|\sigma|=|\underset{\sim}{\mathrm{T}} \times \hat{v}|, \quad$ as $|\hat{v}|=\mathbf{1}$.

This $\tau$ in magnitude is given by the magnitude of vector product of $\frac{\hat{v}}{\mathrm{v}}$ and $\hat{v}$.

### 2.5 STRESS COMPONENTS

Let $P\left(x_{i}\right)$ be any point of the elastic medium whose coordinates are $\left(x_{1}, x_{2}\right.$, $x_{3}$ ) relative to rectangular cartesian system $0 x_{1} X_{2} X_{3}$.


Fig. (2.4)

Let $\stackrel{1}{T}$ denote the stress vector on the plane, with normal along $\mathbf{x}_{1}$ - axis, at the point $P$. Let the stress vector $\underset{\sim}{T}$ has components $\tau_{11}, \tau_{12}, \tau_{13}$, i.e.,

$$
\begin{align*}
\underset{\sim}{T} & =\tau_{\mathbf{1 1}} \hat{e}_{1}+\tau_{\mathbf{1 2}} \hat{e}_{2}+\tau_{13} \hat{e}_{3} \\
& =\tau_{\mathbf{1}} \hat{e}_{j} . \tag{1}
\end{align*}
$$

Let ${\underset{\sim}{\sim}}_{\underset{\sim}{2}}$ be the stress vector acting on the plane $\|$ to $\mathbf{x}_{1} \mathbf{x}_{3}$ - plane at P. Let

$$
\begin{align*}
\stackrel{2}{T} & =\tau_{21} \hat{e}_{1}+\tau_{22} \hat{e}_{2}+\tau_{23} \hat{e}_{3} \\
& =\tau_{2 \mathbf{j}} \hat{e}_{j} . \tag{2}
\end{align*}
$$

Similarly

$$
\begin{align*}
\stackrel{3}{T} & =\tau_{31} \hat{e}_{1}+\tau_{32} \hat{e}_{2}+\tau_{33} \hat{e}_{3} \\
& =\tau_{3 \mathrm{j}} \hat{e}_{j} . \tag{3}
\end{align*}
$$

Equations (1) - (3) can be condensed in the following form

$$
\begin{equation*}
\stackrel{i}{T}=\tau_{\mathrm{ij}} \hat{e}_{j} \tag{4}
\end{equation*}
$$

Then

$$
\begin{align*}
\stackrel{i}{T} \cdot \hat{e}_{k} & =\left(\tau_{\mathbf{i j}} \hat{e}_{j}\right) \cdot \hat{e}_{k} \\
& =\tau_{\mathbf{i j}} \delta_{\mathbf{j k}} \\
& =\tau_{\mathbf{i k}} . \tag{5}
\end{align*}
$$

Thus, for given $\mathbf{i} \& \mathbf{j}$, the quantity $\tau_{i j}$ represent the $\mathbf{j t h}$ components of the stress vector $\stackrel{i}{T}$ acting on a plane having $\hat{e}_{i}$ as the unit normal. Here, the first suffix $i$ indicates the direction of the normal to the plane through $P$ and the second suffix $j$ indicates the direction of the stress component. In all, we have 9 components $\tau_{i j}$ at the point $P\left(x_{i}\right)$ in the $o x_{1} x_{2} x_{3}$ system. These quantities are called stress - components. The matrix

$$
\left(\tau_{\mathbf{i j}}\right)=\left[\begin{array}{lll}
\tau_{11} & \tau_{12} & \tau_{13} \\
\tau_{21} & \tau_{22} & \tau_{23} \\
\tau_{31} & \tau_{32} & \tau_{33}
\end{array}\right]
$$

whose rows are the components of the three stress vectors, is called the matrix of the state of stress at $\mathbf{P}$. The dimensions of stress components are $\frac{\text { force }}{(\text { Length })^{2}}=\mathbf{M} \mathbf{L}^{-1} \mathbf{T}^{-2}$.

The stress components $\tau_{11}, \tau_{22}, \tau_{33}$ are called normal stresses and other components $\tau_{12}, \tau_{13}, \tau_{21}, \tau_{23}, \tau_{31}, \tau_{32}$ are called shearing stresses $\left(\because \frac{1}{T} . \hat{e}_{1}=\right.$ $\mathbf{e}_{11}, \underset{\sim}{T} . \hat{e}_{2}=\mathbf{e}_{12}$, etc)

In CGS system, the stress is measured in dyne per square centimeter.
In English system , it measured in pounds per square inch or tons per square inch.

## Dyadic Representation of Stress

It may be helpful to consider the stress tensor as a vector - like quantity having a magnitude and associated direction (s), specified by unit vector. The dyadic is such a representation. We write the stress tensor or stress dyadic as

$$
\begin{align*}
\overline{\bar{\tau}}= & \tau_{\mathrm{ij}} \hat{e}_{i} \hat{e}_{j} \\
= & \tau_{11} \hat{e}_{1} \hat{e}_{1}+\tau_{12} \hat{e}_{1} \hat{e}_{2}+\tau_{13} \hat{e}_{1} \hat{e}_{3}+\tau_{21} \hat{e}_{2} \hat{e}_{1}+\tau_{22} \hat{e}_{2} \hat{e}_{2} \\
& +\tau_{23} \hat{e}_{2} \hat{e}_{3}+\tau_{31} \hat{e}_{3} \hat{e}_{1}+\tau_{32} \hat{e}_{3} \hat{e}_{2}+\tau_{33} \hat{e}_{3} \hat{e}_{3} \tag{1}
\end{align*}
$$

where the juxtaposed double vectors are called dyads.
The stress vector $\underset{\sim}{i}$ acting on a plane having normal along $\hat{e}_{i}$ is evaluated as follows :

$$
\begin{aligned}
\underset{\sim}{T} & =\overline{\bar{\sigma}} \cdot \hat{e}_{i} \\
& =\left(\tau_{\mathbf{j k}} \hat{e}_{j} \hat{e}_{k}\right) \cdot \hat{e}_{i}
\end{aligned}
$$

$$
\begin{aligned}
& =\tau_{\mathbf{j k}} \hat{e}_{j}\left(\delta_{\mathbf{k} \mathbf{i}}\right) \\
& =\tau_{\mathbf{j i}} \hat{e}_{j} \\
& =\tau_{\mathbf{i j}} \hat{e}_{j} .
\end{aligned}
$$

### 2.6 STATE OF STRESS AT A POINT-THE STRESS TENSOR

We shall show that the state of stress at any point of an elastic medium on an oblique plane is completely characterized by the stress components at P.

Let $\underset{\sim}{v}$ be the stress vector acting on an oblique plane at the material point $P$, the unit normal to this plane being $\hat{v}=v_{i}$.

Through the point $P$, we draw three planar elements parallel to the coordinate planes. A fourth plane $A B C$ at a distance $h$ from the point $P$ and parallel to the given oblique plane at $P$ is also drawn. Now, the tetrahedron PABC contains the elastic material.


Fig. (2.5)
Let $\tau_{\mathrm{ij}}$ be the components of stress at the point P . Regarding the signs (negative or positive) of scalar quantities $\tau_{\mathrm{ij}}$, we adopt the following convention.

If one draws an exterior normal (outside the medium) to a given face of the tetrahedron PABC ,then the positive values of components $\tau_{i j}$ are associated with forces acting in the positive directions of the coordinate axes. On the other hand, if the exterior normal to a given face is pointing in a direction opposite to that of the coordinate axes, then the positive values of $\tau_{\mathrm{ij}}$ are associated with forces directed oppositely to the positive directions of the coordinate axes.

Let $\sigma$ be the area of the face ABC of the tetrahedron in figure. Let $\sigma_{1}, \sigma_{2}$, $\sigma_{3}$ be the areas of the plane faces PBC, PCA and PAB (having normals along $\quad x_{1}-, x_{2}-\& x_{3}-$ axes) respectively.

Then

$$
\begin{equation*}
\sigma_{i}=\sigma \cos \left(x_{i}, \hat{v}\right)=\sigma v_{i} \tag{1}
\end{equation*}
$$

The volume of the tetrahedron is

$$
\begin{equation*}
\mathbf{v}=\frac{1}{3} \mathbf{h} \sigma . \tag{2}
\end{equation*}
$$

Assuming the continuity of the stress vector $\stackrel{v}{T}=\stackrel{v}{T_{i}}$, the $\mathbf{x}_{\mathbf{i}}-$ component of the stress force acting on the face ABC of the tetrahedron PABC (made of elastic material) is

$$
\begin{equation*}
\left(\stackrel{v}{T_{i}}+\epsilon_{\mathfrak{i}}\right) \sigma \tag{3}
\end{equation*}
$$

provided $\quad \lim _{h \rightarrow 0} \epsilon_{i}=\mathbf{0}$.
Here, $\epsilon_{i}$ 's are inserted because the stress force act at points of the oblique plane $A B C$ and not on the given oblique plane through $P$. Under the assumption of continuing of stress field, quantities $\in_{i}$ 's are infinitesimals.
We note that the plane element PBC is a part of the boundary surface of the material contained in the tetrahedron. As such , the unit outward normal to PBC is $-\hat{e}_{1}$. Therefore, the $\mathbf{x}_{\mathbf{i}}$ - component of force due to stress acting on the face PBC of area $\sigma_{1}$ is

$$
\begin{equation*}
\left(-\tau_{1 i}+\epsilon_{1 i}\right) \sigma_{1} \tag{4a}
\end{equation*}
$$

where

$$
\lim _{h \rightarrow 0} \in \in_{1 i}=\mathbf{0} .
$$

Similarly forces on the face PCA and PAB are

$$
\left.\begin{array}{l}
\left(-\tau_{2 i}+\epsilon_{2 i}\right) \sigma_{2} \\
\left(-\tau_{3 i}+\epsilon_{3 i}\right) \sigma_{3}
\end{array}\right]
$$

with $\lim$

$$
\begin{equation*}
\in_{2 \mathbf{i}}=\lim _{\mathrm{h} \rightarrow 0} \in_{3 \mathrm{i}}=\mathbf{0} . \tag{4b}
\end{equation*}
$$

On combining (4a) and (4b), we write

$$
\begin{equation*}
\left(-\tau_{\mathrm{ji}}+\in_{\mathrm{ji}}\right) \sigma_{\mathrm{j}} \tag{5}
\end{equation*}
$$

as the $x_{i}$ - component of stress force acting on the face of area $\sigma_{j}$ provided

$$
\lim _{\mathrm{h} \rightarrow 0} \quad \in_{\mathbf{j i}}=\mathbf{0}
$$

In equation (5), the stress components $\tau_{i j}$ are taken with the negative sign as the exterior normal to a face of area $\sigma_{j}$ is in the negative direction of the $\mathbf{x}_{\mathrm{j}}$ - axis.

Let $F_{i}$ be the body force per unit volume at the point $P$. Then the $x_{i}-$ component of the body force acting on the volume of tetrahedron PABC is

$$
\begin{equation*}
\frac{1}{3} \mathbf{h} \sigma\left(\mathbf{F}_{\mathbf{i}}+\epsilon_{\mathbf{i}}^{\prime}\right) \tag{6}
\end{equation*}
$$

where $\epsilon_{i}$ 's are infinitesimal and

$$
\lim _{\mathrm{h} \rightarrow 0} \in_{\mathbf{i}}^{\prime}=\mathbf{0}
$$

Since the tetrahedral element PABC of the elastic body is in equilibrium, therefore, the resultant force acting on the material contained in PABC must be zero. Thus

$$
\left(\stackrel{v}{T_{i}}+\epsilon_{\mathbf{i}}\right) \sigma+\left(-\tau_{\mathbf{j i}}+\epsilon_{\mathbf{j i}}\right) \sigma_{\mathbf{j}}+\frac{1}{3}\left(\mathbf{F}_{\mathbf{i}}+\epsilon_{\mathbf{i}}^{\prime}\right) \mathbf{h} \sigma=\mathbf{0} .
$$

Using (1), above equation (after cancellation of $\sigma$ ) becomes

$$
\begin{equation*}
\left(\stackrel{v}{T}_{i}+\epsilon_{\mathbf{i}}\right)+\left(-\tau_{\mathbf{j i}}+\epsilon_{\mathbf{j i}}\right) \mathbf{v}_{\mathbf{j}}+\frac{1}{3}\left(\mathbf{F}_{\mathbf{i}}+\epsilon_{\mathbf{i}}^{\prime}\right) \mathbf{h}=\mathbf{0} . \tag{7}
\end{equation*}
$$

As we take the limit $h \rightarrow 0$ in (7), the oblique face ABC tends to the given oblique plane at $P$. Therefore, this limit gives

$$
\stackrel{v}{T_{i}}-\tau_{\mathrm{ji}} v_{\mathbf{j}}=\mathbf{0}
$$

or

$$
\begin{equation*}
\stackrel{v}{T_{i}}=\tau_{\mathrm{j} \mathrm{i}} v_{\mathrm{j}} \tag{8}
\end{equation*}
$$

This relation connecting the stress vector $\stackrel{v}{\sim}$ and the stress components $\tau_{i \mathbf{j}}$ is known as Cauchy's law or formula.
It is convenient to express the equation (8) in the matrix notation. This has the form

$$
\left[\begin{array}{c}
\stackrel{v}{T_{1}}  \tag{8a}\\
\stackrel{v}{T_{2}} \\
\stackrel{v}{T_{3}}
\end{array}\right]=\left[\begin{array}{lll}
\tau_{11} & \tau_{21} & \tau_{31} \\
\tau_{12} & \tau_{22} & \tau_{32} \\
\tau_{13} & \tau_{23} & \tau_{33}
\end{array}\right]\left[\begin{array}{c}
v_{1} \\
v_{2} \\
v_{3}
\end{array}\right]
$$

As $\stackrel{v}{T_{i}}$ and $v_{\mathrm{i}}$ are vectors. Equation (8) shows, by quotient law for tensors, that new components $\tau_{\mathrm{ij}}$ form a second order tensor.
This stress tensor is called the CAUCHY'S STRESS TENSOR.

We note that, through a given point, there exists infinitely many surface plane elements. On every one of these elements we can define a stress vector. The totality of all these stress vectors is called the state of stress at the point. The relation (8) enables us to find the stress vector on any surface element at a point by knowing the stress tensor at that point. As such, the state of stress at a point is completely determined by the stress tensor at the point.

Note: In the above, we have assumed, first , that stress can be defined everywhere in a body, and, second, that the stress field is continuous. These are the basic assumptions of continuum mechanics. Without these assumptions, we can do very little. However, in the further development of the theory , certain mathematical discontinuities will be permitted / allowed.

### 2.7 BASIC BALANCE LAWS

(A) Balance of Linear Momentum :

So far, we have discussed the state of stress at a point. If it is desired to move from one point to another, the stress components will change. Therefore, it is necessary to investigate the equations / conditions which control the way in which they change.

While the strain tensor $\mathrm{e}_{\mathrm{ij}}$ has to satisfy six compatibility conditions, the components of stress tensor must satisfy three linear partial differential equations of the first order. The principle of balance of linear momentum gives us these differential equations. This law , consistent with the Newton's second law of motion, states that the time rate of change of linear momentum is equal to the resultant force on the elastic body.

Consider a continuous medium in equilibrium with volume $\tau$ and bounded by a closed surface $\sigma$. Let $F_{i}$ be the components of the body force per unit volume and $\stackrel{v}{T_{i}}$ be the component of the surface force in the $\mathbf{x}_{\mathbf{i}}$ - direction.

For equilibrium of the medium , the resultant force acting on the matter within $\tau$ must vanish. That is

$$
\begin{equation*}
\int_{\tau} \mathbf{F}_{\mathbf{i}} \mathbf{d} \tau+\int_{\sigma} \stackrel{v}{i}_{i} \mathbf{d} \sigma=\mathbf{0}, \quad \text { for } \mathbf{i}=\mathbf{1}, \mathbf{2}, \mathbf{3} . \tag{1}
\end{equation*}
$$

We know the following Cauchy's formula

$$
\begin{equation*}
\stackrel{v}{T_{i}}=\tau_{\mathrm{j} \mathbf{i}} \boldsymbol{v}_{\mathrm{j}}, \quad(\mathbf{i}=1,2,3) \tag{2}
\end{equation*}
$$

where $\tau_{\mathrm{ij}}$ is the stress tensor and $v_{\mathrm{j}}$ is the unit normal to the surface. Using (2) into equation (1), we obtain

$$
\begin{equation*}
\int_{\tau} \mathbf{F}_{\mathbf{i}} \mathbf{d} \tau+\int_{\sigma} \tau_{\mathrm{ji}} v_{\mathrm{j}} \mathbf{d} \sigma=\mathbf{0}, \quad(\mathbf{i}=\mathbf{1}, \mathbf{2}, \mathbf{3}) \tag{3}
\end{equation*}
$$

We assume that stresses $\tau_{\mathrm{ij}}$ and their first order partial derivatives are also continuous and single valued in the region $\tau$. Under these assumptions, Gauss - divergence theorem can be applied to the surface integral in (3) and we find

$$
\begin{equation*}
\int_{\sigma} \tau_{\mathrm{ji}} v_{\mathbf{j}} \mathbf{d} \sigma=\int_{\tau} \tau_{\mathrm{ji}, \mathbf{j}} \mathbf{d} \tau . \tag{4}
\end{equation*}
$$

From equations (3) and (4), we write

$$
\begin{equation*}
\int_{\tau}\left(\tau_{\mathrm{j} i, \mathrm{j}}+\mathbf{F}_{\mathbf{i}}\right) \mathbf{d} \tau=\mathbf{0} \tag{5}
\end{equation*}
$$

for each $i=1,2,3$. Since the region $\tau$ of integration is arbitrary (every part of the medium is in equilibrium) and the integrand is continuous, so , we must have

$$
\begin{equation*}
\tau_{\mathbf{j i}, \mathrm{j}}+\mathbf{F}_{\mathrm{i}}=\mathbf{0} \tag{6}
\end{equation*}
$$

for each $i=1,2,3$ and at every interior point of the continuous elastic body. These equations are

$$
\begin{align*}
& \frac{\partial \tau_{11}}{\partial x_{1}}+\frac{\partial \tau_{21}}{\partial x_{2}}+\frac{\partial \tau_{31}}{\partial x_{3}}+F_{1}=0 \\
& \frac{\partial \tau_{12}}{\partial x_{1}}+\frac{\partial \tau_{22}}{\partial x_{2}}+\frac{\partial \tau_{32}}{\partial x_{3}}+F_{2}=0  \tag{7}\\
& \frac{\partial \tau_{13}}{\partial x_{1}}+\frac{\partial \tau_{23}}{\partial x_{2}}+\frac{\partial \tau_{33}}{\partial x_{3}}+F_{3}=0
\end{align*}
$$

These equations are referred to as Cauchy's equations of equilibrium. These equations are also called stress equilibrium equations. These equations are associated with undeformed cartesian coordinates.

These equations were obtained by Cauchy in 1827.
Note 1: In the case of motion of an elastic body, these equations (due to balance of linear momentum) take the form

$$
\begin{equation*}
\tau_{j i, j}+F_{i}=\rho \ddot{u}_{i} \tag{8}
\end{equation*}
$$

where $\ddot{u}_{i}$ is the acceleration vector and $\rho$ is the density (mass per unit volume) of the body.

Note 2: When body force $F_{i}$ is absent (or negligible), equations of equilibrium reduce to

$$
\begin{equation*}
\tau_{\mathrm{j}, \mathrm{j}, \mathrm{j}}=\mathbf{0} \tag{9}
\end{equation*}
$$

Example: Show that for zero body force, the state of stress for an elastic body given by

$$
\begin{aligned}
& \tau_{11}=x^{2}+y+3 z^{2}, \tau_{22}=2 x+y^{2}+2 \mathrm{z}, \tau_{33}=-2 x+y+z^{2} \\
& \tau_{12}=\tau_{21}=-x y+z^{3}, \tau_{13}=\tau_{31}=y^{2}-x z, \tau_{23}=\tau_{32}=x^{2}-y z
\end{aligned}
$$

is possible.
Example: Determine the body forces for which the following stress field describes a state of equilibrium

$$
\begin{aligned}
& \tau_{11}=-2 x^{2}-3 y^{2}-5 z, \tau_{22}=-2 y^{2}+7, \tau_{33}=4 x+y+3 z-5 \\
& \tau_{12}=\tau_{21}=z+4 x y-6, \tau_{13}=\tau_{31}=-3 x+2 y+1, \tau_{23}=\tau_{32}=0 .
\end{aligned}
$$

Example: Determine whether the following stress field is admissible in an elastic body when body forces are negligible.

$$
\left[\tau_{\mathrm{ij}}\right]=\left[\begin{array}{ccc}
y z+4 & \mathrm{z}^{2}+2 \mathrm{x} & 5 \mathrm{y}+\mathrm{z} \\
\cdot & \mathrm{xz}+3 \mathrm{y} & 8 \mathrm{x}^{3} \\
\cdot & \cdot & 2 \mathrm{xyz}
\end{array}\right]
$$

(B) Balance of Angular momentum

The principle of balance of angular momentum for an elastic solid is -
"The time rate of change of angular momentum about the origin is equal to the resultant moment about of origin of body and surface forces."

This law assures the symmetry of the stress tensor $\tau_{\mathrm{ij}}$.
Let a continuous elastic body in equilibrium occupies the region $\tau$ bounded by surface $\sigma$. Let $F_{i}$ be the body force acting at a point $P\left(x_{i}\right)$ of the body. Let the position vector of the point $P$ relative to the origin be $\vec{r}=$
$\mathbf{x}_{\mathbf{i}} \hat{e}_{i}$. Then, the moment of force $\overrightarrow{\mathrm{F}}$ is $\overrightarrow{\mathrm{r}} \times \overrightarrow{\mathrm{F}}=\epsilon_{\mathbf{i j k}} \mathbf{x}_{\mathbf{j}} \mathbf{F}_{\mathrm{k}}$, where $\epsilon_{\mathrm{ijk}}$ is the alternating tensor.

As the elastic body is in equilibrium, the resultant moment due to body and surface forces must be zero. So ,

$$
\begin{equation*}
\int_{\tau} \epsilon_{i j k} \mathbf{x}_{\mathbf{j}} \mathbf{F}_{\mathbf{k}} \mathbf{d} \tau+\int_{\sigma} \epsilon_{i j k} \mathbf{x}_{\mathbf{j}} \stackrel{v}{T_{k}} \mathbf{d} \sigma=\mathbf{0} \tag{1}
\end{equation*}
$$

for each $i=1,2,3$.
Since, the body is in equilibrium, so the Cauchy's equilibrium equations gives

$$
\begin{equation*}
\mathbf{F}_{\mathbf{k}}=-\tau_{\mathbf{l k}, \mathrm{l}} \tag{2a}
\end{equation*}
$$

The stress vector $\stackrel{v}{k}_{k}$ in terms of stress components is given by

$$
\begin{equation*}
\stackrel{v}{T_{k}}=\tau_{l \mathbf{k}} \boldsymbol{v}_{l} \tag{2b}
\end{equation*}
$$

The Gauss - divergence theorem gives us

$$
\begin{align*}
\int_{\sigma} \epsilon_{i j k} \mathbf{x}_{\mathbf{j}} \tau_{l \mathbf{k}} \boldsymbol{v}_{l} \mathbf{d} \sigma & =\int_{\tau}\left[\epsilon_{i j k} \mathbf{x}_{\mathbf{j}} \tau_{l \mathbf{k}, l}\right]_{\mathbf{d}} \tau \\
& =\int_{\tau} \epsilon_{i j k}\left[\mathbf{x}_{\mathbf{j}} \tau_{\mathbf{k}, l}+\delta_{\mathbf{j} l} \tau_{l \mathbf{k}}\right] \mathbf{d} \tau \\
& =\int_{\tau} \epsilon_{i j k}\left[\mathbf{x}_{\mathbf{j}} \tau_{l \mathbf{k}, l}+\tau_{\mathbf{j k}}\right] \mathbf{d} \tau . \tag{3}
\end{align*}
$$

From equations (1), (2a) and (3); we write

$$
\int_{\tau} \epsilon_{i j k} \mathbf{x}_{\mathbf{j}}\left(-\tau_{l \mathbf{k}, l}\right) \mathbf{d} \tau+\int_{\tau} \epsilon_{i j k}\left[\mathbf{x}_{\mathbf{j}} \tau_{\mathrm{k}, l}+\tau_{\mathbf{j k}}\right] \mathbf{d} \tau=\mathbf{0} .
$$

This gives

$$
\begin{equation*}
\int_{\tau} \epsilon_{i j k} \tau_{\mathbf{j k}} \mathbf{d} \tau=\mathbf{0}, \tag{4}
\end{equation*}
$$

for $i=1,2,3$. Since the integrand is continuous and the volume is arbitrary, so

$$
\begin{equation*}
\in_{\mathrm{ijk}} \tau_{\mathrm{jk}}=\mathbf{0}, \tag{5}
\end{equation*}
$$

for $i=1,2,3$ and at each point of the elastic body. Expanding (5), we write

$$
\begin{array}{ll} 
& \in_{123} \tau_{23}+\in_{132} \tau_{32}=\mathbf{0} \\
\Rightarrow & \tau_{23}-\tau_{32}=\mathbf{0}, \\
\Rightarrow & \in_{213} \tau_{13}+\in_{231} \tau_{31}=\mathbf{0} \\
& \tau_{13}=\tau_{31},  \tag{6}\\
\Rightarrow & \in_{312} \tau_{12}+\in_{321} \tau_{21}=\mathbf{0} \\
& \tau_{12}=\tau_{21} .
\end{array}
$$

That is

$$
\begin{equation*}
\tau_{i \mathbf{j}}=\tau_{\mathbf{j i}} \quad \text { for } \mathbf{i} \neq \mathbf{j} \tag{7}
\end{equation*}
$$

at every point of the medium.
This proves the symmetry of stress tensor.
This law is also referred to as Cauchy's second law. It is due to Cauchy in 1827.

Note 1: On account of this symmetry, the state of stress at every point is specified by six instead of nine functions of position.

Note 2: In summary, the six components of the state of the stress must satisfy three partial differential equations $\left(\tau_{\mathrm{j}, \mathrm{j}, \mathrm{j}}+\mathrm{F}_{\mathrm{i}}=0\right)$ within the body and the three relations $\left({ }_{T}^{v}=\tau_{\mathrm{j} \mathbf{i}} v_{\mathbf{j}}\right.$ ) on the bounding surface. The equations $\stackrel{v}{T}=\tau_{\mathrm{j} \mathbf{i}} \boldsymbol{v}_{\mathbf{j}}$ are called the boundary conditions.

Note 3 : Because of symmetry of the stress - tensor, the equilibrium equations may be written as

$$
\tau_{\mathrm{i}, \mathrm{j}, \mathrm{j}}+\mathrm{F}_{\mathrm{i}}=\mathbf{0}
$$

Note 4 : Since $\stackrel{i}{T}_{j}=\tau_{\mathbf{i j}}$, equations of equilibrium (using symmetry of $\tau_{\mathbf{i j}}$ ) may also be expressed as

$$
T_{j, j}^{i}=-\mathbf{F}_{\mathbf{i}}
$$

or

$$
\operatorname{div} \underset{\sim}{T}=-\mathbf{F}_{\mathbf{i}}
$$

Note 5 : Because of the symmetry of $\tau_{\mathrm{ij}}$, the boundary conditions can be expressed as

$$
\stackrel{v}{T_{i}}=\tau_{\mathrm{ij}} v_{\mathrm{j}}
$$

Remark : It is obvious that the three equations of equilibrium do not suffice for the determination of the six functions that specify the stress field. This may be expressed by the statement that the stress field is statistically indeterminate. To determine the stress field, the equations of equilibrium must be supplemented by other relations that can't be obtained from static considerations.

### 2.8 TRANSFORMATION OF COORDINATES

We have defined earlier the components of stress $\tau_{i j}$ with respect to cartesian system $0 \mathbf{x}_{1} \mathbf{x}_{2} \mathbf{x}_{3}$. Let $\mathrm{Ox}_{1}{ }^{\prime} \mathbf{x}_{\mathbf{2}}{ }^{\prime} \mathbf{x}_{3}{ }^{\prime}$ be any other cartesian system with the same origin but oriented differently. Let these coordinates be connected by the linear relations

$$
\begin{equation*}
\mathbf{x}_{\mathbf{p}}^{\prime}=\mathbf{a}_{\mathbf{p i}} \mathbf{x}_{\mathbf{i}} \tag{1}
\end{equation*}
$$

where $a_{p i}$ are the direction cosines of the $x_{p}{ }^{\prime}-\operatorname{axis}$ with respect to the $x_{i}$ axis. That is ,

$$
\begin{equation*}
\mathrm{a}_{\mathrm{pi}}=\cos \left(\mathrm{x}_{\mathrm{p}}^{\prime}, \mathrm{x}_{\mathrm{i}}\right) . \tag{2}
\end{equation*}
$$

Let $\tau^{\prime}{ }_{\mathrm{pq}}$ be the components of stress in the new reference system (Fig.).


Given stresses
Fig. (2.6)


## Desired stresses

Fig. (2.7) Transformation of stress components under rotation of co-ordinates system.

We shall now obtain a general formula, in the form of the theorem given below, which enables one to compute the component in any direction $\hat{v}$ of the stress vector acting on any given element with $\hat{v}^{\prime}$.

Theorem: let the surface element $\Delta \sigma$ and $\Delta \sigma^{\prime}$, with unit normals $\hat{v}$ and $\hat{v}^{\prime}$, pass through the point $P$. Show that the component of the stress vector ${ }_{T}^{\hat{v}}$ acting on $\Delta \sigma$ in the direction of $\hat{v}^{\prime}$ is equal to the component of the stress vector ${ }_{\mathrm{T}} \hat{\mathrm{v}}^{\prime}$ acting on $\Delta \sigma^{\prime}$ in the direction of $\hat{v}$.

Proof: In this theorem , it is required to show that

$$
\begin{equation*}
\underset{\sim}{\hat{v}} \cdot \hat{v}^{\prime}=\hat{\mathrm{T}}^{\hat{\mathrm{v}}^{\prime}} \cdot \hat{v} \tag{3}
\end{equation*}
$$

The Cauchy's formula gives us

$$
\begin{equation*}
\stackrel{v}{T_{i}}=\tau_{\mathrm{j} \mathrm{i}} v_{\mathrm{j}} \tag{4}
\end{equation*}
$$

and

$$
\begin{equation*}
\stackrel{v^{\prime}}{T_{i}}=\tau_{\mathrm{ij}} \mathbf{v}_{\mathbf{j}}^{\prime} \tag{5}
\end{equation*}
$$

due to symmetry of stress tensors as with

$$
\hat{v}=v_{\mathrm{j}} \text { and } \hat{v}^{\prime}=v_{\mathrm{j}}^{\prime} .
$$

Now

$$
\begin{aligned}
\stackrel{v^{\prime}}{T} \cdot \hat{v} & =\stackrel{v^{\prime}}{T_{i}} v_{\mathbf{i}} \\
& =\left(\tau_{\mathrm{ij}} v_{\mathrm{j}}^{\prime}\right) v_{\mathrm{i}} \\
& =\tau_{\mathrm{j} \boldsymbol{i}} v_{\mathrm{i}}^{\prime} v_{\mathbf{j}}
\end{aligned}
$$

$$
\begin{align*}
& =\left(\tau_{\mathrm{ij}} v_{\mathrm{j}}\right) v_{\mathrm{i}}^{\prime} \\
& =\stackrel{v}{T}_{i} v_{\mathrm{i}}^{\prime} \\
& =\stackrel{v}{T} \cdot \hat{v}^{\prime} . \tag{6}
\end{align*}
$$

This completes the proof of the theorem.
Article: Use the formula (3) to derive the formulas of transformation of the components of the stress tensor $\tau_{\mathrm{ij}}$.

Solution : Since the stress components $\tau^{\prime}{ }_{p q}$ is the projection on the $\mathbf{x}^{\prime}{ }_{\mathbf{q}}$ - axis of the stress vector acting on a surface element normal to the $\mathbf{x}^{\prime}{ }_{p}-\mathbf{a x i s}$ (by definition), we can write

$$
\begin{equation*}
\tau_{\mathbf{p q}}^{\prime}=\stackrel{p}{T}_{q}=\stackrel{\hat{v}^{\prime}}{T}, \hat{v} \tag{7}
\end{equation*}
$$

where

$$
\left.\begin{array}{l}
\hat{v}^{\prime} \text { is parallel to the } \mathbf{x}_{\mathbf{p}}^{\prime}-\text { axis } \\
\hat{v} \text { is parallel to the } \mathbf{x}_{\mathbf{q}}^{\prime}-\text { axis. } \tag{8}
\end{array}\right\}
$$

Equations (6) and (7) imply

$$
\begin{equation*}
\tau_{\mathrm{pq}}^{\prime}=\tau_{\mathrm{ij}} v_{\mathrm{i}}^{\prime} v_{\mathrm{j} .} \tag{9}
\end{equation*}
$$

Since

$$
\left.\begin{array}{l}
v_{\mathrm{i}}^{\prime}=\cos \left(\mathbf{x}_{\mathrm{p}}^{\prime}, \mathbf{x}_{\mathrm{i}}\right)=\mathbf{a}_{\mathrm{pi}}  \tag{10}\\
v_{\mathrm{j}}=\cos \left(\mathbf{x}_{\mathrm{q}}^{\prime}, \mathrm{x}_{\mathrm{j}}\right)=\mathbf{a}_{\mathbf{q} \cdot} \cdot
\end{array}\right\}
$$

Equation (9) becomes

$$
\begin{equation*}
\tau_{\mathrm{pq}}^{\prime}=\mathbf{a}_{\mathrm{pi}} \mathbf{a}_{\mathbf{q j}} \tau_{\mathrm{ij} j} \tag{11}
\end{equation*}
$$

Equation (11) and definition of a tensor of order 2 show that the stress components $\tau_{\mathrm{ij}}$ transform like a cartesian tensor of order 2 . Thus, the physical concept of stress which is described by $\tau_{\mathrm{ij}}$ agrees with the mathematical definition of a tensor of order 2 in a Euclidean space.

Theorem: Show that the quantity

$$
\theta=\tau_{11}+\tau_{22}+\tau_{33}
$$

is invariant relative to an orthogonal transformation of cartesian coordinates.

Proof: Let $\tau_{\mathrm{ij}}$ be the tensor relative to the cartesian system o $\mathbf{x}_{1} \mathbf{x}_{\mathbf{2}} \mathbf{x}_{\mathbf{3}}$. Let these axes be transformed to 0 $x_{1}{ }^{\prime} \mathbf{x}_{2}{ }^{\prime} \mathbf{x}_{3}{ }^{\prime}$ under the orthogonal transformation

$$
\begin{equation*}
\mathbf{x}_{\mathrm{p}}^{\prime}=\mathbf{a}_{\mathrm{pi}} \mathbf{x}_{\mathrm{i}} \tag{1}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathbf{a}_{\mathrm{pi}}=\cos \left(\mathrm{x}_{\mathrm{p}}^{\prime}, \mathrm{x}_{\mathrm{i}}\right) \tag{2}
\end{equation*}
$$

Let $\tau^{\prime}{ }^{\prime}{ }^{\text {q }}$ be the stress components relative to new axes. Then these components are given by the rule for second order tensors ,

$$
\begin{equation*}
\tau_{\mathbf{p q}}^{\prime}=\mathbf{a}_{\mathbf{p i}} \mathbf{a}_{\mathbf{q j}} \tau_{\mathrm{ij}} \tag{3}
\end{equation*}
$$

Putting $q=p$ and taking summation over the common suffix , we write

$$
\begin{aligned}
\tau_{\mathrm{pp}}^{\prime} & =\mathbf{a}_{\mathrm{pi}} \mathbf{a}_{\mathrm{pj}} \tau_{\mathrm{ij}} \\
& =\delta_{\mathrm{ij}} \tau_{\mathrm{ij}} \\
& =\tau_{\mathrm{ij}}
\end{aligned}
$$

This implies

$$
\begin{equation*}
\tau_{11}^{\prime}+\tau_{22}^{\prime}+\tau_{33}^{\prime}=\tau_{11}+\tau_{22}+\tau_{33}=\theta \tag{4}
\end{equation*}
$$

This proves the theorem.
Remark: This theorem shows that whatever be the orientation of three mutually orthogonal planes passing through a given point , the sum of the normal stresses is independent of the orientation of these planes.

Exercise 1 : Prove that the tangential traction, parallel to a line 1 , across a plane at right angles to a line $l^{\prime}$, the two lines being at right angles to each other, is equal to the tangential traction, parallel to the line $l^{\prime}$, across a plane at right angles to 1 .

Exercise 2 : Show that the following two statements are equivalent.
(a) The components of the stress are symmetric.
(b) Let the surface elements $\Delta \sigma$ and $\Delta \sigma^{\prime}$ with respective normals $\hat{v}$ and $\hat{v}^{\prime}$ pass through a point $P$. Then

$$
\stackrel{v}{T} \cdot \hat{v}^{\prime}=\stackrel{v^{\prime}}{\underset{\sim}{T}} \cdot \hat{v}
$$

Hint : (b) $\Rightarrow \mathbf{( a ) . ~ L e t ~} \hat{v}=\hat{i}$ and $\hat{v}^{\prime}=\hat{j}$

Then

$$
\underset{\sim}{T} \cdot \hat{v}^{\prime}=\stackrel{i}{T} \cdot \hat{\sim}={\underset{T}{T}}_{j}^{i}=\tau_{\mathbf{i j}}
$$

and

$$
\stackrel{v_{\sim}^{\prime}}{\underset{\sim}{T}} \cdot \hat{v}=\stackrel{\underset{\sim}{\mathrm{j}}}{\sim} \cdot \hat{i}=\stackrel{j}{T_{i}}=\tau_{\mathrm{ji}}
$$

by assumption,$\stackrel{v}{\sim} \cdot \hat{v}^{\prime}=\stackrel{v^{\prime}}{T} \cdot \hat{v}$, therefore $\tau_{\mathrm{ij}}=\tau_{\mathbf{j i}}$

This shows that $\tau_{\mathrm{ij}}$ is symmetric.
Example 1: The stress matrix at a point $\mathbf{P}$ in a material is given as

$$
\left[\tau_{\mathrm{ij}}\right]=\left[\begin{array}{ccc}
3 & 1 & 4 \\
1 & 2 & -5 \\
4 & -5 & 0
\end{array}\right]
$$

Find
(i) the stress vector on a plane element through $P$ and parallel to the plane $2 \mathrm{x}_{1}+\mathrm{x}_{2}-\mathrm{x}_{3}=\mathbf{1}$,
(ii) the magnitude of the stress vector, normal stress and the shear stress, (iii) the angle that the stress vector makes with normal to the plane.

Solution: (i) The plane element on which the stress - vector is required is parallel to the plane $\quad 2 x_{1}+x_{2}-x_{3}=1$. Therefore, direction ratios of the normal to the required plane at $P$ are $<2,1,-1>$. So, the d.c.'s of the unit normal $\hat{v}=v_{i}$ to the required plane at $P$ are

$$
v_{1}=\frac{2}{\sqrt{6}}, v_{2}=\frac{1}{\sqrt{6}}, v_{3}=-\frac{1}{\sqrt{6}} .
$$

Let $\underset{\sim}{v}=\stackrel{v}{T}_{i}$ be the required stress vector. Then, Cauchy's formula gives

$$
\left[\begin{array}{c}
T_{1} \\
T_{1} \\
T_{2} \\
T_{3}
\end{array}\right]\left[\begin{array}{ccc}
3 & 1 & 4 \\
1 & 2 & -5 \\
4 & -5 & 0
\end{array}\right]\left[\begin{array}{c}
2 / \sqrt{6} \\
1 / \sqrt{6} \\
-1 / \sqrt{6}
\end{array}\right]
$$

or

$$
\stackrel{v}{T_{1}}=\sqrt{3 / 2}, \stackrel{v}{T_{2}}=3 \sqrt{3 / 2}, \stackrel{v}{T_{3}}=\sqrt{3 / 2} .
$$

So , the required stress vector at $P$ is

$$
\stackrel{v}{\sim}=\sqrt{\frac{3}{2}}\left(\hat{e}_{1}+3 \hat{e}_{2}+\hat{e}_{3}\right)
$$

and

$$
|\stackrel{v}{T}|=\sqrt{\frac{33}{2}}
$$

(ii) The normal stress $\sigma$ is given by

$$
\sigma=\stackrel{v}{\sim} \cdot \hat{v}=\sqrt{\frac{3}{2}} \cdot \frac{1}{\sqrt{6}}(\mathbf{2}+\mathbf{3}-\mathbf{1})=\frac{1}{2} \times \mathbf{4}=\mathbf{2}
$$

and the shear stress $\tau$ is given by

$$
\tau=\sqrt{|\underset{\sim}{T}|^{2}-\sigma^{2}}=\sqrt{\frac{33}{2}-4}=\frac{5}{\sqrt{2}} .
$$

(As $\tau \neq 0$, so the stress vector $\stackrel{v}{T}$ need not be along the normal to the plane element)
(iii) Let $\theta$ be the angle between the stress vector ${ }^{v}$ and normal $\hat{v}$. Then

$$
\cos \theta=\frac{\stackrel{v}{\sim} \cdot \hat{v}}{|\stackrel{v}{T}||\hat{v}|}=\frac{2}{\sqrt{\frac{33}{2}}}=\sqrt{\frac{8}{33}}
$$

This determines the required inclination.
Example 2: The stress matrix at a point $\mathbf{P}\left(\mathbf{x}_{\mathbf{i}}\right)$ in a material is given by

$$
\left[\tau_{\mathbf{i j}}\right]=\left[\begin{array}{ccc}
x_{3} x_{1} & x_{3}{ }^{2} & 0 \\
x_{3}{ }^{2} & 0 & -x_{2} \\
0 & -x_{2} & 0
\end{array}\right] .
$$

Find the stress vector at the point $Q(1,0,-1)$ on the surface $x_{2}{ }^{2}+x_{3}{ }^{2}=x_{1}$. Solution: The stress vector $\underset{\sim}{v}$ is required on the surface element

$$
\mathbf{f}\left(\mathbf{x}_{1}, \mathbf{x}_{2}, \mathbf{x}_{3}\right)=\mathrm{x}_{1}-\mathbf{x}_{2}{ }^{2}-\mathrm{x}_{3}{ }^{2}=\mathbf{0},
$$

at the point $Q(1,0,-1)$.
We find $\nabla \mathbf{f}=\hat{e}_{1}+2 \hat{e}_{3}$ and $|\nabla \mathbf{f}|=\sqrt{5}$ at the point $\mathbf{Q}$.

Hence, the unit outward normal $\hat{v}=\boldsymbol{v}_{\mathbf{i}}$ to the surface $\mathbf{f}=\mathbf{0}$ at the point $Q(1,0,-1)$ is

$$
\hat{v}=\frac{\nabla f}{|\nabla f|}=\frac{1}{\sqrt{5}}\left(\hat{e}_{1}+2 \hat{e}_{3}\right) .
$$

giving

$$
v_{1}=\frac{1}{\sqrt{5}}, v_{2}=0, v_{3}=\frac{2}{\sqrt{5}} .
$$

The stress matrix at the point $Q(1,0,-1)$ is

$$
\left[\tau_{\mathrm{ij}}\right]=\left[\begin{array}{rrr}
-1 & 1 & 0 \\
1 & 0 & 0 \\
0 & 0 & 0
\end{array}\right]
$$

Let $\underset{\sim}{v}=\stackrel{v}{T}$ be the required stress vector at the point $\mathbf{Q}$. Then, by Cauchy's law

$$
\left[\begin{array}{c}
\stackrel{v}{T_{1}} \\
\stackrel{v}{T_{2}} \\
\stackrel{v}{T_{3}}
\end{array}\right]=\left[\begin{array}{rrr}
-1 & 1 & 0 \\
1 & 0 & 0 \\
0 & 0 & 0
\end{array}\right]\left[\begin{array}{c}
1 / \sqrt{5} \\
0 \\
2 / \sqrt{5}
\end{array}\right]
$$

We find

$$
\stackrel{v}{T_{1}}=\frac{-1}{\sqrt{5}}, \quad \stackrel{v}{T_{2}}=\frac{1}{\sqrt{5}}, \quad \stackrel{\mathrm{v}}{\mathrm{~T}_{3}}=0 .
$$

Hence, the required stress vector at $\mathbf{Q}$ is

$$
\underset{\sim}{v}=\frac{1}{\sqrt{5}}\left(-\hat{e}_{1}+\hat{e}_{2}\right) .
$$

Example 3: The stress matrix at a certain point in a material is given by

$$
\left[\tau_{\mathrm{ij}}\right]=\left[\begin{array}{lll}
3 & 1 & 1 \\
1 & 0 & 2 \\
1 & 2 & 0
\end{array}\right]
$$

Find the normal stress and the shear stress on the octahedral plane element through the point.

Solution: An octahedral plane is a plane whose normal makes equal angles with positive directions of the coordinate axes. Hence, the components of the unit normal $\hat{v}=\hat{v}_{\mathbf{i}}$ are

$$
v_{1}=v_{2}=v_{3}=\frac{1}{\sqrt{3}} .
$$

Let $\stackrel{v}{T}=\stackrel{v}{T}$ be the stress vector through the specified point. Then, Cauchy's formula gives

$$
\left[\begin{array}{c}
\stackrel{v}{T_{1}} \\
\stackrel{v}{T_{2}} \\
\stackrel{v}{T_{3}}
\end{array}\right]=\frac{1}{\sqrt{3}}\left[\begin{array}{lll}
3 & 1 & 1 \\
1 & 0 & 2 \\
1 & 2 & 0
\end{array}\right]\left[\begin{array}{l}
1 \\
1 \\
1
\end{array}\right]=\frac{1}{\sqrt{3}}\left[\begin{array}{l}
5 \\
3 \\
3
\end{array}\right]
$$

The magnitude of this stress vector is

$$
|\stackrel{v}{\underset{\sim}{v}}|=\sqrt{\frac{43}{3}}
$$

Let $\sigma$ be the normal stress and $\tau$ be the shear stress. Then

$$
\sigma=\underset{\sim}{\underset{\sim}{v}} \cdot \hat{v}=\frac{1}{3}(\mathbf{5}+\mathbf{3}+\mathbf{3})=\frac{11}{3},
$$

and

$$
\tau=\sqrt{\frac{43}{3}-\frac{121}{9}}=\sqrt{\frac{129=121}{9}}=\sqrt{\frac{8}{9}}=\frac{2 \sqrt{2}}{3} .
$$

Since $\sigma>0$, the normal stress on the octahedral plane is tensile.
Example 4: The state of stress at a point $P$ in cartesian coordinates is given by

$$
\begin{aligned}
& \tau_{11}=500, \tau_{12}=\tau_{21}=500, \tau_{13}=\tau_{31}=800, \tau_{22}=1000, \\
& \tau_{33}=-300, \tau_{23}=\tau_{32}=-750 .
\end{aligned}
$$

Compute the stress vector $T$ and the normal and tangential components on the plane passing through $P$ whose outward normal unit vector is

$$
\hat{v}=\frac{1}{2} \hat{e}_{1}+\frac{1}{2} \hat{e}_{2}+\frac{1}{\sqrt{2}} \hat{e}_{3} .
$$

Solution: The stress vector

$$
T=\mathbf{T}_{\mathbf{i}} \hat{e}_{i}
$$

is given by

$$
\mathbf{T}_{\mathrm{i}}=\tau_{\mathrm{ji}} v_{\mathrm{j}}
$$

## We find

$$
\begin{aligned}
& \mathrm{T}_{1}=\tau_{11} v_{1}+\tau_{21} v_{2}+\tau_{31} v_{3}=250+250+400 \sqrt{2} \\
&=500+400 \times(1.41) \\
&=500+564=1064, \text { approx. } \\
& \mathrm{T}_{2}=\tau_{12} v_{1}+\tau_{22} v_{2}+\tau_{32} v_{3}=250+250+\frac{750}{\sqrt{2}} \\
&=221, \text { App. } \\
& \begin{aligned}
& \mathbf{T}_{3}=\tau_{13} v_{1}+\tau_{23} v_{2}+\tau_{33} v_{3}=400-375-150 \sqrt{2}=25- \\
& 150(1.41)
\end{aligned} \\
&=-187, \text { app. }
\end{aligned}
$$

### 2.9 STRESS QUADRIC

In a trirectangular cartesian coordinate system $0 x_{1} \mathbf{x}_{2} \mathbf{x}_{3}$, consider the equation

$$
\begin{equation*}
\tau_{\mathrm{ij}} \mathbf{x}_{\mathrm{i}} \mathbf{x}_{\mathrm{j}}= \pm \mathbf{k}^{2} \tag{1}
\end{equation*}
$$

where ( $x_{1}, x_{2}, x_{3}$ ) are the coordinates a point $P$ relative to the point $P^{0}$ whose coordinates relative to origin $O$ are $\left(x_{1}{ }^{0}, x_{2}{ }^{0}, x_{3}{ }^{0}\right), \tau_{\mathrm{ij}}$ is the stress tensor at the point $\mathbf{P}^{\mathbf{0}}\left(\mathbf{x}_{\mathrm{i}}{ }^{\mathbf{0}}\right)$, and k is a real constant .

The sign + or - is so chosen that the quadric surface (1) is real.
The quadric surface (1) is known as the stress quadric of Cauchy with its centre at the point $P^{0}\left(x_{i}{ }^{0}\right)$.


Fig. (2.8)
Let $A_{i}$ be the radius vector, of magnitude $A$, on this stress quadric surface which is normal on the plane $\pi$ through the point $\mathbf{P}^{\mathbf{0}}$ having stress tensor $\tau_{\mathrm{ij}}$. Let $\hat{v}$ be the unit vector along the vector $\mathrm{A}_{\mathrm{i}}$. Then

$$
\begin{equation*}
v_{i}=\mathbf{A}_{\mathbf{i}} / \mathbf{A}=\mathbf{x}_{\mathbf{i}} / \mathbf{A} . \tag{2}
\end{equation*}
$$

Let $\underset{\sim}{v}$ denote the stress vector on the plane $\pi$ at the point $\mathbf{P}^{\mathbf{0}}$. Then , the normal stress $\mathbf{N}$ on the plane $\pi$ is given by

$$
\begin{equation*}
\mathbf{N}=\stackrel{v}{T} \cdot \hat{v}=\stackrel{T}{T}_{i}^{v} v_{\mathrm{i}}=\tau_{\mathrm{i} \mathrm{j}} v_{\mathrm{j}} v_{\mathrm{i}}=\tau_{\mathrm{i} \mathrm{j}} v_{\mathrm{i}} v_{\mathrm{j}} \tag{3}
\end{equation*}
$$

From equations (1) and (2), we obtain

$$
\begin{aligned}
& \tau_{\mathrm{ij}}\left(\mathrm{~A} v_{\mathrm{i}}\right)\left(\mathrm{A} v_{\mathrm{j}}\right)= \pm \mathrm{k}^{2} \\
& \tau_{\mathrm{ij}} v_{\mathrm{i}} v_{\mathrm{j}}= \pm \mathrm{k}^{2} / \mathrm{A}^{2}
\end{aligned}
$$

$$
\begin{equation*}
\mathrm{N}= \pm \mathrm{k}^{2} / \mathrm{A}^{2} . \tag{4}
\end{equation*}
$$

This gives the normal stress acting on the plane $\pi$ with orientation $\hat{v}=v_{\mathrm{i}}$ in terms of the length of the radius vector of the stress quadric from the point (centre) $P^{0}$ along the vector $v_{i}$.

The relation (4) shows that the normal stress $N$ on the plane $\pi$ through $P^{0}$ with orientation along $A_{i}$ is inversely proportional to the square of that radius vector $\mathbf{A}_{\mathbf{i}}=\overline{P^{o} P}$ of the stress quadric.

The positive sign in (1) or (4) is chosen whenever the normal stress $N$ represents tension (i.e., $N>0$ ) and negative sign when $N$ represents compression (i.e. , $\mathrm{N}<\mathbf{0}$ ).

The Cauchy's stress quadric (1) possesses another interesting property. This property is
"The normal to the quadric surface at the end of the radius vector $\mathrm{A}_{\mathrm{i}}$ is parallel to the stress vector $\underset{\sim}{T}$ acting on the plane $\pi$ at $\mathrm{P}^{\mathrm{o}}$."

To prove this property, let us write equation (1) in the form

$$
\begin{equation*}
\mathbf{G}\left(\mathbf{x}_{1}, \mathbf{x}_{2}, \mathbf{x}_{3}\right) \equiv \tau_{\mathrm{ij}} \mathbf{x}_{\mathrm{i}} \mathbf{x}_{\mathbf{j}} \mp \mathbf{k}^{2}=\mathbf{0} \tag{5}
\end{equation*}
$$

Then the direction of the normal to the stress quadric surface is given by the gradient of the scalar point function G. The components of gradient are

$$
\begin{align*}
\frac{\partial G}{\partial x_{n}}=\tau_{\mathrm{ij}}\left(\delta_{\mathrm{in}}\right) \mathbf{x}_{\mathbf{j}}+\tau_{\mathrm{ij}} \mathbf{x}_{\mathbf{i}}\left(\delta_{\mathrm{jn}}\right) & =\mathbf{2} \tau_{\mathrm{nj}} \mathbf{x}_{\mathbf{j}} \\
& =2 \mathbf{A} \tau_{\mathrm{nj}} \mathbf{v}_{\mathbf{j}} \\
& =\mathbf{2} \mathbf{A} \stackrel{v}{n}_{n}^{v} \tag{6}
\end{align*}
$$



Fig. (2.9)

Equation (6) shows that vectors $\stackrel{v}{T}_{n}$ and $\frac{\partial G}{\partial x_{n}}$ are parallel. Hence the stress vector $\underset{\sim}{v}$ on the plane $\pi$ at $\mathbf{P}^{\mathbf{0}}$ is directed along the normal to the stress quadric at $\mathrm{P}, \mathrm{P}$ being the end point of the radius vector $\quad \mathbf{A}_{\mathbf{i}}=\overline{\mathrm{P}^{\mathrm{O}} \mathrm{P}}$.

Remark 1: Equation (6) can be rewritten as

$$
\begin{equation*}
\underset{\sim}{\underset{\sim}{v}}=\frac{1}{2 A}(\underline{\nabla} G) . \tag{7}
\end{equation*}
$$

This relations gives an easy way of constructing the stress vector $\underset{\sim}{v}$ from the knowledge of the quadric surface $G\left(x_{1}, x_{2}, x_{3}\right)=$ constant and the magnitude $A$ of the radius vector $A_{i}$.

Remark 2: Taking principal axes along the coordinate axes, the stress quadric of Cauchy assumes the form

$$
\begin{equation*}
\tau_{1} \mathbf{x}_{1}{ }^{2}+\tau_{2} \mathbf{x}_{2}{ }^{2}+\tau_{3} \mathbf{x}_{3}{ }^{2}= \pm \mathbf{k}^{2} \tag{8}
\end{equation*}
$$

Here the coefficients $\tau_{1}, \tau_{2}, \tau_{3}$ are the principal stresses. Let the axes be so numbered that $\tau_{1} \geq \tau_{2} \geq \tau_{3}$.

If $\tau_{1}>\tau_{2}>\tau_{3}>0$, then equation (8) represents an ellipsoid with plus sign. Then, the relation $N=k^{2} / A^{2}$ implies that the force acting on every surface element through $\mathbf{P}^{\mathbf{0}}$ is tensile (as $\mathbf{N}<0$ ).

If $0>\tau_{1}>\tau_{2}>\tau_{3}$, then equation (8) represents an ellipsoid with a negative sign on the right and $N=-k^{2} / A^{2}$ indicates that the normal stress is compressive ( $\mathbf{N}>\mathbf{0}$ ).If $\tau_{1}=\tau_{2} \neq \tau_{3}$ or $\tau_{1} \neq \tau_{2}=\tau_{3}$ or $\tau_{1}=\tau_{3} \neq \tau_{2}$, then the Cauchy's stress quadric is an ellipsoid of revolution.

If $\tau_{1}=\tau_{2}=\tau_{3}$, then the stress quadric is a sphere.

### 2.10 PRINCIPAL STRESSES

In a general state of stress, the stress vector ${ }_{T}^{v}$ acting on a surface with outer normal $\hat{v}$ depends on the direction of $\hat{v}$.

Let us see in what direction $\hat{v}$ the stress vector ${ }_{T}^{v}$ becomes normal to the surface, on which the shearing stress is zero. Such a surface shall be called a principal plane, its normal a principal axis, and the value of normal stress acting on the principal plane shall be called a principal stress.

Let $\hat{v}$ define a principal axis at the point $\mathbf{P}^{0}\left(\mathbf{x}_{i}{ }^{0}\right)$ and let $\tau$ be the corresponding principal stress and $\tau_{\mathrm{ij}}$ be the stress tensor at that point. Let $\stackrel{v}{T}$ be the stress vector. Then

$$
\stackrel{v}{T}=\tau \hat{v},
$$

or

$$
\begin{equation*}
\stackrel{v}{T_{i}}=\tau v_{\mathrm{i}} \tag{1}
\end{equation*}
$$

Also

$$
\begin{equation*}
\stackrel{v}{T_{i}}=\tau_{\mathrm{ij}} v_{\mathrm{j}} \tag{2}
\end{equation*}
$$

Therefore

$$
\tau_{\mathrm{ij}} v_{\mathrm{j}}=\tau v_{\mathrm{i}}=\tau \delta_{\mathrm{ij}} v_{\mathrm{j}}
$$

or

$$
\begin{equation*}
\left(\tau_{\mathrm{ij}}-\tau \delta_{\mathrm{ij}}\right) v_{\mathrm{j}}=\mathbf{0} \tag{3}
\end{equation*}
$$

The three equations, $i=1,2,3$, are to be solved for $v_{1}, v_{2}, v_{3}$. Since $\hat{v}$ is a unit vector, we must find a set of non - trivial solutions for which

$$
v_{1}^{2}+v_{2}^{2}+v_{3}^{2}=1
$$

Thus, equation (3) poses an eigenvalue problem. Equation (3) has a set of non - vanishing solutions $v_{1}, v_{2}, v_{3}$ iff the determinant of the coefficients vanishes, i.e.,
$\begin{array}{cc} & \left|\tau_{i j}-\tau \delta_{\mathrm{ij}}\right|=\mathbf{0}, \\ \text { or } & {\left[\begin{array}{ccc}\tau_{11}-\tau & \tau_{12} & \tau_{13} \\ \tau_{12} & \tau_{22}-\tau & \tau_{23} \\ \tau_{13} & \tau_{23} & \tau_{33}-\tau\end{array}\right]=\mathbf{0 .}}\end{array}$
On expanding (2), we find

$$
\begin{equation*}
-\tau^{3}+\theta_{1} \tau^{2}-\theta_{2} \tau+\theta_{3}=\mathbf{0} \tag{3b}
\end{equation*}
$$

where

$$
\theta_{1}=\tau_{11}+\tau_{22}+\tau_{33}
$$

(4a)

$$
\theta_{2}=\left|\begin{array}{cc}
\tau_{22} & \tau_{23}  \tag{4b}\\
\tau_{23} & \tau_{33}
\end{array}\right|+\left|\begin{array}{cc}
\tau_{11} & \tau_{13} \\
\tau_{31} & \tau_{33}
\end{array}\right|+\left|\begin{array}{cc}
\tau_{11} & \tau_{12} \\
\tau_{21} & \tau_{22}
\end{array}\right|,
$$

$$
\begin{equation*}
\theta_{3}=\epsilon_{\mathrm{ijk}} \tau_{1 \mathrm{i}} \tau_{2 \mathrm{j}} \tau_{3 \mathrm{k}}=\operatorname{det} .\left(\tau_{\mathrm{ij}}\right) . \tag{4c}
\end{equation*}
$$

Equation (3) is a cubic equation in $\tau$. Let its roots be $\tau_{1}, \tau_{2}, \tau_{3}$. Since the matrix of stress, $\left(\tau_{\mathrm{ij}}\right)$ is real and symmetric , the roots $\tau_{\mathrm{i}}$ of $(3)$ are all real. Thus, $\tau_{1}, \tau_{2}, \tau_{3}$ are the principal stresses.

For each value of the principal stress, a unit normal vector $\hat{v}$ can be determined.

Case I : When $\tau_{1} \neq \tau_{2} \neq \tau_{3}$, let $\stackrel{1}{v_{i}}, \stackrel{2}{v_{i}}, v_{i}$ be the unit principal axes corresponding to the principal stresses $\tau_{1}, \tau_{2}, \tau_{3}$, respectively. Then principal axes are mutually orthogonal to each other.

Case II : If $\tau_{\mathbf{1}}=\boldsymbol{\tau}_{\mathbf{2}} \neq \boldsymbol{\tau}_{\mathbf{3}}$ are the principal stresses, then the direction $\stackrel{3}{v_{i}}$ corresponding to principal stress $\tau_{3}$ is a principal direction and any two mutually perpendicular lines in a plane with normal $\stackrel{3}{v_{i}}$ may be chosen as the other two principal direction of stress.

Case III : If $\tau_{1}=\tau_{2}=\tau_{3}$, then any set of orthogonal axes through $\mathbf{P}^{0}$ may be taken as the principal axes.

Remark: Thus, for a symmetric real stress tensor $\tau_{\mathrm{ij}}$, there are three principal stresses which are real and a set of three mutually orthogonal principal directions.

If the reference axes $x_{1}, x_{2}, x_{3}$ are chosen to coincide with the principal axes, then the matrix of stress components becomes

$$
\left[\tau_{\mathbf{i j}}\right]=\left[\begin{array}{ccc}
\tau_{1} & 0 & 0  \tag{5}\\
0 & \tau_{2} & 0 \\
0 & 0 & \tau_{3}
\end{array}\right]
$$

Invariants of the stress - tensor :
Equation (3) can be written as

$$
\begin{equation*}
\left(\tau-\tau_{1}\right)\left(\tau-\tau_{2}\right)\left(\tau-\tau_{3}\right)=0, \tag{6}
\end{equation*}
$$

and we find

$$
\left.\begin{array}{l}
\theta_{1}=\tau_{1}+\tau_{2}+\tau_{3} \\
\theta_{2}=\tau_{1} \tau_{2}+\tau_{2} \tau_{3}+\tau_{3} \tau_{1}
\end{array}\right\}
$$

$$
\begin{equation*}
\theta_{3}=\tau_{1} \tau_{2} \tau_{3} \tag{7}
\end{equation*}
$$

Since the principal stress $\tau_{1}, \tau_{2}, \tau_{3}$ characterize the physical state of stress at point, they are independent of any coordinates of reference.

Hence, coefficients $\theta_{1}, \theta_{2}, \theta_{3}$ of equation (3) are invariant w.r.t. the coordinate transformation. Thus $\theta_{1}, \theta_{2}, \theta_{3}$ are the three scalar invariants of the stress tensor $\tau_{\mathrm{ij}}$.

These scalar invariants are called the fundamental stress invariants.
Components of stress $\tau_{\mathrm{ij}}$ in terms of $\tau_{\alpha}{ }^{\prime} \mathrm{s}$
Let $X_{\alpha}$ be the principal axes. The transformation law for axes is

$$
\mathbf{X}_{\alpha}=\mathbf{a}_{\mathbf{i} \alpha} \mathbf{x}_{\mathbf{i}}
$$

or

$$
\begin{equation*}
\mathbf{x}_{\mathbf{i}}=\mathbf{a}_{\mathbf{i} \alpha} \mathbf{X}_{\alpha} \tag{8}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathbf{a}_{i \alpha}=\cos \left(\mathrm{x}_{\mathrm{i}}, X_{\alpha}\right) . \tag{9}
\end{equation*}
$$

The stress - matrix relative to axes $X_{\alpha}$ is

$$
\begin{equation*}
\tau_{\alpha \beta}^{\prime}=\operatorname{diag}\left(\tau_{1}, \tau_{2}, \tau_{3}\right) \tag{10}
\end{equation*}
$$

Let $\tau_{\mathrm{ij}}$ be the stress - matrix relative to $\mathrm{x}_{\mathrm{i}}$ - axis. Then transformation rule for second order tensor is

$$
\begin{aligned}
\tau_{\mathbf{i j}} & =\mathbf{a}_{\mathbf{i} \alpha} \mathbf{a}_{\mathbf{i \beta}} \tau_{\alpha \beta}^{\prime} \\
& =\sum_{\alpha=1}^{3} a_{i \alpha}\left(\mathbf{a}_{\mathbf{j} \alpha} \tau_{\alpha}\right) .
\end{aligned}
$$

This gives

$$
\begin{equation*}
\tau_{\mathbf{i j}}=\sum_{\alpha=1}^{3} a_{i \alpha} \mathbf{a}_{\mathbf{j} \alpha} \tau_{\alpha} \tag{11}
\end{equation*}
$$

Definition (Principal axes of Stress)
A system of coordinate axes chosen along the principal directions of stress is referred to as principal axes of stress.

Question: Show that, as the orientation of a surface element at a point $P$ varies, the normal stress on the surface element assumes an extreme value
when the element is a principal plane of stress at $P$ and that this extremum value is a principal stress.

Solution: Let $\tau_{\mathrm{ij}}$ be the stress tensor at the point P . Let $\tau$ be the normal stress on a surface element at $P$ having normal in the direction of unit vector $\hat{v}=v_{i}$. Then, we know that

$$
\begin{equation*}
\tau=\tau_{\mathrm{ij}} v_{\mathrm{i}} v_{\mathrm{j}} \tag{1}
\end{equation*}
$$

we have to find $\hat{v}=\boldsymbol{v}_{\mathbf{i}}$ for which $\tau$ is an extremum. Since $\hat{v}=\boldsymbol{v}_{\mathrm{i}}$ is a unit vector, we have the restriction

$$
\begin{equation*}
v_{k} v_{k}-1=0 \tag{2}
\end{equation*}
$$

We use the method of lagrange's multiplier to find the extremum values of $\tau$. The extreme values are given by

$$
\begin{equation*}
\frac{\partial}{\partial v_{i}}\left\{\tau_{\mathrm{ij}} v_{\mathrm{i}} v_{\mathrm{j}}-\lambda\left(v_{\mathrm{k}} v_{\mathrm{k}}-\mathbf{1}\right)\right\}=\mathbf{0} \tag{3}
\end{equation*}
$$

where $\lambda$ is a Lagrange's multiplier. From (3), we find

$$
\begin{array}{ll} 
& \tau_{\mathrm{ij}}\left\{v_{\mathrm{j}}+\delta_{\mathrm{ij}} v_{\mathrm{i}}\right\}-\lambda\left\{2 v_{\mathrm{k}} \delta_{\mathrm{ik}}\right\}=0 \\
\Rightarrow & 2 \tau_{\mathrm{ij}} v_{\mathrm{j}}-2 \lambda v_{\mathrm{i}}=0 \\
\Rightarrow & \tau_{\mathrm{ij}} v_{\mathrm{j}}-\lambda \delta_{\mathrm{ij}} v_{\mathrm{j}}=0 \\
\Rightarrow & \left(\tau_{\mathrm{ij}}-\lambda \delta_{\mathrm{ij}}\right) v_{\mathrm{j}}=0 . \tag{4}
\end{array}
$$

These conditions are satisfied iff $\hat{v}=\boldsymbol{v}_{\mathbf{j}}$ is a principal direction of stress and $\tau=\lambda$ is the corresponding principal stress.

Thus, $\tau$ assumes an extreme value on a principal plane of stress and a principal stress is an extreme value of $\tau$ given by (1).

### 2.11 MAXIMUM NORMAL AND SHEAR STRESSES

Let the co-ordinate axes at a point $P^{0}$ be taken along the principle directions of stress. Let $\tau_{1}, \tau_{2}, \tau_{3}$ be the principal stresses as $\mathbf{P}^{0}$. Then

$$
\begin{aligned}
& \tau_{11}=\tau_{1}, \tau_{22}=\tau_{2}, \tau_{33}=\tau_{3}, \\
& \tau_{12}=\tau_{23}=\tau_{31}=0 .
\end{aligned}
$$

Let $\stackrel{v}{T}$ be the stress vector on a planar element at $\mathbf{P}^{\mathbf{0}}$ having the normal $\hat{v}$ $=v_{\mathrm{i}}$. Let N be the normal stress and S be the shearing stress. Then

$$
\begin{equation*}
|\stackrel{v}{T}|^{2}=\mathbf{N}^{2}+\mathbf{S}^{2} \tag{1}
\end{equation*}
$$

The relation

$$
\stackrel{v}{T_{i}}=\tau_{\mathrm{ij}} \mathbf{v}_{\mathbf{j}}
$$

gives

$$
\begin{equation*}
\stackrel{v}{T_{1}}=\tau_{\mathbf{1}} v_{\mathbf{1}}, \stackrel{v}{T_{2}}=\tau_{\mathbf{2}} \boldsymbol{v}_{\mathbf{2}}=\stackrel{v}{T_{3}}=\tau_{\mathbf{3}} \boldsymbol{v}_{\mathbf{3}}, \tag{1a}
\end{equation*}
$$

so that

$$
\begin{equation*}
\mathbf{N}=\stackrel{v}{T} \cdot \hat{v}=\stackrel{v}{T}_{i} v_{\mathrm{i}}=\tau_{1} v_{1}^{2}+\tau_{2} v_{2}^{2}+\tau_{3}^{2} . \tag{1b}
\end{equation*}
$$

N is a function of three variables $v_{1}, v_{2}, v_{3}$ connected by the relation

$$
\begin{equation*}
v_{\mathrm{k}} v_{\mathrm{k}}-1=0 \tag{2}
\end{equation*}
$$

From (1) \& (2), we write

$$
\begin{align*}
\mathrm{N} & =\tau_{1}\left(1-v_{2}^{2}-v_{3}^{2}\right)+\tau_{2} v_{2}^{2}+\tau_{3} v_{3}^{2} \\
& =\tau_{1}+\left(\tau_{2}-\tau_{1}\right) v_{2}^{2}+\left(\tau_{3}-\tau_{1}\right) v_{3}^{2} \tag{3}
\end{align*}
$$

The extreme value of $\mathbf{N}$ are given by

$$
\frac{\partial N}{\partial v_{2}}=0 \quad, \quad \frac{\partial N}{\partial v_{3}}=\mathbf{0}
$$

which yield

$$
v_{2}=0, v_{3}=0 \text { for } \tau_{2} \neq \tau_{1} \& \tau_{3} \neq \tau_{1} .
$$

Hence

$$
v_{1}= \pm 1, v_{2}=v_{3}=0 \& \mathbf{N}=\tau_{1} .
$$

Similarly , we can find other two directions

$$
\begin{aligned}
& v_{1}=0, v_{2}= \pm 1, v_{1}=0, \mathrm{~N}=\tau_{2} \\
& v_{1}=0, v_{2}=0, v_{3}= \pm 1, \mathrm{~N}=\tau_{3}
\end{aligned}
$$

Thus, we find that the extreme values of the Normal stress $\mathbf{N}$ are along the principal directions of stress and the extreme values are themselves principal stresses. So , the absolute maximum normal stress is the maximum of the set $\left\{\tau_{1}, \tau_{2}, \tau_{3}\right\}$. Along the principal directions, the shearing stress is zero (i.e. , the minimum).

Now $\quad S^{2}=\left(\tau_{1}{ }^{2} v_{1}{ }^{2}+\tau_{2}{ }^{2} v_{2}{ }^{2}+\tau_{3}{ }^{2} v_{3}{ }^{2}\right)-\left(\tau_{1} v_{1}{ }^{2}+\tau_{2} v_{2}{ }^{2}+\tau_{3} v_{3}{ }^{2}\right)^{2}$
To determine the directions associated with the maximum values of $|S|$. We maximize the function $\mathrm{S}\left(v_{1}, v_{2}, v_{3}\right)$ in (3) subject to the constraining relation $v_{i} v_{i}=1$.

For this , we use the method of Lagrange multipliers to find the free extremum of the functions

$$
\begin{equation*}
\mathrm{F}\left(v_{1}, v_{2}, v_{3}\right)=\mathrm{S}^{2}-\lambda\left(v_{\mathrm{i}} v_{\mathrm{i}}-1\right) \tag{4}
\end{equation*}
$$

For extreme values of $F$. We must have

$$
\begin{equation*}
\frac{\partial F}{\partial v_{1}}=\frac{\partial F}{\partial v_{2}}=\frac{\partial F}{\partial v_{3}}=\mathbf{0} . \tag{5}
\end{equation*}
$$

The equation $\frac{\partial F}{\partial v_{1}}=\mathbf{0}$ gives
or

$$
\begin{align*}
& 2 \tau_{1}^{2} v_{1}-4 \tau_{1} v_{1}\left(\tau_{1} v_{1}^{2}+\tau_{2} v_{2}^{2}+\tau_{3} v_{3}^{2}\right)-2 \lambda v_{1}=0 \\
& \lambda=\tau_{1}^{2}-2 \tau_{1}\left(\tau_{1} v_{1}^{2}+\tau_{2} v_{2}^{2}+\tau_{3} v_{3}^{2}\right) \tag{6}
\end{align*}
$$

Similarly from other equations, we obtain

$$
\begin{align*}
& \lambda=\tau_{2}^{2}-2 \tau_{2}\left(\tau_{1} v_{1}^{2}+\tau_{2} v_{2}^{2}+\tau_{3} v_{3}^{2}\right),  \tag{7}\\
& \lambda=\tau_{3}^{2}-2 \tau_{3}\left(\tau_{1} v_{1}^{2}+\tau_{2} v_{2}^{2}+\tau_{3} v_{3}^{2}\right) . \tag{8}
\end{align*}
$$

Equations (6) \& (7) yield

$$
\tau_{2}^{2}-\tau_{1}^{2}=2\left(\tau_{2}-\tau_{1}\right)\left(\tau_{1} v_{1}^{2}+\tau_{2} v_{2}^{2}+\tau_{3} v_{3}^{2}\right)
$$

For $\tau_{1} \neq \tau_{2}$, this leads to
or

$$
\begin{aligned}
& \tau_{2}+\tau_{1}=2\left(\tau_{1} v_{1}^{2}+\tau_{2} v_{2}^{2}+\tau_{3} v_{3}^{2}\right) \\
& \left(2 v_{1}^{2}-1\right) \tau_{1}+\left(2 v_{2}^{2}-1\right) \tau_{2}+2 v_{3}^{2} \tau_{3}=0 .
\end{aligned}
$$

This relation is identically satisfied if

$$
\begin{equation*}
v_{1}= \pm \frac{1}{\sqrt{2}}, v_{2}=\frac{1}{\sqrt{2}}, v_{3}=0 \tag{9}
\end{equation*}
$$

From equations (1b), (3a) and (9), the corresponding maximum value of | $S$ |is
and

$$
\left.\begin{array}{l}
|\mathbf{S}|_{\max }=\frac{1}{2}\left|\tau_{2}-\tau_{1}\right| \\
|\mathbf{N}|=\frac{1}{2}\left|\tau_{1}+\tau_{2}\right| \tag{10}
\end{array}\right\}
$$

For this direction

$$
\begin{aligned}
\mathbf{S}^{2} & =\frac{1}{2}\left(\tau_{1}{ }^{2}+\tau_{2}^{2}\right)-\left(\frac{\tau_{1}+\tau_{2}}{2}\right)^{2} \\
& =\frac{1}{4}\left[2 \tau_{1}^{2}+2 \tau_{2}^{2}-\left(\tau_{1}^{2}+\tau_{2}^{2}+2 \tau_{1} \tau_{2}\right)\right] \\
& =\frac{1}{4}\left(\tau_{1}^{2}+\tau_{2}^{2}-2 \tau_{1} \tau_{2}\right) .
\end{aligned}
$$

This implies

$$
|S|_{\max }=\frac{1}{2}\left|\tau_{1}-\tau_{2}\right| .
$$

Similarly , for the directions

$$
v_{1}= \pm \frac{1}{\sqrt{2}}, v_{2}=0, v_{3}= \pm \frac{1}{\sqrt{2}}
$$

we have

$$
\begin{array}{r}
|\mathbf{S}|_{\max }=\frac{1}{2}\left|\tau_{3}-\tau_{1}\right| \\
|\mathbf{N}|=\frac{1}{2}\left|\tau_{3}+\tau_{1}\right|
\end{array}
$$

Also, for the direction

$$
v_{1}=0, v_{2}= \pm \frac{1}{\sqrt{2}}, v_{3}= \pm \frac{1}{\sqrt{2}}
$$

the corresponding values of $|S|_{\text {max }}$ and $|N|$ are , respectively ,

$$
\frac{1}{2}\left|\tau_{2}-\tau_{3}\right| \text { and } \frac{1}{2}\left|\tau_{2}+\tau_{3}\right| .
$$

These results can recorded in the following table

| $\boldsymbol{v}_{1}$ | $\boldsymbol{V}_{2}$ | $\boldsymbol{v}_{3}$ | $\|\mathbf{S}\|_{\max / \min }$ | $\|\mathbf{N}\|$ |
| :---: | :---: | :---: | :---: | :---: |
| $\mathbf{0}$ | $\mathbf{0}$ | $\pm \mathbf{1}$ | $\mathbf{M i n} \mathbf{S}=\mathbf{0}$ | $\left\|\tau_{3}\right\|=$ Max. |
| $\mathbf{0}$ | $\pm \mathbf{1}$ | $\mathbf{0}$ | $\mathbf{0}$ (Min.) | $\left\|\tau_{2}\right\|=$ Max. |
| $\pm \mathbf{1}$ | $\mathbf{0}$ | $\mathbf{0}$ | $\mathbf{0}$ (Min.) | $\left\|\tau_{1}\right\|=$ Max. |
| $\mathbf{0}$ | $\pm \frac{1}{\sqrt{2}}$ | $\pm \frac{1}{\sqrt{2}}$ | $\frac{1}{2}\left\|\tau_{2}-\tau_{3}\right\|$ Max. | $\frac{1}{2}\left\|\tau_{2}+\tau_{3}\right\|$ |
| Min. |  |  |  |  |

If $\tau_{1}>\tau_{2}>\tau_{3}$, then $\tau_{1}$ is the absolute maximum values of $\mathbf{N}$ and $\tau_{3}$ is its minimum value, and the maximum value of $|S|$ is

$$
|\mathbf{S}|_{\max }=\frac{1}{2}\left(\tau_{3}-\tau_{1}\right)
$$

and the maximum shearing stress acts on the surface element containing the $x_{2}$ principal axis and bisecting the angle between the $x_{1}-$ and $x_{3}$-axes. Hence the following theorem is proved.

Theorem : Show that the maximum shearing stress is equal to one - half the difference between the greatest and least normal stress and acts on the plane that bisects the angle between the directions of the largest and smallest principal stresses.
2.12 MOHR'S CIRCLE
(GEOMETRICAL PROOF OF THE THEOREM AS PROPOSED BY O.(OTTO) MOHR(1882))

## We know that

$$
\begin{array}{ll} 
& \mathrm{N}=\tau_{1} v_{1}^{2}+\tau_{2} v_{2}^{2}+\tau_{3} v_{3}^{2} \\
\text { and } & \mathrm{S}^{2}+\mathrm{N}^{2}=\tau_{1}^{2} v_{1}^{2}+\tau_{2}^{2} v_{2}^{2}+\tau_{3}^{2} v_{3}^{2} . \\
\text { Also } & v_{1}^{2}+v_{2}^{2}+v_{3}^{2}=1 .
\end{array}
$$

Solving equations (1) - (3), by Cramer's rule, for $v_{1}{ }^{2}, v_{2}{ }^{2}, v_{3}{ }^{2}$; we find

$$
\begin{align*}
& v_{1}^{2}=\frac{S^{2}+\left(N-\tau_{2}\right)\left(N-\tau_{3}\right)}{\left(\tau_{1}-\tau_{2}\right)\left(\tau_{1}-\tau_{3}\right)},  \tag{4}\\
& v_{2}^{2}=\frac{S^{2}+\left(N-\tau_{3}\right)\left(N-\tau_{1}\right)}{\left(\tau_{2}-\tau_{1}\right)\left(\tau_{2}-\tau_{3}\right)},  \tag{5}\\
& v_{3}^{2}=\frac{S^{2}+\left(N-\tau_{1}\right)\left(N-\tau_{2}\right)}{\left(\tau_{3}-\tau_{1}\right)\left(\tau_{3}-\tau_{2}\right)}, \tag{6}
\end{align*}
$$

Assume that $\tau_{1}>\tau_{2}>\tau_{3}$ so that $\tau_{1}-\tau_{2}>0$ and $\tau_{1}-\tau_{3}>0$. Since $v_{1}{ }^{2}$ is non negative. We conclude from equation (4) that
or

$$
\mathbf{S}^{2}+\left(\mathbf{N}-\tau_{2}\right)\left(\mathbf{N}-\tau_{3}\right) \geq \mathbf{0} .
$$

$$
\mathbf{S}^{2}+\mathbf{N}^{2}-\mathbf{N}\left(\tau_{2}+\tau_{3}\right)+\tau_{2} \tau_{3} \geq \mathbf{0}
$$

or

$$
\begin{equation*}
\mathbf{S}^{2}+\left(N-\frac{\tau_{2}+\tau_{3}}{2}\right)^{2} \geq\left(\frac{\tau_{2}-\tau_{3}}{2}\right)^{2} \tag{7}
\end{equation*}
$$

This represents a region outside the circle

$$
\mathbf{S}^{2}+\left(N-\frac{\tau_{2}+\tau_{3}}{2}\right)^{2}=\left(\frac{\tau_{2}-\tau_{3}}{2}\right)^{2}
$$

in the ( $\mathrm{N}, \mathrm{S}$ ) plane.


## Fig. (2.10) Mohr’s Circles

This circle, say $C_{1}$, has centre $\left(\frac{\tau_{2}+\tau_{3}}{2}, 0\right)$ and radius $\frac{\tau_{2}-\tau_{3}}{2}$ in the cartesian SN -plane with the values of $\mathbf{N}$ as abscissas and those of S as ordinates.

Since $\tau_{2}-\tau_{3}>0$ and $\tau_{2}-\tau_{1}<0$, we conclude from (5) that

$$
\begin{equation*}
\mathbf{S}^{2}+\left(\mathbf{N}-\tau_{3}\right)\left(\mathbf{N}-\tau_{1}\right) \leq \mathbf{0} . \tag{8}
\end{equation*}
$$

Thus, the region defined by (8) is a closed region, interior to the circle $\mathbf{c}_{\mathbf{2}}$, whose equation is

$$
\begin{equation*}
\mathbf{S}^{2}+\left(\mathbf{N}-\tau_{3}\right)\left(\mathbf{N}-\tau_{1}\right)=\mathbf{0} . \tag{8a}
\end{equation*}
$$

The circle $C_{2}$ passed through the points $\left(\tau_{3}, 0\right),\left(\tau_{1}, 0\right)$ and have centre on the N - axis.

Finally, equation (6) yields

$$
\begin{equation*}
\mathbf{S}^{2}+\left(\mathbf{N}-\tau_{1}\right)\left(\mathbf{N}-\tau_{2}\right) \geq \mathbf{0}, \tag{9}
\end{equation*}
$$

since

$$
\tau_{3}-\tau_{1}<0 \text { and } \tau_{3}-\tau_{2}<0 .
$$

The region defined by (9) is exterior to the circle $c_{3}$, with centre on the N -axis and passing through the points $\left(\tau_{1}, 0\right),\left(\tau_{2}, 0\right)$.

It follows from inequalities (7) to (9) that the admissible values of $S$ and $\mathbf{N}$ lie in the shaded region bounded by the circles as shown in the figure.

From figure, it is clear that the maximum value of shearing stress $S$ is represented by the greatest ordinate $\mathrm{O}^{1} \mathrm{Q}$ of the circle $\mathrm{C}_{2}$.

Hence

$$
\begin{equation*}
\mathbf{S}_{\max }=\frac{\tau_{1}-\tau_{3}}{2} \tag{10a}
\end{equation*}
$$

The value of N , corresponding to $\mathrm{S}_{\text {max }}$ is $\mathrm{OO}^{\prime}$ where

$$
\begin{equation*}
\mathbf{O O}^{\prime}=\tau_{3}+\frac{\tau_{1}-\tau_{3}}{2}=\frac{\tau_{1}+\tau_{3}}{2} \tag{10b}
\end{equation*}
$$

Putting the values of $S \& N$ from equations (10a,b) into equations (4) (6). We find

$$
v_{1}^{2}=v_{3}^{2}=\frac{1}{2}, v_{2}^{2}=0
$$

or

$$
\begin{equation*}
v_{1}= \pm \frac{1}{\sqrt{2}}, v_{3}= \pm \frac{1}{\sqrt{2}}, v_{2}=0 . \tag{11}
\end{equation*}
$$

Equation (11) determines the direction of the maximum shearing stress. Equation (11) shows that the maximum shearing stress acts on the plane that bisects the angle between the directions of the largest and smallest principal stresses.

### 2.13 OCTAHEDRAL STRESSES

Consider a plane which is equally inclined to the principal directions of stress. Stresses acting on such a plane are known as octahedral stresses. Assume that coordinate axes coincide with the principal directions of stress. Let $\tau_{1}, \tau_{2}, \tau_{3}$ be the principal stresses. Then the stress matrix is

$$
\left[\begin{array}{lll}
\tau_{1} & 0 & 0 \\
0 & \tau_{2} & 0 \\
0 & 0 & \tau_{3}
\end{array}\right] .
$$

A unit normal $\hat{v}=\boldsymbol{v}_{\mathbf{i}}$ to this plane is

$$
v_{1}=v_{2}=v_{3}=\frac{1}{\sqrt{3}}
$$

Then the stress vector $\underset{\sim}{v}$ on a plane element with normal $\hat{v}$ is given by

$$
\stackrel{v}{T}_{i}=\tau_{\mathrm{ij}} v_{\mathrm{j}} .
$$

This gives

$$
\stackrel{v}{T_{1}}=\tau_{1} \boldsymbol{v}_{\mathbf{1}}, \stackrel{v}{T_{2}}=\tau_{\mathbf{2}} \boldsymbol{v}_{\mathbf{2}}, \stackrel{v}{T_{3}}=\tau_{3} \boldsymbol{v}_{\mathbf{3}} .
$$

Let $\mathbf{N}$ be the normal stress and $S$ be the shear stress. Then

$$
\mathbf{N}=\underset{\sim}{\underset{\sim}{v}} \cdot \hat{v}=\tau_{1} v_{1}^{2}+\tau_{2} v_{2}^{2}+\tau_{3} v_{3}^{2}=\frac{1}{3}\left(\tau_{1}+\tau_{2}+\tau_{3}\right),
$$

and

$$
\begin{aligned}
\mathbf{S}^{2} & =|\underset{\sim}{T}|^{2}-\mathbf{N}^{2} \\
& =\left(\tau_{1}{ }^{2} v_{1}{ }^{2}+\tau_{2}{ }^{2} v_{2}{ }^{2}+\tau_{3}{ }^{2} v_{3}{ }^{2}\right)-\frac{1}{9}\left(\tau_{1}+\tau_{2}+\tau_{3}\right)^{2} \\
& =\frac{1}{3}\left(\tau_{1}{ }^{2}+\tau_{2}{ }^{2}+\tau_{3}{ }^{2}\right)-\frac{1}{9}\left(\tau_{1}+\tau_{2}+\tau_{3}\right)^{2} \\
& =\frac{1}{9}\left[3\left(\tau_{1}{ }^{2}+\tau_{2}{ }^{2}+\tau_{3}{ }^{2}\right)-\left(\tau_{1}{ }^{2}+\tau_{2}{ }^{2}+\tau_{3}{ }^{2}+2 \tau_{1} \tau_{2}+\tau_{2} \tau_{3}+\tau_{3} \tau_{1}\right)\right] \\
& =\frac{1}{9}\left[\left(\tau_{1}{ }^{2}+\tau_{2}{ }^{2}-2 \tau_{1} \tau_{2}\right)+\left(\tau_{2}{ }^{2}+\tau_{3}{ }^{2}-2 \tau_{2} \tau_{3}\right)+\left(\tau_{3}{ }^{2}+\tau_{1}{ }^{2}-2 \tau_{1} \tau_{3}\right)\right] \\
& =\frac{1}{9}\left[\left(\tau_{1}-\tau_{2}\right)^{2}+\left(\tau_{2}-\tau_{3}\right)^{2}+\left(\tau_{3}-\tau_{1}\right)^{2}\right],
\end{aligned}
$$

giving

$$
\mathbf{S}=\frac{1}{3} \sqrt{\left(\tau_{1}-\tau_{2}\right)^{2}+\left(\tau_{2}-\tau_{3}\right)^{2}+\left(\tau_{3}-\tau_{1}\right)^{2}} .
$$

Example : At a point $\mathbf{P}$, the principal stresses are $\tau_{1}=4, \tau_{2}=1, \tau_{3}=\mathbf{- 2}$. Find the stress vector, the normal stress and the shear stress on the octahedral plane at $P$.
[Hint : N = $\mathbf{1}, \mathbf{S}=\sqrt{6}, \underset{\sim}{\underset{\sim}{v}}=\frac{1}{\sqrt{3}}\left(\mathbf{4} \hat{e}_{1}+\hat{e}_{2}-2 \hat{e}_{3}\right)$.]

### 2.14. STRESS DEVIATOR TENSOR

Let $\tau_{\mathrm{ij}}$ be the stress tensor. Let

$$
\sigma_{0}=\frac{1}{3}\left(\tau_{11}+\tau_{22}+\tau_{33}\right)=\frac{1}{3}\left(\tau_{1}+\tau_{2}+\tau_{3}\right)
$$

Then the tensor

$$
\tau^{(\mathrm{d})}{ }_{\mathrm{ij}}=\tau_{\mathrm{ij}}-\sigma_{0} \delta_{\mathrm{ij}}
$$

is called the stress deviator tensor. It specifies the deviation of the state of stress from the mean stress $\sigma_{0}$.

## Chapter-3 <br> Analysis of Strain

### 3.1 INTRODUCTION

## Rigid Body

A rigid body is an ideal body such that the distance between every pair of its points remains unchanged under the action of external forces.

The possible displacements in a rigid body are translation and rotation. These displacements are called rigid displacements. In translation, each point of the rigid body moves a fixed distance in a fixed direction. In rotation about a line, every point of the body (rigid) moves in a circular path about the line in a plane perpendicular to the line.


Fig. (3.1)
In a rigid body motion, there is a uniform motion throughout the body.

## Elastic Body

A body is called elastic if it possesses the property of recovering its original shape and size when the forces causing deformation are removed.

## Continuous Body

In a continuous body, the atomistic structure of matter can be disregarded and the body is replaced by a continuous mathematical region of the space whose geometrical points are identified with material points of the body.

The mechanics of such continuous elastic bodies is called mechanics of continuum. This branch covers a vast range of problems of elasticity , hydromechanics , aerodynamics , plasticity , and electrodynamics , seismology , etc.

## Deformation of Elastic Bodies

The change in the relative position of points in a continuous is called deformation, and the body itself is then called a strained body. The study of
deformation of an elastic body is known as the analysis of strain. The deformation of the body is due to relative movements or distortions within the body.

### 3.2 TRANSFORMATION OF AN ELASTIC BODY

We consider the undeformed and deformed both positions of an elastic body. Let $\mathrm{Ox}_{1} \mathrm{x}_{2} \mathrm{x}_{3}$ be mutually orthogonal cartesian coordinates fixed in space. Let a continuous body $B$, referred to system $\mathrm{Ox}_{1} \mathrm{x}_{2} \mathrm{x}_{3}$, occupies the region R in the undeformed state. In the deformed state, the points of the body $B$ will occupy some region ,say R'.


Fig. (3.2)
Let $\mathrm{P}\left(\mathrm{x}_{1}, \mathrm{x}_{2}, \mathrm{x}_{3}\right)$ be the coordinates of a material point P of the elastic body in the initial or unstained state. In the transformation or deformed state, let this material point occupies the geometric point $\mathrm{P}^{\prime}\left(\xi_{1}, \xi_{2}, \xi_{3}\right)$. We shall be concerned only with continuous deformations of the body from region R into the region $\mathrm{R}^{\prime}$ and we assume that the deformation is given by the equation

$$
\begin{align*}
& \xi_{1}=\xi_{1}\left(x_{1}, x_{2}, x_{3}\right), \\
& \xi_{2}=\xi_{2}\left(x_{1}, x_{2}, x_{3}\right), \\
& \xi_{3}=\xi_{3}\left(x_{1}, x_{2}, x_{3}\right) . \tag{1}
\end{align*}
$$

The vector $\overline{P P^{1}}$ is called the displacement vector of the point P and is denoted by $u_{i}$.
Thus

$$
\begin{equation*}
\mathrm{u}_{\mathrm{i}}=\xi_{\mathrm{i}}-\mathrm{x}_{\mathrm{i}} \quad ; \quad \mathrm{i}=1,2,3 \tag{2}
\end{equation*}
$$

or

$$
\begin{equation*}
\xi_{i}=x_{i}+u_{i} . \quad i=1,2,3 \tag{3}
\end{equation*}
$$

Equation (1) expresses the coordinates of the points of the body in the transformed state in terms of their coordinates in the initial undeformed state. This type of description of deformation is known as the Lagrangian method of describing the transformation of a continuous medium.

Another method, known as Euler's method expresses the coordinates in the undeformed state in terms of the coordinates in the deformed state.

## The transformation (1) is invertible when

$$
\mathbf{J} \neq \mathbf{0} .
$$

## Then, we may write

$$
\begin{equation*}
\mathrm{x}_{\mathrm{i}}=\mathrm{x}_{\mathrm{i}}\left(\xi_{1}, \xi_{2}, \xi_{3}\right) \quad, \quad \mathrm{i}=1,2,3 . \tag{4}
\end{equation*}
$$

In this case, the transformation from the region $R$ into region $R^{\prime}$ is one to - one.

Each of the above description of deformation of the body has its own advantages. It is however, more convenient in the study of the mechanics of solids to use Lagrangian approach because the undeformed state of the body often possesses certain symmetries which make it convenient to use a simple system of coordinates.

A part of the transformation defined by equation (1) may represent rigid body motions
(i.e., translations and rotations) of the body as a whole. This part of the deformation leaves unchanged the length of every vector joining a pair of points within the body and is of no interest in the analysis of strain.

The remaining part of transformation (1) will be called pure deformation.
Now, we shall learn how to distinguish between pure deformation and rigid body motions when the latter are present in the transformation equations (1).
3.3 LINEAR TRANSFORMATION OR AFFINE TRANSFORMATION

Definition: The transformation

$$
\xi_{\mathrm{i}}=\xi_{\mathrm{i}}\left(\mathrm{x}_{1}, \mathrm{x}_{2}, \mathrm{x}_{3}\right)
$$

is called a linear transformation or affine transformation when the functions $\xi_{\mathrm{i}}$ are linear functions of the coordinates $\mathrm{x}_{1}, \mathrm{x}_{2}, \mathrm{x}_{3}$.

In order to distinguish between rigid motion and pure deformation, we consider the simple case in which the transformation (1) is linear.

We assume that the general form of the linear transformation (1) is of the type

$$
\left.\begin{array}{l}
\xi_{1}=\alpha_{10}+\left(\alpha_{11}+1\right) x_{1}+\alpha_{12} x_{2}+\alpha_{13} x_{3}  \tag{5}\\
\xi_{2}=\alpha_{20}+\alpha_{21} x_{1}+\left(1+\alpha_{22}\right) x_{2}+\alpha_{23} x_{3} \\
\xi_{3}=\alpha_{30}+\alpha_{31} x_{1}+\alpha_{32} x_{2}+\left(1+\alpha_{33}\right) x_{3}
\end{array}\right\}
$$

or

$$
\begin{equation*}
\xi_{\mathrm{i}}=\alpha_{\mathrm{i} 0}+\left(\alpha_{\mathrm{ij}}+\delta_{\mathrm{ij}}\right) \mathrm{x}_{\mathrm{j}}, \tag{6}
\end{equation*}
$$

where the coefficients $\alpha_{\mathrm{ij}}$ are constants and are well known.
Equation (5) can written in the matrix form as

$$
\left[\begin{array}{l}
\xi_{1}-\alpha_{10}  \tag{7}\\
\xi_{2}-\alpha_{20} \\
\xi_{3}-\alpha_{30}
\end{array}\right]=\left[\begin{array}{ccc}
1+\alpha_{11} & \alpha_{12} & \alpha_{13} \\
\alpha_{21} & 1+\alpha_{22} & \alpha_{23} \\
\alpha_{31} & \alpha_{32} & 1+\alpha_{33}
\end{array}\right]\left[\begin{array}{l}
x_{1} \\
x_{2} \\
x_{3}
\end{array}\right],
$$

or

$$
\left[\begin{array}{l}
u_{1}-\alpha_{10}  \tag{8}\\
u_{2}-\alpha_{20} \\
u_{3}-\alpha_{30}
\end{array}\right]=\left[\begin{array}{lll}
\alpha_{11} & \alpha_{12} & \alpha_{13} \\
\alpha_{21} & \alpha_{22} & \alpha_{23} \\
\alpha_{31} & \alpha_{32} & \alpha_{33}
\end{array}\right]\left[\begin{array}{c}
x_{1} \\
x_{2} \\
x_{3}
\end{array}\right] .
$$

We can look upon the matrix $\left(\alpha_{\mathrm{ij}}+\delta_{\mathrm{ij}}\right)$ as an operator acting on the vector $\bar{x}=x_{i}$ to give the
vector $\alpha_{i 0}$.
If the matrix $\left(\alpha_{\mathrm{ij}}+\delta_{\mathrm{ij}}\right)$ is non - singular , then we obtain

$$
\left[\begin{array}{c}
x_{1}  \tag{9}\\
x_{2} \\
x_{3}
\end{array}\right]=\left(\alpha_{\mathrm{ij}}+\delta_{\mathrm{ij}}\right)^{-1}\left[\begin{array}{l}
\xi_{1}-\alpha_{10} \\
\xi_{2}-\alpha_{20} \\
\xi_{3}-\alpha_{30}
\end{array}\right],
$$

which is also linear as inverse of a linear transformation is linear.
Infact, matrix algebra was developed basically to express linear transformations in a concise and lucid manner.

Result (1) : Sum of two linear transformations is a linear transformation.

Result (2) : Product of two linear transformations is a linear transformation which is not commutative.

Result (3) : Under a linear transformation, a plane is transformed into a plane.
Proof of (3) : Let

$$
a x_{1}+b x_{2}+c x_{3}+d=0,
$$

be an equation of a plane in the undeformed state. Let

$$
\left[\begin{array}{l}
\xi_{1} \\
\xi_{2} \\
\xi_{3}
\end{array}\right]=\left[\begin{array}{lll}
a_{1} & b_{1} & c_{1} \\
a_{2} & b_{2} & c_{2} \\
a_{3} & b_{3} & c_{3}
\end{array}\right]\left[\begin{array}{l}
x_{1} \\
x_{2} \\
x_{3}
\end{array}\right]
$$

be the linear transformation of points. Let its inverse be

$$
\left[\begin{array}{l}
x_{1} \\
x_{2} \\
x_{3}
\end{array}\right]=\left[\begin{array}{lll}
A_{1} & B_{1} & C_{1} \\
A_{2} & B_{2} & C_{2} \\
A_{3} & B_{3} & C_{3}
\end{array}\right]\left[\begin{array}{l}
\xi_{1} \\
\xi_{2} \\
\xi_{3}
\end{array}\right] .
$$

Then the equation of the plane is transformed to

$$
\begin{gathered}
\mathrm{a}\left(\mathrm{~A}_{1} \xi_{1}+\mathrm{B}_{1} \xi_{2}+\mathrm{C}_{1} \xi_{3}\right)+\mathrm{b}\left(\mathrm{~A}_{2} \xi_{1}+\mathrm{B}_{2} \xi_{2}+\mathrm{C}_{2} \xi_{3}\right) \\
+\mathrm{c}\left(\mathrm{~A}_{3} \xi_{1}+\mathrm{B}_{3} \xi_{2}+\mathrm{C}_{3} \xi_{3}\right)+\mathrm{d}=0
\end{gathered}
$$

or

$$
\alpha_{1} \xi_{1}+\beta_{1} \xi_{2}+\gamma_{1} \xi_{3}+d=0
$$

which is again an equation of a plane in terms of new coordinates $\left(\xi_{1}, \xi_{2}, \xi_{3}\right)$.
Hence the result.
Result (4) : A linear transformation carries line segments into line segments.
Thus, it is the linear transformation that allows us to assume that a line segment is transformed to a line segment and not to a curve.

### 3.4 SMALL/ INFINITESIMAL LINEAR DEFORMATIONS

## Definition (Small / Infinitesimal Deformations)

A linear transformation of the type

$$
\xi_{\mathrm{i}}=\alpha_{\mathrm{i} 0}+\left(\alpha_{\mathrm{ij}}+\delta_{\mathrm{ij}}\right) \mathrm{x}_{\mathrm{j}}
$$

is said to be a small linear transformation of the coefficients $\alpha_{\mathrm{ij}}$ are so small that their products can be neglected in comparison with the linear terms.

Result (1): The product of two small linear transformations is a small linear transformation which is COMMUTATIVE and the product transformation is obtained by superposition of the original transformations and the result is independent of the order in which the transformations are performed.

Result (2) : In the study of finite deformations (as compared to the infinitesimal affine deformation), the principle of superposition of effects and the independence of the order of transformations are no longer valid.

Note : If a body is subjected to large linear transformations, a straight line element seldom remains straight. A curved element is more likely to result. The linear transformation then expresses the transformation of element $\mathrm{P}_{1} \mathrm{P}_{2}$ to the tangent $\mathrm{P}_{1}{ }^{\prime} \mathrm{T}_{1}{ }^{\prime}$ to the curve at $\mathrm{P}_{1}{ }^{\prime}$ for the curve itself.


Fig. (3.3)
For this reason, a linear transformation is sometimes called linear tangent transformation.
It is obvious that the smaller the element $\mathrm{P}_{1} \mathrm{P}_{2}$, the better approximation of $\mathrm{P}_{1}{ }^{\prime} \mathrm{P}_{2}{ }^{\prime}$ by its tangent $\mathrm{P}_{1}{ }^{\prime} \mathrm{T}_{1}{ }^{\prime}$.

### 3.5 HOMOGENEOUS DEFORMATION

Suppose that a body B, occupying the region $R$ in the undeformed state, is transformed to the region $R^{\prime}$ under the linear transformation.

$$
\begin{equation*}
\xi_{\mathrm{i}}=\alpha_{\mathrm{i} 0}+\left(\alpha_{\mathrm{ij}}+\delta_{\mathrm{ij}}\right) \mathrm{x}_{\mathrm{j}} \tag{1}
\end{equation*}
$$

referred to orthogonal cartesian system $\mathrm{Ox}_{1} \mathrm{x}_{2} \mathrm{x}_{3}$. Let $\hat{e}_{1}, \hat{e}_{2}, \hat{e}_{3}$ be the unit base vectors directed along the coordinate axes $\mathrm{x}_{1}, \mathrm{x}_{2}, \mathrm{x}_{3}$ (Fig. 3.4)


Fig. (3.5)
Let $\mathrm{P}_{1}\left(\begin{array}{cc}\mathrm{x}_{1} & 1 \\ \mathrm{x}_{2}, & \mathrm{x}_{3}\end{array}\right)$ and $\mathrm{P}_{2}\left(\begin{array}{c}2 \\ \mathrm{x}_{1},\end{array} \mathrm{x}_{2},{ }_{2}^{2}, \mathrm{x}_{3}\right)$ be two points of the elastic body in the initial state. Let the positions of these points in the deformed state, due to linear transformation (1), be $\mathrm{P}_{1}{ }^{\prime}\left(\xi_{1}, \xi_{2}, \xi_{3}\right)$ and $\mathrm{P}_{2}{ }^{\prime}\left(\xi_{1}^{2}, \xi_{2}^{2}, \xi_{3}^{2}\right)$.

Since the transformation (1) is linear, so the line segment $\overline{P_{1} P_{2}}$ is transformed into a line segment $\overline{P_{1}{ }^{\prime} P_{2}{ }^{\prime}}$.

Let the vector $\overline{P_{1} P_{2}}$ has components $\mathrm{A}_{\mathrm{i}}$ and vector $\overline{P_{1}{ }^{\prime} P_{2}{ }^{\prime}}$ has components $\mathrm{A}_{\mathrm{i}}{ }^{\prime}$. Then

$$
\begin{equation*}
\overline{P_{1} P_{2}}=\mathrm{A}_{\mathrm{i}} \hat{e}_{i}, \quad \mathrm{~A}_{\mathrm{i}}=\stackrel{2}{x_{i}}-\stackrel{1}{x_{i}}, \tag{2}
\end{equation*}
$$

and

$$
\begin{equation*}
\overline{P_{1}^{\prime} P_{2}^{\prime}}=\mathrm{A}_{\mathrm{i}}{ }^{\prime} \hat{e}_{i}, \quad \mathrm{~A}_{\mathrm{i}}^{\prime}=\stackrel{2}{\xi}_{i}-\xi_{i} \tag{3}
\end{equation*}
$$

Let

$$
\begin{equation*}
\delta \mathrm{A}_{\mathrm{i}}=\mathrm{A}_{\mathrm{i}}{ }^{\prime}-\mathrm{A}_{\mathrm{i}}, \tag{4}
\end{equation*}
$$

be the change in vector $A_{i}$.
The vectors $A_{i}$ and $A_{i}{ }^{\prime}$, in general, differ in direction and magnitude.
From equations (1), (2) and (3), we write

$$
\begin{aligned}
\mathrm{A}_{\mathrm{i}}^{\prime} & =\stackrel{2}{\xi_{i}}-\stackrel{1}{\xi_{i}} \\
& =\left[\alpha_{\mathrm{i} 0}+\left(\alpha_{\mathrm{ij}}+\delta_{\mathrm{ij}}\right)\right. \\
\left.\stackrel{2}{x}_{j}\right]-\left[\alpha_{\mathrm{i} 0}+\left(\alpha_{\mathrm{ij}}+\delta_{\mathrm{ij}}\right)\right. & \left.\stackrel{1}{x_{j}}\right]
\end{aligned}
$$

$$
\begin{aligned}
& =\left(\stackrel{2}{x}_{x_{i}}-x_{i}^{1}\right)+\alpha_{\mathrm{ij}}\left(\stackrel{2}{x}_{j}-\stackrel{1}{x}_{j}\right) \\
& =\mathrm{A}_{\mathrm{i}}+\alpha_{\mathrm{ij}} \mathrm{~A}_{\mathrm{j}} .
\end{aligned}
$$

This implies

$$
\begin{align*}
& \mathrm{A}_{\mathrm{i}}^{\prime}-\mathrm{A}_{\mathrm{i}}=\alpha_{\mathrm{ij}} \mathrm{~A}_{\mathrm{j}} \\
& \delta \mathbf{A}_{\mathbf{i}}=\alpha_{\mathrm{ij}} \mathbf{A}_{\mathbf{j}} . \tag{5}
\end{align*}
$$

Thus, the linear transformation (1) changes the vector $\mathrm{A}_{\mathrm{i}}$ into vector $\mathrm{A}_{\mathrm{i}}{ }^{\prime}$ where

$$
\left[\begin{array}{c}
A_{1}^{\prime}  \tag{6}\\
A_{2}^{\prime} \\
A_{3}^{\prime}
\end{array}\right]=\left[\begin{array}{ccc}
1+\alpha_{11} & \alpha_{12} & \alpha_{13} \\
\alpha_{21} & 1+\alpha_{22} & \alpha_{23} \\
\alpha_{31} & \alpha_{32} & 1+\alpha_{33}
\end{array}\right]\left[\begin{array}{c}
A_{1} \\
A_{2} \\
A_{3}
\end{array}\right],
$$

or

$$
\left[\begin{array}{l}
\delta A_{1}  \tag{7}\\
\delta A_{2} \\
\delta A_{3}
\end{array}\right]=\left[\begin{array}{lll}
\alpha_{11} & \alpha_{12} & \alpha_{13} \\
\alpha_{21} & \alpha_{22} & \alpha_{23} \\
\alpha_{31} & \alpha_{32} & \alpha_{33}
\end{array}\right]\left[\begin{array}{c}
A_{1} \\
A_{2} \\
A_{3}
\end{array}\right]
$$

Thus, the linear transformation (1) or (6) or (7) are all equivalent.
From equation (6), it is clear that two vectors $A_{i}$ and $B_{i}$ whose components are equal transform into two vectors $\mathrm{A}_{\mathrm{i}}{ }^{\prime}$ and $\mathrm{B}_{\mathrm{i}}{ }^{\prime}$ whose components are again equal. Also two parallel vectors transformation into parallel vectors.

Hence,two equal and similarly oriented rectilinear polygons located in different parts of the region R will be transformed into equal and similarly oriented polygons in the transformed region $\mathrm{R}^{\prime}$ under the linear transformation (1).

Thus, the different parts of the body B, when the latter is subjected to the linear transformation (1), experience the same deformation independent of the position of the parts of the body.

For this reason , the linear deformation (1) is called a homogeneous deformation.

Theorem: Prove that the necessary and sufficient condition for an infinitesimal affine transform

$$
\xi_{\mathrm{ij}}=\alpha_{\mathrm{i} 0}+\left(\alpha_{\mathrm{ij}}+\delta_{\mathrm{ij}}\right) \mathrm{x}_{\mathrm{j}}
$$

to represent a rigid body motion is that the matrix $\alpha_{\mathrm{ij}}$ is skew - symmetric

Proof: With reference to an orthogonal cartesian system o $\mathrm{x}_{1} \mathrm{x}_{2} \mathrm{x}_{3}$ fixed in space, let the line segment $\overrightarrow{P^{o} P}$ of the body in the undeformed state be transferred to the line segment $\overline{P^{0 \prime} P^{\prime}}$ in the deformed state due to infinitesimal affine transformation

$$
\begin{equation*}
\xi_{\mathrm{i}}=\alpha_{\mathrm{i} 0}+\left(\alpha_{\mathrm{ij}}+\delta_{\mathrm{ij}}\right) \mathrm{x}_{\mathrm{j}}, \tag{1}
\end{equation*}
$$

in which $\alpha_{\mathrm{ij}}$ are known constants. Let $\mathrm{A}_{\mathrm{i}}$ be vector $\overline{P^{o} P}$ and $\mathrm{A}_{\mathrm{i}}{ }^{\prime}$ be the vector $\overline{P^{0} P^{\prime}}$


Fig. (3.6)
Then

$$
\begin{equation*}
A_{i}=x_{i}-x_{i}^{\circ} \quad, \quad A_{i}^{\prime}=\xi_{i}-\xi_{i}^{\circ} . \tag{2}
\end{equation*}
$$

Let

$$
\begin{equation*}
\delta \mathrm{A}_{\mathrm{i}}=\mathrm{A}_{\mathrm{i}}{ }^{\prime}-\mathrm{A}_{\mathrm{i}} . \tag{3}
\end{equation*}
$$

From (1) and (2), we find

$$
\begin{aligned}
\mathrm{A}_{\mathrm{i}}^{\prime} & =\xi_{\mathrm{i}}-\xi_{\mathrm{i}}^{\circ} \\
& =\left(\alpha_{\mathrm{i} 0}+\alpha_{\mathrm{ij}} x_{j}+\mathrm{x}_{\mathrm{i}}\right)-\left(\alpha_{\mathrm{i} 0}+\alpha_{\mathrm{ij}} \mathrm{x}_{\mathrm{j}}^{\circ}+\mathrm{x}_{\mathrm{i}}^{\circ}\right) \\
& =\left(\mathrm{x}_{\mathrm{i}}-\mathrm{x}_{\mathrm{i}}^{\circ}\right)+\alpha_{\mathrm{ij}}\left(\mathrm{x}_{\mathrm{j}}-\mathrm{x}_{\mathrm{j}}^{\circ}\right) \\
& =\mathrm{A}_{\mathrm{i}}+\alpha_{\mathrm{ij}} A_{\mathrm{j}} .
\end{aligned}
$$

This gives

$$
\begin{equation*}
\delta \mathrm{A}_{\mathrm{i}}=\mathrm{A}_{\mathrm{i}}{ }^{\prime}-\mathrm{A}_{\mathrm{i}}=\alpha_{\mathrm{ij}} \mathrm{~A}_{\mathrm{j}} . \tag{4}
\end{equation*}
$$

Let A denote the length of the vector. Then

$$
\begin{equation*}
\mathbf{A}=\left|\mathbf{A}_{\mathbf{i}}\right|=\sqrt{A_{i} A_{i}}=\sqrt{A_{1}^{2}+A_{2}^{2}+A_{3}^{2}} . \tag{5}
\end{equation*}
$$

Let $\delta \mathrm{A}$ denote the change in length A due to deformation. Then

$$
\begin{equation*}
\delta \mathbf{A}=\left|\mathrm{A}_{\mathrm{i}}{ }^{\prime}\right|-\left|\mathrm{A}_{\mathrm{i}}\right| . \tag{6}
\end{equation*}
$$

It is obvious hat $\delta \mathrm{A} \neq\left|\delta \mathrm{A}_{\mathrm{i}}\right|$, but

$$
\delta \mathrm{A}=\sqrt{\left(A_{i}+\delta A_{i}\right)\left(A_{i}+\delta A_{i}\right)}-\sqrt{A_{i} A_{i}} .
$$

This imply

$$
(\mathrm{A}+\delta \mathrm{A})^{2}=\left(\mathrm{A}_{\mathrm{i}}+\delta \mathrm{A}_{\mathrm{i}}\right)\left(\mathrm{A}_{\mathrm{i}}+\delta \mathrm{A}_{\mathrm{i}}\right),
$$

or

$$
\begin{equation*}
(\delta \mathrm{A})^{2}+2 \mathrm{~A} \delta \mathrm{~A}=\left(\delta \mathrm{A}_{\mathrm{i}}\right)\left(\delta \mathrm{A}_{\mathrm{i}}\right)+2 \mathrm{~A}_{\mathrm{i}}\left(\delta \mathrm{~A}_{\mathrm{i}}\right) \tag{7}
\end{equation*}
$$

Since the linear transformation (1) or (4) is small, the terms $(\delta A)^{2}$ and $\left(\delta \mathrm{A}_{\mathrm{i}}\right)$ $\left(\delta A_{i}\right)$ are to be neglected in (7). Therefore, after neglecting these terms in (7), we write

$$
2 \mathrm{~A} \delta \mathrm{~A}=2 \mathrm{~A}_{\mathrm{i}} \delta \mathrm{~A}_{\mathrm{i}},
$$

or

$$
\begin{equation*}
\mathbf{A} \delta \mathbf{A}=\mathbf{A}_{\mathbf{i}} \delta \mathbf{A}_{\mathbf{i}}=\mathrm{A}_{1} \delta \mathrm{~A}_{1}+\mathrm{A}_{2} \delta \mathrm{~A}_{2}+\mathrm{A}_{3} \delta \mathrm{~A}_{3} \tag{8}
\end{equation*}
$$

Using (4), equation (8) becomes

$$
\begin{align*}
& \quad \mathrm{A} \delta \mathrm{~A}=\mathrm{A}_{\mathrm{i}}\left(\alpha_{\mathrm{ij}} \mathrm{~A}_{\mathrm{j}}\right) \\
& =\alpha_{\mathrm{ij}} \mathrm{~A}_{\mathrm{i}} \mathrm{~A}_{\mathrm{j}} \\
& =\alpha_{11} \mathrm{~A}_{1}^{2}+\alpha_{22} \mathrm{~A}_{2}^{2}+\alpha_{33} \mathrm{~A}_{3}^{2}+\left(\alpha_{12}+\alpha_{21}\right) \mathrm{A}_{1} \mathrm{~A}_{2} \\
& \quad+\left(\alpha_{13}+\alpha_{31}\right) \mathrm{A}_{3} \mathrm{~A}_{1}+\left(\alpha_{23}+\alpha_{32}\right) \mathrm{A}_{2} \mathrm{~A}_{3} . \tag{9}
\end{align*}
$$

Case 1: Suppose that the infinitesimal linear transformation (1) represents a rigid body motion.

Then, the length of the vector $\mathrm{A}_{\mathrm{i}}$ before deformation and after deformation remains unchanged.

That is

$$
\begin{equation*}
\delta \mathrm{A}=0, \tag{10}
\end{equation*}
$$

for all vectors $\mathrm{A}_{\mathrm{i}}$.
Using (9) , we then get

$$
\alpha_{11} \mathrm{~A}_{1}^{2}+\alpha_{22} \mathrm{~A}_{2}^{2}+\alpha_{33} \mathrm{~A}_{3}^{2}+\left(\alpha_{12}+\alpha_{21}\right) \mathrm{A}_{1} \mathrm{~A}_{2}+\left(\alpha_{23}+\alpha_{32}\right) \mathrm{A}_{2} \mathrm{~A}_{3}
$$

$$
\begin{equation*}
+\left(\alpha_{13}+\alpha_{31}\right) \mathrm{A}_{1} \mathrm{~A}_{3}=0, \tag{11}
\end{equation*}
$$

for all vectors $\mathrm{A}_{\mathrm{i}}$.
This is possible only when
i.e.,

$$
\begin{align*}
& \alpha_{11}=\alpha_{22}=\alpha_{33}=0, \\
& \alpha_{12}+\alpha_{21}=\alpha_{13}+\alpha_{31}=\alpha_{23}+\alpha_{32}=0, \\
& \alpha_{\mathrm{ij}}=-\alpha_{\mathrm{ji}}, \quad \text { for all i \& } \mathrm{j} \tag{12}
\end{align*}
$$

i.e. , the matrix $\alpha_{\mathrm{ij}}$ is skew - symmetric.

Case 2: Suppose $\alpha_{\mathrm{ij}}$ is skew - symmetric. Then, equation (9) shows that

$$
\begin{equation*}
\mathrm{A} \delta \mathrm{~A}=0 \tag{13}
\end{equation*}
$$

for all vectors $\mathrm{A}_{\mathrm{i}}$. This implies

$$
\begin{equation*}
\delta \mathrm{A}=0 \tag{14}
\end{equation*}
$$

for all vectors $\mathrm{A}_{\mathrm{i}}$
This shows that the transformation (1) represents a rigid body linear small transformation.

This completes the proof of the theorem.
Remark: When the quantities $\alpha_{\mathrm{ij}}$ are skew - symmetric, then the linear infinitesimal transformation

$$
\delta A_{i}=\alpha_{i j} A_{j}
$$

equation (11) takes the form

$$
\begin{align*}
& \delta \mathrm{A}_{1}=-\alpha_{21} \mathrm{~A}_{2}+\alpha_{13} \mathrm{~A}_{3}, \\
& \delta \mathrm{~A}_{2}=\alpha_{21} \mathrm{~A}_{1}-\alpha_{32} \mathrm{~A}_{3}, \\
& \delta \mathrm{~A}_{3}=-\alpha_{13} \mathrm{~A}_{1}+\alpha_{32} \mathrm{~A}_{2} . \tag{15}
\end{align*}
$$

Let

$$
\begin{align*}
& \mathrm{w}_{1}=\alpha_{32}=-\alpha_{23}, \\
& \mathrm{w}_{2}=\alpha_{13}=-\alpha_{31}, \\
& \mathrm{w}_{3}=\alpha_{21}=-\alpha_{12} . \tag{16}
\end{align*}
$$

Then , the transformation (15) can be written as the vector product

$$
\begin{equation*}
\overline{\delta \mathrm{A}}=\overline{\mathrm{w}} \times \overline{\mathrm{A}}, \tag{17}
\end{equation*}
$$

where $\overline{\mathrm{w}}=\mathrm{w}_{\mathrm{i}}$ is the infinitesimal rotation vector. Further

$$
\begin{align*}
\delta \mathrm{A}_{\mathrm{i}} & =\mathrm{A}_{\mathrm{i}}^{\prime}-\mathrm{A}_{\mathrm{i}} \\
& =\left(\xi_{\mathrm{i}}-\xi_{\mathrm{i}}^{\circ}\right)-\left(\mathrm{x}_{\mathrm{i}}-\mathrm{x}_{\mathrm{i}}^{\circ}\right) \\
& =\left(\xi_{\mathrm{i}}-\mathrm{x}_{\mathrm{i}}\right)-\left(\xi_{\mathrm{i}}^{\circ}-\mathrm{x}_{\mathrm{i}}{ }^{\circ}\right) \\
& =\delta \mathrm{x}_{\mathrm{i}}-\delta \mathrm{x}_{\mathrm{i}}^{\circ} . \tag{18}
\end{align*}
$$

This yields

$$
\delta \mathrm{x}_{\mathrm{i}}=\delta \mathrm{x}_{\mathrm{i}}{ }^{\circ}+\delta \mathrm{A}_{\mathrm{i}},
$$

or

$$
\begin{equation*}
\delta \mathrm{x}_{\mathrm{i}}=\delta \mathrm{x}_{\mathrm{i}}{ }^{\circ}+(\overline{\mathrm{w}} \times \overline{\mathrm{A}}) . \tag{19}
\end{equation*}
$$

Here, the quantities

$$
\delta x_{i}{ }^{\circ}=\xi_{i}{ }^{\circ}-x_{i}{ }^{\circ}
$$

are the components of the displacement vector representing the translation of the point $\mathrm{P}^{\circ}$ and the remaining terms of (19) represent rotation of the body about the point $\mathrm{P}^{\circ}$.

### 3.6 PURE DEFORMATION AND COMPONENTS OF STRAIN TENSOR <br> We consider the infinitesimal linear transformation

$$
\begin{equation*}
\delta \mathrm{A}_{\mathrm{i}}=\alpha_{\mathrm{ij}} \mathrm{~A}_{\mathrm{j}} \tag{1}
\end{equation*}
$$

Let

$$
\begin{equation*}
\mathrm{w}_{\mathrm{ij}}=\frac{1}{2}\left(\alpha_{\mathrm{ij}}-\alpha_{\mathrm{ij}}\right), \tag{2}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathrm{e}_{\mathrm{ij}}=\frac{1}{2}\left(\alpha_{\mathrm{ij}}+\alpha_{\mathrm{ji}}\right) . \tag{3}
\end{equation*}
$$

Then the matrix $\mathrm{w}_{\mathrm{ij}}$ is antisymmetric while $\mathrm{e}_{\mathrm{ij}}$ is symmetric .
Moreover,

$$
\begin{equation*}
\alpha_{\mathrm{ij}}=\mathrm{e}_{\mathrm{ij}}+\mathrm{w}_{\mathrm{ij}}, \tag{4}
\end{equation*}
$$

and this decomposition of $\alpha_{\mathrm{ij}}$ as a sum of a symmetric and skew - symmetric matrices is unique.
From (1) and (4), we write

$$
\begin{equation*}
\delta \mathrm{A}_{\mathrm{i}}=\mathrm{e}_{\mathrm{ij}} \mathrm{~A}_{\mathrm{j}}+\mathrm{w}_{\mathrm{ij}} \mathrm{~A}_{\mathrm{j}} \tag{5}
\end{equation*}
$$

This shows that the transformation of the components of a vector $\mathrm{A}_{\mathrm{i}}$ given by

$$
\begin{equation*}
\delta \mathrm{A}_{\mathrm{i}}=\mathrm{w}_{\mathrm{ij}} \mathrm{~A}_{\mathrm{j}} \tag{6}
\end{equation*}
$$

represents rigid body motion with the components of rotation vector $\mathrm{w}_{\mathrm{i}}$ given by

$$
\begin{equation*}
\mathrm{w}_{1}=\mathrm{w}_{32}, \mathrm{w}_{2}=\mathrm{w}_{13}, \mathrm{w}_{3}=\mathrm{w}_{21}, \tag{7}
\end{equation*}
$$

and the transformation

$$
\begin{equation*}
\delta \mathbf{A}_{\mathbf{i}}=\mathbf{e}_{\mathbf{i j}} \mathbf{A}_{\mathbf{j}} \tag{8}
\end{equation*}
$$

with

$$
\begin{equation*}
\mathbf{e}_{\mathrm{ij}}=\mathbf{e}_{\mathrm{ji}}, \tag{9}
\end{equation*}
$$

represents a pure deformation.

## Strain Components

The symmetric coefficients, $\mathrm{e}_{\mathrm{ij}}$, in the pure deformation

$$
\delta \mathrm{A}_{\mathrm{i}}=\mathrm{e}_{\mathrm{ij}} \mathrm{~A}_{\mathrm{j}}
$$

are called the strain components.
Note (1): These components of strain characterize pure deformation of the elastic body. Since $A_{j}$ and $\delta A_{i}$ are vectors (each is a tensor of order 1), therefore, by quotient law, the strain components $\mathbf{e}_{\mathbf{i j}}$ form a tensor of order 2.

Note (2): For most materials / structures, the strains are of the order of $\mathbf{1 0}^{\mathbf{- 3}}$. Such strains certainly deserve to be called small.

Note (3) : The strain components $e_{11}, e_{22}, e_{33}$ are called normal strain components while $e_{12}, e_{13}, e_{23}, e_{21}, e_{31}, e_{32}$ are called shear strain components.

Example : For the deformation defined by the linear transformation

$$
\xi_{1}=x_{1}+x_{2}, \xi_{2}=x_{1}-2 x_{2}, \xi_{3}=x_{1}+x_{2}-x_{3},
$$

find the inverse transformation, components of rotation and strain tensor, and axis of rotation.

Solution : The given transformation is expresses as

$$
\left[\begin{array}{l}
\xi_{1}  \tag{1}\\
\xi_{2} \\
\xi_{3}
\end{array}\right]=\left[\begin{array}{rrr}
1 & 1 & 0 \\
1 & -2 & 0 \\
1 & 1 & -1
\end{array}\right]\left[\begin{array}{l}
x_{1} \\
x_{2} \\
x_{3}
\end{array}\right],
$$

and its inverse transformation is

$$
\begin{align*}
{\left[\begin{array}{l}
x_{1} \\
x_{2} \\
x_{3}
\end{array}\right] } & =\left[\begin{array}{rrr}
1 & 1 & 0 \\
1 & -2 & 0 \\
1 & 1 & -1
\end{array}\right]^{-1}\left[\begin{array}{l}
\xi_{1} \\
\xi_{2} \\
\xi_{3}
\end{array}\right] \\
& =\frac{1}{3}\left[\begin{array}{rrr}
2 & 1 & 0 \\
1 & -1 & 0 \\
1 & 0 & -1
\end{array}\right]\left[\begin{array}{l}
\xi_{1} \\
\xi_{2} \\
\xi_{3}
\end{array}\right], \tag{2}
\end{align*}
$$

giving

$$
\begin{align*}
& \mathrm{x}_{1}=\frac{1}{3}\left(2 \xi_{1}+\xi_{2}\right), \\
& \mathrm{x}_{2}=\frac{1}{3}\left(\xi_{1}-\xi_{2}\right), \\
& \mathrm{x}_{3}=\xi_{1}-\xi_{3} \tag{3}
\end{align*}
$$

Comparing (1) with

$$
\begin{equation*}
\xi_{\mathrm{i}}=\left(\alpha_{\mathrm{ij}}+\delta_{\mathrm{ij}}\right) \mathrm{x}_{\mathrm{j}} \tag{4}
\end{equation*}
$$

we find

$$
\left(\alpha_{\mathrm{ij}}\right)=\left[\begin{array}{rrr}
0 & 1 & 0  \tag{5}\\
1 & -3 & 0 \\
1 & 1 & -2
\end{array}\right]
$$

Then

$$
\mathrm{w}_{\mathrm{ij}}=\frac{1}{2}\left(\alpha_{\mathrm{ij}}-\alpha_{\mathrm{ji}}\right)=\frac{1}{2}\left[\begin{array}{rrr}
0 & 0 & -1  \tag{6}\\
0 & 0 & -1 \\
1 & 1 & 0
\end{array}\right],
$$

and

$$
\begin{align*}
\mathrm{e}_{\mathrm{ij}} & =\frac{1}{2}\left(\alpha_{\mathrm{ij}}+\alpha_{\mathrm{ji}}\right) \\
& =\left[\begin{array}{ccc}
0 & 1 & \frac{1}{2} \\
1 & -3 & \frac{1}{2} \\
\frac{1}{2} & \frac{1}{2} & -2
\end{array}\right], \tag{7}
\end{align*}
$$

and

$$
\begin{equation*}
\alpha_{i j}=w_{i j}+e_{i j}, \tag{8}
\end{equation*}
$$

The axis of rotation is

$$
\overline{\mathrm{w}}=\mathrm{w}_{\mathrm{i}} \hat{e}_{i}
$$

where

$$
\begin{align*}
& \mathrm{w}_{1}=\mathrm{w}_{32}=\frac{1}{2}, \\
& \mathrm{w}_{2}=\mathrm{w}_{13}=-\frac{1}{2}, \\
& \mathrm{w}_{3}=\mathrm{w}_{21}=0 . \tag{9}
\end{align*}
$$

### 3.7 GEOMETRICAL INTERPRETATION OF THE COMPONENTS OF STRAIN

## Normal Strain Component $\mathbf{e}_{11}$

Let $\mathrm{e}_{\mathrm{ij}}$ be the components of strains. The pure infinitesimal linear deformation of a vector $\mathrm{A}_{\mathrm{i}}$ is given by

$$
\begin{equation*}
\delta \mathrm{A}_{\mathrm{i}}=\mathrm{e}_{\mathrm{ij}} \mathrm{~A}_{\mathrm{j}}, \tag{1}
\end{equation*}
$$

with $\mathrm{e}_{\mathrm{ij}}=\mathrm{e}_{\mathrm{j} \mathrm{j}}$.
Let e denote the extension (or change) in length per unit length of the vector $\mathrm{A}_{\mathrm{i}}$ with magnitude A. Then, by definition,

$$
\begin{equation*}
\mathrm{e}=\frac{\delta A}{A} . \tag{2}
\end{equation*}
$$

we note that e is positive or negative depending upon whether the material line element $\mathrm{A}_{\mathrm{i}}$ experiences an extension or a contraction.

Also $\mathrm{e}=0$ iff the vector A retains its length during a deformation.
This number $e$ is referred to as the normal strain of the vector $\mathbf{A}_{\mathbf{i}}$.
Since the deformation is linear and infinitesimal, we have (proved earlier)

$$
\begin{equation*}
\mathrm{A} \delta \mathrm{~A}=\mathrm{A}_{\mathrm{i}} \delta \mathrm{~A}_{\mathrm{i}} \tag{3}
\end{equation*}
$$

or

$$
\frac{\delta A}{A}=\frac{A_{i} \delta A_{i}}{A^{2}} .
$$

Now from (1) - (3), we write

$$
\begin{aligned}
\mathrm{e} & =\frac{\delta A}{A}=\frac{A_{i} \delta A_{i}}{A^{2}} \\
& =\frac{1}{A^{2}} \mathrm{~A}_{\mathrm{i}} \mathrm{e}_{\mathrm{ij}} \mathrm{~A}_{\mathrm{j}} .
\end{aligned}
$$

This implies
$\mathrm{e}=\frac{1}{A^{2}}\left[\mathrm{e}_{11} \mathrm{~A}_{1}^{2}+\mathrm{e}_{22} \mathrm{~A}_{2}^{2}+\mathrm{e}_{33} \mathrm{~A}_{3}^{2}+2 \mathrm{e}_{12} \mathrm{~A}_{1} \mathrm{~A}_{2}+2 \mathrm{e}_{13} \mathrm{~A}_{1} \mathrm{~A}_{3}+2 \mathrm{e}_{23} \mathrm{~A}_{2} \mathrm{~A}_{3}\right] \backslash$
since $\mathrm{e}_{\mathrm{ij}}=\mathrm{e}_{\mathrm{j} \mathrm{i}}$.
In particular, we consider the case in which the vector $\mathrm{A}_{\mathrm{i}}$ in the undeformed state is parallel to the $\mathrm{x}_{1}$ - axis. Then

$$
\begin{equation*}
\mathrm{A}_{1}=\mathrm{A}, \mathrm{~A}_{2}=\mathrm{A}_{3}=0 \tag{5}
\end{equation*}
$$

Using (5) , equation (4) gives

$$
\begin{equation*}
\mathrm{e}=\mathrm{e}_{11} . \tag{6}
\end{equation*}
$$

Thus, the component $\mathrm{e}_{11}$ of the strain tensor represents, to a good approximation the extension or change in length per unit initial length of a material line segment (or fibre of the material) originally placed parallel to the $x_{1}$ - axis in the undeformed state.

Similarly, normal strains $\mathrm{e}_{22}$ and $\mathrm{e}_{33}$ are to be interpreted.

Illustration : Let $\mathrm{e}_{\mathrm{ij}}=\left[\begin{array}{lll}e_{11} & 0 & 0 \\ 0 & 0 & 0 \\ 0 & 0 & 0\end{array}\right]$.
Then all unit vectors parallel to the $\mathrm{x}_{1}-$ axis will be extended by an amount $\mathrm{e}_{11}$. In this case, one has a homogeneous deformation of material in the direction of the $\mathrm{x}_{1}$ - axis. A cube of material whose edges before deformation are ' $l$ ' units long will become (after deformation due to $\mathrm{e}_{\mathrm{ij}}$ ) a rectangular parallelopiped whose dimension in the $\mathrm{x}_{1}$ - direction is $l\left(1+\mathrm{e}_{11}\right)$ units and whose dimensions in the direction of the $x_{2}-$ and $x_{3}-$ axes are unchanged.

Remark: The vector

$$
\overline{\mathrm{A}}=\mathrm{A}_{\mathrm{i}}=(\mathrm{A}, 0,0)
$$

is changed to (due to deformation)

$$
\overline{\mathrm{A}}^{\prime}=\left(\mathrm{A}+\delta \mathrm{A}_{1}\right) \hat{e}_{1}+\delta A_{2} \hat{e}_{2}+\delta A_{3} \hat{e}_{3}
$$

in which

$$
\delta \mathrm{A}_{\mathrm{i}}=\mathrm{e}_{\mathrm{ij}} \mathrm{~A}_{\mathrm{j}}=\mathrm{e}_{\mathrm{il}} \mathrm{~A}_{1}
$$

give

$$
\delta \mathrm{A}_{1}=\mathrm{e}_{11} \mathrm{~A}, \delta \mathrm{~A}_{2}=\mathrm{e}_{12} \mathrm{~A}, \delta \mathrm{~A}_{3}=\mathrm{e}_{13} \mathrm{~A} .
$$



Fig. (3.7)
Thus

$$
\overline{\mathrm{A}}^{\prime}=\left(\mathrm{A}+\mathrm{e}_{11} \mathrm{~A}, \mathrm{e}_{12} \mathrm{~A}, \mathrm{e}_{13} \mathrm{~A}\right) .
$$

This indicates that vector $\mathrm{A}_{\mathrm{i}}=(\mathrm{A}, 0,0)$ upon deformation, in general, changes its orientation also.

The length of the vector due to deformation becomes $\left(1+e_{11}\right)$ A .
Question : From the relation $\delta A_{i}=e_{i j} A_{j}$, find $\delta A$ and $\delta A_{i}$ for a vector lying initially along the x - axis (i.e., $\overline{\mathrm{A}}=\mathrm{A} \hat{e}_{1}$ ) and justify the fact that $\frac{\delta A}{A}=\mathrm{e}_{11}$. Does $\delta \mathrm{A}_{\mathrm{i}}$ lie along the $\mathrm{x}-$ axis ?

Answer : It is given that $A_{i}=(A, 0,0)$. The given relation

$$
\begin{equation*}
\delta \mathrm{A}_{\mathrm{i}}=\mathrm{e}_{\mathrm{ij}} \mathrm{~A}_{\mathrm{j}} \tag{1}
\end{equation*}
$$

gives

$$
\begin{equation*}
\delta \mathrm{A}_{1}=\mathrm{e}_{11} \mathrm{~A}, \delta \mathrm{~A}_{2}=\mathrm{e}_{12} \mathrm{~A}, \delta \mathrm{~A}_{3}=\mathrm{e}_{13} \mathrm{~A} . \tag{2}
\end{equation*}
$$

Thus, in general, the vector $\delta A_{i}$ does not lie along the $x$-axis.

## Further

$$
\begin{align*}
(\mathrm{A}+\delta \mathrm{A}) & =\sqrt{\left[A\left(1+e_{11}\right)^{2}+\left(e_{12} A\right)^{2}+\left(e_{13} A\right)^{2}\right.} \\
& =\mathrm{A} \sqrt{1+2 e_{11}+{e_{11}}^{2}+{e_{12}}^{2}+e_{13}{ }^{2}} . \tag{3}
\end{align*}
$$

Neglecting square terms as deformation is small, equation (3) gives

$$
\begin{gather*}
(\mathrm{A}+\delta \mathrm{A})^{2}=\mathrm{A}^{2}\left(1+2 \mathrm{e}_{11}\right), \\
\Rightarrow \quad A^{\not 2}+2 \mathrm{~A} \delta \mathrm{~A}=\not A^{2}+2 \mathrm{~A}^{2} \mathrm{e}_{11}, \\
2 \mathrm{~A} \delta \mathrm{~A}=2 \mathrm{~A}^{2} \mathrm{e}_{11} \\
\frac{\delta A}{A}=\mathrm{e}_{11} . \tag{4}
\end{gather*}
$$

This shows that $e_{11}$ gives the extension of a vector (A, 0,0 ) per unit length due to deformation.

Remark : The strain components $\mathrm{e}_{\mathrm{ij}}$ refer to the chosen set of coordinate axes. If the axes are changed, the strain components $\mathrm{e}_{\mathrm{ij}}$ will, in general, change as per tensor transformation laws.

The shearing strain component $\mathrm{e}_{23}$ may be interpreted by considering intersecting vectors initially parallel to two coordinate axes $-x_{2}-$ and $x_{3}-$ axes.

Now, we consider in the undeformed state two vectors.

$$
\begin{align*}
& \overline{\mathrm{A}}=\mathrm{A}_{2} \hat{e}_{2}, \\
& \overline{\mathrm{~B}}=\mathrm{B}_{3} \hat{e}_{3}, \tag{1}
\end{align*}
$$

directed along $x_{2}-$ and $x_{3}-$ axis, respectively.

## The relations of small linear deformation are

$\delta A_{i}=$
$\mathbf{e}_{\mathrm{ij}} \mathrm{A}_{\mathrm{j}}$,

$$
\delta \mathbf{B}_{\mathrm{i} \delta \overline{\mathbb{B}_{3}}} \mathbf{e}_{\mathrm{Q}} \mathbf{B}_{\mathbf{3}},
$$

Further, the vectors $\mathrm{A}_{\mathrm{i}}$ and $\mathrm{B}_{\mathrm{i}} \mathrm{d} \mathrm{Q}^{\prime}$ e to deformation become (Figure)


Fig. (3.8)

$$
\left.\begin{array}{l}
\overline{\mathbf{A}}^{\prime}=\delta \mathbf{A}_{\mathbf{1}} \hat{e}_{1}+\left(\mathbf{A}_{\mathbf{2}}+\delta \mathbf{A}_{\mathbf{2}}\right) \hat{e}_{2}+\delta \mathbf{A}_{\mathbf{3}} \hat{e}_{3}  \tag{3}\\
\overline{\mathbf{B}}^{\prime}=\delta \mathbf{B}_{1} \hat{e}_{1}+\delta \mathbf{B}_{\mathbf{2}} \hat{e}_{2}+\left(\mathbf{B}_{\mathbf{3}}+\delta \mathbf{B}_{\mathbf{3}}\right) \hat{e}_{3} .
\end{array}\right\}
$$

Deformed vectors $\overline{\mathbf{A}}^{\prime}$ and $\overline{\mathbf{B}}^{\prime}$ need not lie in the $\mathbf{x}_{2} \mathbf{x}_{3}-$ plane. Let $\theta$ be the angle between $\overline{\mathbf{A}}^{\prime}$ and $\overline{\mathbf{B}}^{\prime}$. Then
$\cos \theta=\frac{\bar{A}^{\prime} \cdot \bar{B}^{\prime}}{A^{\prime} B^{\prime}}=\frac{\delta A_{1} \delta B_{1}+\left(A_{2}+\delta A_{2}\right) \delta B_{2}+\delta A_{3}\left(B_{3}+\delta B_{3}\right)}{\sqrt{\left(\delta A_{1}\right)^{2}+\left(A_{2}+\delta A_{2}\right)^{2}+\left(\delta A_{3}\right)^{2}} \sqrt{\left(\delta B_{1}\right)^{2}+\left(\delta B_{2}\right)^{2}+\left(B_{3}+\delta B_{3}\right)^{2}}}$.

Since, the deformation is small, we may neglect the products of the changes in the components of the vector $A_{i}$ and $B_{i}$. Neglecting these products, equation (4) gives

$$
\begin{aligned}
\cos \theta & =\left(\mathbf{A}_{2} \delta \mathbf{B}_{2}+\mathbf{B}_{3} \delta \mathbf{A}_{\mathbf{3}}\right)\left(\mathbf{A}_{\mathbf{2}}+\delta \mathbf{A}_{2}\right)^{-1}\left(\mathbf{B}_{3}+\delta \mathbf{B}_{3}\right)^{-1} \\
& =\frac{\mathbf{A}_{2} \delta \mathrm{~B}_{2}+\mathrm{B}_{3} \delta \mathrm{~A}_{3}}{\mathrm{~A}_{2} \mathrm{~B}_{3}}\left(1+\frac{\delta \mathrm{A}_{2}}{\mathrm{~A}_{2}}\right)^{-1}\left(1+\frac{\delta \mathrm{B}_{3}}{\mathrm{~B}_{3}}\right)^{-1} \\
& =\left(\frac{\delta B_{2}}{B_{3}}+\frac{\delta A_{3}}{A_{2}}\right)\left(1-\frac{\delta A_{2}}{A_{2}}\right)\left(1-\frac{\delta B_{3}}{B_{3}}\right),
\end{aligned}
$$

neglecting other terms. This gives

$$
\begin{equation*}
\cos \theta=\frac{\delta B_{2}}{B_{3}}+\frac{\delta A_{3}}{A_{2}} \tag{5}
\end{equation*}
$$

neglecting the product terms involving changes in the components of the vectors $A_{i}$ and $B_{i}$.

Since in formula (5), all increments in the components of initial vectors $A_{i}$ and $B_{i}$ have been neglected except $\delta A_{3}$ and $\delta B_{2}$, the deformation of these vectors on assuming (w.l.o.g)

$$
\delta \mathbf{A}_{\mathbf{1}}=\delta \mathbf{A}_{\mathbf{2}} \equiv \mathbf{0},
$$

and

$$
\delta \mathbf{B}_{1}=\delta \mathbf{B}_{3} \equiv \mathbf{0},
$$

can be represented as shown in the figure below (It shows that vectors $\mathbf{A}_{\mathbf{i}}{ }^{\prime}$ and $B_{i}{ }^{\prime}$ lie in the
$\mathbf{x}_{2} \mathbf{x}_{3}$ - plane). We call that equations (3)
now may be taken as

$$
\begin{align*}
& \overline{\mathbf{A}}^{\prime}=\mathbf{A}_{2} \hat{e}_{2}+\delta \mathbf{A}_{\mathbf{3}} \hat{e}_{3}, \\
& \overline{\mathbf{B}}^{\prime}=\delta \mathbf{B}_{2} \hat{e}_{2}+\mathbf{B}_{3} \hat{e}_{3} . \tag{6}
\end{align*}
$$



Fig. (3.9)
From equations (1) and (2), we obtain

$$
\begin{align*}
& \delta \mathbf{A}_{3}=\mathbf{e}_{32} \mathbf{A}_{2}, \\
& \delta \mathbf{B}_{2}=\mathbf{e}_{23} \mathbf{B}_{3} \tag{7}
\end{align*}
$$

This gives

$$
\begin{align*}
& \mathbf{e}_{32}=\frac{\delta A_{3}}{A_{2}}=\tan \left\lfloor\mathbf{P}^{\prime} \mathbf{O} \mathbf{P}\right.  \tag{8}\\
& \mathbf{e}_{23}=\frac{\delta \mathbf{B}_{2}}{\mathbf{B}_{3}}=\tan \mathbf{Q}^{\prime} \mathbf{O} \mathbf{Q} . \tag{9}
\end{align*}
$$

Since strains $\mathbf{e}_{23}=\mathbf{e}_{32}$ are small , so

$$
\angle \mathbf{P}^{\prime} \mathbf{O} \mathbf{P}=\angle \mathbf{Q}^{\prime} \mathbf{O} \mathbf{Q} \cong \mathbf{e}_{23},
$$

and hence

$$
\begin{equation*}
\mathbf{2} \mathrm{e}_{23} \cong 90^{\circ}-\theta=\pi / 2-\theta . \tag{10}
\end{equation*}
$$

Thus, a positive value of $2 \mathrm{e}_{23}$ represents a decrease in the right angle between the vectors $A_{i}$ and $B_{i}$ due to small linear deformation which were initially directed along the positive $x_{2}$ - and $x_{3}$-axes. The quantity /strain component $\mathbf{e}_{23}$ is called the shearing strain.

A similar interpretation can be made for the shear strain components $\mathbf{e}_{12}$ and $\mathrm{e}_{13}$.

Shear strain components represent the changes in the relative orientations of material arcs.

Remark 1: By rotating the parallelogram $\mathbf{R}^{\prime} \mathbf{O P}^{\prime} \mathbf{Q}^{\prime}$ through an angle $\mathbf{e}_{23}$ about the origin (in the $x_{2} x_{3}$ - plane), we obtain the following configurations (Figure)


Fig. (3.10)
This figure shows a slide or a shear of planar elements parallel to the $\mathbf{x}_{1} \mathbf{x}_{\mathbf{2}}$ - plane.

Remark 2: Figure shows that the areas of the rectangle OQRP and the parallelogram $O Q^{\prime} R^{\prime} P^{\prime}$ are equal as they have the same height and same base in the $x_{2} \mathbf{x}_{3}$ - plane.

Remark 3: For the strain tensor $\left[\begin{array}{ccc}0 & 0 & 0 \\ 0 & 0 & e_{23} \\ 0 & e_{32} & 0\end{array}\right]$,
a cubical element is deformed into a parallelopiped and the volumes of the cube and parallelopiped remain the same.

Such a small linear deformation is called a pure shear.

### 3.8 NORMAL AND TANGENTIAL DISPLACEMENTS

Consider a point $P\left(x_{1}, x_{2}, x_{3}\right)$ of the material. Let it be moved to $Q$ under a small linear transformation. Let the components of the displacement vector $\overline{P Q}$ be $\mathbf{u}_{1}, \mathbf{u}_{\mathbf{2}}, \mathbf{u}_{\mathbf{3}}$. In the plane OPQ , let $\overline{P N}=\overline{\mathbf{n}}$ be the projection of $\overline{P Q}$ on the line OPN and let $\overline{P T}=\overline{\mathbf{t}}$ be the tangential component of $\overline{P Q}$ in the plane of OPQ or PQN.

Fig. (3.11)


Definition: Vectors $\overline{\mathbf{n}}$ and $\overline{\mathbf{t}}$ are, respectively, called the normal and the tangential components of the displacement of $P$.

Note: The magnitude $n$ of normal displacement $\bar{n}$ is given by the dot product of vectors

$$
\overline{O P}=\left(\mathbf{x}_{1}, \mathbf{x}_{2}, \mathbf{x}_{3}\right) \text { and } \overline{P Q}=\left(\mathbf{u}_{1}, \mathbf{u}_{2}, \mathbf{u}_{3}\right) .
$$

The magnitude $\mathbf{t}$ of tangential vector $\overline{\mathbf{t}}$ is given the vector product of vectors $\overline{O P}$ and $\overline{P Q}$ (This does not give the direction of $\overline{\mathbf{t}}$ ).

Thus

$$
\begin{aligned}
& \mathbf{n}=\mathbf{P Q} \cos \mathbf{N P Q}=\frac{\overline{O P} \cdot \overline{P Q}}{|\overline{O P}|}, \\
& \mathbf{t}=\mathbf{P Q} \sin \left\lvert\, \mathbf{N P Q}=\frac{(O P)(P Q) \sin (N P Q)}{O P}=\frac{\mid \overline{O P} \times \overline{P Q \mid}}{|\overline{O P}|}\right.,
\end{aligned}
$$

and

$$
\mathbf{n}^{2}+\mathbf{t}^{2}=\mathbf{u}_{1}{ }^{2}+\mathbf{u}_{2}{ }^{2}+\mathbf{u}_{3}{ }^{2} .
$$

### 3.9 STRAIN QUADRIC CAUCHY

Let $\mathbf{P}^{\circ}\left(\mathbf{x}_{1}{ }^{\circ}, \mathbf{x}_{2}{ }^{\circ}, \mathbf{x}_{3}{ }^{\circ}\right)$ be any but fixed point of a continuous medium, with reference axes $0 x_{1} X_{2} x_{3}$ fixed in space.

We introduce a local system of axes with origin at the point $\mathbf{P}^{\circ}$ and with axes parallel to the fixed axes (figure).


Fig. (3.12)

With reference to these local axes, consider the equation

$$
\begin{equation*}
\mathbf{e}_{\mathrm{ij}} \mathbf{x}_{\mathbf{i}} \mathbf{x}_{\mathbf{j}}= \pm \mathbf{k}^{2} \tag{1}
\end{equation*}
$$

where $k$ is a real constant and $e_{i j}$ is the strain tensor at $P^{\circ}$. This equation represents a quadric surface with its centre at $\mathbf{P}^{\circ}$.

This quadric is called the quadric surface of deformation or strain quadric or strain quadric of Cauchy.

The sign + or - in equation (1) be chosen so that the quadric surface (1) becomes a real one.

The nature of this quadric surface depends on the value of the strains $\mathrm{e}_{\mathrm{ij}}$.
If $\left|\mathrm{e}_{\mathrm{i} j}\right| \neq \mathbf{0}$, the quadric is either an ellipsoid or a hyperboloid.
If $\left|e_{i j}\right|=0$, the quadric surface degenerates into a cylinder of the elliptic or hyperbolic type or else into two parallel planes symmetrically situated with respect to the origin $\mathrm{P}^{\circ}$ of the quadric surface.

This strain quadric is completely determined once the strain components $\mathrm{e}_{\mathrm{ij}}$ at point $\mathrm{P}^{\circ}$ are known.

Let $\overline{P^{\circ} P}$ be the radius vector $A_{i}$ of magnitude $A$ to any point $P\left(x_{1}, x_{2}, x_{3}\right)$, referred to local axis, on the strain quadric surface (1). Let $e$ be the extension of the vector $A_{i}$ due to some linear deformation characterized by

$$
\begin{equation*}
\delta \mathbf{A}_{\mathbf{i}}=\mathbf{e}_{\mathbf{i j}} \mathbf{A}_{\mathbf{j}} . \tag{2}
\end{equation*}
$$

Then, by definition ,

$$
\mathbf{e}=\frac{\delta A}{A}=\frac{A \delta A}{A^{2}}=\frac{A_{i} \delta A_{i}}{A^{2}} .
$$

This gives

$$
\begin{equation*}
\mathbf{e}=\frac{e_{i j} A_{i} A_{j}}{A^{2}} \tag{3}
\end{equation*}
$$

using (2).
Since $\overline{P^{\circ} P}=\mathbf{A}_{\mathbf{i}}$ and the coordinate of the point $\mathbf{P}$, on the surface (1), relative to $\mathrm{P}^{\circ}$ are ( $\mathrm{x}_{1}, \mathrm{x}_{2}, \mathrm{x}_{3}$ ), it follows that

$$
\begin{equation*}
\mathbf{A}_{\mathbf{i}}=\mathbf{x}_{\mathbf{i}} . \tag{4}
\end{equation*}
$$

From equations (1), (3) and (4) ; we obtain

$$
\mathbf{e} \mathbf{A}^{2}=\mathbf{e}_{\mathrm{ij}} \mathbf{A}_{\mathbf{i}} \mathbf{A}_{\mathbf{j}}=\mathbf{e}_{\mathrm{ij}} \mathbf{x}_{\mathbf{i}} \mathbf{x}_{\mathrm{j}}= \pm \mathbf{k}^{2}
$$

$$
e A^{2}= \pm \mathbf{k}^{2}
$$

or

$$
\begin{equation*}
\mathrm{e}= \pm \frac{k^{2}}{A^{2}} \tag{5}
\end{equation*}
$$

Result (I) : Relation (5) shows that the extension or elongation of any radius vector $A_{i}$ of the strain quadric of Cauchy, given by equation (1), is inversely proportional to the square of the length $A$ of that radius vector. This determines the elongation of any radius vector of the strain quadric at the point $\mathrm{P}^{\circ}\left(\mathrm{x}_{\mathrm{i}}{ }^{\circ}\right)$.

Result (II) : We know that the length $\mathbf{A}$ of the radius vector $\mathbf{A}_{\mathbf{i}}$ of strain quadric (1) at the point $P^{\circ}\left(\mathbf{x}_{i}{ }^{\circ}\right)$ has maximum and minimum values along the axes of the quadric. In general , axes of the strain quadric (1) differs from the coordinate axes through $P^{\circ}\left(x_{i}{ }^{\circ}\right)$.

Therefore, the maximum and minimum extensions /elongations of radius vectors of strain quadric (1) will be along its axes.

Result (III) : Another interesting property of the strain quadric (1) is that normal $v_{i}$ to this surface at the end point $\mathbf{P}$ of the vector $\overline{P^{\circ} P}=\mathbf{A}_{\mathbf{i}}$ is parallel to the displacement vector $\delta \mathbf{A}_{\mathbf{i}}$.

To prove this property, let us write equation (1) in the form

$$
\begin{equation*}
G=e_{i j} \mathbf{x}_{j} \mathbf{x}_{\mathrm{i}} \mp \mathbf{k}^{2}=\mathbf{0} \tag{6}
\end{equation*}
$$

Then the direction of the normal $\hat{v}$ to the strain quadric (6) is given by the gradient of the scalar function $G$. The components of the gradient are

$$
\begin{aligned}
\frac{\partial \mathrm{G}}{\partial \mathrm{x}_{\mathrm{k}}} & =\mathbf{e}_{\mathbf{i j}} \delta_{\mathbf{i k}} \mathbf{x}_{\mathbf{j}}+\mathbf{e}_{\mathbf{i j}} \mathbf{x}_{\mathbf{i}} \delta_{\mathbf{k j}} \\
& =\mathbf{e}_{\mathbf{k j}} \mathbf{x}_{\mathbf{j}}+\mathbf{e}_{\mathbf{i k}} \mathbf{x}_{\mathbf{i}} \\
& =\mathbf{2} \mathbf{e}_{\mathbf{k j}} \mathbf{x}_{\mathbf{j}}
\end{aligned}
$$

or

$$
\begin{equation*}
\frac{\partial \mathrm{G}}{\partial \mathrm{x}_{\mathrm{k}}}=\mathbf{2} \delta \mathbf{A}_{\mathrm{k}} \tag{7}
\end{equation*}
$$

This shows that vector $\frac{\delta G}{\delta x_{k}}$ and $\delta \mathbf{A}_{\mathbf{k}}$ are parallel.

$\delta \mathbf{A}_{\mathbf{i}} \| \mathbf{v}_{\mathbf{i}}$
Fig. (3.13)
Hence, the vector $\overline{\delta A}$ is directed along the normal at $\mathbf{P}$ to the strain quadric of Cauchy.

### 3.10 STRAIN COMPONENTS AT A POINT IN A ROTATION OF COORDINATE AXES

Let new axes $0 \mathbf{x}_{1}{ }^{\prime} \mathbf{x}^{\prime}{ }^{\prime} \mathbf{x}_{3}{ }^{\prime}$ be obtained from the old reference system $0 \mathbf{x}_{1} \mathbf{x}_{\mathbf{2}}$ $x_{3}$ by a rotation (figure).


Fig. (3.14)
Let the directions of the new axes $x_{i}{ }^{\prime}$ be specified relative to the old system $\mathbf{x}_{\mathrm{i}}$ by the following table of direction cosines in which $\mathrm{a}_{\mathrm{pi}}$ is the cosine of the angle between the $x_{p}{ }^{\prime}$ - and $x_{i}-$ axis.

That is,

$$
\mathbf{a}_{\mathrm{pi}}=\cos \left(\mathrm{x}_{\mathrm{p}}{ }^{\prime}, \mathrm{x}_{\mathrm{i}}\right) .
$$

Thus

|  | $\mathbf{x}_{1}$ | $\mathbf{x}_{2}$ | $\mathbf{\mathbf { x } _ { 3 }}$ |
| :--- | :---: | :---: | :---: |
| $\mathbf{x}_{1}{ }^{\prime}$ | $\mathbf{a}_{11}$ | $\mathbf{a}_{12}$ | $\mathbf{a}_{13}$ |
| $\mathbf{x}_{2}{ }^{\prime}$ | $\mathbf{a}_{21}$ | $\mathbf{a}_{22}$ | $\mathbf{a}_{23}$ |
| $\mathbf{x}_{3^{\prime}}$ | $\mathbf{a}_{31}$ | $\mathbf{a}_{32}$ | $\mathbf{a}_{33}$ |

Then the transformation law for coordinates is

$$
\mathbf{x}_{\mathrm{i}}=\mathbf{a}_{\mathbf{p i}} \mathbf{x}_{\mathbf{p}}{ }^{\prime},
$$

or

$$
\begin{equation*}
\mathbf{x}_{\mathbf{p}}^{\prime}=\mathbf{a}_{\mathbf{p i}} \mathbf{x}_{\mathbf{i}} \tag{2}
\end{equation*}
$$

The well - known orthogonality relation are

$$
\begin{align*}
& \mathbf{a}_{\mathbf{p i}} \mathbf{a}_{\mathbf{q i}}=\delta_{p q},  \tag{3}\\
& \mathbf{a}_{\mathbf{p i}} \mathbf{a}_{\mathbf{p j}}=\delta_{\mathrm{ij}}, \tag{4}
\end{align*}
$$

with reference to new $\mathbf{x}_{\mathrm{p}}{ }^{\prime}$ system, a new set of strain components $\mathbf{e}^{\prime}{ }_{\mathrm{pq}}$ is determined at the point $O$ while $\mathrm{e}_{\mathrm{ij}}$ are the components of strain at O relative to old axes $O x_{1} \mathbf{x}_{2} x_{3}$.

Let

$$
\begin{equation*}
\mathbf{e}_{\mathrm{ij}} \mathbf{x}_{\mathbf{i}} \mathbf{x}_{\mathbf{j}}= \pm \mathbf{k}^{2} \tag{5}
\end{equation*}
$$

be the equation of the strain quadric surface relative to old axis. The equation of quadric surface with reference to new prime system becomes

$$
\begin{equation*}
\mathbf{e}_{\mathbf{p q}}^{\prime} \mathbf{x}_{\mathbf{p}}^{\prime} \mathbf{x}_{\mathbf{q}}^{\prime}= \pm \mathbf{k}^{2} \tag{6}
\end{equation*}
$$

as we know that quadric form is invariant w.r.t. an orthogonal transformation of coordinates.

Further, equations (2) to (6) together yield

$$
\begin{aligned}
\mathbf{e}_{\mathrm{pq}}^{\prime} \mathbf{x}_{\mathrm{p}}^{\prime} \mathbf{x}_{\mathrm{q}}^{\prime} & =\mathrm{e}_{\mathrm{ij}} \mathbf{x}_{\mathrm{i}} \mathbf{x}_{\mathrm{j}} \\
& =\mathbf{e}_{\mathrm{ij}}\left(\mathbf{a}_{\mathrm{pi}} \mathbf{x}_{\mathrm{p}}^{\prime}\right)\left(\mathbf{a}_{\mathrm{qj}} \mathbf{x}_{q}^{\prime}\right) \\
& =\left(\mathbf{e}_{\mathrm{ij}} \mathbf{a}_{\mathrm{pi}} \mathbf{a}_{\mathrm{qj}}\right) \mathbf{x}_{\mathrm{p}}^{\prime} \mathbf{x}_{\mathrm{q}}^{\prime}
\end{aligned}
$$

or

$$
\begin{equation*}
\left(\mathbf{e}_{\mathrm{pq}}^{\prime}-\mathbf{a}_{\mathrm{pi}} \mathbf{a}_{\mathbf{q j}} \mathbf{e}_{\mathrm{ij}}\right) \mathbf{x}_{\mathbf{p}}^{\prime} \mathbf{x}_{\mathbf{q}}^{\prime}=\mathbf{0} . \tag{7}
\end{equation*}
$$

Since equation (7) is satisfied for arbitrary vector $x_{p}^{\prime}$, we must have

$$
\begin{equation*}
\mathbf{e}_{\mathbf{p q}}^{\prime}=\mathbf{a}_{\mathbf{p i}} \mathbf{a}_{\mathbf{q j}} \mathbf{e}_{\mathrm{ij}} \tag{8}
\end{equation*}
$$

Equation (8) is the law of transformation for a second order tensor.

We, therefore, conclude that the components of strain form a second order tensor.

Similarly, it can be verified that

$$
\begin{equation*}
\mathbf{e}_{\mathrm{ij}}=\mathbf{a}_{\mathrm{pi}} \mathbf{a}_{\mathbf{q j}} \mathbf{e}_{\mathrm{pq}}^{\prime} . \tag{9}
\end{equation*}
$$

Question : Assuming that $\mathbf{e}_{\mathbf{i j}}$ is a tensor of order 2, show that quadratic form $e_{i j} x_{i} x_{j}$ is an variant.

Solution : We have
so

$$
\begin{aligned}
& \mathbf{e}_{\mathrm{ij}}=\mathbf{a}_{\mathrm{p} i} \mathbf{a}_{\mathbf{q j}} \mathbf{e}_{\mathrm{pq}}^{\prime}, \\
& \begin{aligned}
\mathbf{e}_{\mathrm{ij}} \mathbf{x}_{\mathbf{i}} \mathbf{x}_{\mathbf{j}} & =\mathbf{a}_{\mathrm{pi}} \mathbf{a}_{\mathbf{q j}} \mathbf{e}_{\mathrm{pq}}^{\prime} \mathbf{x}_{\mathrm{i}} \mathbf{x}_{\mathbf{j}} \\
& =\mathbf{e}_{\mathrm{pq}}^{\prime}\left(\mathbf{a}_{\mathrm{pi}} \mathbf{x}_{\mathbf{i}}\right)\left(\mathbf{a}_{\mathbf{q j}} \mathbf{x}_{\mathbf{j}}\right) \\
& =\mathbf{e}_{\mathrm{pq}}^{\prime} \mathbf{x}_{\mathbf{p}}^{\prime} \mathbf{x}_{\mathbf{q}}^{\prime} .
\end{aligned}
\end{aligned}
$$

Hence the result.

### 3.11 PRINCIPAL STRAINS AND INVARIANTS

From a material point $\mathbf{P}^{\circ}\left(\mathbf{x}_{\mathrm{i}}{ }^{\circ}\right)$ there emerges infinitely many material arcs /filaments, and each of these arcs generally changes in length and orientation under a deformation.

We seek now the lines through $\mathbf{P}^{\circ}\left(\mathbf{x}_{\mathbf{i}}{ }^{\circ}\right)$ whose orientation is left unchanged by the small linear deformation given by

$$
\begin{equation*}
\delta \mathbf{A}_{\mathbf{i}}=\mathbf{e}_{\mathbf{i j}} \mathbf{A}_{\mathbf{j}} \tag{1}
\end{equation*}
$$

where the strain components $\mathrm{e}_{\mathrm{ij}}$ are small and constant.
In this situation, vectors $\mathrm{A}_{\mathrm{i}}$ and $\delta \mathrm{A}_{\mathrm{i}}$ are parallel and, therefore,

$$
\begin{equation*}
\delta \mathbf{A}_{\mathbf{i}}=\mathbf{e} \mathbf{A}_{\mathbf{i}} \tag{2}
\end{equation*}
$$

for some constant e .
Equation (2) shows that the constant e represents the extension

$$
\left(e=\frac{\left|\delta A_{i}\right|}{\left|A_{i}\right|}=\frac{\delta A}{A}\right)
$$

of vector $\mathrm{A}_{\mathrm{i}}$.
From equations (1) and (2), we write

$$
\begin{align*}
\mathbf{e}_{\mathbf{i j}} \mathbf{A}_{\mathbf{j}} & =\mathbf{e} \mathbf{A}_{\mathbf{i}} \\
& =\mathbf{e} \delta_{\mathrm{ij}} \mathbf{A}_{\mathbf{j}} . \tag{3}
\end{align*}
$$

This implies

$$
\begin{equation*}
\left(\mathbf{e}_{\mathrm{ij}}-\mathbf{e} \delta_{\mathrm{ij}}\right) \mathbf{A}_{\mathbf{j}}=\mathbf{0} \tag{4}
\end{equation*}
$$

We know that $\mathrm{e}_{\mathrm{ij}}$ is a real symmetric tensor of order 2. The equation (3) shows that the scalar $e$ is an eigenvalue of the real symmetric tensor $\mathbf{e}_{\mathrm{ij}}$ with corresponding eigenvector $\mathbf{A}_{i}$. Therefore, we conclude that there are precisely three mutually orthogonal directions whose orientations are not changed on account of deformation and these directions coincide with the three eigenvectors of the strain tensor $\mathrm{e}_{\mathrm{ij}}$.

These directions are known as principal directions or invariant directions of strain.

Equation (4) gives us a system of three homogeneous equations in the unknowns $\mathbf{A}_{1}, \mathbf{A}_{2}, \mathbf{A}_{3}$. This system possesses a non - trivial solution iff the determinant of the coefficients of the $A_{1}, A_{2}, A_{3}$ is equal to zero, i.e.,

$$
\left|\begin{array}{ccc}
e_{11}-e & e_{12} & e_{13}  \tag{5}\\
e_{21} & e_{22}-e & e_{23} \\
e_{31} & e_{32} & e_{33}-e
\end{array}\right|=\mathbf{0}
$$

which is a cubic equation in $e$.
Let $e_{1}, e_{2}, e_{3}$ be the three roots of equation (5). These are known as principal strains.

Evidently, the principal strains are the eigenvalues of the second order real symmetric strain tensor $\mathrm{e}_{\mathrm{ij}}$. Consequently, these principal strains are all real (not necessarily distinct).

Physically , the principal strains $e_{1}, e_{2}, e_{3}$ (all different) are the extensions of the vectors, say $\frac{i}{A}$, in the principal /invariant directions of strain. So, vectors $\stackrel{i}{\mathrm{~A}}, \stackrel{i}{\delta} \mathrm{~A}_{\sim}, \stackrel{i}{\mathrm{~A}}+\delta{\underset{\sim}{\mathrm{A}}}_{\sim}^{\mathrm{A}}$ are collinear.

Fig. (3.15)


At the point $\mathbf{P}^{\circ}$, consider the strain quadric

$$
\begin{equation*}
\mathbf{e}_{\mathrm{ij}} \mathbf{x}_{\mathrm{i}} \mathbf{x}_{\mathrm{j}}= \pm \mathbf{k}^{2} \tag{6}
\end{equation*}
$$

For every principal direction of strain $\stackrel{i}{A}$, we know that $\delta \stackrel{i}{A}$ is normal to the quadric surface (6). Therefore, the principal directions of strain are also normal to the strain quadric of Cauchy. Hence, principal direction of strain must be the three principal axes of the strain quadric of Cauchy.

If some of the principal strains $e_{i}$ are equal, then the associated directions become indeterminate but one can always select three directions that all mutually orthogonal (already proved).

If $\mathbf{e}_{1} \neq \mathbf{e}_{2}=\mathbf{e}_{3}$, then the quadric surface of Cauchy is a surface revolution and our principal direction,
say $\underset{\sim}{A}$, will be directed along the axis of revolution.


Fig. (3.16)
In this case, any two mutually perpendicular vectors lying in the plane normal to $\stackrel{1}{A}$ may be taken as the other two principal directions of strain.

If $\mathbf{e}_{1}=\mathbf{e}_{2}=\mathbf{e}_{3}$, then strain quadric of Cauchy becomes a sphere and any three orthogonal directions may be chosen as the principal directions of strain.

Result 1 : If the principal directions of strain are taken as the coordinate axes, then

$$
\mathbf{e}_{11}=\mathbf{e}_{1}, e_{22}=\mathbf{e}_{2}, e_{33}=\mathbf{e}_{3}
$$

and

$$
\mathbf{e}_{12}=\mathbf{e}_{13}=\mathbf{e}_{23}=\mathbf{0},
$$

as a vector initially along an axis remains in the same direction after deformation (so changes in right angles are zero). In this case, the strain quadric of Cauchy has the equation

$$
\begin{equation*}
e_{1} \mathbf{x}_{1}^{2}+e_{2} \mathbf{x}_{2}^{2}+e_{3}^{2} \mathbf{x}_{3}^{2}= \pm k^{2} \tag{7}
\end{equation*}
$$

Result 2: Expanding the cubic equation (5), we write

$$
\text { where } \quad \begin{align*}
&-\mathbf{e}^{3}+v_{1} \mathbf{e}^{2}-v_{2} \mathbf{e}+v_{3}=0 . \\
& v_{1}=\mathbf{e}_{11}+\mathbf{e}_{22}+\mathbf{e}_{33} \\
&=\mathbf{e}_{i \mathrm{i}}=\operatorname{tr}(\mathbf{E}) \quad, \\
& v_{2}=\mathbf{e}_{11} \mathbf{e}_{22}+\mathbf{e}_{22} \mathbf{e}_{33}+\mathbf{e}_{33} \mathbf{e}_{11}-\mathbf{e}_{23}^{2}-\mathbf{e}_{13}{ }^{2}-\mathbf{e}_{12}^{2}  \tag{8}\\
&=\operatorname{tr}\left(\mathbf{E}^{2}\right)=\frac{1}{2}\left(\mathbf{e}_{\mathrm{ii}} \mathbf{e}_{\mathrm{ij}}-\mathbf{e}_{\mathrm{ij}} \mathbf{e}_{\mathrm{jij}}\right) \\
& v_{3}=\epsilon_{\mathrm{ijk}} \mathbf{e}_{1 \mathrm{i}} \mathbf{e}_{2 \mathrm{j}} \mathbf{e}_{3 \mathrm{k}}  \tag{9}\\
&=\left|\mathbf{e}_{\mathrm{ij}}\right|=\operatorname{tr}\left(\mathbf{E}^{3}\right) .
\end{align*}
$$

Also , $e_{1}, e_{2}, e_{3}$ are roots of the cubic equation (8), so

$$
\left.\begin{array}{l}
v_{1}=e_{1}+e_{2}+e_{3}  \tag{11}\\
v_{2}=e_{1} e_{2}+e_{2} e_{3}+e_{3} e_{1} \\
v_{3}=e_{1} e_{2} e_{3}
\end{array}\right\}
$$

We know that eigenvalues of a second order real symmetric tensor are independent of the choice of the coordinate system.

It follows that $v_{1}, v_{2}, v_{3}$ are, as given by (10), three invariants of the strain tensor $\mathrm{e}_{\mathrm{ij}}$ with respect to an orthogonal transformation of coordinates.

Geometrical Meaning of the First Strain Invariant $v=\mathrm{e}_{\mathrm{ii}}$
The quantity $v=\mathrm{e}_{\mathrm{ii}}$ has a simple geometrical meaning. Consider a volume element in the form of a rectangular parallelepiped whose edges of length $l_{1}, l_{2}, l_{3}$ are parallel to the principal directions of strain.

Due to small linear transformation /deformation, this volume element becomes again a rectangular parallelepiped with edges of length $\boldsymbol{l}_{1}\left(1+e_{1}\right)$, $l_{2}\left(1+e_{2}\right), l_{3}\left(1+e_{3}\right)$, where $e_{1}, e_{2}, e_{3}$ are principal strains.

Hence, the change $\delta V$ in the volume $V$ of the element is

$$
\delta \mathrm{V}=l_{1} l_{2} l_{3}\left(1+\mathrm{e}_{1}\right)\left(1+\mathrm{e}_{2}\right)\left(1+\mathrm{e}_{3}\right)-l_{1} l_{2} l_{3}
$$

$$
\begin{aligned}
& =l_{1} l_{2} l_{3}\left(1+\mathrm{e}_{1}+\mathrm{e}_{2}+\mathrm{e}_{3}\right)-l_{1} l_{2} l_{3}, \quad \text { ignoring small strains } \mathrm{e}_{\mathrm{i}} . \\
& =l_{1} l_{2} l_{3}\left(\mathrm{e}_{1}+\mathrm{e}_{2}+\mathrm{e}_{3}\right)
\end{aligned}
$$

This implies

$$
\frac{\delta V}{V}=\mathbf{e}_{1}+\mathbf{e}_{2}+\mathbf{e}_{3}=v_{1}
$$

Thus, the first strain invariant $v_{\mathbf{1}}$ represents the change in volume per unit initial volume due to strain produced in the medium.

The quantity $v_{1}$ is called the cubical dilatation or simply the dilatation.
Note : If $\mathbf{e}_{\mathbf{1}}>\mathbf{e}_{\mathbf{2}}>\mathbf{e}_{3}$ then $\mathbf{e}_{\mathbf{3}}$ is called the minor principal strain , $\mathbf{e}_{\mathbf{2}}$ is called the intermediate principal strain, and $\mathbf{e}_{1}$ is called the major principal strain.

Question : For small linear deformation , the strains $\mathbf{e}_{\mathrm{ij}}$ are given by

$$
\left(\mathbf{e}_{\mathbf{i j}}\right)=\alpha\left[\begin{array}{ccc}
x_{2} & \left(x_{1}+x_{2}\right) / 2 & x_{3} \\
\left(x_{1}+x_{2}\right) / 2 & x_{1} & x_{3} \\
x_{3} & x_{3} & 2\left(x_{1}+x_{2}\right)
\end{array}\right], \quad \alpha
$$

$=$ constant.
Find the strain invariants, principal strains and principal directions of strain at the point $\mathrm{P}(\mathbf{1}, \mathbf{1}, \mathbf{0})$.

Solution : The strain matrix at the point $\mathbf{P}(1,1,0)$ becomes

$$
\left(\mathbf{e}_{\mathrm{ij}}\right)=\left[\begin{array}{ccc}
\alpha & \alpha & 0 \\
\alpha & \alpha & 0 \\
0 & 0 & 4 \alpha
\end{array}\right]
$$

whose characteristic equation becomes

$$
e(e-2 \alpha)(e-4 \alpha)=0
$$

Hence, the principle strains are

$$
\mathrm{e}_{1}=0, \mathrm{e}_{2}=2 \alpha, \mathrm{e}_{3}=4 \alpha .
$$

The three scalar invariants are

$$
v_{1}=e_{1}+e_{2}+e_{3}=6 \alpha, v_{2}=8 \alpha^{2}, v_{3}=0
$$

The three principal unit directions are found to be

$$
\stackrel{1}{A_{i}}=\left(\frac{1}{\sqrt{2}}, \frac{-1}{\sqrt{2}}, 0\right), \quad \stackrel{2}{A_{i}}=\left(\frac{1}{\sqrt{2}}, \frac{1}{\sqrt{2}}, 0\right)
$$

$$
\stackrel{3}{A}_{i}=0,0,1_{-}^{-}
$$

Exercise : The strain field at a point $\mathbf{P}(\mathbf{x}, \mathbf{y}, \mathbf{z})$ in an elastic body is given by

$$
\mathbf{e}_{\mathrm{ij}}=\left[\begin{array}{ccr}
20 & 3 & 2 \\
3 & -10 & 5 \\
2 & 5 & -8
\end{array}\right] \times \mathbf{1 0}^{-6}
$$

Determine the strain invariants and the principal strains.
Question : Find the principal directions of strain by finding the extremal value of the extension $e$.

## OR

Find the directions in which the extension e is stationary.
Solution: Let $\mathbf{e}$ be the extension of a vector $\mathbf{A}_{\mathbf{i}}$ due to small linear deformation

$$
\begin{equation*}
\delta \mathbf{A}_{\mathbf{i}}=\mathbf{e}_{\mathbf{i j}} \mathbf{A}_{\mathbf{j}} \tag{1}
\end{equation*}
$$

Then

$$
\begin{equation*}
\mathbf{e}=\frac{\delta A}{A} . \tag{2}
\end{equation*}
$$

We know that for an infinitesimal linear deformation (1), we have

Thus

$$
\begin{equation*}
\mathbf{A} \delta \mathbf{A}=\mathbf{A}_{\mathbf{i}} \delta \mathbf{A}_{\mathbf{i}} \tag{3}
\end{equation*}
$$

Let $\quad \frac{A_{i}}{A}=\mathbf{a}_{\mathbf{i}}$.
Then

$$
\begin{equation*}
\mathbf{a}_{\mathbf{i}} \mathbf{a}_{\mathbf{i}}=\mathbf{1} \tag{6}
\end{equation*}
$$

and equation (4) then gives

$$
\begin{equation*}
\mathbf{e}\left(\mathbf{a}_{1}, \mathbf{a}_{2}, \mathbf{a}_{3}\right)=\mathbf{e}_{\mathrm{ij}} \mathbf{a}_{\mathrm{i}} \mathbf{a}_{\mathbf{j}} . \tag{7}
\end{equation*}
$$

Thus, the extension $e$ is a function of $a_{1}, a_{2}, a_{3}$ which are not independent because of relation (6). The extreme/stationary (or Max/Min) values of the extension $e$ are to be found by making use of Lagrange's method of multipliers.

For this purpose, we consider the auxiliary function

$$
\begin{equation*}
\mathbf{F}\left(\mathbf{a}_{1}, \mathbf{a}_{2}, \mathbf{a}_{3}\right)=\mathbf{e}_{\mathbf{i j}} \mathbf{a}_{\mathbf{i}} \mathbf{a}_{\mathbf{j}}-\lambda\left(\mathbf{a}_{\mathrm{i}} \mathbf{a}_{\mathbf{i}}-\mathbf{1}\right), \tag{8}
\end{equation*}
$$

where $\lambda$ is a constant.
In order to find the values of $\mathbf{a}_{1}, a_{2}, a_{3}$ for which the function (7) may have a maximum or minimum, we solve the equations

$$
\begin{equation*}
\frac{\partial \mathrm{F}}{\partial \mathrm{a}_{\mathrm{k}}}=0 \quad, \mathbf{k}=\mathbf{1}, \mathbf{2}, \mathbf{3} \tag{9}
\end{equation*}
$$

Thus, the stationary values of $e$ are given by

$$
\mathbf{e}_{i j}\left(\delta_{i k} \mathbf{a}_{\mathrm{j}}+\mathbf{a}_{\mathrm{i}} \delta_{\mathrm{jk}}\right)-\lambda 2 \mathbf{a}_{\mathrm{i}} \delta_{\mathrm{ik}}=\mathbf{0}
$$

or

$$
\mathbf{e}_{\mathrm{kj}} \mathbf{a}_{\mathrm{j}}+\mathbf{e}_{\mathrm{ik}} \mathbf{a}_{\mathrm{i}}-2 \lambda \mathbf{a}_{\mathrm{k}}=\mathbf{0}
$$

$$
2 e_{k i} a_{i}-2 \lambda a_{k}=0
$$

or

$$
\begin{equation*}
\mathbf{e}_{\mathrm{ki}} \mathbf{a}_{\mathrm{i}}=\lambda \mathbf{a}_{\mathrm{k}} . \tag{10}
\end{equation*}
$$

This shows that $\lambda$ is an eigenvalue of the strain tensor $e_{i j}$ and $\mathbf{a}_{i}$ is the corresponding eigenvector. Therefore, equations in (10) determine the principal strains and principal directions of strain.

Thus, the extension e assumes the stationary values along the principal directions of strain and the stationary/extreme values are precisely the principal strains.

Remark : Let $M$ be the square matrix with eigenvectors of the strain tensor $\mathbf{e}_{\mathrm{ij}}$ as columns. That is

$$
\mathbf{M}=\left[\begin{array}{ccc}
1 & 2 & 3 \\
A_{1} & A_{1} & A_{1} \\
1 & 2 & 3 \\
A_{2} & A_{2} & A_{2} \\
1 & 2 & 3 \\
A_{3} & A_{3} & A_{3}
\end{array}\right]
$$

Then

$$
\begin{aligned}
& \mathbf{e}_{\mathbf{i j}} \stackrel{1}{A}_{j}=\mathbf{e}_{\mathbf{1}} \stackrel{1}{A_{i}}, \\
& \mathbf{e}_{\mathbf{i j}} \stackrel{2}{A}_{j}^{2}=\mathbf{e}_{\mathbf{2}} \stackrel{2}{A_{i}}, \\
& \mathbf{e}_{\mathbf{i j}} \stackrel{3}{A_{j}}=\mathbf{e}_{\mathbf{3}} \stackrel{3}{A_{i}}
\end{aligned}
$$

The matrix $M$ is called the modal matrix of strain tensor $\mathbf{e}_{\mathrm{ij}}$.
Let

$$
E=\left(e_{i j}\right), D=\operatorname{dia}\left(e_{1}, e_{2}, e_{3}\right) .
$$

Then, we find
or

$$
\mathbf{E} \mathbf{M}=\mathbf{M} \mathbf{D}
$$

$\mathbf{M}^{-1} \mathbf{E M}=\mathbf{D}$.
This shows that the matrices $E$ and $D$ are similar.
We know that two similar matrices have the same eigenvalues. Therefore, the characteristic equation associated with $\mathbf{M}^{-1} E M$ is the same as the one associated with $E$. Consequently, eigenvalues of $E$ and $D$ are identical.

Question : Show that, in general, at any point of the elastic body there exists (at least) three mutually perpendicular principal directions of strain due to an infinitesimal linear deformation.

Solution: Let $\mathbf{e}_{1}, \mathbf{e}_{\mathbf{2}}, \mathbf{e}_{3}$ be the three principal strains of the strain tensor $\mathrm{e}_{\mathrm{ij} \text {. }}$ Then, they are the roots of the cubic equation

$$
\left(\mathbf{e}-\mathbf{e}_{1}\right)\left(\mathbf{e}-\mathbf{e}_{2}\right)\left(\mathbf{e}-\mathbf{e}_{3}\right)=0,
$$

and

$$
\begin{aligned}
\mathbf{e}_{1}+\mathbf{e}_{2}+\mathbf{e}_{3}=\mathbf{e}_{11}+\mathbf{e}_{22}+\mathbf{e}_{33}=\mathbf{e}_{\mathrm{ii}}, \\
\mathbf{e}_{1} \mathbf{e}_{2}+\mathbf{e}_{2} \mathbf{e}_{3}+\mathbf{e}_{3} \mathbf{e}_{1}=\frac{1}{2}\left(\mathbf{e}_{\mathrm{ii}} \mathbf{e}_{\mathrm{ij}}-\mathbf{e}_{\mathrm{ij}} \mathbf{e}_{\mathrm{ij}}\right), \\
\mathbf{e}_{1} \mathbf{e}_{2} \mathbf{e}_{3}=\left|\mathbf{e}_{\mathrm{ij}}\right|=\epsilon_{\mathrm{ijk}} \mathbf{e}_{1 \mathrm{i}} \mathbf{e}_{2 \mathrm{j}} \mathbf{e}_{3 \mathrm{k}} .
\end{aligned}
$$

We further assume that coordinate axes coincide with the principal directions of strain. Then, the strain components are given by

$$
\begin{aligned}
& \mathbf{e}_{11}=\mathbf{e}_{1}, e_{22}=\mathbf{e}_{2}, e_{33}=e_{3}, \\
& \mathbf{e}_{12}=e_{13}=\mathbf{e}_{23}=0,
\end{aligned}
$$

and the strain quadric of Cauchy becomes

$$
\begin{equation*}
e_{1} \mathbf{x}_{1}^{2}+e_{2} \mathbf{x}_{2}^{2}+e_{3} \mathbf{x}_{3}{ }^{2}= \pm k^{2} . \tag{1}
\end{equation*}
$$

Now, we consider the following three possible cases for principal strains.
Case 1: When $\mathbf{e}_{1} \neq \mathbf{e}_{2} \neq \mathbf{e}_{3}$. In this case, it is obvious that there exists three mutually orthogonal eigenvectors of the second order real symmetric strain tensor $\mathbf{e}_{\mathrm{ij}}$. These eigenvectors are precisely the three principal directions that are mutually orthogonal.

Case 2: When $\mathbf{e}_{\mathbf{1}} \neq \mathbf{e}_{\mathbf{2}}=\mathbf{e}_{\mathbf{3}}$.
Let $\stackrel{1}{A_{i}}$ and $\stackrel{2}{A_{i}}$ be the corresponding principal orthogonal directions corresponding to strains (distinct) $e_{1}$ and $e_{2}$, respectively. Then

$$
\begin{align*}
& \mathbf{e}_{\mathbf{i j}} \stackrel{1}{A}_{j}=\mathbf{e}_{\mathbf{1}} \stackrel{1}{A_{i}}, \\
& \mathbf{e}_{\mathbf{i j}} \stackrel{2}{A}_{j}^{2}=\mathbf{e}_{\mathbf{2}} \stackrel{2}{A}_{i} \tag{2}
\end{align*}
$$

Let $\mathbf{p}_{\mathbf{i}}$ be a vector orthogonal to both $\stackrel{1}{A_{i}}$ and $\stackrel{2}{A_{i}}$. Then

$$
\begin{equation*}
\mathbf{p}_{\mathbf{i}} \stackrel{1}{A}_{i}=\mathbf{p}_{\mathbf{i}} \stackrel{2}{A}_{i}=\mathbf{0} . \tag{3}
\end{equation*}
$$

Let

$$
\begin{equation*}
\mathbf{e}_{\mathrm{ij}} \mathbf{p}_{\mathrm{i}}=\mathbf{q}_{\mathbf{j}} . \tag{4}
\end{equation*}
$$

Then

$$
\begin{equation*}
\mathbf{q}_{\mathbf{j}} \stackrel{1}{A_{j}}=\left(\mathbf{e}_{\mathbf{i j}} \mathbf{p}_{\mathbf{i}}\right) \stackrel{1}{A_{j}}=\left(\mathbf{e}_{\mathbf{i j}} \stackrel{1}{A_{j}}\right) \mathbf{p}_{\mathbf{i}}=\mathbf{e}_{\mathbf{1}} \stackrel{1}{A_{i}} \mathbf{p}_{\mathbf{i}}=\mathbf{0} \tag{5a}
\end{equation*}
$$

Similarly

$$
\begin{equation*}
\mathbf{q}_{\mathbf{j}} \stackrel{2}{A}_{j}=\mathbf{0} \tag{5b}
\end{equation*}
$$

This shows that the vector $\mathbf{q}_{\mathbf{j}}$ is orthogonal to both $\stackrel{1}{A_{j}}$ and $\stackrel{2}{A}_{j}$. Hence, the vectors $q_{i}$ and $p_{i}$ must be parallel. Let

$$
\begin{equation*}
\mathbf{q}_{i}=\alpha \mathbf{p}_{\mathbf{i}} \tag{6}
\end{equation*}
$$

for some scalar $\alpha$. From equations (4) and (6), we write

$$
\begin{equation*}
\mathbf{e}_{\mathrm{ij}} \mathbf{p}_{\mathrm{j}}=\mathbf{q}_{\mathrm{i}}=\alpha \mathbf{p}_{\mathbf{i}}, \tag{7}
\end{equation*}
$$

which shows that the scalar $\alpha$ is an eigenvalue/principal strain of the strain tensor $\mathrm{e}_{\mathrm{ij}}$ with corresponding principal direction $\mathrm{p}_{\mathrm{i}}$.

Since $\mathbf{e}_{\mathbf{i j}}$ has only three principal strains $\mathbf{e}_{1}, \mathbf{e}_{2}, \alpha$ and two of these are equal, so $\alpha$ must be equal to $e_{2}=e_{3}$.

We denote the normalized form of $\mathbf{p}_{\mathbf{i}}$ by $\stackrel{3}{A_{i}}$.
This shows the existence of three mutually orthogonal principal directions in this case.

Further, let $\mathbf{v}_{\mathbf{i}}$ be any vector normal to $\stackrel{1_{i}}{A_{i}}$. Then $\mathbf{v}_{\mathbf{i}}$ lies in the plane containing principal directions $\stackrel{2}{A_{i}}$ and $\stackrel{3}{A_{i}}$. Let

$$
\begin{equation*}
\mathbf{v}_{\mathbf{i}}=\mathbf{k}_{\mathbf{1}} \stackrel{2}{A_{i}}+\mathbf{k}_{\mathbf{2}} \stackrel{3}{A_{i}} \quad \text { for some scalars } \mathbf{k}_{\mathbf{1}} \text { and } \mathbf{k}_{\mathbf{2}} \tag{8}
\end{equation*}
$$

Now

$$
\begin{aligned}
\mathbf{e}_{\mathbf{i j}} \mathbf{v}_{\mathbf{j}} & =\mathbf{e}_{\mathbf{i j}}\left(\mathbf{k}_{\mathbf{1}} \stackrel{2}{A_{j}}+\mathbf{k}_{\mathbf{2}} \stackrel{3}{A}_{j}\right) \\
& =\mathbf{k}_{\mathbf{1}}\left(\mathbf{e}_{\mathbf{i j}} \stackrel{2}{A}_{j}\right)+\mathbf{k}_{\mathbf{2}}\left(\mathbf{e}_{\mathbf{i j}} \stackrel{3}{A_{j}}\right) \\
& =\mathbf{k}_{\mathbf{1}}\left(\mathbf{e}_{\mathbf{2}} \stackrel{2}{A_{i}}\right)+\mathbf{k}_{\mathbf{2}}\left(\mathbf{e}_{\mathbf{3}} \stackrel{3}{A}_{i}\right) \\
& =\mathbf{e}_{\mathbf{2}}\left[\mathbf{k}_{\mathbf{1}} \stackrel{2}{A_{i}}+\mathbf{k}_{\mathbf{2}} \stackrel{3}{A}_{i}\right] \quad\left(\because \mathbf{e}_{\mathbf{2}}=\mathbf{e}_{\mathbf{3}}\right) \\
& =\mathbf{e}_{\mathbf{2}} \mathbf{v}_{\mathbf{i}}
\end{aligned}
$$

This shows that the direction $\mathbf{v}_{\mathbf{i}}$ is also a principal direction corresponding to principal strain $e_{2}$. Thus, in this case, any two orthogonal(mutually) vectors lying on the plane normal to $\stackrel{1}{A_{i}}$ can be chosen as the other two principal directions. In this case, the strain quadric surface is a surface of revolution.

Case 3: When $\mathbf{e}_{1}=\mathbf{e}_{2}=\mathbf{e}_{3}$, then the strain quadric of Cauchy is a sphere with equation
or

$$
\begin{aligned}
& \mathbf{e}_{1}\left(\mathbf{x}_{1}{ }^{2}+\mathbf{x}_{2}{ }^{2}+\mathbf{x}_{3}{ }^{2}\right)= \pm \mathbf{k}^{2} \\
& \mathbf{x}_{1}{ }^{2}+{\mathbf{x}_{2}}^{2}+\mathbf{x}_{3}{ }^{2}= \pm \frac{k^{2}}{e_{1}}
\end{aligned}
$$

and any three mutually orthogonal directions can be taken as the coordinate axes which are coincident with principal directions of strain.

Hence, the result.

### 3.12 GENERAL INFINITESIMAL DEFORMATION

Now we consider the general functional transformation and its relation to the linear deformation. Consider an arbitrary material point $\mathbf{P}^{\circ}\left(\mathbf{x}_{i}{ }^{\circ}\right)$ in a continuous medium. Let the same material point assume after deformation the point $\mathbf{Q}^{\circ}\left(\xi_{i}{ }^{\circ}\right)$. Then

$$
\begin{equation*}
\xi_{i}{ }^{\circ}=\mathbf{x}_{\mathbf{i}}{ }^{\circ}+\mathbf{u}_{i}\left(\mathbf{x}_{1}{ }^{\circ}, \mathbf{x}_{2}{ }^{\circ}, \mathbf{x}_{3}{ }^{\circ}\right), \tag{1}
\end{equation*}
$$

where $u_{i}$ are the components of the displacement vector $\overline{P^{\circ} Q^{\circ}}$. We assume that $\mathbf{u}_{1}, \mathbf{u}_{\mathbf{2}}, \mathbf{u}_{3}$, as well as their partial derivatives are continuous functions.


Fig. (3.17)
The nature of the deformation in the neighbourhood of the point $\mathrm{P}^{\circ}$ can be determined by considering the change in the vector $\overline{P^{\circ} P}=\mathbf{A}_{\mathbf{i}}$ in the undeformed state, where $P\left(x_{1}, x_{2}, x_{3}\right)$ is an arbitrary neighbouring point of $\mathbf{P}^{\circ}$.

Let $Q\left(\xi_{1}, \xi_{2}, \xi_{3}\right)$ be the deformed position of $P$. Then the displacement $u_{i}$ at the point $P$ is

$$
\begin{equation*}
\mathbf{u}_{\mathbf{i}}\left(\mathbf{x}_{1}, \mathbf{x}_{2}, \mathbf{x}_{3}\right)=\xi_{\mathrm{i}}-\mathbf{x}_{\mathbf{i}} . \tag{2}
\end{equation*}
$$

The vector

$$
\begin{equation*}
\mathbf{A}_{i}=\mathbf{x}_{\mathrm{i}}-\mathbf{x}_{\mathrm{i}}{ }^{\circ}, \tag{3}
\end{equation*}
$$

has now deformed to the vector

$$
\begin{equation*}
\xi_{i}-\xi_{i}^{\circ}=\mathbf{A}_{i}^{\prime}(\mathbf{s a y}) \tag{4}
\end{equation*}
$$

Therefore,

$$
\begin{align*}
& \quad \delta \mathbf{A}_{\mathbf{i}}=\mathbf{A}_{\mathbf{i}}^{\prime}-\mathbf{A}_{\mathbf{i}} \\
& \quad=\left(\xi_{\mathrm{i}}-\xi_{\mathrm{i}}^{\circ}\right)-\left(\mathbf{x}_{\mathbf{i}}-\mathbf{x}_{\mathbf{i}}^{\circ}\right) \\
& =\left(\xi_{\mathrm{i}}-\mathbf{x}_{\mathbf{i}}\right)-\left(\xi_{\mathrm{i}}{ }^{\circ}-\mathbf{x}_{\mathbf{i}}^{\circ}\right) \\
& =\mathbf{u}_{\mathbf{i}}\left(\mathbf{x}_{1}, \mathbf{x}_{\mathbf{2}}, \mathbf{x}_{3}\right)-\mathbf{u}_{\mathbf{i}}\left(\mathbf{x}_{1}{ }^{\circ}, \mathbf{x}_{2}^{\circ}, \mathbf{x}_{3}{ }^{\circ}\right) \\
& =\mathbf{u}_{\mathbf{i}}\left(\mathbf{x}_{1}{ }^{\circ}+\mathbf{A}_{\mathbf{1}}, \mathbf{x}_{2}{ }^{\circ}+\mathbf{A}_{\mathbf{2}}, \mathbf{x}_{3}{ }^{\circ}+\mathbf{A}_{3}\right)-\mathbf{u}_{\mathbf{i}}\left(\mathbf{x}_{1}{ }^{\circ}, \mathbf{x}_{2}^{\circ}, \mathbf{x}_{3}{ }^{\circ}\right) \\
&  \tag{5}\\
& =\left(\frac{\partial u_{i}}{\partial x_{j}}\right)_{0} A_{j},
\end{align*}
$$

plus the higher order terms of Taylor's series. The subscript $\circ$ indicates that the derivative is to be evaluated at the point $\mathrm{P}^{\circ}$.

If the region in the nbd. of $\mathbf{P}^{\circ}$ is chosen sufficiently small, i.e., if the vector $\mathbf{A}_{\mathbf{i}}$ is sufficiently small, then the product terms like $\mathbf{A}_{\mathbf{i}} \mathbf{A}_{\mathbf{j}}$ may be ignored. Ignoring the product terms and dropping the subscript oin (5), we write

$$
\begin{equation*}
\delta \mathbf{A}_{\mathbf{i}}=\mathbf{u}_{\mathbf{i}, \mathbf{j}} \mathbf{A}_{\mathbf{j}} \tag{6}
\end{equation*}
$$

where the symbol $\mathbf{u}_{\mathrm{i}, \mathrm{j}}$ has been used for $\frac{\partial u_{i}}{\partial x_{j}}$.
Result (6) holds for small vectors $\mathbf{A}_{\mathbf{i}}$.
If we further assume that the displacements $u_{i}$ as well as their partial derivatives are so small that their products can be neglected, then the transformation (which is linear) given by (4) becomes infinitesimal in the nbd of the point $\left(\mathrm{P}^{\circ}\right)$ under consideration and

$$
\begin{equation*}
\delta \mathbf{A}_{\mathbf{i}}=\alpha_{\mathrm{ij}} \mathbf{A}_{\mathbf{j}} \tag{7}
\end{equation*}
$$

with

$$
\begin{equation*}
\alpha_{i j}=\mathbf{u}_{\mathrm{i}, \mathrm{j}} \tag{8}
\end{equation*}
$$

Hence, all results discussed earlier are immediately applicable.
The transformation (6) can be splitted into pure deformation and rigid body motion as

$$
\begin{align*}
\delta \mathbf{A}_{\mathbf{i}}=\mathbf{u}_{\mathbf{i}, \mathbf{j}} \mathbf{A}_{\mathbf{j}} & =\left(\frac{u_{i, j}+u_{j, i}}{2}+\frac{u_{i, j}-u_{j, i}}{2}\right) \mathbf{A}_{\mathbf{j}} \\
& =\mathbf{e}_{\mathbf{i j}} \mathbf{A}_{\mathbf{j}}+\mathbf{w}_{\mathbf{i j}} \mathbf{A}_{\mathbf{j}}, \tag{9}
\end{align*}
$$

where $\quad \mathbf{e}_{\mathrm{ij}}=\frac{1}{2}\left(\mathbf{u}_{\mathrm{i}, \mathrm{j}}+\mathbf{u}_{\mathrm{j}, \mathrm{i}}\right)$,

$$
\begin{equation*}
\mathbf{w}_{\mathbf{i j}}=\frac{1}{2}\left(\mathbf{u}_{\mathbf{i}, \mathbf{j}}-\mathbf{u}_{\mathbf{j}, \mathbf{i}}\right) . \tag{11}
\end{equation*}
$$

The transformation

$$
\begin{equation*}
\delta \mathbf{A}_{\mathbf{i}}=\mathbf{e}_{\mathbf{i j}} \mathbf{A}_{\mathbf{j}} \tag{12}
\end{equation*}
$$

represents pure deformation and

$$
\begin{equation*}
\delta \mathbf{A}_{\mathbf{i}}=\mathbf{w}_{\mathrm{ij}} \mathbf{A}_{\mathbf{j}}, \tag{13}
\end{equation*}
$$

represents rotation. In general , the transformation (9) is no longer homogeneous as both the strain components $e_{i j}$ and components of rotation $\omega_{i j}$ are functions of the coordinates. We find

$$
\begin{equation*}
v=\mathrm{e}_{\mathrm{ii}}=\frac{\partial u_{i}}{\partial x_{i}}=\mathbf{u}_{\mathrm{i}, \mathrm{i}}=\operatorname{div} \overline{\mathbf{u}} . \tag{14}
\end{equation*}
$$

That is , the cubic dilatation is the divergence of the displacement vector $\overline{\mathbf{u}}$ and it differs,
in general , from point to point of the body.
The rotation vector $w_{i}$ is given by

$$
\begin{equation*}
\mathbf{w}_{1}=\mathbf{w}_{32}, w_{2}=w_{13}, w_{3}=w_{21} . \tag{15}
\end{equation*}
$$

Question : For the small linear deformation given by

$$
\overline{\mathbf{u}}=\alpha \mathbf{x}_{1} \mathbf{x}_{2}\left(\hat{e}_{1}+\hat{e}_{2}\right)+2 \alpha\left(\mathbf{x}_{1}+\mathbf{x}_{2}\right) \mathbf{x}_{3} \hat{e}_{3}, \alpha=\text { constant. }
$$

find the strain tensor , the rotation and the rotation vector.
Solution: We find

$$
\begin{aligned}
& \mathbf{u}_{1}=\alpha \mathbf{x}_{1} \mathbf{x}_{2}, \mathbf{u}_{2}=\alpha \mathbf{x}_{1} \mathbf{x}_{2}, \mathbf{u}_{3}=\mathbf{2} \alpha\left(\mathbf{x}_{1}+\mathbf{x}_{2}\right) \mathbf{x}_{3}, \\
& \mathbf{e}_{11}=\frac{\partial u_{1}}{\partial x_{1}}=\alpha \mathbf{x}_{2}, \mathbf{e}_{22}=\frac{\partial u_{2}}{\partial x_{2}}=\alpha \mathbf{x}_{1}, \mathbf{e}_{33}=\frac{\partial u_{3}}{\partial x_{3}}=\mathbf{2} \alpha\left(\mathbf{x}_{1}+\mathbf{x}_{2}\right)
\end{aligned}
$$

$$
\begin{aligned}
& \mathbf{e}_{12}=\frac{1}{2}\left(\frac{\partial u_{1}}{\partial x_{1}}+\frac{\partial u_{2}}{\partial x_{2}}\right)=\frac{\alpha}{2}\left(\mathbf{x}_{1}+\mathbf{x}_{2}\right), \\
& \mathbf{e}_{13}=\frac{1}{2}\left(\frac{\partial u_{1}}{\partial x_{3}}+\frac{\partial u_{3}}{\partial x_{1}}\right)=\alpha \mathbf{x}_{3}, \mathbf{e}_{23}=\alpha \mathbf{x}_{3}
\end{aligned}
$$

## Hence

$$
\left(\mathbf{e}_{\mathbf{i j}}\right)=\alpha\left[\begin{array}{ccc}
x_{2} & \left(x_{1}+x_{2}\right) / 2 & x_{3} \\
\left(x_{1}+x_{2}\right) / 2 & x_{1} & x_{3} \\
x_{3} & x_{3} & 2\left(x_{1}+x_{2}\right)
\end{array}\right] .
$$

## We know that

$$
\mathbf{w}_{\mathrm{ij}}=\frac{1}{2}\left(\frac{\partial \mathrm{u}_{\mathrm{i}}}{\partial \mathrm{x}_{\mathrm{j}}}-\frac{\partial \mathrm{u}_{\mathrm{j}}}{\partial \mathrm{x}_{\mathrm{i}}}\right)
$$

## We find

$$
\begin{aligned}
& \mathbf{w}_{11}=\mathbf{w}_{22}=\mathbf{w}_{33}=\mathbf{0}, \\
& \mathbf{w}_{12}=\frac{\alpha}{2}\left[\mathbf{x}_{1}-\mathbf{x}_{2}\right]=-\mathbf{w}_{21}, \\
& \mathbf{w}_{13}=-\alpha \mathbf{x}_{3}=-\mathbf{w}_{31}, \\
& \mathbf{w}_{23}=-\alpha \mathbf{x}_{3}=-\mathbf{w}_{32},
\end{aligned}
$$

Therefore

$$
\left(\mathbf{w}_{\mathbf{i j}}\right)=\alpha\left[\begin{array}{ccc}
0 & \left(x_{1}-x_{2}\right) / 2 & -x_{3}  \tag{2}\\
-\left(x_{1}-x_{2}\right) / 2 & 0 & -x_{3} \\
x_{3} & x_{3} & 0
\end{array}\right]
$$

The rotation vector $\overline{\mathbf{w}}=\mathbf{w}_{\mathbf{i}}$ is given by $\mathbf{w}_{\mathbf{i}}=\epsilon_{\mathrm{ijk}} \mathbf{u}_{\mathbf{k}, \mathbf{j}}$. We find

$$
\begin{gathered}
\mathbf{w}_{1}=\mathbf{w}_{32}=\alpha \mathbf{x}_{3}, \mathbf{w}_{2}=\mathbf{w}_{13}=-\alpha \mathbf{x}_{3}, \mathbf{w}_{3}=\mathbf{w}_{21}=\frac{\alpha}{2}\left(\mathbf{x}_{2}-\mathbf{x}_{1}\right) . \\
\overline{\mathbf{w}}=\alpha \mathbf{x}_{3}\left(\hat{e}_{1}-\hat{e}_{2}\right)+\frac{\alpha}{2}\left(\mathbf{x}_{2}-\mathbf{x}_{1}\right) \hat{e}_{3} .
\end{gathered}
$$

So

Exercise 1: For small deformation defined by the following displacements, find the strain tensor, rotation tensor and rotation vector.
(i) $\mathbf{u}_{1}=-\alpha \mathbf{x}_{2} \mathbf{x}_{3}, \mathbf{u}_{2}=\alpha \mathbf{x}_{1} \mathbf{x}_{2}, \mathbf{u}_{3}=\mathbf{0}$.
(ii) $\mathbf{u}_{1}=\alpha^{2}\left(\mathbf{x}_{1}-\mathbf{x}_{3}\right)^{2}, \mathbf{u}_{2}=\alpha^{2}\left(\mathbf{x}_{2}+\mathbf{x}_{3}\right)^{2}, u_{3}=-\alpha \mathbf{x}_{1} \mathbf{x}_{2}, \alpha=$ constant.

Exercise 2: The displacement components are given by

$$
\mathbf{u}=-\mathbf{y z}, \mathbf{v}=\mathrm{xz}, \mathbf{w}=\phi(\mathbf{x}, \mathbf{y})
$$

calculate the strain components.
Exercise 3: Given the displacements

$$
u=3 x^{2} y, v=y^{2}+6 x z, w=6 z^{2}+2 y z,
$$

calculate the strain components at the point $(1,0,2)$. What is the extension of a line element (parallel to the $x$-axis) at this point?

Exercise 4: Find the strain components and rotation components for the small displacement components given below
(a) Uniform dilatation -

$$
\begin{aligned}
& \mathbf{u}=\mathbf{e} \mathbf{x}, \mathbf{v}=\mathbf{e} \mathbf{y}, \mathbf{w}=\mathbf{e} \mathbf{z} \\
& \mathbf{u}=\mathbf{e} \mathbf{x}, \mathbf{v}=\mathbf{w}=\mathbf{0}, \\
& \mathbf{u}=\mathbf{2} \mathbf{s y}, \mathbf{v}=\mathbf{w}=\mathbf{0}, \\
& \mathbf{u}=\mathbf{u}(\mathbf{x}, \mathbf{y}), \mathbf{v}=\mathbf{v}(\mathbf{x}, \mathbf{y}), \mathbf{w}=\mathbf{0} .
\end{aligned}
$$

(c) Shearing strain -
(d) Plane strain -

### 3.13 SAINT-VENANT'S EQUATIONS OF COMPATIBILITY

By definition , the strain components $\mathrm{e}_{\mathrm{ij}}$ in terms of displacement components $\mathbf{u}_{\mathbf{i}}$ are given by

$$
\begin{equation*}
\mathrm{e}_{\mathrm{ij}}=\frac{1}{2}\left[\mathrm{u}_{\mathrm{i}, \mathrm{j}}+\mathrm{u}_{\mathrm{j}, \mathrm{i}}\right] \tag{1}
\end{equation*}
$$

Equation (1) is used to find the components of strain if the components of displacement are given. However, if the components of strain, $\mathrm{e}_{\mathrm{ij}}$, are given then equation (1) is a set of 6 partial differential equations in the three unknowns $\mathbf{u}_{\mathbf{1}}, \mathbf{u}_{\mathbf{2}}, \mathbf{u}_{\mathbf{3}}$. Therefore, the system (1) will not have a single - valued solution for $u_{i}$ unless given strains $e_{i j}$ satisfy certain conditions which are known as the conditions of compatibility or equations of compatibility.

Geometrical meaning of Conditions of Compatibility

$$
\mathrm{R} \quad \mathrm{R}^{\prime}
$$



Fig (1) shows a portion of the material of the body in the undeformed state in the form of a continuous triangle ABC. If we deform the body by an arbitrarily specified strain field then we might end up at the points $C^{\prime}$ and $D^{\prime}$ with a gap between them, after deformation, as shown in fig(2a) or with overlapping material as shown in fig(2b).

For a single valued continuous solution to exist the points $C^{\prime}$ and $D^{\prime}$ must be the same in the strained state. This cannot be guaranteed unless the specified strain components satisfy certain conditions, known as the conditions (or relations or equations) of compatibility.

Equations of Compatibility
$\begin{array}{ll}\text { We have } & \mathbf{e}_{\mathrm{ij}}=\frac{1}{2}\left(\mathbf{u}_{\mathrm{i}, \mathrm{j}}+\mathbf{u}_{\mathrm{j}, \mathrm{j}}\right) . \\ \text { So , } & \mathbf{e}_{\mathrm{ij}, \mathrm{k} l}=\frac{1}{2}\left(\mathbf{u}_{\mathrm{i}, \mathrm{j} k l}+\mathbf{u}_{\mathrm{j}, \mathrm{i} k l}\right) .\end{array}$
Interchanging i with k and j with $l$ in equation (2), we write

$$
\begin{equation*}
\mathbf{e}_{\mathrm{k} l, \mathrm{ij}}=\frac{1}{2}\left(\mathbf{u}_{\mathrm{k}, \mathrm{ij}}+\mathbf{u}_{l, \mathrm{kj}}\right) . \tag{3}
\end{equation*}
$$

Adding (2) and (3), we get

$$
\begin{equation*}
\mathbf{e}_{\mathrm{ij}, \mathrm{k} l}+\mathbf{e}_{\mathrm{k} l, \mathrm{j}}=\frac{1}{2}\left[\mathbf{u}_{\mathrm{i}, \mathrm{j} l}+\mathbf{u}_{\mathrm{j}, \mathrm{k} l}+\mathbf{u}_{\mathrm{k}, \mathrm{ij}}+\mathbf{u}_{l, \mathrm{kj}}\right] \tag{4}
\end{equation*}
$$

Interchanging i and $l$ in (4), we get

From (4) and (5), we obtain
or

$$
\begin{align*}
& \mathbf{e}_{\mathrm{ij}, \mathrm{k} l}+\mathbf{e}_{\mathrm{k} l, \mathrm{j}}=\mathbf{e}_{l \mathrm{j}, \mathrm{ki}}+\mathbf{e}_{\mathrm{k}, \mathrm{l}, \mathrm{j}}, \\
& \mathbf{e}_{\mathrm{ij}, \mathrm{k} l}+\mathbf{e}_{\mathrm{k} l, \mathrm{j}}-\mathbf{e}_{\mathrm{ik}, \mathrm{j} l}-\mathbf{e}_{\mathrm{j} l, \mathrm{i}}=\mathbf{0} . \tag{6}
\end{align*}
$$

These equations are known as equations of compatibility.
These equations are necessary conditions for the existence of a single valued continuous displacement field. These are 81 equations in number. Because of symmetry in indicies $i$, $j$, and $k, l$; some of these equations are identically satisfied and some are repetitions.

Only 6 out of these 81 equations are essential. These equations were first obtained by Saint - Venant in 1860.

A strain tensor $\mathrm{e}_{\mathrm{ij}}$ that satisfies these conditions is referred to as a possible strain tensor.

Show that the conditions of compatibility are sufficient for the existence of a single valued continuous displacement field.
Let $P^{\circ}\left(x_{i}{ }^{\circ}\right)$ be some point of a simply connected region at which the displacements $u_{i}{ }^{\circ}$ and rotations $w_{i j}{ }^{\circ}$ are known. The displacements $u_{i}$ of an arbitrary point $P^{\prime}\left(x_{i}{ }^{\prime}\right)$ can be obtained in terms of the known functions $\mathrm{e}_{\mathrm{ij}}$ by means of a line integral along a continuous curve $\mathbf{C}$ joining the points $\mathbf{P}_{0} \quad$ and $\mathbf{P}^{\prime}$.


Fig. (3.18)

$$
\begin{equation*}
\mathbf{u}_{\mathrm{j}}\left(\mathbf{x}_{1} \boldsymbol{1}^{\prime}, \mathbf{x}_{2}{ }^{\prime}, \mathbf{x}_{3}{ }^{\prime}\right)=\mathbf{u}_{\mathrm{j}}{ }^{\circ}\left(\mathbf{x}_{1}{ }^{\circ}, \mathbf{x}_{2}{ }^{\circ}, \mathbf{x}_{3}{ }^{\circ}\right)+\int_{P^{\circ}}^{P^{\prime}} d u_{j} \cdot \tag{7}
\end{equation*}
$$

If the process of deformation does not create cracks or holes, i.e., if the body remains continuous, the displacements $\mathbf{u}_{j}{ }^{\prime}$ should be independent of the path of integration.

That is , $u_{j}{ }^{\prime}$ should have the same value regardless of whether the integration is along curve $C$ or any other curve. We write

$$
\begin{equation*}
\mathbf{d} \mathbf{u}_{\mathbf{j}}=\frac{\partial u_{j}}{\partial x_{k}} \mathbf{d} \mathbf{x}_{k}=\mathbf{u}_{\mathbf{j}, \mathbf{k}} \mathbf{d} \mathbf{x}_{\mathrm{k}}=\left(\mathbf{e}_{\mathbf{j k}}+\mathbf{w}_{\mathbf{j k}}\right) \mathbf{d} \mathbf{x}_{\mathbf{k}} . \tag{8}
\end{equation*}
$$

Therefore,
$\mathbf{u}_{\mathbf{j}}{ }^{\prime}=\mathbf{u}_{\mathbf{j}}{ }^{\circ}+\int_{P^{\circ}}^{P^{\prime}} \mathbf{e}_{\mathbf{j} \mathbf{k}} \mathbf{d} \mathbf{x}_{\mathbf{k}}+\quad \int_{P^{\circ}}^{P^{\prime}} \mathbf{w}_{\mathbf{j k}} \mathbf{d} \mathbf{x}_{\mathbf{k}}, \quad\left(\mathbf{P}\left(\mathbf{x}_{\mathbf{k}}\right)\right.$ being point the joining curve).

Integrating by parts the second integral, we write

$$
\begin{align*}
\int_{P^{\circ}}^{P^{\prime}} \mathbf{w}_{\mathbf{j k}} \mathbf{d} \mathbf{x}_{\mathbf{k}} & =\int_{P^{\circ}}^{P^{\prime}} \mathbf{w}_{\mathbf{j k}} \mathbf{d}\left(\mathbf{x}_{\mathbf{k}}-\mathbf{x}_{\mathbf{k}}{ }^{\prime}\right) \quad\left(\text { The point } \mathbf{P}^{\prime}\left(\mathbf{x}_{\mathbf{k}}{ }^{\prime}\right) \text { being fixed so } \mathbf{d}_{\mathbf{x}_{\mathbf{k}}}=\mathbf{0}\right) \\
& =\left\{\left(\mathbf{x}_{\mathbf{k}}-\mathbf{x}_{\mathbf{k}}{ }^{\prime}\right) \mathbf{w}_{\mathbf{j k}}{ }_{P_{-}^{\circ}}^{P^{\circ}}-\int_{P^{\circ}}^{P^{\prime}}\left(\mathbf{x}_{\mathbf{k}}-\mathbf{x}_{\mathbf{k}}{ }^{\prime}\right) \mathbf{w}_{\mathbf{j} \mathbf{k}, l} \mathbf{d x}_{l}\right. \\
& =\left(\mathbf{x}_{\mathbf{k}}{ }^{\prime}-\mathbf{x}_{\mathbf{k}}{ }^{\circ}\right) \mathbf{w}_{\mathbf{j k}}{ }^{\circ}+\int_{P^{\circ}}^{P^{\prime}}\left(\mathbf{x}_{\mathbf{k}}^{\prime}-\mathbf{x}_{\mathbf{k}}\right) \mathbf{w}_{\mathbf{j k}, l} \mathbf{d} \mathbf{x}_{l} \tag{10}
\end{align*}
$$

From equations (9) and (10), we write

$$
\begin{gather*}
\mathbf{u}_{\mathbf{j}}\left(\mathbf{x}_{\mathbf{1}}{ }^{\prime}, \mathbf{x}_{\mathbf{2}}{ }^{\prime}, \mathbf{x}_{\mathbf{3}}{ }^{\prime}\right)=\mathbf{u}_{\mathbf{j}}{ }^{\circ}+\left(\mathbf{x}_{\mathbf{k}}{ }^{\prime}-\mathbf{x}_{\mathbf{k}}{ }^{\circ}\right) \mathbf{w}_{\mathbf{j} \mathbf{k}}{ }^{\mathrm{o}}+\int_{P^{\circ}}^{P^{\prime}} \mathbf{e}_{\mathbf{j k}} \mathbf{d} \mathbf{x}_{\mathbf{k}}+\int_{P^{\circ}}^{P^{\prime}}\left(\mathbf{x}_{\mathbf{k}}{ }^{\prime}-\mathbf{x}_{\mathbf{k}}\right) \mathbf{w}_{\mathbf{j k}, l} \mathbf{d} \mathbf{x}_{l} \\
=\mathbf{u}_{\mathbf{j}}{ }^{\circ}+\left(\mathbf{x}_{\mathbf{k}}{ }^{\prime}-\mathbf{x}_{\mathbf{k}}{ }^{\circ}\right) \mathbf{w}_{\mathbf{j k}}{ }^{\circ}+\int_{P^{\circ}}^{P^{\prime}}\left[\mathbf{e}_{\mathbf{j} \mathbf{l}}+\left(\mathbf{x}_{\mathbf{k}}{ }^{\prime}-\mathbf{x}_{\mathbf{k}}\right) \mathbf{w}_{\mathbf{j k}, l}\right] \mathbf{d} \mathbf{x}_{l} \tag{11}
\end{gather*}
$$

where the dummy index k of $\mathrm{e}_{\mathbf{j k}}$ has been changed to $\boldsymbol{l}$.
But $\quad \mathbf{w}_{\mathbf{j k}, l}=\frac{1}{2} \frac{\partial}{\partial x_{l}}\left[\mathbf{u}_{\mathbf{j}, \mathbf{k}}-\mathbf{u}_{\mathbf{k}, \mathbf{j}}\right]$

$$
=\frac{1}{2}\left[\mathbf{u}_{\mathrm{j}, \mathrm{k} l}-\mathbf{u}_{\mathrm{k}, \mathrm{j} l}\right]
$$

$$
=\frac{1}{2}\left[\mathbf{u}_{\mathrm{j}, \mathbf{k} l}+\mathbf{u}_{l, \mathbf{j} \mathbf{k}}\right]-\frac{1}{2}\left[\mathbf{u}_{l, \mathrm{jk}}+\mathbf{u}_{\mathrm{k}, \mathrm{j} l}\right]
$$

$$
\begin{equation*}
=\mathbf{e}_{\mathrm{j} l, \mathbf{k}}-\mathbf{e}_{l \mathbf{k}, \mathbf{j}} \tag{12}
\end{equation*}
$$

Using (12), equation (11) becomes
$\mathbf{u}_{\mathbf{j}}\left(\mathbf{x}_{1}{ }^{\prime}, \mathbf{x}_{2}{ }^{\prime}, \mathbf{x}_{3}{ }^{\prime}\right)=\mathbf{u}_{\mathbf{j}}{ }^{\circ}+\left(\mathbf{x}_{\mathbf{k}}{ }^{\prime}-\mathbf{x}_{\mathbf{k}}{ }^{\circ}\right) \mathbf{w}_{\mathbf{j} \mathbf{k}}{ }^{\circ}+\int_{P^{\circ}}^{P^{\prime}}\left[\mathbf{e}_{\mathrm{j} l}+\left\{\mathbf{x}_{\mathbf{k}}{ }^{\prime}-\mathbf{x}_{\mathbf{k}}\right\}\left\{\mathbf{e}_{\mathrm{j} l, \mathbf{k}}-\mathbf{e}_{\mathrm{k} l \mathrm{j}}\right\}\right] \mathbf{d} \mathbf{x}_{l}$

$$
\begin{equation*}
=\mathbf{u}_{\mathbf{j}}{ }^{\circ}+\left(\mathbf{x}_{\mathbf{k}}{ }^{\prime}-\mathbf{x}_{\mathbf{k}}{ }^{\circ}\right) \mathbf{w}_{\mathbf{j} \mathbf{k}}{ }^{\circ}+\int_{P^{\circ}}^{P^{\prime}} \mathbf{U}_{\mathbf{j} l} \mathbf{d} \mathbf{x}_{l} \tag{13}
\end{equation*}
$$

where for convenience we have set

$$
\begin{equation*}
\mathbf{U}_{\mathbf{j} l}=\mathbf{e}_{j l}+\left(\mathbf{x}_{\mathbf{k}}{ }^{\prime}-\mathbf{x}_{\mathbf{k}}\right)\left(\mathbf{e}_{\mathrm{j} l, \mathrm{k}}-\mathbf{e}_{\mathbf{k} l \mathbf{j}}\right) \tag{14}
\end{equation*}
$$

which is a known function as $\mathrm{e}_{\mathrm{ij}}$ are known.
The first two terms in the right side of equation (13) are independent of the path of integration. From the theory of line integrals, the third term become independent of the path of integration when the integrands $U_{j 1} d_{x_{1}}$ must be exact differentials.

Therefore, if the displacements $\mathbf{u}_{\mathrm{i}}\left(\mathrm{x}_{1}{ }^{\prime}, \mathrm{x}_{2}{ }^{\prime}, \mathrm{x}_{3}{ }^{\prime}\right)$ are to be independent of the path of integration, we must have

$$
\begin{equation*}
\frac{\partial U_{j l}}{\partial x_{l}}=\frac{\partial U_{j i}}{\partial x_{l}}, \quad \text { for } \mathbf{i}, \mathbf{j}, l=\mathbf{1}, \mathbf{2}, \mathbf{3} . \tag{15}
\end{equation*}
$$

Now

$$
\begin{align*}
& \mathbf{U}_{j l, i}=\mathbf{e}_{j l, i}+\left(\mathbf{x}_{\mathbf{k}}{ }^{\prime}-\mathbf{x}_{\mathbf{k}}\right)\left(\mathbf{e}_{\mathrm{j} l, \mathrm{ki}}-\mathbf{e}_{\mathrm{k} l, \mathrm{j} \boldsymbol{i}}\right)-\delta_{\mathrm{ki}}\left(\mathbf{e}_{\mathrm{j} l, \mathrm{k}}-\mathbf{e}_{\mathrm{k} l, \mathrm{j}}\right) \\
& =\mathbf{e}_{j l, i}-\mathbf{e}_{j l, i}+\mathbf{e}_{\mathrm{i} l, \mathrm{j}}+\left(\mathbf{x}_{\mathbf{k}}{ }^{\prime}-\mathbf{x}_{\mathrm{k}}\right)\left(\mathbf{e}_{\mathrm{j}, \mathrm{ki}}-\mathbf{e}_{\mathrm{k} l, \mathrm{j}}\right), \tag{16a}
\end{align*}
$$

and

$$
\begin{align*}
\mathbf{U}_{\mathbf{j}, l} & =\mathbf{e}_{\mathbf{j}, l}+\left(\mathbf{x}_{\mathbf{k}}{ }^{\prime}-\mathbf{x}_{\mathbf{k}}\right)\left(\mathbf{e}_{\mathbf{j}, \mathrm{k}, \mathrm{l}}-\mathbf{e}_{\mathbf{k}, j l}\right)-\delta_{\mathbf{k} l}\left(\mathbf{e}_{\mathbf{j}, \mathrm{k}}-\mathbf{e}_{\mathbf{k}, \mathrm{j}}\right) \\
& =\mathbf{e}_{\mathbf{j}, l}-\mathbf{e}_{\mathbf{j}, l}+\mathbf{e}_{l i, j}+\left(\mathbf{x}_{\mathbf{k}}{ }^{\prime}-\mathbf{x}_{\mathbf{k}}\right)\left(\mathbf{e}_{\mathbf{j} \mathbf{j}, k l}-\mathbf{e}_{\mathbf{k}, \mathrm{j}, l}\right) \tag{16b}
\end{align*}
$$

Therefore, equations (15) and (16 a , b) yields

$$
\left(\mathbf{x}_{\mathbf{k}}^{\prime}-\mathbf{x}_{\mathbf{k}}\right)\left[\mathbf{e}_{\mathrm{j} l, \mathrm{ki}}-\mathbf{e}_{\mathbf{k}, \mathrm{ji}}-\mathbf{e}_{\mathrm{j}, \mathrm{k}, \mathrm{l}}+\mathbf{e}_{\mathbf{k}, \mathrm{j}, \mathrm{l}}\right]=\mathbf{0} .
$$

Since this is true for an arbitrary choice of $\mathbf{x}_{\mathbf{k}}{ }^{\prime}-\mathrm{x}_{\mathrm{k}}$ (as $\mathbf{P}^{\prime}$ is arbitrary), it follows that

$$
\begin{equation*}
\mathbf{e}_{\mathbf{i j}, k l}+\mathbf{e}_{\mathbf{k} l, \mathrm{j}}-\mathbf{e}_{\mathrm{ik}, \mathbf{j} l}-\mathbf{e}_{\mathbf{j} l, \mathbf{k}}=\mathbf{0} \tag{17}
\end{equation*}
$$

which is true as these are the compatibility relations.
Hence, the displacement (7) independent of the path of Integration. Thus , the compatibility conditions (6) are sufficient also.

Remark 1: The compatibility conditions (6) are necessary and sufficient for the existence of a single valued continuous displacement field when the strain components are prescribed.

In detailed form, these 6 conditions are

$$
\begin{aligned}
& \frac{\partial^{2} e_{11}}{\partial x_{2} \partial x_{3}}=\frac{\partial}{\partial x_{1}}\left(\frac{-\partial e_{23}}{\partial x_{1}}+\frac{\partial e_{31}}{\partial x_{2}}+\frac{\partial e_{12}}{\partial x_{3}}\right) \\
& \frac{\partial^{2} e_{22}}{\partial x_{3} \partial x_{1}}=\frac{\partial}{\partial x_{2}}\left(\frac{-\partial e_{31}}{\partial x_{2}}+\frac{\partial e_{12}}{\partial x_{3}}+\frac{\partial e_{23}}{\partial x_{1}}\right) \\
& \frac{\partial^{2} e_{33}}{\partial x_{1} \partial x_{2}}=\frac{\partial}{\partial x_{3}}\left(\frac{-\partial e_{12}}{\partial x_{3}}+\frac{\partial e_{23}}{\partial x_{1}}+\frac{\partial e_{31}}{\partial x_{2}}\right) \\
& \frac{2 \partial^{2} e_{12}}{\partial x_{1} \partial x_{2}}=\frac{\partial^{2} e_{11}}{\partial x_{2}^{2}}+\frac{\partial^{2} e_{22}}{\partial x_{1}^{2}} \\
& \frac{2 \partial^{2} e_{23}}{\partial x_{2} \partial x_{3}}=\frac{\partial^{2} e_{22}}{\partial x_{3}^{2}}+\frac{\partial^{2} e_{33}}{\partial x_{2}^{2}} \\
& \frac{2 \partial^{2} e_{31}}{\partial x_{3} \partial x_{1}}=\frac{\partial^{2} e_{33}}{\partial x_{1}^{2}}+\frac{\partial^{2} e_{11}}{\partial x_{3}^{2}}
\end{aligned}
$$

These are the necessary and sufficient conditions for the strain components $e_{i j}$ to give single valued displacements $u_{i}$ for a simply connected region.

Definition : A region space is said to be simply connected if an arbitrary closed curve lying in the region can be shrunk to a point, by continuous deformation, without passing outside of the boundaries.

Remark 2: The specification of the strains $\mathbf{e}_{\mathbf{i j}}$ only does not determine the displacements $u_{i}$ uniquely because the strains $\mathrm{e}_{\mathrm{ij}}$ characterize only the pure deformation of an elastic medium in the neighbourhood of the point $x_{i}$.

The displacements $\mathbf{u}_{\mathrm{i}}$ may involve rigid body motions which do not affect $\mathbf{e}_{\mathrm{ij}}$.

Example 1 : (i) Find the compatibility condition for the strain tensor $\mathbf{e}_{\mathbf{i j}}$ if $e_{11}, e_{22}, e_{12}$ are independent of $x_{3}$ and $e_{31}=e_{32}=e_{33}=0$.
(ii) Find the condition under which the following are possible strain components.

$$
\begin{aligned}
& e_{11}=k\left(x_{1}{ }^{2}-x_{2}^{2}\right), e_{12}=k^{\prime} x_{1} x_{2}, e_{22}=k x_{1} x_{2}, \\
& e_{31}=e_{32}=e_{33}=0, k \& k^{\prime} \text { are constants }
\end{aligned}
$$

(iii) when $\mathrm{e}_{\mathrm{ij}}$ given above are possible strain components, find the corresponding displacements, given that $\mathbf{u}_{3}=0$.

Solution: (i) We verify that all the compatibility conditions except one are obviously satisfied. The only compatibility condition to be satisfied by $\mathbf{e}_{\mathbf{i j}}$ is

$$
\begin{equation*}
\mathbf{e}_{11,22}+\mathbf{e}_{22,11}=2 \mathbf{e}_{12,12} . \tag{1}
\end{equation*}
$$

(ii) Five conditions are trivially satisfied. The remaining condition (1) is satisfied iff
as

$$
\begin{aligned}
& k^{\prime}=k \\
& e_{1122}=-2 k, e_{12,12}=k^{\prime}, e_{22,11}=0 .
\end{aligned}
$$

(iii) We find

$$
\begin{aligned}
& \mathbf{e}_{11}=\mathbf{u}_{1,1}=\mathbf{k}\left(\mathbf{x}_{1}{ }^{2}-\mathbf{x}_{2}{ }^{2}\right), \mathbf{u}_{2,2}=\mathbf{k} \mathbf{x}_{1} \mathbf{x}_{2}, \mathbf{u}_{1,2}+\mathbf{u}_{2,1}=-2 \mathbf{k} \mathbf{x}_{1} \mathbf{x}_{2}, \\
& \left(\because \mathbf{k}^{\prime}=-\mathbf{k}\right) \\
& \qquad \mathbf{u}_{2,3}=\mathbf{u}_{1,3}=\mathbf{0} .
\end{aligned}
$$

This shows that the displacement components $\mathbf{u}_{1}$ and $\mathbf{u}_{\mathbf{2}}$ are independent of $x_{3}$. We find (exercise)

$$
\begin{aligned}
& \mathbf{u}_{1}=\frac{1}{6}\left(2 \mathrm{x}_{1}{ }^{3}-6 \mathrm{x}_{1}{\left.\mathbf{x}_{2}{ }^{2}+\mathrm{x}_{2}{ }^{3}\right)-\mathbf{c} \mathbf{x}_{2}+\mathrm{c}_{1},}^{\mathbf{u}_{2}=\frac{1}{2} k \mathbf{x}_{1} \mathbf{x}_{2}{ }^{2}+\mathbf{c} \mathbf{x}_{1}+\mathrm{c}_{2}},\right.
\end{aligned}
$$

where $c_{1}, c_{2}$ and $c$ are constants.
Example 2: Show that the following are not possible strain components

$$
\begin{aligned}
& \mathbf{e}_{11}=k\left(\mathbf{x}_{1}{ }^{2}+x_{2}{ }^{2}\right), e_{22}=k\left(\mathbf{x}_{2}{ }^{2}+\mathbf{x}_{3}{ }^{2}\right), e_{33}=0 \\
& \mathbf{e}_{12}=k^{\prime} \mathbf{x}_{1} \mathbf{x}_{2} \mathbf{x}_{3}, \mathbf{e}_{13}=\mathbf{e}_{23}=0, k{ }^{\&} k^{\prime} \text { being constants. }
\end{aligned}
$$

Solution : The given components $\mathbf{e}_{\mathbf{i j}}$ are possible strain components if each of the six compatibility conditions is satisfied. On substitution, we find

$$
2 k=2 k^{\prime} x_{3} .
$$

This can't be satisfied for $x_{3} \neq 0$.
For $\mathbf{x}_{3}=\mathbf{0}$, this gives $\mathrm{k}=0$ and then all $\mathrm{e}_{\mathrm{ij}}$ vanish.
Hence, the given $\mathrm{e}_{\mathrm{ij}}$ are not possible strain components.
Exercise 1: Consider a linear strain field associated with a simply connected region $R$ such that $\quad e_{11}=A x_{2}{ }^{2}, e_{22}=A x_{1}{ }^{2}, e_{12} B \mathbf{x}_{1} \mathbf{x}_{2}$, $\mathbf{e}_{13}=\mathbf{e}_{23}=\mathbf{e}_{33}=0$. Find the relationship between constant $A$ and $B$ such that it is possible to obtain a single - valued continuous displacement field which corresponds to the given strain field.

Exercise 2: Show by differentiation of the strain displacement relations that the compatibility conditions are necessary conditions for the existence of continuous single - valued displacements.

Exercise 3: Is the following state of strain possible? (c = constant)
$e_{11}=c\left(x_{1}{ }^{2}+x_{2}{ }^{2}\right) x_{3}, e_{22}=c{x_{2}}^{2} x_{3}, e_{12}=2 c x_{1} x_{2} x_{3}, e_{31}=e_{32}=e_{33}=0$.
Exercise 4: Show that the equations of compatibility represents a set of necessary and sufficient conditions for the existence of single - valued displacements. Derive the equations of compatibility for plane strain.

Exercise 5: If $\mathbf{e}_{\mathbf{1 1}}=\mathbf{e}_{\mathbf{2 2}}=\mathbf{e}_{\mathbf{1 2}}=\mathbf{e}_{33}=\mathbf{0}, \mathbf{e}_{13}=\phi_{, 2}$ and $\mathbf{e}_{23}=\phi_{, 1}$; where $\phi$ is a function of $x_{1}$ and $x_{2}$, show that $\phi$ must satisfy the equation

$$
\nabla^{2} \phi=\text { constant }
$$

Exercise 6: If $\mathbf{e}_{13}$ and $\mathbf{e}_{23}$ are the only non - zero strain components and $\mathbf{e}_{\mathbf{1 3}}$ and $e_{23}$ are independent of $x_{3}$, show that the compatibility conditions may be reduced to the following single condition

$$
\mathbf{e}_{13,2}-\mathbf{e}_{23,1}=\text { constant. }
$$

Exercise 7: Find which of the following values of $\mathbf{e}_{\mathrm{ij}}$ are possible linear strains.
(i) $\mathbf{e}_{11}=\alpha\left(\mathbf{x}_{1}{ }^{2}+\mathbf{x}_{2}{ }^{2}\right), \mathbf{e}_{22}=\alpha \mathbf{x}_{2}{ }^{2}, \mathbf{e}_{12}=2 \alpha \mathbf{x}_{1} \mathbf{x}_{2}, \mathbf{e}_{31}=\mathbf{e}_{32}=\mathbf{e}_{33}=0, \alpha=$ constant.

$$
\left(\mathbf{e}_{\mathbf{i j}}\right)=\left[\begin{array}{ccc}
x_{1}+x_{2} & x_{1} & x_{2}  \tag{ii}\\
x_{1} & x_{2}+x_{3} & x_{3} \\
x_{2} & x_{3} & x_{1}+x_{3}
\end{array}\right]
$$

compute the displacements in the case (i).

### 3.14 FINITE DEFORMATIONS

All the results reported in the preceding sections of this chapter were that of the classical theory of infinitesimal strains. Infinitesimal transformations permits the application of the principle of superposition of effects.

Finite deformations are those deformations in which the displacements $\mathbf{u}_{\mathbf{i}}$ together with their derivatives are no longer small.

Consider an aggregate of particles in a continuous medium. We shall use the same reference frame for the location of particles in the deformed and undeformed states.

Let the coordinates of a particle lying on a curve $\mathbf{C}_{0}$, before deformation , be denoted by ( $a_{1}, a_{2}, a_{3}$ ), and let the coordinates of the same particle after deformation (now lying on some curve $C$ ) be ( $x_{1}, x_{2}$, $\mathrm{x}_{3}$ ).

Then the elements of arc of the curve $\mathbf{C}_{0}$ and $\mathbf{C}$ are given, respectively, by

$$
\begin{equation*}
\mathbf{d} \mathrm{s}_{\mathrm{o}}^{2}=\mathbf{d} \mathrm{a}_{\mathrm{i}} \mathbf{d} \mathrm{a}_{\mathrm{i}} \tag{1}
\end{equation*}
$$

and

$$
\begin{equation*}
d s^{2}=d x_{i} d x_{i} . \tag{2}
\end{equation*}
$$

We consider first the Eulerian description of the strain and write

$$
\begin{equation*}
\mathbf{a}_{\mathrm{i}}=\mathbf{a}_{\mathrm{i}}\left(\mathbf{x}_{1}, \mathbf{x}_{2}, \mathbf{x}_{3}\right) . \tag{3}
\end{equation*}
$$

Then

$$
\begin{equation*}
\mathbf{d a}_{\mathrm{i}}=\mathbf{a}_{\mathrm{i}, \mathrm{j}} \mathbf{d} \mathbf{x}_{\mathrm{j}}=\mathbf{a}_{\mathrm{i}, \mathrm{k}} \mathbf{d x _ { k }} \tag{4}
\end{equation*}
$$

Substituting from (4) into (1), we write

$$
\begin{equation*}
\mathbf{d s}_{\mathbf{o}}{ }^{2}=\mathbf{a}_{\mathrm{i}, \mathrm{j}} \mathbf{a}_{\mathrm{i}, \mathrm{k}} \mathbf{d x}_{\mathbf{j}} \mathbf{d} \mathbf{x}_{\mathbf{k}} . \tag{5}
\end{equation*}
$$

Using the substitution tensor, equation (2) can be rewritten as

$$
\begin{equation*}
\mathbf{d ~ s} \mathbf{s}^{2}=\delta_{j k} \mathbf{d} \mathbf{x}_{j} \mathbf{d} \mathbf{x}_{k} \tag{6}
\end{equation*}
$$

we know that the measure of the strain is the difference $\mathrm{ds}^{\mathbf{2}}-\mathrm{ds}_{\mathbf{0}}{ }^{\mathbf{}}$.
From equations (5) and (6), we get

$$
\mathbf{d s ^ { 2 }}-\mathbf{d s}_{\mathbf{0}}{ }^{2}=\left(\delta_{j \mathbf{k}}-\mathbf{a}_{\mathrm{i}, \mathbf{j}} \mathbf{a}_{\mathrm{i}, \mathrm{k}}\right) \mathbf{d x _ { j }} \mathbf{d} \mathbf{x}_{\mathbf{k}}
$$

$$
\begin{equation*}
=2 \eta_{j k} d_{x_{j}} d x_{k} \tag{7}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathbf{2} \eta_{\mathbf{j k}}=\delta_{\mathbf{j k}}-\mathbf{a}_{\mathbf{i}, \mathbf{j}} \mathbf{a}_{\mathrm{i}, \mathbf{k}} \tag{8}
\end{equation*}
$$

We now write the strain components $\eta_{j k}$ in terms of the displacement components $u_{i}$, where

$$
\begin{equation*}
\mathbf{u}_{i}=\mathbf{x}_{\mathbf{i}}-\mathbf{a}_{\mathbf{i}} . \tag{9}
\end{equation*}
$$

This gives

$$
\mathbf{a}_{i}=\mathbf{x}_{\mathbf{i}}-\mathbf{u}_{\mathbf{i}} .
$$

Hence

$$
\begin{align*}
& \mathbf{a}_{\mathbf{i}, \mathrm{j}}=\delta_{\mathrm{ij}}-\mathbf{u}_{\mathbf{i}, \mathbf{j}}  \tag{10}\\
& \mathbf{a}_{\mathbf{i}, \mathbf{k}}=\delta_{\mathbf{i k}}-\mathbf{u}_{\mathbf{i}, \mathbf{k}} \tag{11}
\end{align*}
$$

Equations (8) , (10) and (11) yield

$$
\begin{align*}
& \mathbf{2} \eta_{j \mathbf{j k}}=\delta_{\mathbf{j k}}-\left(\delta_{\mathrm{ij}}-\mathbf{u}_{\mathrm{i}, \mathbf{j}}\right)\left(\delta_{\mathbf{i k}}-\mathbf{u}_{\mathbf{i}, \mathbf{k}}\right) \\
& =\delta_{\mathbf{j k}}-\left[\delta_{\mathbf{j k}}-\mathbf{u}_{\mathbf{k}, \mathbf{j}}-\mathbf{u}_{\mathbf{j}, \mathbf{k}}+\mathbf{u}_{\mathbf{i}, \mathbf{j}} \mathbf{u}_{\mathbf{i}, \mathbf{k}}\right] \\
& =\left(\mathbf{u}_{\mathbf{j}, \mathbf{k}}+\mathbf{u}_{\mathbf{k}, \mathbf{j}}\right)-\mathbf{u}_{\mathbf{i}, \mathbf{j}} \mathbf{u}_{\mathbf{i}, \mathbf{k}} . \tag{12}
\end{align*}
$$

The quantities $\eta_{\mathrm{jk}}$ are called the Eulerian strain components.
If, on the other hand, Lagrangian coordinates are used, and equations of transformation are of the form

$$
\begin{equation*}
\mathbf{x}_{\mathrm{i}}=\mathbf{x}_{\mathrm{i}}\left(\mathbf{a}_{1}, \mathbf{a}_{2}, \mathbf{a}_{3}\right) \tag{13}
\end{equation*}
$$

then

$$
\begin{equation*}
\mathbf{d} \mathbf{x}_{\mathrm{i}}=\mathrm{x}_{\mathrm{i}, \mathrm{j}} \mathbf{d a _ { i }}=\mathrm{x}_{\mathrm{i}, \mathrm{k}} \mathbf{d a _ { k }} \tag{14}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathbf{d s}^{2}=\mathbf{x}_{\mathrm{i}, \mathrm{j}} \mathrm{x}_{\mathrm{i}, \mathrm{k}} \mathbf{d} \mathrm{a}_{\mathrm{j}} \mathbf{d} \mathrm{a}_{\mathrm{k}}, \tag{15}
\end{equation*}
$$

while

$$
\begin{equation*}
\mathbf{d s}_{\mathbf{o}}{ }^{2}=\delta_{\mathrm{jk}} \mathbf{d a _ { j }} \mathbf{d a _ { k }} . \tag{16}
\end{equation*}
$$

The Lagrangian components of strain $\epsilon_{\mathrm{jk}}$ are defined by

$$
\begin{equation*}
d s^{2}-d s_{0}{ }^{2}=2 \in_{j k} d a_{j} d a_{k} . \tag{17}
\end{equation*}
$$

Since

$$
\begin{equation*}
\mathbf{x}_{\mathbf{i}}=\mathbf{a}_{\mathbf{i}}+\mathbf{u}_{\mathrm{i}} \tag{18}
\end{equation*}
$$

therefore,

$$
\begin{aligned}
& \mathbf{x}_{\mathrm{i}, \mathrm{j}}=\delta_{\mathrm{ij}}+\mathbf{u}_{\mathrm{i}, \mathrm{j}}, \\
& \mathbf{x}_{\mathrm{i}, \mathrm{k}}=\delta_{\mathrm{i} \mathbf{k}}+\mathbf{u}_{\mathrm{i}, \mathrm{k}} .
\end{aligned}
$$

Now

$$
\begin{align*}
& \mathbf{d s}{ }^{2}-\mathbf{d s}_{0}{ }^{2}=\left(\mathbf{x}_{\mathrm{i}, \mathrm{j}} \mathrm{X}_{\mathrm{i}, \mathrm{k}}-\delta_{\mathrm{jk}}\right) \mathbf{d a _ { j }} \mathrm{da}_{\mathrm{k}} \\
& =\left[\left(\delta_{i j}+\mathbf{u}_{\mathrm{i}, \mathrm{j}}\right)\left(\delta_{\mathrm{ik}}+\mathbf{u}_{\mathrm{i}, \mathrm{k}}\right)-\delta_{\mathrm{jk}}\right] \mathrm{da}_{\mathbf{j}} \mathbf{d} \mathbf{a}_{\mathrm{k}=}= \\
& =\left(\mathbf{u}_{\mathrm{j}, \mathrm{k}}+\mathbf{u}_{\mathrm{k}, \mathrm{j}}+\mathbf{u}_{\mathrm{i}, \mathrm{j}} \mathbf{u}_{\mathbf{i}, \mathbf{k}}\right) \mathrm{da}_{\mathrm{j}} \mathbf{d} \mathbf{a}_{\mathbf{k}} . \tag{19}
\end{align*}
$$

Equations (17) and (19) give

$$
\begin{equation*}
\mathbf{2}_{\in_{\mathbf{j k}}}=\mathbf{u}_{\mathbf{j}, \mathbf{k}}+\mathbf{u}_{\mathbf{k}, \mathbf{j}}+\mathbf{u}_{\mathbf{i}, \mathbf{j}} \mathbf{u}_{\mathbf{i}, \mathbf{k} \cdot} \tag{20}
\end{equation*}
$$

It is mentioned here that the differentiation in (12) is carried out with respect to the variables $x_{i}$, while in (20) the $a_{i}$ are regarded as the independent variables.

To make the difference explicitly clear, we write out the typical expressions $\eta_{\mathbf{j k}}$ and $\epsilon_{\mathbf{j k}}$ in unabridged notation,

$$
\begin{align*}
& \boldsymbol{\eta}_{\mathrm{xx}}=\frac{\partial u}{\partial x}-\frac{1}{2}\left[\left(\frac{\partial u}{\partial x}\right)^{2}+\left(\frac{\partial v}{\partial x}\right)^{2}+\left(\frac{\partial w}{\partial x}\right)^{2}\right]  \tag{21}\\
& \mathbf{2}_{\mathbf{x y}}=\left(\frac{\partial u}{\partial y}+\frac{\partial v}{\partial x}\right)-\left(\frac{\partial u}{\partial x} \frac{\partial u}{\partial y}+\frac{\partial v}{\partial x} \frac{\partial v}{\partial y}+\frac{\partial w}{\partial x} \frac{\partial w}{\partial y}\right),  \tag{22}\\
& \epsilon_{\mathrm{xx}}=\frac{\partial u}{\partial a}+\frac{1}{2}\left[\left(\frac{\partial u}{\partial a}\right)^{2}+\left(\frac{\partial v}{\partial b}\right)^{2}+\left(\frac{\partial w}{\partial a}\right)^{2}\right]  \tag{23}\\
& \mathbf{2} \in_{\mathrm{xy}}=\left(\frac{\partial u}{\partial b}+\frac{\partial v}{\partial a}\right)+\left(\frac{\partial u}{\partial a} \frac{\partial u}{\partial b}+\frac{\partial v}{\partial a} \frac{\partial v}{\partial b}+\frac{\partial w}{\partial a} \frac{\partial w}{\partial b}\right) \tag{24}
\end{align*}
$$

When the strain components are large, it is no longer possible to give simple geometrical interpretations of the strains $\epsilon_{j k}$ and $\eta_{j k}$.

Now, we consider some particular cases.
Case I: Consider a line element with

$$
\begin{equation*}
d s_{0}=d a_{1}, d a_{2}=0, d a_{3}=0 . \tag{25}
\end{equation*}
$$

Define the extension $E_{1}$ of this element by

$$
\mathbf{E}_{\mathbf{1}}=\frac{d s-d s_{o}}{d s_{o}}
$$

Then

$$
\begin{equation*}
d s=\left(1+E_{1}\right) d s_{0}, \tag{26}
\end{equation*}
$$

and consequently

$$
\begin{align*}
\mathbf{d s} \mathbf{s}^{2}-\mathbf{d s}_{\mathbf{o}}{ }^{2} & =2 \in_{\mathbf{j k}} \mathbf{d} \mathbf{a}_{\mathbf{j}} \mathbf{d} \mathbf{a}_{\mathbf{k}} \\
& =2 \in_{11} \mathbf{d} \mathbf{a}_{\mathbf{j}}^{2} . \tag{27}
\end{align*}
$$

Equations (25) to (27) yield

$$
\left(1+E_{1}\right)^{2}-1=2 \in_{11}
$$

or

$$
\begin{equation*}
\mathbf{E}_{\mathbf{1}}=\sqrt{1+2 \epsilon_{11}}-\mathbf{1} \tag{28}
\end{equation*}
$$

When the strain $\in_{11}$ is small, (28) reduces to

$$
\mathbf{E}_{1} \cong \epsilon_{11},
$$

as was shown in the discussion of infinitesimal strains.

## Case II: Consider next two line elements

$$
\begin{equation*}
\mathrm{ds}_{0}=\mathrm{da}_{2}, \mathrm{da}_{1}=\mathbf{0}, \mathrm{da}_{3}=\mathbf{0}, \tag{29}
\end{equation*}
$$

and

$$
\begin{equation*}
d \bar{s}_{\mathbf{0}}=\mathbf{d} \overline{\mathbf{a}}_{3}, \mathbf{d} \overline{\mathbf{a}}_{\mathbf{1}}=\mathbf{d} \overline{\mathbf{a}}_{2}=\mathbf{0} . \tag{30}
\end{equation*}
$$

These two elements lie initially along the $a_{2}$ - and $a_{3}$ - axes.
Let $\theta$ denote the angle between the corresponding deformed elements $\mathbf{d} x_{i}$ and $d \bar{x}_{i}$, of lengths ds and $d \bar{s}$, respectively. Then

$$
\begin{align*}
& \mathbf{d s} d \bar{s} \boldsymbol{\operatorname { c o s }} \theta=\mathbf{d} x_{i} d \overline{x_{i}} \\
& =\boldsymbol{x}_{i, \alpha} \quad \bar{x}_{i, \beta} \mathbf{d a}_{\alpha} d \overline{a_{\beta}} \\
& =\boldsymbol{x}_{i, 2} \quad \bar{x}_{i, 3} \mathbf{d a}_{2} d \overline{a_{3}} \\
& =\mathbf{2} \in{ }_{23} \mathbf{d a}_{2} d \overline{a_{3}} . \tag{31}
\end{align*}
$$

Let

$$
\begin{equation*}
\alpha_{23}=\frac{\pi}{2}-\theta \tag{32}
\end{equation*}
$$

denotes the change in the right angle between the line elements in the initial state. Then , we have

$$
\begin{align*}
\sin \alpha_{23} & =\mathbf{2} \in_{23}\left(\frac{d a_{2}}{d s}\right)\left(\frac{d \overline{a_{3}}}{d \bar{s}}\right) .  \tag{33}\\
& =\frac{2 \in_{23}}{\sqrt{1+2 \epsilon_{22}} \sqrt{1+2 \epsilon_{33}}}, \tag{34}
\end{align*}
$$

using relations in (26) and (28).
Again, if the strains $\in_{\mathrm{ij}}$ are so small that their products can be neglected, then

$$
\begin{equation*}
\alpha_{23} \cong 2 \in_{23}, \tag{35}
\end{equation*}
$$

as proved earlier for infinitesimal strains.
Remark: If the displacements and their derivatives are small, then it is immaterial whether the derivatives of the displacements are calculated at the position of a point before or after deformation. In this case, we may neglect the nonlinear terms in the partial derivatives in (12) and (20) and reduce both sets of formulas to

$$
2 \eta_{j k}=\mathbf{u}_{\mathbf{j}, \mathrm{k}}+\mathbf{u}_{\mathrm{k}, \mathrm{j}}=\mathbf{2}_{\in_{\mathbf{j k}}},
$$

which were obtained for an infinitesimal transformation.
It should be emphasized that the transformations of finite homogeneous strain are not in general commutative and that the simple superposition of effects is no longer applicable to finite deformation.

## Chapter-4

## Constitutive Equations of Linear Elasticity

### 4.1. INTRODUCTION <br> It is a fact of experience that deformation of a solid body induces stresses within. The relationship between stress and deformation is expressed as a constitutive relation for the material and depends on the material properties and also on other physical observables like temperative and, perhaps, the electromagnetic field.

An elastic deformation is defined to be one in which the stress is determined by the current value of the strain only, and not on rate of strain or strain history : $\tau=\tau(\mathrm{e})$.

An elastic solid that undergoes only an infinitesimal deformation and for which the governing material is linear is called a linear elastic solid or Hookean solid.

From experimental observations, it is known that , under normal loadings , many structural materials such as metals, concrete, wood and rocks behave as linear elastic solids.

The classical theory of elasticity (or linear theory) serves as an excellent model for studying the mechanical behaviour of a wide variety of such solid materials.

Hook's law : In 1678 , Robert Hook, on experimental grounds, stated that the extension is proportional to the force. Cauchy in $\mathbf{1 8 2 2}$ generalized Hook law for the deformation of elastic solids. According to Cauchy, " Each component of stress at any point of an elastic body is a linear function of the components of strain at the point".

This law is now known as Generalized Hooke's Law. Here, linearity means that stress - strain relations are linear.

### 4.2. GENERALIZED HOOKE'S LAW

## In general, we write the following set of linear

## relations

$$
\begin{aligned}
& \tau_{11}=c_{1111} e_{11}+c_{1112} e_{12}+\ldots \ldots \ldots . .+c_{1133} e_{33}, \\
& \tau_{12}=c_{1211} e_{11}+c_{1212} e_{12}+\ldots \ldots \ldots . .+c_{1233} e_{33}, \\
& \tau_{33}=c_{3311} e_{11}+c_{3312} e_{12}+\ldots \ldots \ldots \ldots .+c_{3333} e_{33},
\end{aligned}
$$

or

$$
\begin{equation*}
\tau_{\mathrm{ij}}=\mathrm{c}_{\mathrm{ijk} \mathrm{l}} \mathrm{e}_{\mathrm{k} 1}, \tag{1}
\end{equation*}
$$

where $\tau_{\mathrm{ij}}$ is the stress tensor and $\mathrm{e}_{\mathrm{k} 1}$ is the strain tensor. The coefficients, which are $81=3^{4}$ in number, are called elastic moduli.

In general, these coefficients depend on the physical properties of the medium and are independent of the strain components $\mathrm{e}_{\mathrm{i} j}$.

We suppose that relations (1) hold at every point of the medium and at every instant of time and are solvable for $\mathrm{e}_{\mathrm{ij}}$ in terms of $\tau_{\mathrm{ij}}$.

From (1), it follows that $\tau_{\mathrm{ij}}$ are all zero whenever all $\mathrm{e}_{\mathrm{ij}}$ are 0 .
It means that in the initial unstrained state the body is unstressed. From quotient law for tensors, relation (1) shows that $\mathrm{c}_{\mathrm{ijk} l}$ are components of a fourth - order tensor.

This tensor is called elasticity tensor. Since $\mathbf{e}_{\mathrm{ij}}$ are dimensionless quantities, it follows that elastic moduli $\mathrm{c}_{\mathrm{ijk} l}$ have the same dimensions as the stresses (force/Area).

If, however, $\mathrm{c}_{\mathrm{ijkl}}$ do not change throughout the medium for all time, we say that the medium is (elastically) homogeneous.

Thus , for a homogeneous elastic solid, the elastic moduli are constants so that the mechanical properties remain the same throughout the solid for all times. The tensor equation (1) represents the generalized Hooke's law in the $\mathbf{x}_{\mathbf{i}}$ - system.

Since $\tau_{\mathrm{ij}}$ is symmetric and $\mathrm{e}_{\mathrm{kl}}$ is symmetric , there are left 6 independent equations in relation (1) and each equation contains 6 independent elastic moduli.

So, the number of independent elastic coefficients are, in fact, 36 for a generalized an anisotropic medium.

For simplicity, we introduce the following (engineering) notations
$\left.\begin{array}{l}\tau_{11}=\tau_{1}, \tau_{22}=\tau_{2}, \tau_{33}=\tau_{3}, \tau_{23}=\tau_{4}, \tau_{31}=\tau_{5}, \tau_{12}=\tau_{6}, \\ \mathbf{e}_{11}=\mathbf{e}_{1}, \mathbf{e}_{22}=\mathbf{e}_{2}, \mathbf{e}_{33}=\mathbf{e}_{3}, \mathbf{2} \mathbf{e}_{23}=\mathbf{e}_{4}, \mathbf{2} \mathbf{e}_{31}=\mathbf{e}_{5}, \mathbf{2} \mathbf{e}_{12}=\mathbf{e}_{6} .\end{array}\right\}$
Then, the generalized Hooke's law may be written in the form

$$
\begin{equation*}
\tau_{\mathrm{i}}=\mathbf{c}_{\mathrm{ij}} \mathbf{e}_{\mathrm{j}} \quad ; \quad \mathbf{i}, \mathbf{j}=\mathbf{1}, \mathbf{2}, \ldots \ldots, \mathbf{6} \tag{3}
\end{equation*}
$$

or in the matrix form

$$
\left[\begin{array}{c}
\tau_{1}  \tag{4}\\
\tau_{2} \\
\tau_{3} \\
\tau_{4} \\
\tau_{5} \\
\tau_{6}
\end{array}\right]=\left[\begin{array}{ll}
c_{11} & c_{12} \ldots \ldots \ldots \ldots c_{16} \\
c_{21} & c_{22} \ldots \ldots \ldots \ldots c_{26} \\
c_{31} \ldots \ldots \ldots \ldots \ldots \ldots \ldots \\
c_{41} \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \\
c_{51} \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \\
c_{61} & c_{62} \ldots \ldots \ldots \ldots . c_{66}
\end{array}\right]\left[\begin{array}{c}
e_{1} \\
e_{2} \\
e_{3} \\
e_{4} \\
e_{5} \\
e_{6}
\end{array}\right] .
$$

If the elastic properties (or mechanical properties) of a medium at a point are independent of the orientation (i.e. $c_{i j}^{\prime}=c_{i j}$ ) of the coordinate axes, then we say that the medium under consideration is isotropic.

If at a point of the medium, properties of medium (i.e. , $\mathrm{c}_{\mathrm{ij}}$ 's) depend upon the orientation, then medium is called an Anisotropic or Aelotropic medium.

The $6 \times 6$ matrix $\left(c_{i j}\right)$ in (4) is called stiffness matrix.

### 4.3. HOMOGENEOUS ISOTROPIC MEDIA

When the elastic coefficients $c_{i j k}$ in the generalized Hooke's law are constants throughout the medium and they are independent of the orientation of the coordinate axes , the elastic media is termed as homogeneous isotropic media.

We know that the generalized Hooke's law is

$$
\begin{equation*}
\tau_{\mathrm{ij}}=\mathbf{c}_{\mathrm{ijk} l} \mathbf{e}_{k l} \tag{1}
\end{equation*}
$$

where $\tau_{i j}$ and $\mathrm{e}_{\mathrm{k} l}$ are symmetric strain and stress tensor, respectively, and $c_{i j k l}$ are the components of a tensor of order 4.

Since the media is isotropic, therefore, the tensor $\mathrm{c}_{\mathrm{ijk} l}$ is an isotropic tensor. Hence, it can be represented in the form

$$
\begin{equation*}
\mathbf{c}_{\mathrm{ij} k l}=\alpha \delta_{\mathrm{ij}} \delta_{\mathrm{k} l}+\beta \delta_{\mathrm{ik}} \delta_{\mathrm{j} l}+\gamma \delta_{\mathrm{il}} \delta_{\mathrm{jk}} \tag{2}
\end{equation*}
$$

where $\alpha, \beta, \gamma$ are some scalars. From equations (1) and (2), we obtain

$$
\begin{align*}
\tau_{\mathrm{ij}} & =\left(\alpha \delta_{\mathrm{ij}} \delta_{\mathrm{k} l}+\beta \delta_{\mathrm{ik}} \delta_{\mathrm{j} l}+\gamma \delta_{\mathrm{il}} \delta_{\mathrm{jk}}\right) \mathbf{e}_{\mathrm{k} l} \\
& =\alpha \delta_{\mathrm{ij}} \mathbf{e}_{\mathrm{kk}}+\beta \delta_{\mathrm{ik}} \mathbf{e}_{\mathrm{kj}}+\gamma \delta_{\mathrm{i} l} \mathbf{e}_{\mathbf{j} l} \\
& =\alpha \delta_{\mathrm{ij}} \mathbf{e}_{\mathrm{kk}}+\beta \mathbf{e}_{\mathrm{ij}}+\gamma \mathbf{e}_{\mathrm{ji}} \\
& =\alpha \delta_{\mathrm{ij}} \mathbf{e}_{\mathrm{kk}}+(\beta+\gamma) \mathbf{e}_{\mathrm{ij}}, \tag{3}
\end{align*}
$$

since $\mathrm{e}_{\mathrm{ij}}=\mathrm{e}_{\mathrm{ji}}$. On redesignating $\alpha$ by $\lambda$ and $(\beta+\gamma)$ by $2 \mu$, relation (3) yields

$$
\begin{equation*}
\tau_{\mathrm{ij}}=\lambda \delta_{\mathrm{ij}} \mathrm{e}_{\mathrm{kk}}+2 \mu \mathrm{e}_{\mathrm{ij}} \tag{4}
\end{equation*}
$$

The two elastic coefficients $\lambda$ and $\mu$ are known as Lame constants. $\delta_{\mathrm{ij}}$ is the substitution tensor.

Let

$$
\begin{equation*}
v=\mathbf{e}_{\mathrm{kk}}, \theta=\tau_{\mathrm{ij}} . \tag{5}
\end{equation*}
$$

Taking $j=i$ in (4) and using summation convention according, we find

$$
\begin{align*}
\theta & =3 \lambda v+2 \mu v \\
& =(3 \lambda+2 \mu) v \\
& =3 k v \tag{6}
\end{align*}
$$

where

$$
\begin{equation*}
\mathbf{k}=\lambda+\frac{2}{3} \mu \tag{7}
\end{equation*}
$$

is the bulk modulus.
From (4), we write

$$
\begin{align*}
\mathbf{e}_{\mathrm{ij}} & =\frac{\lambda}{2 \mu} \delta_{\mathrm{ij}} \mathbf{v}+\frac{1}{2 \mu} \tau_{\mathrm{ij}} \\
& =\frac{-\lambda}{2 \mu(3 \lambda+2 \mu)} \delta_{\mathrm{ij}} \tau_{\mathrm{kk}}+\frac{1}{2 \mu} \tau_{\mathrm{ij}} \tag{8}
\end{align*}
$$

This relation expresses the strain components as a linear functions of components of stress tensor.

Question : Show that if the medium is isotropic , the principal axes of stress are coincident with the principal axes of strain.

Solution : Let the $\mathbf{x}_{\mathrm{i}}$-axes be directed along the principal axes of strain.
Then

$$
\begin{equation*}
\mathbf{e}_{12}=\mathbf{e}_{13}=\mathbf{e}_{23}=\mathbf{0} \tag{1}
\end{equation*}
$$

The stress - strain relations for an isotropic medium are

$$
\begin{equation*}
\tau_{\mathrm{ij}}=\lambda \delta_{\mathrm{ij}} \mathbf{e}_{\mathrm{kk}}+2 \mu \mathbf{e}_{\mathrm{ij}} \tag{2}
\end{equation*}
$$

Combining (1) $\&(2)$, we find

$$
\begin{equation*}
\tau_{12}=\tau_{13}=\tau_{23}=0 \tag{3}
\end{equation*}
$$

This shows that the coordinates axes $x_{i}$ are also the principal axes of stress.
This proves the result. Thus, there is no distinction between the principal axes of stress and of strain for isotropic media.

### 4.4. PHYSICAL MEANINGS OF ELASTIC MODULI FOR AN ISOTROPIC MEDIUM

We have already introduced two elastic moduli $\lambda$ and $\mu$ in the generalized Hooke's law for an isotropic medium. We introduce three more elastic moduli defined below

$$
\begin{equation*}
\mathbf{E}=\frac{\mu(3 \lambda+2 \mu)}{\lambda+\mu}, \quad \sigma=\frac{\lambda}{2(\lambda+\mu)} \quad, \quad \mathbf{k}=\lambda+\frac{2}{3} \mu \tag{1}
\end{equation*}
$$

The quantity $\sigma$ is dimensionless and is called the Poisson ratio. It was introduced by

Simon D. Poisson in 1829.
The quantity $\mathbf{E}$ is called Young's modulus after Thomas Young who introduced it in the early $19^{\text {th }}$ century, probably in 1807. Its dimension is that of a stress (force/area).

The elastic modulus $k$ is called the modulus of compression or the bulk modulus.

Solving the first two equations for $\lambda$ and $\mu$ (in terms $\sigma$ and $E$ ), we find

$$
\begin{equation*}
\lambda=\frac{E \sigma}{(1+\sigma)(1-2 \sigma)} \quad, \quad \mu=\frac{E}{2(1+\sigma)} . \tag{2}
\end{equation*}
$$

From (2), we find the following relations

$$
\left.\begin{array}{l}
\lambda+2 \mu=\frac{E(1-\sigma)}{(1+\sigma)(1-2 \sigma)}, \frac{\lambda+\mu}{\mu}=\frac{1}{1-2 \sigma}  \tag{3}\\
\frac{\lambda+2 \mu}{\mu}=\frac{2-\sigma}{1-2 \sigma}, \frac{\lambda}{\lambda+2 \mu}=\frac{\sigma}{1-\sigma} \cdot
\end{array}\right\}
$$

Generalized Hooke's Law in terms of Elastic Moduli $\sigma$ and E
We know that the generalized Hooke's law (giving strain components in terms of stresses) in terms of Lame's constants $\lambda$ and $\mu$ is

$$
\begin{equation*}
\mathbf{e}_{\mathrm{ij}}=\frac{-\lambda}{2 \mu(3 \lambda+2 \mu)} \delta_{\mathrm{ij}} \tau_{\mathrm{kk}}+\frac{1}{2 \mu} \tau_{\mathrm{ij}} \tag{4}
\end{equation*}
$$

Substituting the values of $\lambda$ and $\mu$ in terms of $E$ and $\sigma$ from (2) into (4), we find

$$
\begin{equation*}
\mathbf{e}_{\mathrm{ij}}=\frac{-\sigma}{E} \delta_{\mathrm{ij}} \tau_{\mathrm{kk}}+\frac{1+\sigma}{E} \tau_{\mathrm{ij}} . \tag{5}
\end{equation*}
$$

Note : (1) Out of five elastic moduli (namely ; $\lambda, \mu, \mathbf{E}, \sigma, \mathbf{k}$ ) only two are independent.

Note : (2) The Hooke's law, given in (5), is frequently used in engineering problems.

Remark : The following three experiments give some insight into the physical significance of various elastic moduli for isotropic media.
(I) Simple Tension :

Consider a right cylinder with its axis parallel to the $\mathbf{x}_{1}$ - axis which is subjected to longitudinal forces applied to the ends of the cylinder. These applied forces give rise to a uniform tension $T$ in every cross - section of the cylinder so that the stress tensor $\tau_{\mathrm{ij}}$ has only one non - zero component $\tau_{11}=\mathrm{T}$.

That is

$$
\begin{equation*}
\tau_{11}=\mathbf{T}, \tau_{22}=\tau_{33}=\tau_{12}=\tau_{23}=\tau_{31}=\mathbf{0} . \tag{1}
\end{equation*}
$$



Fig. (4.1)
Since the body forces are absent $\left(f_{i}=0\right)$, the state of stress given by (1) satisfies the equilibrium equation $\tau_{\mathrm{i}, \mathrm{j}}=0$ in the interior of the cylinder.

A normal $\hat{v}$ to the lateral surface lies in the plane parallel to $\mathbf{x}_{\mathbf{2}} \mathbf{x}_{\mathbf{3}}$ plane, so $\hat{v}=\left(0, v_{2}, v_{3}\right)$.

The relation $\stackrel{v}{T_{i}}=\tau_{\mathrm{ij}} \boldsymbol{v}_{\mathbf{j}}$ implies that $\stackrel{v}{T_{1}}=\stackrel{v}{T_{2}}=\stackrel{v}{T}=\mathbf{0}$.

Hence $\stackrel{v}{T}=\mathbf{0}$.

This shows that the lateral surface of the cylinder is free from tractions.
The generalized Hooke's law giving strains in terms of stresses is

$$
\begin{equation*}
\mathbf{e}_{\mathrm{ij}}=\frac{-\lambda}{2 \mu(3 \lambda+2 \mu)} \delta_{\mathrm{ij}} \tau_{\mathrm{kk}}+\frac{1}{2 \mu} \tau_{\mathrm{ij}} \tag{2}
\end{equation*}
$$

We find from equations (1) \& (2) that

$$
\begin{align*}
& \mathbf{e}_{11}=\frac{\lambda+\mu}{\mu(3 \lambda+2 \mu)} \mathbf{T}, \\
& \mathbf{e}_{22}=\mathbf{e}_{33}=\frac{-\lambda}{2 \mu(3 \lambda+2 \mu)} \mathbf{T}, \\
& \mathbf{e}_{12}=\mathbf{e}_{23}=\mathbf{e}_{31}=\mathbf{0} . \tag{3}
\end{align*}
$$

Since $\frac{\lambda+\mu}{\mu(3 \lambda+2 \mu)}=\frac{1}{E} \quad$ and $\quad \frac{\sigma}{E}=\frac{\lambda}{2 \mu(3 \lambda+2 \mu)}$,
Therefore

$$
\mathbf{e}_{11}=\frac{T}{E},
$$

$$
\begin{equation*}
\mathbf{e}_{22}=\mathbf{e}_{33}=-\frac{\sigma}{E} \mathbf{T}=-\sigma \mathbf{e}_{11} \tag{5}
\end{equation*}
$$

and

$$
\mathbf{e}_{12}=\mathbf{e}_{13}=\mathbf{e}_{23}=\mathbf{0} .
$$

These strain components obviously satisfies the compatibility equations

$$
\mathbf{e}_{\mathrm{ij}, \mathrm{kl}}+\mathbf{e}_{\mathrm{k}, \mathrm{ij}}-\mathbf{e}_{\mathrm{ik}, \mathrm{j} \mathbf{l}}-\mathbf{e}_{\mathbf{j}, \mathrm{l}, \mathrm{ik}}=\mathbf{0},
$$

and therefore, the state of stress given in (1) actually corresponds to one which can exist in a deformed elastic body. From equation (5), we write

$$
\begin{equation*}
\frac{\tau_{11}}{e_{11}}=\mathbf{E} \quad, \quad \frac{e_{22}}{e_{11}}=\frac{e_{33}}{e_{11}}=-\sigma . \tag{6}
\end{equation*}
$$

Experiments conducted on most naturally occurring elastic media show that a tensile longitudinal stress produces a longitudinal extension together with a contraction in a transverse directions. According for $\tau_{11}=$ T>0, we take

$$
\mathbf{e}_{11}>0 \text { and } e_{22}<0, e_{33}<0 .
$$

It then follows from (6) that

$$
\begin{equation*}
E>0 \text { and } \sigma>0 . \tag{7}
\end{equation*}
$$

From equation (6), we see that $E$ represents the ratio of the longitudinal stress $\tau_{11}$ to the corresponding longitudinal strain $e_{11}$ produced by the stress $\tau_{11}$.

From equation (6), we get

$$
\begin{equation*}
\left|\frac{e_{22}}{e_{11}}\right|=\left|\frac{e_{33}}{e_{11}}\right|=\sigma . \tag{8}
\end{equation*}
$$

Thus, the Poisson's ratio $\sigma$ represents the numerical value of the ratio of the contraction $\mathrm{e}_{22}$ (or $\mathbf{e}_{33}$ ) in a transverse direction to the corresponding extension $\mathrm{e}_{11}$ in the longitudinal direction.
(II) Pure Shear

From generalized Hooke's law for an isotropic medium, we write

$$
\begin{equation*}
\mathbf{2} \mu=\frac{\tau_{12}}{e_{12}}=\frac{\tau_{13}}{e_{13}}=\frac{\tau_{23}}{e_{23}} . \tag{9}
\end{equation*}
$$

The constant $2 \mu$ is thus the ratio of a shear stress component to the corresponding shear strain component. It is, therefore, related to the rigidity of the elastic material.

For this reason, the coefficient $\mu$ is called the modulus of rigidity or the shear modulus.

The other lame constant $\lambda$ has no direct physical meaning.
The value of $\mu$ in terms of Young's modulus $E$ and Poisson ratio $\sigma$ is given by

$$
\begin{equation*}
\mu=\frac{E}{2(1+\sigma)} \tag{10}
\end{equation*}
$$

Since $\mathbf{E}>\mathbf{0}, \sigma>0$, it follows that

$$
\begin{equation*}
\mu>0 \tag{11}
\end{equation*}
$$

(III) Hydrostatic Pressure

Consider an elastic body of arbitrary shape which is put in a large vessel containing a liquid. A hydrostatic pressure $p$ is exerted on it by the liquid and the elastic body experience all around pressure. The stress tensor is given by $\tau_{\mathrm{ij}}=-\mathbf{p} \delta_{\mathrm{ij}}$. That is ,

$$
\tau_{11}=\tau_{22}=\tau_{33}=-\mathbf{p}, \tau_{12}=\tau_{23}=\tau_{31}=\mathbf{0} .
$$



Fig. (4.2)
These stress components satisfy the equilibrium equations for zero body force. We find

$$
\tau_{k k}=-3 p
$$

and the generalized Hooke's law giving strains in terms of stresses

$$
\begin{equation*}
\mathbf{e}_{\mathrm{ij}}=\frac{1}{2 \mu}\left[\frac{-\lambda}{3 \lambda+2 \mu} \delta_{i j} \tau_{k k}+\tau_{i j}\right] \tag{13}
\end{equation*}
$$

gives

$$
\begin{align*}
& \mathbf{e}_{12}=\mathbf{e}_{23}=\mathbf{e}_{31}=\mathbf{0}, \\
& \mathbf{e}_{11}=\mathbf{e}_{22}=\mathbf{e}_{33}=\frac{1}{2 \mu}\left[\frac{3 \lambda p}{3 \lambda+2 \mu}+(-p)\right]=-\frac{p}{3 \lambda+2 \mu}, \tag{14}
\end{align*}
$$

which obviously satisfy the compatibility equations. We find

$$
\mathbf{e}_{\mathbf{k k}}=\frac{-3 p}{3 \lambda+2 \mu}=\frac{-p}{\lambda+\frac{2}{3} \mu}=\frac{-p}{k} .
$$

That is ,

$$
\begin{equation*}
v(\text { cubical dilatation })=\frac{-p}{k} \tag{15}
\end{equation*}
$$

From experiments, it has been found that a hydrostatic pressure tends to reduce the volume of the elastic material. That is , if $\mathbf{p}>\mathbf{0}$, then

$$
\mathrm{e}_{\mathrm{kk}}=v<0 .
$$

Consequently, it follows from (15) that $k>0$.
Relation (15) also shows that the constant $k$ represents the numerical value of the ratio of the compressive stress to the dilatation.

Substituting the value of $\lambda$ and $\mu$ in terms of $E$ and $\sigma$, we find

$$
\begin{equation*}
\mathbf{k}=\frac{E}{3(1-2 \sigma)} . \tag{16}
\end{equation*}
$$

Since $k>0$ and $E>0$, it follows that $o<\sigma<\frac{1}{2}$ for all physical substances.

Since

$$
\begin{equation*}
\lambda=\frac{E \sigma}{(1+\sigma)(1-2 \sigma)}, \tag{17}
\end{equation*}
$$

and $\mathrm{E}>0,0<\sigma<\frac{1}{2}$, it follows that $\lambda>0$.
Remark : The solutions of many problems in elasticity are either exactly or approximately independent of the value chosen for Poisson's ratio. This
fact suggests that approximate solutions may be found by so choosing Poisson's ratio as to simplify the problem. Show that , if one take $\sigma=0$, then

$$
\lambda=\mathbf{0}, \mu=\frac{E}{2}, \mathbf{k}=\frac{1}{3} \mathbf{E},
$$

and Hooke's law is expressed by

$$
\tau_{\mathrm{ij}}=\mathbf{E} \mathrm{e}_{\mathrm{ij}}=\frac{1}{2} \mathbf{E}\left(\mathbf{u}_{\mathrm{i}, \mathrm{j}}+\mathbf{u}_{\mathrm{j}, \mathrm{i}}\right)
$$

Note 1: The elastic constants $\mu, \mathbf{E}, \sigma, \mathbf{k}$ have definite physical meanings. These constants are called engineering elastic modulus.

Note 2: The material such as steel, brass, copper, lead, glass, etc. are isotropic elastic medium.

Note 3: We find

$$
\mathbf{e}_{\mathbf{k k}}=\frac{\tau_{k k}}{3 k}=\frac{1-2 \sigma}{E} \tau_{\mathbf{k k}}
$$

Thus $\mathbf{e}_{\mathbf{k k}}=0$ iff $\sigma=\frac{1}{2}$, provided $E$ and $\tau_{\mathbf{k k}}$ remain finite.
When $\sigma \rightarrow \frac{1}{2}, \lambda \rightarrow \infty, \mathbf{k} \rightarrow \infty, \mu=\frac{E}{3}, v=\mathbf{e}_{\mathrm{ii}}=\mathbf{u}_{\mathrm{i}, \mathrm{i}}=\mathbf{0}$,
This limiting case corresponds to which is called an incompressible elastic body.

Question : In an elastic beam placed along the $x_{3}$ - axis and bent by a couple about the $x_{2}$-axis, the stresses are found to be

$$
\tau_{33}=-\frac{E}{R} \mathbf{x}_{1}, \tau_{11}=\tau_{22}=\tau_{12}=\tau_{13}=\tau_{23}=\mathbf{0}, \mathbf{R}=\mathbf{c o n s t a n t}
$$

Find the corresponding strains.
Solution : The strains in terms of stresses $\&$ elastic moduli $E$ and $\sigma$ are given by the Hooke's law

$$
\begin{equation*}
\mathbf{e}_{\mathrm{ij}}=-\frac{\sigma}{E} \delta_{\mathrm{ij}} \tau_{\mathrm{kk}}+\frac{1+\sigma}{E} \tau_{\mathrm{ij}} \tag{1}
\end{equation*}
$$

Here

$$
\tau_{\mathrm{kk}}=-\frac{E}{R} \mathbf{x}_{1} .
$$

Hence, (1) becomes

$$
\begin{equation*}
\mathbf{e}_{\mathrm{ij}}=\frac{\sigma}{R} \mathbf{x}_{\mathbf{1}} \delta_{\mathrm{ij}}+\frac{1+\sigma}{E} \tau_{\mathrm{ij}} . \tag{2}
\end{equation*}
$$

This gives $\mathbf{e}_{11}=\mathbf{e}_{22}=\frac{\sigma}{R} \mathbf{x}_{1}, \mathbf{e}_{33}=-\frac{1}{R} \mathbf{x}_{1}, \mathbf{e}_{12}=\mathbf{e}_{23}=\mathbf{e}_{13}=\mathbf{0}$.
Question : A beam placed along the $\mathrm{x}_{1}$ - axis and subjected to a longitudinal stress $\tau_{11}$ at every point is so constrained that $\mathbf{e}_{22}=\mathbf{e}_{33}=\mathbf{0}$ at every point. Show that $\tau_{22}=\sigma \tau_{11}, \mathbf{e}_{11}=\frac{1-\sigma^{2}}{E} \tau_{11}$, $\frac{-\sigma \mathbf{1}+\sigma^{-}}{E} \tau_{11}$.

Solution : The Hooke's law giving the strains in terms of stresses is

$$
\begin{equation*}
\mathbf{e}_{\mathrm{ij}}=-\frac{\sigma}{E} \delta_{\mathrm{ij}} \tau_{\mathrm{kk}}+\frac{1+\sigma}{E} \tau_{\mathrm{ij}} \tag{1}
\end{equation*}
$$

It gives

$$
\begin{align*}
\mathbf{e}_{22} & =-\frac{\sigma}{E}\left(\tau_{11}+\tau_{22}+\tau_{33}\right)+\frac{1+\sigma}{E} \tau_{22} \\
& =\frac{1}{E} \tau_{22}-\frac{\sigma}{E}\left(\tau_{11}+\tau_{33}\right) . \tag{2}
\end{align*}
$$

Putting $e_{22}=e_{33}=0$ in (2), we get

$$
\begin{equation*}
\tau_{22}=\sigma \tau_{11} \tag{3}
\end{equation*}
$$

Also, from (1), we find

$$
\begin{align*}
\mathbf{e}_{11} & =-\frac{\sigma}{E}\left(\tau_{11}+\tau_{22}+\tau_{33}\right)+\frac{1+\sigma}{E} \tau_{11} \\
& =-\frac{\sigma}{E}\left(\tau_{11}+\sigma \tau_{11}\right)+\frac{1+\sigma}{E} \tau_{11} \\
& =\frac{1}{E}\left[-\sigma-\sigma^{2}+1+\sigma\right] \tau_{11}=\frac{1-\sigma^{2}}{E} \tau_{11} . \tag{4}
\end{align*}
$$

Also , from (1), we get

$$
\begin{align*}
\mathbf{e}_{33} & =-\frac{\sigma}{E}\left(\tau_{11}+\tau_{22}\right)+\frac{1+\sigma}{E} \tau_{33} \\
& =-\frac{\sigma}{E}\left(\tau_{11}+\sigma \tau_{11}\right)=\frac{-\sigma 1+\sigma^{-}}{E} \tau_{11} . \tag{5}
\end{align*}
$$

Exercise 1: Find the stresses with the following displacement fields :

$$
\begin{align*}
& \text { (i) } \mathrm{u}=\mathrm{kyz}, \mathrm{v}=\mathrm{kzx}, \mathrm{w}=\mathrm{kxy}  \tag{i}\\
& \text { (ii) } \mathrm{u}=\mathrm{kyz}, \mathrm{v}=\mathrm{kzx}, \mathrm{w}=\mathrm{k}\left(\mathrm{x}^{2}-\mathrm{y}^{2}\right)
\end{align*}
$$

where $k=$ constant .
Exercise 2:A rod placed along the $\mathbf{x}_{\mathbf{1}}$ - axis and subjected to a longitudinal stress $\tau_{11}$ is so constrained that there is no lateral contraction. Show that

$$
\tau_{11}=\frac{1-\sigma \underline{E}}{(1+\sigma)(1-2 \sigma)} e_{11}
$$

### 4.5. EQUILIBRIUM AND DYNAMIC EQUATIONS <br> FOR AN ISOTROPIC ELASTIC SOLID

We know that Cauchy's equation's of equilibrium in term of stress components are

$$
\begin{equation*}
\tau_{i \mathrm{i}, \mathrm{j}}+\mathbf{F}_{\mathrm{i}}=\mathbf{0}, \tag{1}
\end{equation*}
$$

where $F_{i}$ is the body force per unit volume and $i, j=1,2,3$.
The generalized Hooke's law for a homogeneous isotropic elastic body is

$$
\begin{align*}
\tau_{\mathrm{ij}} & =\lambda \delta_{\mathrm{ij}} \mathbf{e}_{\mathrm{k} \mathbf{k}}+2 \mu \mathbf{e}_{\mathrm{i} j} \\
& =\lambda \delta_{\mathrm{ij}} \mathbf{u}_{\mathrm{k}, \mathbf{k}}+\mu\left(\mathbf{u}_{\mathrm{i}, \mathrm{j}}+\mathbf{u}_{\mathrm{j}, \mathrm{i}}\right), \tag{2}
\end{align*}
$$

where $\lambda$ and $\mu$ are Lame constants. Putting the value of $\tau_{\mathrm{ij}}$ from (2) into equation (1), we find

$$
\begin{align*}
& \lambda \delta_{i j} \mathbf{u}_{k, k j}+\mu\left(\mathbf{u}_{i, j j}+\mathbf{u}_{\mathbf{j}, \mathbf{i}}\right)+\mathbf{F}_{\mathbf{i}}=\mathbf{0} \\
& \lambda \mathbf{u}_{\mathrm{k}, \mathbf{k}}+\mu \nabla^{2} \mathbf{u}_{\mathrm{i}}+\mu \mathbf{u}_{\mathrm{k}, \mathbf{k i}}+\mathbf{F}_{\mathbf{i}}=\mathbf{0} \\
& (\lambda+\mu) \frac{\partial \theta}{\partial x_{i}}+\mu \nabla^{2} \mathbf{u}_{\mathrm{i}}+\mathbf{F}_{\mathrm{i}}=\mathbf{0} \tag{3}
\end{align*}
$$

where $\theta=\mathbf{u}_{\mathrm{k}, \mathrm{k}}=\operatorname{div} \overline{\mathbf{u}}=$ cubical dilatation and $\mathbf{i}=1,2,3$.
Equations in (3) form a synthesis of the analysis of strain, analysis of stress and the stress - strain relation.

These fundamental partial differential equations of the elasticity theory are known as Navier's equations of equilibrium , after Navier (1821).

Equation (3) can be put in several different forms.
Form (A) : In vector form , equation (3) can be written as

$$
\begin{equation*}
(\lambda+\mu) \text { grad div } \overline{\mathbf{u}}+\mu \nabla^{2} \overline{\mathbf{u}}+\overline{\mathbf{F}}=\overline{\mathbf{0}} \tag{4}
\end{equation*}
$$

Form (B) :We know the following vector identity :

$$
\begin{equation*}
\text { curl curl } \overline{\mathbf{u}}=\operatorname{grad} \operatorname{div} \overline{\mathbf{u}}-\nabla^{2} \overline{\mathbf{u}} . \tag{5}
\end{equation*}
$$

Putting the value of $\nabla^{2} \overline{\mathbf{u}}$ from (5) into (4), we obtain

$$
(\lambda+\mu) \text { grad div } \overline{\mathbf{u}}+\mathbf{u}[\operatorname{grad} \operatorname{div} \overline{\mathbf{u}}-\operatorname{curl} \operatorname{curl} \overline{\mathbf{u}}]+\overline{\mathbf{F}}=\overline{\mathbf{0}}
$$

or

$$
\begin{equation*}
(\lambda+2 \mu) \text { grad div } \overline{\mathbf{u}}-\mu \operatorname{curl} \operatorname{curl} \overline{\mathbf{u}}+\overline{\mathbf{F}}=0 . \tag{6}
\end{equation*}
$$

Form (C) : Putting the value of grad div $\overline{\mathbf{u}}$ from (5) into (4), we get

$$
(\lambda+\mu)\left[\nabla^{2} \overline{\mathbf{u}}+\operatorname{curl} \operatorname{curl} \overline{\mathbf{u}}\right]+\mu \nabla^{2} \overline{\mathbf{u}}+\overline{\mathbf{F}}=\overline{\mathbf{0}}
$$

or

$$
\begin{equation*}
(\lambda+2 \mu) \nabla^{2} \overline{\mathbf{u}}+(\lambda+\mu) \text { curl curl } \overline{\mathbf{u}}+\overline{\mathbf{F}}=\overline{\mathbf{0}} \tag{7}
\end{equation*}
$$

Form D : We know that

$$
\begin{equation*}
\frac{\lambda+\mu}{\mu}=\frac{1}{1-2 \sigma} \tag{8}
\end{equation*}
$$

From (8) and (4), we find

$$
\begin{equation*}
\nabla^{2} \overline{\mathbf{u}}+\frac{1}{1-2 \sigma} \operatorname{grad} \operatorname{div} \overline{\mathbf{u}}+\frac{1}{\mu} \overline{\mathbf{F}}=\overline{\mathbf{0}} \tag{9}
\end{equation*}
$$

Dynamical Equations for an Isotropic Elastic Solid

Let $\rho$ be the density of the medium. The components of the force (mass $\times$ acceleration/volume) per unit volume are $\rho \frac{\partial^{2} u_{i}}{\partial t^{2}}$.

Hence, the dynamical equations in terms of the displacements $\mathbf{u}_{i}$ become

$$
(\lambda+\mu) \frac{\partial \theta}{\partial x_{i}}+\mu \nabla^{2} \mathbf{u}_{\mathbf{i}}+\mathbf{F}_{\mathbf{i}}=\rho \frac{\partial^{2} u_{i}}{\partial t^{2}}
$$

for $\mathbf{i}=\mathbf{1 , 2 , 3}$.
Various form of it can be obtained as above for equilibrium equations.
Question : In an isotropic elastic body in equilibrium under the body force $\underline{\mathbf{f}}=\mathbf{a} \mathbf{x}_{1} \mathbf{x}_{2} \hat{e}_{3}$, where $\mathbf{a}$ is a constant, the displacements are of the form

$$
\mathbf{u}_{1}=A \mathbf{x}_{1}^{2} \mathbf{x}_{2} \mathbf{x}_{3}, \mathbf{u}_{2}=\mathbf{B} \mathbf{x}_{1} \mathbf{x}_{2}^{2} \mathbf{x}_{3}, \mathbf{u}_{3}=\mathbf{C} \mathbf{x}_{1} \mathbf{x}_{2} \mathbf{x}_{3}^{2}
$$

where A, B, C are constants. Find A , B , C. Evaluate the corresponding stresses.

### 4.6. BELTRAMI-MICHELL COMPATIBILITY EQUATIONS IN TERMS OF THE STRESSES FOR AN ISOTROPIC SOLID

The strain - stress relations for an isotropic elastic solid are

$$
\begin{equation*}
\mathbf{e}_{\mathrm{ij}}=\frac{1+\sigma}{E} \tau_{\mathrm{ij}}-\frac{\sigma}{E} \delta_{\mathrm{ij}} \theta \quad, \quad \theta=\tau_{\mathrm{ij}} \tag{1}
\end{equation*}
$$

in which $\sigma$ is the Poisson's ratio and $E$ is the Young's modulus.
The Saint - Venant's compatibility equations in terms of strain components are

$$
\begin{equation*}
\mathbf{e}_{\mathbf{i j}, \mathrm{kl}}+\mathbf{e}_{\mathbf{k}, \mathrm{ij}}-\mathbf{e}_{\mathbf{i k}, \mathbf{j} \mathbf{l}}-\mathbf{e}_{\mathbf{j}, \mathrm{l}, \mathrm{ik}}=\mathbf{0}, \tag{2}
\end{equation*}
$$

which impose restrictions on the strain components to ensure that given $\mathbf{e}_{\mathbf{i j}}$ yield single - valued continuous displacements $u_{i}$.

When the region $\tau$ is simply connected, using (1) in (2), we find

$$
\begin{align*}
& \frac{1+\sigma}{E}\left\{\tau_{\mathrm{ij}, \mathrm{k} l}+\tau_{\mathrm{k}, \mathrm{ij}}-\tau_{\mathrm{i}, \mathrm{j} l}-\tau_{\mathrm{j} l, \mathrm{k} \mathrm{i}}\right\}=\frac{\sigma}{E}\left\{\delta_{\mathrm{ij}} \theta_{, \mathrm{k} l}+\delta_{\mathrm{k} l} \theta_{, \mathrm{ij}}-\delta_{\mathrm{ik}} \theta_{\mathrm{j} l}-\delta_{\mathrm{j} l} \theta_{, \mathbf{i k}}\right\} \\
& \tau_{\mathrm{ij}, \mathrm{k} l}+\tau_{\mathrm{k} l, \mathrm{ij}}-\tau_{\mathrm{ik}, \mathrm{j} l}-\tau_{\mathrm{j} l}, \mathrm{ki}=\frac{\sigma}{1+\sigma}\left(\delta_{\mathrm{ij}} \theta, \mathrm{k} l+\delta_{\mathrm{k} l} \theta, \mathrm{ij}-\delta_{\mathrm{ik}} \theta, \mathrm{j} l-\delta_{\mathrm{j} l} \theta, \mathrm{ik}\right) \tag{3}
\end{align*}
$$

with

$$
\tau_{\mathrm{ij}, \mathrm{k} l}=\frac{\partial^{2} \tau_{i j}}{\partial x_{k} \partial x_{l}}, \quad \theta_{, \mathrm{ij}}=\frac{\partial^{2} \theta}{\partial x_{i} \partial x_{j}} .
$$

These are equations of compatibility in stress components. These are $81\left(=3^{4}\right)$ in number but all of them are not independent. If i \& j or $k \& l$ are interchanged, we get same equations. Similarly for $\mathbf{i}=\mathbf{j}=\mathbf{k}=l$, equations are identically satisfied. Actually, the set of equations (3) contains only six independent equations obtained by setting

$$
\begin{aligned}
& \mathbf{k}=\boldsymbol{l}=\mathbf{1} \quad, \quad \mathbf{i}=\mathbf{j}=\mathbf{2} \\
& \mathbf{k}=\boldsymbol{l}=\mathbf{2} \quad, \quad \mathbf{i}=\mathbf{j}=\mathbf{3} \\
& \mathbf{k}=\boldsymbol{l}=\mathbf{3} \quad, \quad \mathbf{i}=\mathbf{j}=\mathbf{1} \\
& \mathrm{k}=l=\mathbf{1} \quad, \quad \mathbf{i}=\mathbf{2}, \mathbf{j}=\mathbf{3} \\
& \mathrm{k}=l=\mathbf{2} \quad, \quad \mathbf{i}=\mathbf{3}, \mathbf{j}=\mathbf{1} \\
& \mathrm{k}=l=\mathbf{3} \quad, \quad \mathbf{i}=\mathbf{1}, \mathbf{j}=\mathbf{2} .
\end{aligned}
$$

Setting $\mathrm{k}=l$ in (3) and then taking summation over the common index, we get
$\tau_{\mathbf{i j}, \mathbf{k k}}+\tau_{\mathbf{k k}, \mathbf{i j}}-\tau_{\mathbf{i k}, \mathbf{j k}}-\tau_{\mathbf{j k}, \mathbf{i k}}=\frac{\sigma}{1+\sigma}\left(\delta_{\mathbf{i j}} \theta_{, \mathbf{k k}}+\delta_{\mathbf{k k}} \theta_{, \mathbf{i j}}-\delta_{\mathbf{i k}} \theta_{, \mathbf{j k}}-\delta_{\mathbf{j k}} \theta_{, \mathbf{i k}}\right)$,
Since

$$
\begin{aligned}
& \theta_{, \mathrm{kk}}=\nabla^{2} \theta, \tau_{\mathrm{ij}, \mathrm{kk}}=\nabla^{2} \tau_{\mathrm{ij}}, \\
& \tau_{\mathrm{kk}, \mathrm{ij}}=\theta_{, \mathrm{ij}}, \delta_{\mathrm{kk}}=\mathbf{3},
\end{aligned}
$$

therefore, above equations become

$$
\nabla^{2} \tau_{\mathrm{ij}}+\theta_{, \mathrm{ij}}-\tau_{\mathrm{ik}, \mathbf{j k}}-\tau_{\mathrm{jk}, \mathbf{i k}}=\frac{\sigma}{1+\sigma}\left[\delta_{\mathrm{ij}} \nabla^{2} \theta+3 \theta_{, \mathrm{ij}}-2 \theta, \mathrm{ij}\right]
$$

or

$$
\begin{equation*}
\nabla^{2} \tau_{\mathrm{ij}}+\frac{1}{1+\sigma} \theta_{, \mathrm{ij}}-\tau_{\mathrm{ik}, \mathbf{j k}}-\tau_{\mathrm{jk}, \mathbf{i k}}=\frac{\sigma}{1+\sigma} \delta_{\mathrm{ij}} \nabla^{2} \theta \tag{4}
\end{equation*}
$$

This is a set of 9 equations and out of which only 6 are independent due to the symmetry of i \& j . In combining equations (3) linearly, the number of independent equations is not reduced.

Hence the resultant set of equations in (4) is equivalent to the original equations in (3).

Equilibrium equations are

$$
\tau_{\mathrm{ik}, \mathbf{k}}+\mathbf{F}_{\mathrm{i}}=\mathbf{0},
$$

where $F_{i}$ is the body force per unit volume.
Differentiating these equations with respect to $\mathrm{x}_{\mathrm{j}}$, we get

$$
\begin{equation*}
\tau_{i k, k j}=-\mathbf{F}_{\mathbf{i}, \mathbf{j}} . \tag{5}
\end{equation*}
$$

Using (5), equation (4) can be rewritten in the form

$$
\begin{equation*}
\nabla^{2} \tau_{\mathrm{ij}}+\frac{1}{1+\sigma} \theta, \mathrm{ij}-\frac{\sigma}{1+\sigma} \delta_{\mathrm{ij}} \nabla^{2} \theta=-\left(\mathbf{F}_{\mathbf{i}, \mathbf{j}}+\mathbf{F}_{\mathbf{j}, \mathbf{i}}\right) . \tag{6}
\end{equation*}
$$

Setting $j=i$ in (6) and adding accordingly, we write

$$
\begin{array}{r}
\nabla^{2} \theta+\frac{1}{1+\sigma} \nabla^{2} \theta-\frac{3 \sigma}{1+\sigma} \nabla^{2} \theta=-2 \mathbf{F}_{\mathrm{i}, \mathrm{i}} \\
\left(1+\frac{1}{1+\sigma}-\frac{3 \sigma}{1+\sigma}\right) \nabla^{2} \theta=-2 \mathrm{~F}_{\mathrm{i}, \mathrm{i}} \\
\frac{2(1-\sigma)}{1+\sigma} \nabla^{2} \theta=-2 \mathrm{~F}_{\mathrm{i}, \mathrm{i}}=-\mathbf{2} d i v \vec{F}
\end{array}
$$

giving

$$
\begin{equation*}
\nabla^{2} \theta=-\frac{1+\sigma}{1-\sigma} d i v \vec{F} \tag{7}
\end{equation*}
$$

Using (7) in (6), we find the final form of the compatibility equations in terms of stresses.

We get

$$
\begin{equation*}
\nabla^{2} \tau_{\mathrm{ij}}+\frac{1}{1+\sigma} \theta, \mathrm{ij}=-\frac{\sigma}{1-\sigma} \delta_{\mathrm{ij}} \operatorname{div} \vec{F}-\left(\mathbf{F}_{\mathbf{i}, \mathbf{j}}+\mathbf{F}_{\mathbf{j}, \mathbf{i}}\right) . \tag{8}
\end{equation*}
$$

These equations in cartesian coordinates ( $x, y, z$ ) can be written as

$$
\left.\nabla^{2} \tau_{\mathrm{xx}}+\frac{1}{1+\sigma} \frac{\partial^{2} \theta}{\partial x^{2}}=-\frac{\sigma}{1-\sigma} \operatorname{div} \vec{F}-\mathbf{2} \frac{\partial F_{x}}{\partial x}\right)
$$

$$
\left.\begin{array}{l}
\nabla^{2} \tau_{\mathrm{yy}}+\frac{1}{1+\sigma} \frac{\partial^{2} \theta}{\partial y^{2}}=-\frac{\sigma}{1-\sigma} \operatorname{div} \vec{F}-\mathbf{2} \frac{\partial F_{y}}{\partial y} \\
\nabla^{2} \tau_{\mathrm{zz}}+\frac{1}{1+\sigma} \frac{\partial^{2} \theta}{\partial z^{2}}=-\frac{\sigma}{1-\sigma} \operatorname{div} \vec{F}-\mathbf{2} \frac{\partial F_{z}}{\partial z} \\
\nabla^{2} \tau_{\mathrm{yz}}+\frac{1}{1+\sigma} \frac{\partial^{2} \theta}{\partial y \partial z}=-\left(\frac{\partial F_{y}}{\partial z}+\frac{\partial F_{z}}{\partial y}\right) \\
\nabla^{2} \tau_{\mathrm{zx}}+\frac{1}{1+\sigma} \frac{\partial^{2} \theta}{\partial z \partial x}=-\left(\frac{\partial F_{z}}{\partial x}+\frac{\partial F_{x}}{\partial z}\right)  \tag{9}\\
\nabla^{2} \tau_{\mathrm{xy}}+\frac{1}{1+\sigma} \frac{\partial^{2} \theta}{\partial x \partial y}=-\left(\frac{\partial F_{x}}{\partial y}+\frac{\partial F_{y}}{\partial x}\right)
\end{array}\right\}
$$

In 1892 , Beltrami obtained these equations for $\vec{F}=\overrightarrow{0}$ and in 1900 Michell obtained then in form as given in (9).

These equations in (9) are called the Beltrami - Michell compatibility equations.

Definition : A function $\mathbf{V}$ of class $\mathbf{C}^{\mathbf{4}}$ is called a biharmonic function when

$$
\nabla^{2} \nabla^{2} \mathbf{V}=\mathbf{0}
$$

Theorem : When the components of the body $\vec{F}$ are constants, show that the stress and strain invariants $\theta$ and $v$ are harmonic functions and the stress components $\tau_{\mathrm{ij}}$ and strain components $\mathrm{e}_{\mathrm{ij}}$ are biharmonic functions.

Proof : The Beltrami - Michell compatibility equations in terms of stress are

$$
\begin{equation*}
\nabla^{2} \tau_{\mathrm{ij}}+\frac{1}{1+\sigma} \theta_{, \mathrm{ij}}=-\frac{\sigma}{1-\sigma} \delta_{\mathrm{ij}} \operatorname{div} \vec{F}-\left(\mathbf{F}_{\mathbf{i}, \mathbf{j}}+\mathbf{F}_{\mathbf{j}, \mathbf{i}}\right) \tag{1}
\end{equation*}
$$

in which $\vec{F}$ is the body force per unit volume.
It is given that the vector $\vec{F}$ is constant. In this case, equations in (1) reduce to

$$
\begin{equation*}
\nabla^{2} \tau_{\mathrm{ij}}+\frac{1}{1+\sigma} \theta_{\mathrm{ij}}=\mathbf{0} \tag{2}
\end{equation*}
$$

Setting $j=i$ in (2) and taking summation accordingly, we get

$$
\begin{align*}
& \nabla^{2} \tau_{\mathrm{ii}}+\frac{1}{1+\sigma} \theta_{, \mathrm{ii}}=\mathbf{0} \\
& \nabla^{2} \theta+\frac{1}{1+\sigma} \nabla^{2} \theta=\mathbf{0} \\
& \left(1+\frac{1}{1+\sigma}\right) \nabla^{2} \theta=\mathbf{0} \\
& \nabla^{2} \theta=\mathbf{0} \tag{3}
\end{align*}
$$

This shows that the stress invariant $\theta=\tau_{\mathbf{k k}}$ is a harmonic function.
The standard relation between the invariants $\theta$ and $v$ is

$$
\begin{equation*}
\theta=(3 \lambda+2 \mu) v, \tag{4}
\end{equation*}
$$

and equation (3) implies that

$$
\begin{equation*}
\nabla^{2} \boldsymbol{v}=\mathbf{0} \tag{5}
\end{equation*}
$$

showing that the strain invariant $\mathrm{v}=\mathrm{e}_{\mathrm{kk}}$ is also a harmonic function.
Again

$$
\begin{align*}
\nabla^{2} \nabla^{2} \tau_{\mathrm{ij}} & =\nabla^{2}\left(-\frac{1}{1+\sigma} \theta_{, i j}\right) \\
& =-\frac{1}{1+\sigma} \nabla^{2}(\theta, \mathrm{ij}) \\
& =-\frac{1}{1+\sigma}\left(\nabla^{2} \theta\right)_{, \mathrm{ij}} \\
\nabla^{2} \nabla^{2} \tau_{\mathrm{ij}} & =\mathbf{0} \tag{6}
\end{align*}
$$

giving
This shows that the stress components $\tau_{\mathrm{ij}}$ are biharmonic functions.
The following strain - stress relations

$$
\mathbf{e}_{\mathrm{ij}}=\frac{-\lambda}{2 \mu(3 \lambda+2 \mu)} \delta_{\mathrm{ij}} \theta+\frac{1}{2 \mu} \tau_{\mathrm{ij}}
$$

give

$$
\begin{gather*}
\nabla^{2} \nabla^{2} \mathbf{e}_{\mathrm{ij}}=\frac{-\lambda}{2 \mu(3 \lambda+2 \mu)} \delta_{\mathrm{ij}} \nabla^{2} \nabla^{2} \theta+\frac{1}{2 \mu} \nabla^{2} \nabla^{2} \tau_{\mathrm{ij}} \\
\nabla^{2} \nabla^{2} \mathbf{e}_{\mathrm{ij}}=\mathbf{0} \tag{7}
\end{gather*}
$$

Equation (7) shows that the strain components $\mathbf{e}_{\mathrm{ij}}$ are also biharmonic functions.

Theorem 2: If the body force $\vec{F}$ is derived from a harmonic potential function, show that the strain and stress invariants $\mathbf{e}_{\mathbf{k k}} \& \tau_{\mathbf{k k}}$ are harmonic functions and the strain and stress components are biharmonic function.

Proof : Let $\phi$ be the potential function and $\vec{F}$ is derived from $\phi$ so that

$$
\begin{equation*}
\vec{F}=\underline{\nabla} \phi \quad \text { or } \quad \mathbf{F}_{\mathbf{j}}=\phi, \mathbf{j} \tag{1}
\end{equation*}
$$

Then

$$
\begin{equation*}
\operatorname{div} \vec{F}=\phi_{\mathrm{jj}}=\nabla^{2} \phi=\mathbf{0} \tag{2}
\end{equation*}
$$

since $\phi$ is a harmonic function (given). Further

$$
\begin{equation*}
\mathbf{F}_{\mathbf{i}, \mathbf{j}}=\mathbf{F}_{\mathbf{j}, \mathbf{i}}=\phi, \mathbf{i j} \tag{3}
\end{equation*}
$$

The Beltrami - Michell compatibility equations in terms of stresses, in this case, reduce to

$$
\begin{equation*}
\nabla^{2} \tau_{\mathrm{ij}}+\frac{1}{1+\sigma} \theta_{, \mathrm{ij}}=-\mathbf{2} \phi_{, \mathrm{ij}} . \tag{4}
\end{equation*}
$$

Putting $\mathbf{j}=\mathrm{i}$ in (4) and taking the summation accordingly, we obtain

$$
\begin{aligned}
\nabla^{2} \tau_{\mathrm{ii}}+\frac{1}{1+\sigma} \theta_{\mathrm{ii}} & =-\mathbf{2} \phi_{, \mathrm{ii}} \\
& =-\mathbf{2} \nabla^{2} \phi \\
& =\mathbf{0}
\end{aligned}
$$

giving

$$
\begin{equation*}
\nabla^{2} \theta=\mathbf{0}, \tag{5}
\end{equation*}
$$

This shows that $\theta$ is harmonic.
The relation $\theta=(3 \lambda+2 \mu) v$ immediately shows that $v$ is also harmonic.

From equation (4), we write

$$
\begin{aligned}
\nabla^{2} \nabla^{2} \tau_{\mathrm{ij}}+\frac{1}{1+\sigma} \nabla^{2} \theta_{, \mathrm{ij}} & =-2\left(\nabla^{2} \phi\right)_{, \mathrm{ij}} \\
& =\mathbf{0}
\end{aligned}
$$

This gives
as

$$
\begin{align*}
& \nabla^{2} \nabla^{2} \tau_{\mathrm{ij}}=\mathbf{0}  \tag{6}\\
& \nabla^{2} \theta=\nabla^{2} \phi=\mathbf{0}
\end{align*}
$$

It shows that the components $\tau_{i \mathrm{j}}$ are biharmonic.
The strain - stress relations yields that the strain components are also biharmonic function.

Question : Find whether the following stress system can be a solution of an elastostatic problem in the absence of body forces :

$$
\tau_{11}=\mathbf{x}_{2} \mathbf{x}_{3}, \tau_{22}=\mathbf{x}_{3} \mathbf{x}_{1}, \tau_{12}=\mathbf{x}_{3}{ }^{2}, \tau_{13}=\tau_{33}=\tau_{32}=\mathbf{0}
$$

Solution : In order that the given stress system can be a solution of an elastostatic problem in the absence of body forces, the following equations are to be satisfied :
(i) Cauchy's equations of equilibrium with $f_{i}=0$, i.e.,

$$
\begin{align*}
& \tau_{11,1}+\tau_{12,2}+\tau_{13,3}=0 \\
& \tau_{12,1}+\tau_{22,2}+\tau_{23,3}=0, \\
& \tau_{13,1}+\tau_{23,2}+\tau_{33,3}=\mathbf{0} \tag{1}
\end{align*}
$$

(ii) Beltrami - Michell equations with $f_{i}=0$, i.e.,

$$
\begin{aligned}
& \nabla^{2} \tau_{11}+\frac{1}{1+\sigma}\left(\tau_{11}+\tau_{22}+\tau_{33}\right)_{, 11}=\mathbf{0} \\
& \nabla^{2} \tau_{22}+\frac{1}{1+\sigma}\left(\tau_{11}+\tau_{22}+\tau_{33}\right)_{, 22}=\mathbf{0} \\
& \nabla^{2} \tau_{33}+\frac{1}{1+\sigma}\left(\tau_{11}+\tau_{22}+\tau_{33}\right)_{, 33}=\mathbf{0}
\end{aligned}
$$

$$
\begin{align*}
& \nabla^{2} \tau_{12}+\frac{1}{1+\sigma}\left(\tau_{11}+\tau_{22}+\tau_{33}\right)_{, 12}=0 \\
& \nabla^{2} \tau_{13}+\frac{1}{1+\sigma}\left(\tau_{11}+\tau_{22}+\tau_{33}\right)_{, 13}=0 \\
& \nabla^{2} \tau_{23}+\frac{1}{1+\sigma}\left(\tau_{11}+\tau_{22}+\tau_{33}\right)_{, 23}=0 \tag{2}
\end{align*}
$$

It is easy to check that all the equilibrium equations in (1) are satisfied.
Moreover, all except the fourth one in (2) are satisfied by the given stress system.

Since the given system does not satisfy the Beltrami - Michell equations fully, it can not form a solution of an elastostatic problem.

Remark : The example illustrates the important fact that a stress system may not be a solution of an elasticity problem even though it satisfies Cauchy's equilibrium equations.

Exercise 1: Show that the stress - system $\tau_{11}=\tau_{22}=\tau_{13}=\tau_{23}=\tau_{12}=\mathbf{= 0}, \tau_{33}$ $=\rho g x_{3}$, where $\rho$ and $g$ are constants, satisfies that equations of equilibrium and the equations of compatibility for a suitable body force.

Exercise 2: Show that the following stress system can not be a solution of an elastostatic problem although it satisfies cauchy's equations of equilibrium with zero body forces :

$$
\begin{aligned}
& \tau_{11}=\mathbf{x}_{2}{ }^{2}+\sigma\left(\mathbf{x}_{1}{ }^{2}-\mathbf{x}_{2}{ }^{2}\right), \tau_{22}=\mathbf{x}_{1}{ }^{2}+\sigma\left(\mathbf{x}_{2}{ }^{2}-\mathbf{x}_{1}{ }^{2}\right), \tau_{33}=\sigma\left(\mathbf{x}_{1}{ }^{2}+\mathbf{x}_{2}{ }^{2}\right) \\
& \tau_{12}=-2 \sigma \mathbf{x}_{1} \mathbf{x}_{2}, \tau_{23}=\tau_{31}=0
\end{aligned}
$$

where $\sigma$ is a constant of elasticity.
Exercise 3: Determine whether or not the following stress components are a possible solution in elastostatics in the absence of body forces :

$$
\begin{aligned}
& \tau_{11}=\mathbf{a} \mathbf{x}_{2} \mathbf{x}_{3}, \tau_{22}=\mathbf{b} \mathbf{x}_{3} \mathbf{x}_{1}, \tau_{33}=\mathbf{c} \mathbf{x}_{1} \mathbf{x}_{2}, \tau_{12}=\mathbf{d} \mathbf{x}_{3}{ }^{2} \\
& \tau_{13}=\mathbf{e} \mathbf{x}_{2}^{2}, \tau_{23}=\mathbf{f} \mathbf{x}_{1}^{2},
\end{aligned}
$$

where $a, b, c, d, e, f$ are all constants.
Exercise 4: In an elastic body in equilibrium under the body force $\underline{\mathbf{f}}=\mathbf{a} \mathbf{x}_{\mathbf{1}}$ $\mathbf{x}_{2} \hat{e}_{3}$, where a is a constant, the stresses are of the form

$$
\begin{aligned}
& \tau_{11}=a \mathbf{x}_{1} \mathbf{x}_{2} \mathbf{x}_{3}, \tau_{22}=b \mathbf{x}_{1} \mathbf{x}_{2} \mathbf{x}_{3}, \tau_{33}=\mathbf{c} \mathbf{x}_{1} \mathbf{x}_{2} \mathbf{x}_{3} \\
& \tau_{12}=\left(a \mathbf{x}_{1}{ }^{2}+b \mathbf{x}_{2}{ }^{2}\right) \mathbf{x}_{3}, \tau_{23}=\left(b \mathbf{x}_{2}{ }^{2}+\mathbf{c} \mathbf{x}_{3}{ }^{2}\right) \mathbf{x}_{1}, \tau_{13}=\left(\mathbf{c} \mathbf{x}_{3}{ }^{2}+a \mathbf{x}_{1}{ }^{2}\right) \mathbf{x}_{2}
\end{aligned}
$$

where a, b, c are constants. Find these constants.
Exercise 5: Define the stress function $\mathbf{S}$ by

$$
\tau_{\mathrm{ij}}=\mathbf{S}_{, \mathrm{ij}}=\frac{\partial^{2} S}{\partial x_{i} \partial x_{j}}
$$

and consider the case of zero body force. Show that, if $\sigma=0$, then the equilibrium and compatibility equations reduce to

$$
\nabla^{2} \mathbf{S}=\text { Constant }
$$

### 4.7. UNIQUENESS OF SOLUTION

The most general problem of the elasticity theory is to determine the distribution of stresses and strains as well as displacements at all points of a body when certain boundary conditions and certain initial conditions are specified (under the assumption that the body force $\underline{f}$ is known before hand).
In the linear elasticity, the displacements, strains and stresses are governed by the following equations
(I) $\quad \mathbf{e}_{\mathbf{i j}}=\frac{1}{2}\left(\mathbf{u}_{\mathrm{i}, \mathrm{j}}+\mathbf{u}_{\mathrm{j}, \mathrm{i}}\right)$ strain - displacement relations
(II) $\quad \tau_{i j}=\lambda \delta_{i j} \mathbf{e}_{\mathrm{kk}}+2 \mu \mathbf{e}_{\mathrm{ij}}=\lambda \delta_{\mathrm{ij}} \mathbf{u}_{\mathrm{k}, \mathrm{k}}+\mu\left(\mathbf{u}_{\mathrm{i}, \mathrm{j}}+\mathbf{u}_{\mathrm{j}, \mathrm{i}}\right)$ material law
or

$$
\mathbf{e}_{\mathrm{ij}}=-\frac{\sigma}{E} \delta_{\mathrm{ij}} \tau_{\mathrm{kk}}+\frac{1+\sigma}{E} \tau_{\mathrm{ij}}
$$

$$
\begin{equation*}
\tau_{i, j, j}+f_{i}=0 \text { Cauchy's equation of equilibrium } \tag{III}
\end{equation*}
$$

or
$(\lambda+\mu) \operatorname{grad} \operatorname{div} \overline{\mathbf{u}}+\mu \nabla^{2} \overline{\mathbf{u}}+\overline{\mathbf{f}}=\overline{\mathbf{0}} \quad$ Navier equation of equilibrium Accordingly, solving a problem in linear elasticity generally amounts to solving these equations for $\mathrm{u}_{\mathrm{i}}, \mathrm{e}_{\mathrm{ij}}$ and $\tau_{\mathrm{ij}}$ under certain specified boundary conditions and initial conditions.
Let a body occupying a region $V$ has the boundary $S$.
Initially, It is assumed that the body is in the undeformed state. That is,

$$
\begin{equation*}
\mathrm{u}_{\mathrm{i}}=0 \text { for } \mathrm{e}_{\mathrm{ij}}=0 \quad \text { in } \mathrm{V} . \quad(\text { at time } \mathbf{t}=\mathbf{0}) \tag{4}
\end{equation*}
$$

The boundary conditions specified are usually of one of following three kinds :
(i) The stress vector is specified at every point of boundary $S$ for all times, i.e.,

$$
\stackrel{v}{T}=s \quad \text { on } \mathbf{S}
$$

where $s$ is a known vector point function.
(ii) the displacement vector is specified at every point of $S$ and for all times , i.e. ,

$$
\begin{equation*}
u=u^{*} \quad \text { on S } \tag{6}
\end{equation*}
$$

when $u$ * is a known function.
(ii) the stress vector is specified at every point of a part $S_{\tau}$ of $S$ and the displacement vector is specified at every point of the remaining $\operatorname{part} S_{u}=S-S_{\tau}$; i.e. ,

$$
\left.\begin{array}{ll}
\underset{\sim}{v}=\underset{\sim}{s} &  \tag{7}\\
\underset{\sim}{u}={\underset{\sim}{u}}^{*} & \\
\text { on } \mathbf{S}_{\tau} \\
\mathbf{o n ~}_{\mathbf{u}}=\mathbf{S}-\mathbf{S}_{\tau}
\end{array}\right\}
$$

The problem of solving equations (1) , (2) , (3) under the initial condition (4) and one of the boundary conditions in (5) - (7) to determine $u_{i}$, $\mathbf{e}_{\mathrm{ij}}, \tau_{\mathrm{ij}}$ is known as boundary value problem in elastostatics.

A set $\left\{\mathbf{u}_{\mathrm{i}}, \mathbf{e}_{\mathrm{ij}}, \tau_{\mathrm{ij}}\right\}$ so obtained/determined, if it exists , is called a solution of the problem.

When the boundary condition is of the form (5), the problem is referred to as the traction (or stress) boundary value problem ; and when the boundary condition of the form (6), the problem is referred to as the displacement BV problem ; and when the boundary condition is of the form (7), the problem is referred to as mixed BV problem.

The three problems are together called the fundamental boundary value problems. The boundary conditions valid for all the three problems can be written down in the form of (7). For the traction problem $S_{\tau}=S$ and $S_{u}=\phi$ , for the displacement problem $S=S_{u}$ and $S_{\tau} \neq \phi$, and for mixed problem $\mathbf{S}_{\mathbf{u}} \neq \mathbf{S} \neq \mathbf{S}_{\tau}$.

Uniqueness : The solutions of an elastostatic problem governed by equations (1) - (3) and the boundary condition (7) is unique within a rigid body displacement.

Remark 1: The displacement boundary value problem is completely solved if one obtains a solution of the Navier equation subject to the boundary condition (6). Note that we need not adjoin the compatibility equations

$$
\begin{equation*}
\mathbf{e}_{\mathbf{i j}, \mathbf{k} l}+\mathbf{e}_{\mathbf{k} l, \mathbf{j}}-\mathbf{e}_{\mathbf{i k}, \mathbf{j} l}-\mathbf{e}_{\mathrm{j} l, \mathbf{i k}}=\mathbf{0}, \tag{8}
\end{equation*}
$$

for the only purpose of the latter is to impose restrictions on the strain components that shall ensure that the $\mathrm{e}_{\mathrm{ij}}$ yield single - valued continuous displacements $u_{i}$, when the region is simply connected. From the knowledge of functions $u_{i}$, one can determine the strains, and hence stresses by making use of Hooke's law in (2).

Remark 2: The stress boundary value problem suggests the desirability of expressing all the differential equations entirely in terms of stress. The compatibility equations (Betrami - Michell compatibility equations) in terms of stresses are

$$
\begin{equation*}
\nabla^{2} \tau_{\mathrm{ij}}+\frac{1}{1+\sigma} \tau_{\mathrm{kk}, \mathrm{ij}}=-\frac{\sigma}{1-\sigma} \delta_{\mathrm{ij}} \operatorname{div} \bar{F}-\left(\mathbf{F}_{\mathrm{i}, \mathrm{j}}+\mathbf{F}_{\mathrm{j}, \mathrm{i}}\right) \tag{9}
\end{equation*}
$$

In order to determine the state of stress in the interior of an elastic body, one must solve the system of equations consisting of Cauchy's equations of equilibrium (3) and $B-M$ compatibility equations subject to the boundary conditions in (5).

### 4.8. ST. VENANT'S PRINCIPAL

In the analysis of actual structures subjected to external loads, it can be invariably found that the distributions of surface forces are so complex as to define them more accurately for solving the appropriate governing equations. It is true that the solutions obtained for such problems using the elasticity equations are exact only if the external loads are applied in a specific manner. However, in many cases it is possible to predict the net effect of the external surface tractions without worrying about the precise manner in which they are distributed over the boundary.

In 1853 in his "Memoire Sur la Torsion des Prismes", Saint - Vanant developed solutions for the torsion of prismatic bars which gave the same stress distribution for all cross - sections. He attempted to justify the usefulness of his formulation by the following : " The fact is that the means of application and distribution of the forces towards the extremities of the prisms is immaterial to the perceptible effects produced on the rest of the length, so that one can always, in a sufficiently similar manner, replace the forces applied with equivalent static forces or with those having the same total moments and the same resultant forces".

St. Venant's Principal : If a certain distribution of forces acting on a portion of the surface of a body is replaced by different distribution of forces acting on the same portion of the body, then the effects of the two different distributions on the parts sufficiently far removed (large compared to linear dimensions of the body) from the region of application of forces, are essentially the same, provided that the two distributions of forces are statically equivalent (that is, the same resultant forces and the same resultant moment).

St . Venant principal is profitable when solving problems in rigid - body mechanics to employ the concept of a point force when we had a force distribution over a small area. At other times, we employed the rigid -
body resultant force system of some distribution in the handling of a problem. Such replacements led to reasonably accurate and direct solutions. From the viewpoint of rigid - body mechanics, this principal states that the stresses reasonably distant from an applied load on a boundary are not significantly altered if this load is changed to another load which is equivalent to it. We may call such a second load the statically equivalent load.

This principal is actually summarized by one of its more detailed names, "The principal of the elastic equivalence of statically equipollent systems of load".

St. Venant's principal is very convenient and useful in obtaining solutions to various problems in elasticity. However, the statements are in general vague. They do not specifically state either the extent of the region within which the effects of two different statically equivalent force systems are not quite the same or the magnitude of the error.

Therefore , St. Venant's principal is only qualitative and expresses only a trend.

Nevertheless , St. Venant's principal has many important implications with respect to many practical problems. For instance, in many structures , the overall deflections are not unduly affected by the local changes in the distribution of forces or localized stress concentrations due the holes, cracks, etc. But it should be realized that the presence of defects in a region , or a non - uniform application of load will cause changes in stress distribution.

This principle is mainly used in elasticity to solve the problems of extension/bending/torsion of elastic beams. Under this technique , certain assumptions about the components of stress, strain and displacements are made, while leaving enough degree of freedom, so that the equations of equilibrium and compatibility are satisfied. The solution so obtained will be unique by the uniqueness of solution of the general boundary - value problems of linear elasticity.

Table for various elastic coefficients for an isotropic media

| $\lambda$ | $\frac{2 \mu \sigma}{1-2 \sigma}$ | $\frac{\mu(\mathrm{E}-2 \mu)}{3 \mu-\mathrm{E}}$ | $\mathrm{k}-\frac{2}{3} \mu$ | $\frac{\mathrm{E} \sigma}{(1+\sigma)(1-2 \sigma)}$ | $\frac{3 \mathrm{k} \sigma}{1+\sigma}$ | $\frac{3 \mathrm{k}(3 \mathrm{k}-\mathrm{E})}{(9 \mathrm{k}-\mathrm{E})}$ | - | - | - |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\mu$ | - | - | - | $\frac{E}{2(1+\sigma)}$ | $\frac{3 \mathrm{k}(1-2 \sigma)}{2(1+\sigma)}$ | $\frac{3 \mathrm{kE}}{9 \mathrm{k}-\mathrm{E}}$ | - | $\frac{\lambda(1-2 \sigma)}{2 \sigma}$ | $\frac{3}{2}(\mathbf{k}-\lambda)$ |
| $\sigma$ | - | $\frac{E}{2 \mu}-1$ | $\frac{3 \mathrm{k}-2 \mu}{2(3 \mathrm{k}+\mu)}$ | - | - | $\frac{3 \mathrm{k}-\mathrm{E}}{6 \mathrm{k}}$ | $\frac{\lambda}{2(\lambda+\mu)}$ | - | $\frac{\lambda}{3 k-\lambda}$ |
| E | $2 \mu(1+\sigma)$ | - | $\frac{9 \mathrm{k} \mu}{3 \mathrm{k}+\mu}$ | - | $3 \mathrm{k}(1-2 \sigma)$ | - | $\frac{\mu(3 \lambda+2 \mu)}{\lambda+\mu}$ | $\frac{\lambda(1+\sigma)(1-2 \sigma)}{\sigma}$ | $\frac{9 \mathrm{k}(\mathrm{k}-\lambda)}{3 \mathrm{k}-\lambda}$ |
| k | $\frac{2 \mu+\sigma}{3(1-2 \sigma)}$ | $\frac{\mu \mathrm{E}}{3(3 \mu-\mathrm{E})}$ | - | $\frac{E}{3(1-2 \sigma)}$ | - | - | $\lambda+\frac{2}{3} \mu$ | $\frac{\lambda(1+\sigma)}{3 \sigma}$ | - |

## Chapter-5 <br> Strain-Energy Function

### 5.1. INTRODUCTION

The energy stored in an elastic body by virtue of its deformation is called the strain energy. This energy is acquired by the body when the body forces and surface tractions do some work. This is also termed as internal energy. It depends upon the shape and temperature of the body.

### 5.2. STRAIN - ENERGY FUNCTION

Let $\tau_{\mathrm{ij}}$ be the stress tensor and $\mathrm{e}_{\mathrm{ij}}$ be the strain tensor for an infinitesimal affine deformation of an elastic body. We write

$$
\left.\begin{array}{l}
\tau_{11}=\tau_{1}, \tau_{22}=\tau_{2}, \tau_{33}=\tau_{4}  \tag{1}\\
\tau_{23}=\tau_{4}, \tau_{31}=\tau_{5}, \tau_{12}=\tau_{6}
\end{array}\right\}
$$

and

$$
\left.\begin{array}{l}
\mathrm{e}_{11}=\mathrm{e}_{1}, \mathrm{e}_{22}=\mathrm{e}_{2}, \mathrm{e}_{33}=\mathrm{e}_{3}  \tag{2}\\
2 \mathrm{e}_{23}=\mathrm{e}_{4}, 2 \mathrm{e}_{13}=\mathrm{e}_{5}, 2 \mathrm{e}_{12}=\mathrm{e}_{6}
\end{array}\right\},
$$

in terms of engineering notations.
We assume that the deformation of the elastic body is isothermal or adiabatic. Love(1944) has proved that, under this assumption there exists a function of strains

$$
\begin{equation*}
W=W\left(e_{1}, e_{2}, e_{3}, \ldots \ldots, e_{6}\right), \tag{3}
\end{equation*}
$$

with the property

$$
\begin{equation*}
\frac{\partial \mathrm{W}}{\partial \mathrm{e}_{\mathrm{i}}}=\tau_{\mathrm{i}} \quad, \quad \text { for } \mathrm{i}=1,2, \ldots . ., 6 . \tag{4}
\end{equation*}
$$

## This function $\mathbf{W}$ is called the strain energy function.

W represents strain energy, per unit of undeformed volume, stored up in the body by the strains $\mathrm{e}_{\mathrm{i}}$.

The units of W are $\frac{\text { force } \times L}{L^{3}}=\frac{\text { force }}{L^{2}}$, that of a stress.

The existence of W was first introduced by George Green (1839).
Expanding the strain energy function W , given in (3) in a power series in terms of strains $e_{i}$, we write

$$
\begin{equation*}
2 \mathrm{~W}=\mathrm{d}_{0}+2 \mathrm{~d}_{\mathrm{i}} \mathrm{e}_{\mathrm{i}}+\mathrm{d}_{\mathrm{ij}} \mathrm{e}_{\mathrm{i}} \mathrm{e}_{\mathrm{j}} \quad, \mathrm{i}, \mathrm{j}=1,2, \ldots, 6 \tag{5}
\end{equation*}
$$

after discarding all terms of order 3 and higher in the strains $e_{i}$ as strains $e_{i}$ are assumed to be small. In second term, summation of is to be taken and in $3^{\text {rd }}$ term , summation over dummy sufficies i \& j are to be taken.

In the natural state, $\mathbf{e}_{\mathbf{i}}=\mathbf{0}$, consequently $\mathrm{W}=0$ for $\mathbf{e}_{\mathbf{i}}=\mathbf{0}$.
This gives

$$
\begin{equation*}
\mathrm{d}_{0}=0 . \tag{6}
\end{equation*}
$$

Even otherwise, the constant term in (5) can be neglected since we are interested only in the partial derivatives of W. Therefore ,equations (5) and (6) yield

$$
\begin{equation*}
2 \mathrm{~W}=2 \mathrm{~d}_{\mathrm{i}} \mathrm{e}_{\mathrm{i}}+\mathrm{d}_{\mathrm{ij}} \mathrm{e}_{\mathrm{i}} \mathrm{e}_{\mathrm{j}} \tag{7}
\end{equation*}
$$

This gives

$$
\begin{aligned}
\frac{\partial \mathrm{W}}{\partial \mathrm{e}_{\mathrm{k}}} & =\mathrm{d}_{\mathrm{i}} \delta_{\mathrm{ik}}+\frac{1}{2} \frac{\partial}{\partial e_{k}}\left\{\mathrm{~d}_{\mathrm{ij}} \mathrm{e}_{\mathrm{i}} \mathrm{e}_{\mathrm{j}}\right\} \\
& =\mathrm{d}_{\mathrm{k}}+\frac{1}{2}\left\{\mathrm{~d}_{\mathrm{ij}} \delta_{\mathrm{ki}} \mathrm{e}_{\mathrm{j}}+\mathrm{d}_{\mathrm{ij}} \mathrm{e}_{\mathrm{i}} \delta_{\mathrm{kj}}\right] \\
& =\mathrm{d}_{\mathrm{k}}+\frac{1}{2}\left[\mathrm{~d}_{\mathrm{kj}} \mathrm{e}_{\mathrm{j}}+\mathrm{d}_{\mathrm{ki}} \mathrm{e}_{\mathrm{i}}\right] \\
& =\mathrm{d}_{\mathrm{k}}+\frac{1}{2}\left(\mathrm{~d}_{\mathrm{kj}}+\mathrm{d}_{\mathrm{kj}}\right) \mathrm{e}_{\mathrm{j}} \\
& =\mathrm{d}_{\mathrm{k}}+\left(\mathrm{d}_{\mathrm{kj}} \mathrm{e}_{\mathrm{j}}\right) \mathrm{e}_{\mathrm{j}} .
\end{aligned}
$$

This gives

$$
\begin{equation*}
\tau_{\mathrm{i}}=\mathrm{d}_{\mathrm{i}}+\mathrm{c}_{\mathrm{ij}} \mathrm{e}_{\mathrm{j}} \tag{8}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathrm{c}_{\mathrm{ij}}=\frac{1}{2}\left(\mathrm{~d}_{\mathrm{ij}}+\mathrm{d}_{\mathrm{ji}}\right)=\mathrm{c}_{\mathrm{ji}} . \tag{9}
\end{equation*}
$$

We observe that $\mathbf{c}_{\mathrm{ij}}$ is symmetric.
We further assume that the stresses $\tau_{i}=0$ in the undeformed state, when $\mathrm{e}_{\mathrm{i}}=0$.
This assumption , using equation (8), gives

$$
\begin{equation*}
\mathbf{d}_{\mathbf{i}}=\mathbf{0}, \quad \mathrm{i}=1,2, \ldots . ., 6 . \tag{10}
\end{equation*}
$$

Equations (7) , (8) and (10) give

$$
\begin{equation*}
\tau_{i}=c_{i j} e_{j} \tag{11}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathrm{W}=\frac{1}{2} \mathrm{c}_{\mathrm{ij}} \mathrm{e}_{\mathrm{i}} \mathrm{e}_{\mathrm{j}}=\frac{1}{2} \mathrm{e}_{\mathrm{i}} \tau_{\mathrm{i}} \tag{12}
\end{equation*}
$$

since, two quadric homogeneous forms for W are equal as

$$
\begin{equation*}
\mathrm{d}_{\mathrm{ij}} \mathrm{e}_{\mathrm{i}} \mathrm{e}_{\mathrm{j}}=\mathrm{c}_{\mathrm{ij}} \mathrm{e}_{\mathrm{i}} \mathrm{e}_{\mathrm{j}} . \tag{13}
\end{equation*}
$$

Equation (12) shows that the strain energy function W is a homogeneous function of degree 2 in strains $e_{i}, i=1,2, \ldots ., 6$, and coefficients $c_{i j}$ are symmetric.

The generalized Hooke's law under the conditions of existence of strain energy function is given in equations (9) and (11).

In matrix form , it can be expressed as

$$
\left[\begin{array}{c}
\tau_{11}  \tag{14}\\
\tau_{22} \\
\tau_{33} \\
\tau_{23} \\
\tau_{13} \\
\tau_{12}
\end{array}\right]=\left[\begin{array}{llllll}
c_{11} & c_{12} & c_{13} & c_{14} & c_{15} & c_{16} \\
c_{12} & c_{22} & c_{23} & c_{24} & c_{25} & c_{26} \\
c_{13} & c_{23} & c_{33} & c_{34} & c_{35} & c_{36} \\
c_{14} & c_{24} & c_{34} & c_{44} & c_{45} & c_{46} \\
c_{15} & c_{25} & c_{35} & c_{45} & c_{55} & c_{56} \\
c_{16} & c_{26} & c_{36} & c_{46} & c_{56} & c_{66}
\end{array}\right]\left[\begin{array}{c}
e_{11} \\
e_{22} \\
e_{33} \\
2 e_{23} \\
2 e_{13} \\
2 e_{12}
\end{array}\right] .
$$

This law contains 21 independent elastic constants.
Result 1: From equation (2) ; we write

$$
\mathrm{W}=\frac{1}{2}\left[\tau_{1} \mathrm{e}_{1}+\tau_{2} \mathrm{e}_{2}+\tau_{3} \mathrm{e}_{3}+\tau_{4} \mathrm{e}_{4}+\tau_{5} \mathrm{e}_{5}+\tau_{6} \mathrm{e}_{6}\right]
$$

$$
\begin{align*}
& =\frac{1}{2}\left[\tau_{11} \mathrm{e}_{11}+\tau_{22} \mathrm{e}_{22}+\tau_{33} \mathrm{e}_{33}+2 \tau_{23} \mathrm{e}_{23}+2 \tau_{13} \mathrm{e}_{13}+2 \tau_{12} \mathrm{e}_{12}\right] \\
& =\frac{1}{2} \tau_{\mathrm{ij}} \mathrm{e}_{\mathrm{ij}} \quad, \mathrm{i}, \mathrm{j}=1,2,3 \tag{15}
\end{align*}
$$

## This result in (14) is called Claperon formula.

Result II : For an isotropic elastic medium, the Hooke's law gives

$$
\begin{equation*}
\tau_{\mathrm{ij}}=\lambda \delta_{\mathrm{ij}} \mathrm{e}_{\mathrm{kk}}+2 \mu \mathrm{e}_{\mathrm{ij}}, \quad \mathrm{i}, \mathrm{j}=1,2,3 \tag{16}
\end{equation*}
$$

This gives

$$
\begin{align*}
\mathrm{W} & \left.=\frac{1}{2} \mathrm{e}_{\mathrm{ij}} \lambda \lambda \delta_{\mathrm{ij}} \mathrm{e}_{\mathrm{kk}}+2 \mu \mathrm{e}_{\mathrm{ij}}\right] \\
& =\frac{1}{2} \lambda \mathrm{e}_{\mathrm{kk}} \mathrm{e}_{\mathrm{kk}}+\mu \mathrm{e}_{\mathrm{ij}} \mathrm{e}_{\mathrm{ij}} . \\
& =\frac{1}{2} \lambda \mathrm{e}_{\mathrm{kk}}^{2}+\mu \mathrm{e}_{\mathrm{ij}}^{2} \\
& =\frac{1}{2} \lambda\left(\mathrm{e}_{11}+\mathrm{e}_{22}+\mathrm{e}_{33}\right)^{2}+\mu\left(\mathrm{e}_{11}^{2}+\mathrm{e}_{22}^{2}+\mathrm{e}_{33}^{2}+2 \mathrm{e}_{12}^{2}+2 \mathrm{e}_{13}^{2}+2 \mathrm{e}_{23}^{2}\right) . \tag{17}
\end{align*}
$$

Result 3: Also, we have

$$
\begin{equation*}
\mathrm{e}_{\mathrm{ij}}=-\frac{\sigma}{E} \delta_{\mathrm{ij}} \tau_{\mathrm{kk}}+\frac{1+\sigma}{E} \tau_{\mathrm{ij}} . \tag{18}
\end{equation*}
$$

Hence

$$
\begin{align*}
\mathrm{W} & =\frac{1}{2} \tau_{\mathrm{ij}}\left[-\frac{\sigma}{E} \delta_{\mathrm{ij}} \tau_{\mathrm{kk}}+\frac{1+\sigma}{E} \tau_{\mathrm{ij}}\right] \\
& =-\frac{\sigma}{2 E} \tau_{\mathrm{ii}} \tau_{\mathrm{kk}}+\frac{1+\sigma}{2 E} \tau_{\mathrm{ij}} \tau_{\mathrm{ij}} \tag{19}
\end{align*}
$$

Result 4: From equation (12), we note that in the value of W, we may interchange $e_{i}$ and $\tau_{i}$. Consequently, interchanging $e_{i}$ and $\tau_{\mathrm{i}}$ in equation (4), we obtain

$$
\begin{equation*}
\frac{\partial W}{\partial \tau_{i}}=\mathrm{e}_{\mathrm{i}} \quad, \quad \text { for } \mathrm{i}=1,2,3, \ldots \ldots .6 \tag{20}
\end{equation*}
$$

This result is due to Castigliano (1847-1884).
It follows from the assumed linear assumed linear stress - strain relations.
Result 5: We know that the elastic moduli $\lambda$ and $\mu$ are both positive for all physical elastic solids. The quadratic form on the right side of (17) takes only positive values for every set of values of the strains.

This shows that the strain energy function W is a positive definite form in the strain components $\mathrm{e}_{\mathrm{ij}}$, for an isotropic elastic solid.

Question : Show that the strain - energy function W for an isotropic solid is independent of the choice of coordinate axes.

Solution : We know that the strain energy function W is given by

$$
\begin{align*}
& \mathrm{W}=\frac{1}{2} \tau_{\mathrm{ij}} \mathrm{e}_{\mathrm{ij}} \\
& =\frac{1}{2} \mathrm{e}_{\mathrm{ij}}\left(\lambda \delta_{\mathrm{ij}} \mathrm{e}_{\mathrm{kk}}+2 \mu \mathrm{e}_{\mathrm{ij}}\right) \\
& =\frac{1}{2} \lambda\left(\mathrm{e}_{11}+\mathrm{e}_{22}+\mathrm{e}_{33}\right)^{2}+\mu\left(\mathrm{e}_{11}^{2}+\mathrm{e}_{22}^{2}+\mathrm{e}_{33}^{2}+2 \mathrm{e}_{12}^{2}+2 \mathrm{e}_{13}^{2}+2 \mathrm{e}_{23}^{2}\right] \tag{1}
\end{align*}
$$

Let

$$
\begin{align*}
& I_{1}=e_{i i}=e_{11}+e_{22}+e_{33},  \tag{2}\\
& I_{2}=e_{i i} e_{j j}-e_{i j} e_{j i} . \tag{3}
\end{align*}
$$

be the first and second invariants of the strain tensor $\mathrm{e}_{\mathrm{ij}}$. As the given medium is isotropic, the elastic moduli $\lambda$ and $\mu$ are also independent of the choice of coordinate axes. We write

$$
\begin{aligned}
\mathrm{W}= & \frac{1}{2} \lambda \mathrm{I}_{1}^{2}+\mu\left[\left(\mathrm{e}_{11}+\mathrm{e}_{22}+\mathrm{e}_{33}\right)^{2}-2 \mathrm{e}_{11} \mathrm{e}_{22}-2 \mathrm{e}_{22} \mathrm{e}_{33}\right. \\
& \left.-2 \mathrm{e}_{33} \mathrm{e}_{11}+2 \mathrm{e}_{12}^{2}+2 \mathrm{e}_{13}^{2}+2 \mathrm{e}_{23}^{2}\right] \\
= & \frac{1}{2} \lambda \mathrm{I}^{2}+\mu\left[\mathrm{I}_{1}^{2}-2\left\{\left(\mathrm{e}_{11} \mathrm{e}_{22}-\mathrm{e}_{12}^{2}\right)+\left(\mathrm{e}_{22} \mathrm{e}_{33}-\mathrm{e}_{23}^{2}\right)+\left(\mathrm{e}_{11} \mathrm{e}_{33}-\mathrm{e}_{13}^{2}\right\}\right]\right. \\
= & \frac{1}{2} \lambda \mathrm{I}_{1}^{2}+\mu \mathrm{I}_{1}^{2}-2 \mu \mathrm{I}_{2}
\end{aligned}
$$

$$
\begin{equation*}
=\left(\frac{\lambda}{2}+\mu\right) \mathrm{I}_{1}^{2}-2 \mu \mathrm{I}_{2} . \tag{4}
\end{equation*}
$$

Hence, equation (4) shows that the strain energy function W is invariant relative to all rotations of cartesian axes.

Question : Evaluate W for the stress field (for an isotropic solid)

$$
\begin{aligned}
& \tau_{11}=\tau_{22}=\tau_{33}=\tau_{12}=0 \\
& \tau_{13}=-\mu \alpha x_{2}, \tau_{23}=\mu \alpha x_{1} \quad, \alpha \neq 0 \text { is a constant and } \mu \text { is the }
\end{aligned}
$$

Lame's constant
Solution : We find

$$
\tau_{\mathrm{kk}}=\tau_{11}+\tau_{22}+\tau_{33}=0
$$

Hence, the relation

$$
\mathrm{e}_{\mathrm{ij}}=\frac{1}{2 \mu}\left[\tau_{\mathrm{ij}}-\frac{\lambda}{3 \lambda+2 \mu} \delta_{\mathrm{ij}} \tau_{\mathrm{kk}}\right] \quad ; \quad \mathrm{i}, \mathrm{j}=1,2,3 .
$$

gives

$$
\mathrm{e}_{\mathrm{ij}}=\frac{1}{2 \mu} \tau_{\mathrm{ij}}
$$

That is, $\quad e_{11}=e_{22}=e_{33}=e_{12}=0$,

$$
\begin{equation*}
e_{13}=-\frac{1}{2} \alpha x_{2}, e_{12}=\frac{1}{2} \alpha x_{1} . \tag{1}
\end{equation*}
$$

The energy function W is given by

$$
\begin{aligned}
\mathrm{W} & =\frac{1}{2} \tau_{\mathrm{ij}} \mathrm{e}_{\mathrm{ij}} \\
& =\frac{1}{4 \mu} \tau_{\mathrm{ij}} \tau_{\mathrm{ij}} \\
& =\frac{1}{4 \mu}\left[\tau_{13}^{2}+\tau_{23}^{2}\right] \\
& =\frac{1}{4} \mu \alpha^{2}\left(\mathrm{x}_{1}{ }^{2}+\mathrm{x}_{2}^{2}\right) .
\end{aligned}
$$

Exercise : Show that the strain energy function W is given by

$$
\mathrm{W}=\mathrm{W}_{1}+\mathrm{W}_{2},
$$

where $\mathrm{W}_{1}=\frac{1}{2} \mathrm{k} \mathrm{e}_{\mathrm{ii}} \mathrm{e}_{\mathrm{ii}}=\frac{1}{18 k} \tau_{\mathrm{ii}} \tau_{\mathrm{ii}}, \quad \mathrm{k}=$ bulk modulus ,
and $\quad W_{2}=\frac{1}{3} \mu\left[\left(e_{11}-e_{22}\right)^{2}+\left(e_{22}-e_{33}\right)^{2}+\left(e_{33}-e_{11}\right)^{2}+6\left(e_{12}{ }^{2}+e_{23}{ }^{2}+e_{31}{ }^{2}\right)\right]$

$$
=\frac{1}{12 \mu}\left[\left(\tau_{11}-\tau_{22}\right)^{2}+\left(\tau_{22}-\tau_{33}\right)^{2}+\left(\tau_{33}-\tau_{11}\right)^{2}+6\left(\tau_{12}^{2}+\tau_{23}^{2}+\tau_{31}^{2}\right)\right] .
$$

Question : If $\mathrm{W}=\frac{1}{2}\left[\lambda \mathrm{e}_{\mathrm{kk}}^{2}+2 \mu \mathrm{e}_{\mathrm{ij}} \mathrm{e}_{\mathrm{ij}}\right]$, prove the following,
(i) $\frac{\partial W}{\partial e_{i j}}=\tau_{\mathrm{ij}}$
(ii) $\mathrm{W}=\frac{1}{2} \tau_{\mathrm{ij}} \mathrm{e}_{\mathrm{ij}}$
(iii) W is a scalar invariant.
(iv) $\mathrm{W} \geq 0$ and $\mathrm{W}=0$ iff $\mathrm{e}_{\mathrm{ij}}=0$
(v) $\frac{\partial W}{\partial \tau_{i j}}=\mathrm{e}_{\mathrm{ij}}$.

Solution : We note that W is a function of $\mathrm{e}_{\mathrm{ij}}$. Partial differentiation of this function w.r.t. $\mathrm{e}_{\mathrm{ij}}$ gives

$$
\begin{align*}
\frac{\partial W}{\partial e_{i j}} & =\frac{1}{2}\left[\lambda \cdot 2 \mathrm{e}_{\mathrm{kk}} \frac{\partial e_{k k}}{\partial e_{i j}}+4 \mu \mathrm{e}_{\mathrm{ij}}\right] \\
& =\left[\lambda \mathrm{e}_{\mathrm{kk}} \delta_{\mathrm{ij}}+2 \mu \mathrm{e}_{\mathrm{ij}}\right] \\
& =\tau_{\mathrm{ij}} \tag{1}
\end{align*}
$$

(ii) We write

$$
\begin{aligned}
\mathrm{W} & =\frac{1}{2}\left[\lambda \mathrm{e}_{\mathrm{kk}} \mathrm{e}_{\mathrm{kk}}+2 \mu \mathrm{e}_{\mathrm{ij}} \mathrm{e}_{\mathrm{ij}}\right] \\
& =\frac{1}{2}\left[\lambda \mathrm{e}_{\mathrm{kk}}\left\{\delta_{\mathrm{ij}} \mathrm{e}_{\mathrm{ij}}\right\}+2 \mu \mathrm{e}_{\mathrm{ij}} \mathrm{e}_{\mathrm{ij}}\right] \\
& =\frac{1}{2}\left(\lambda \delta_{\mathrm{ij}} \mathrm{e}_{\mathrm{kk}}+2 \mu \mathrm{e}_{\mathrm{ij}}\right) \mathrm{e}_{\mathrm{ij}}
\end{aligned}
$$

$$
\begin{equation*}
\mathrm{W}=\frac{1}{2} \tau_{\mathrm{ij}} \mathrm{e}_{\mathrm{ij}} \tag{2}
\end{equation*}
$$

(iii) Since $\tau_{\mathrm{ij}}$ and $\mathrm{e}_{\mathrm{ij}}$ are components of tensors, each of order 2, respectively.

So by contraction rule, $\mathrm{W}=\frac{1}{2} \tau_{\mathrm{ij}} \mathrm{e}_{\mathrm{ij}}$ is a scalar invariant.
(iv) Since $\lambda>0, \mu>0, \mathrm{e}_{\mathrm{kk}}^{2} \geq 0$ and $\mathrm{e}_{\mathrm{ij}} . \mathrm{e}_{\mathrm{ij}} \geq 0$, if follows that $\mathrm{W} \geq 0$.

Moreover $\mathrm{W}=0$ iff $\mathrm{e}_{\mathrm{kk}}=0$ and $\mathrm{e}_{\mathrm{ij}}=0$. Since $\mathrm{e}_{\mathrm{ij}}=0$ automatically implies that $\mathrm{e}_{\mathrm{kk}}=0$. Hence $\mathrm{W}=0$ holds iff $\mathrm{e}_{\mathrm{ij}}=0$.
(v) Putting

$$
\mathrm{e}_{\mathrm{ij}}=\frac{1+\sigma}{E} \tau_{\mathrm{ij}}-\frac{\sigma}{E} \tau_{\mathrm{kk}} \delta_{\mathrm{ij}}
$$

into (2), we find

$$
\mathrm{W}=\frac{1}{2}\left[\frac{1+\sigma}{E} \tau_{\mathrm{ij}} \tau_{\mathrm{ij}}-\frac{\sigma}{E} \tau_{\mathrm{kk}} \delta_{\mathrm{ij}} \tau_{\mathrm{ij}}\right]=\frac{1}{2}\left[\frac{1+\sigma}{E} \tau_{\mathrm{ij}} \tau_{\mathrm{ij}}-\frac{\sigma}{E} \tau_{\mathrm{kk}}^{2}\right]
$$

This implies

$$
\begin{equation*}
\frac{\partial W}{\partial \tau_{i j}}=\frac{1+\sigma}{E} \tau_{\mathrm{ij}}-\frac{\sigma}{E} \tau_{\mathrm{kk}} \frac{\partial \tau_{k k}}{\partial \tau_{i j}} \Rightarrow \frac{\partial W}{\partial \tau_{i j}}=\frac{1+\sigma}{E} \tau_{\mathrm{ij}}-\frac{\sigma}{E} \tau_{\mathrm{kk}} \delta_{\mathrm{ij}}=\mathrm{e}_{\mathrm{ij}} \tag{3}
\end{equation*}
$$

Theorem : Show that the total work done by the external forces in altering (changing) the configuration of the natural state to the state at time $\mathbf{t}$ is equal to the sum of the kinetic energy and the strain energy.

Proof : the natural / unstrained state of an elastic body is one in which there is a uniform temperature and zero displacement with reference to which all strains will be specified.
Let the body be in the natural state when $t=0$. Let ( $\mathrm{x}_{1}, \mathrm{x}_{2}, \mathrm{x}_{3}$ ) denote the coordinate of an arbitrary material point of the elastic body in the undeformed / unstrained state.


If the elastic body is subjected to the action of external forces, then it may produce a deformation of the body and at any time ' $t$ ', the coordinate of the same material point will be $\mathrm{x}_{\mathrm{i}}+\mathrm{u}_{\mathrm{i}}\left(\mathrm{x}_{1}, \mathrm{x}_{2}, \mathrm{x}_{3}, \mathrm{t}\right)$.


## The displacement of the point $P$ in the interval of time ( $\mathbf{t}, \mathbf{t}+\mathbf{d t}$ ) is given by

$$
\frac{\partial u_{i}}{\partial t} \mathrm{dt}=\dot{u}_{i} \mathrm{dt}
$$

where

$$
\dot{u}_{i}=\frac{\partial u_{i}}{\partial t} .
$$

The work done by the body forces $\mathrm{F}_{\mathrm{i}}$ acting on the volume element $\mathrm{d} \tau$, in time ' dt ' sec, located at the material point P is

$$
\left(\mathrm{F}_{\mathrm{i}} \mathrm{~d} \tau\right)\left(\dot{u}_{i} \mathrm{dt}\right)=\mathrm{F}_{\mathrm{i}} \dot{u}_{i} \mathrm{~d} \tau \mathrm{dt},
$$

and the work performed by the external surface forces $\stackrel{v}{T}_{i}$ in time interval $(\mathrm{t}, \mathrm{t}$ +dt ) is

$$
\stackrel{\nu}{T_{i}} \dot{u}_{i} \mathrm{~d} \sigma \mathrm{dt},
$$

where $\mathrm{d} \sigma$ is the element of surface.
Let E denote the work done by the body and surface forces acting on the elastic body.

Then, the rate of doing work on the body originally occupying some region $\tau$ (by external forces) is

$$
\begin{equation*}
\frac{d E}{d t}=\int_{\tau} \mathrm{F}_{\mathrm{i}} \dot{u}_{i} \mathrm{~d} \tau+\int_{\Sigma} \stackrel{v}{T_{i}} \dot{u}_{i} \mathrm{~d} \sigma \tag{1}
\end{equation*}
$$

where $\Sigma$ denotes the original surface of the elastic body.
Now $\int_{\Sigma} \stackrel{v}{T_{i}} \dot{u}_{i} \mathrm{~d} \sigma=\int_{\Sigma}\left(\tau_{\mathrm{ij}} v_{j}\right) \dot{u}_{i} \mathrm{~d} \sigma$

$$
\begin{align*}
& =\int_{\Sigma}\left(\tau_{\mathrm{ij}} \dot{u}_{i}\right) v_{j} \mathrm{~d} \sigma \\
& =\int_{\tau}\left(\tau_{\mathrm{ij}} \dot{u}_{i}\right), \mathrm{j} \mathrm{~d} \tau, \\
& =\int_{\tau}\left[\tau_{\mathrm{ij}, \mathrm{j}} \dot{u}_{i}+\tau_{\mathrm{ij}} \dot{u}_{i, j}\right] \mathrm{d} \tau \\
& =\int_{\tau} \tau_{\mathrm{ij}, \mathrm{j}} \dot{u}_{i} \mathrm{~d} \tau+\int_{\tau} \tau_{\mathrm{ij}} \dot{e}_{i j} \mathrm{~d} \tau+\int_{\tau} \tau_{\mathrm{ij}} \dot{w}_{i j} \mathrm{~d} \tau, \tag{2}
\end{align*}
$$

where

$$
\begin{align*}
& \dot{e}_{i j}=\left(\dot{u}_{i, \mathrm{j}}+\dot{u}_{j, i}\right) / 2 \\
& \dot{w}_{i j}=\left(\dot{u}_{i, \mathrm{j}}-\dot{u}_{j, i}\right) / 2 \tag{3}
\end{align*}
$$

Since

$$
\dot{w}_{i j}=-\dot{w}_{j i} \text { and } \tau_{\mathrm{ij}}=\tau_{\mathrm{ji}}
$$

so

$$
\begin{equation*}
\tau_{\mathrm{ij}} \dot{w}_{i j}=0 \tag{4a}
\end{equation*}
$$

From dynamical equations of motion for an isotropic body, we write

$$
\begin{equation*}
\tau_{\mathrm{ij}, \mathrm{j}}=\rho \ddot{u}_{i}-\mathrm{F}_{\mathrm{i}} . \tag{4b}
\end{equation*}
$$

Therefore, $\quad \tau_{\mathrm{ij}, \mathrm{j}} \dot{u}_{i}=\rho \ddot{u}_{i} \dot{u}_{i}-\mathrm{F}_{\mathrm{i}} \dot{u}_{i}$.
Using results (4a,b) ; we write from equations (3) and (1),

$$
\begin{align*}
\frac{d E}{d t} & =\int_{\tau} \mathrm{F}_{\mathrm{i}} \dot{u}_{i} \mathrm{~d} \tau+\int_{\tau}\left[\rho \ddot{u}_{i} \dot{u}_{i}-\mathrm{F}_{\mathrm{i}} \dot{u}_{i}\right] \mathrm{d} \tau+\int_{\tau} \tau_{\mathrm{ij}} \dot{e}_{i j} \mathrm{~d} \tau \\
& =\int_{\tau} \rho \ddot{u}_{i} \dot{u}_{i} \mathrm{~d} \tau+\int_{\tau} \tau_{\mathrm{ij}} \dot{e}_{i j} \mathrm{~d} \tau . \tag{5}
\end{align*}
$$

The kinetic energy K of the body in motion is given by

$$
\begin{equation*}
\mathrm{K}=\frac{1}{2} \int_{\tau} \rho \dot{u}_{i} \dot{u}_{i} \mathrm{~d} \tau \tag{6}
\end{equation*}
$$

Then

$$
\begin{equation*}
\frac{d K}{d t}=\int_{\tau} \rho \ddot{u}_{i} \dot{u}_{i} \mathrm{~d} \tau \tag{7}
\end{equation*}
$$

We define the engineering notation

$$
\left.\begin{array}{l}
\tau_{11}=\tau_{1}, \tau_{22}=\tau_{2}, \tau_{33}=\tau_{3}  \tag{8}\\
\tau_{23}=\tau_{4}, \tau_{13}=\tau_{5}, \tau_{12}=\tau_{6} \\
\mathrm{e}_{11}=\mathrm{e}_{1}, \mathrm{e}_{22}=\mathrm{e}_{2}, \mathrm{e}_{33}=\mathrm{e}_{3} \\
2 \mathrm{e}_{23}=\mathrm{e}_{4}, 2 \mathrm{e}_{13}=\mathrm{e}_{5}, 2 \mathrm{e}_{12}=\mathrm{e}_{6}
\end{array}\right\}
$$

Then

$$
\begin{equation*}
\int_{\tau} \tau_{\mathrm{ij}} \dot{e}_{i j} \mathrm{~d} \tau=\int_{\tau} \tau_{\mathrm{i}} \frac{\partial e_{i}}{\partial t} \mathrm{~d} \tau \tag{9}
\end{equation*}
$$

for $\mathrm{i}=1,2,3, \ldots, 6$, and under isothermal condition, there exists a energy function

$$
\mathrm{W}=\mathrm{W}\left(\mathrm{e}_{1}, \mathrm{e}_{2}, \ldots \ldots, \mathrm{e}_{6}\right)
$$

with the property that

$$
\begin{equation*}
\frac{\partial W}{\partial e_{i}}=\tau_{\mathrm{i}} \tag{10}
\end{equation*}
$$

$1 \leq i \leq 6$. From equations (9) and (10), we write

$$
\begin{align*}
\int_{\tau} \tau_{\mathrm{ij}} \dot{e}_{i j} \mathrm{~d} \tau & =\int_{\tau}\left(\frac{\partial W}{\partial e_{i}} \frac{\partial e_{i}}{\partial t}\right) \mathrm{d} \tau=\frac{d}{d t} \int_{\tau} \mathrm{W} \mathrm{~d} \tau \\
& =\frac{d U}{d t} \tag{11}
\end{align*}
$$

where

$$
\begin{equation*}
\mathrm{U}=\int_{\tau} \mathrm{W} \mathrm{~d} \tau \tag{12}
\end{equation*}
$$

From equations (5), (7) and (11), we write

$$
\begin{equation*}
\frac{d E}{d t}=\frac{d K}{d t}+\frac{d U}{d t} \tag{13}
\end{equation*}
$$

Integrating equation (13) w.r.t. ' $t$ ' between the limits $t=0$ and $t=t$, we obtain

$$
\begin{equation*}
\mathrm{E}=\mathrm{K}+\mathrm{U}, \tag{14}
\end{equation*}
$$

since both E and K are zero at $\mathrm{t}=0$.
The equation (14) proves the required result.
Note 1: If the elastic body is in equilibrium instead of in motion, then $K=0$ and
consequently $\mathrm{E}=\mathrm{U}$.
Note 2: $U$ is called the total strain energy of the deformation.

### 5.3. CLAPEYRON'S THEOREM

Statement. If an elastic body is in equilibrium under a given system of body forces $\mathrm{F}_{\mathrm{i}}$ and surface forces $\stackrel{v}{T}_{i}$, then the strain energy of deformation is equal to one - half the work that would be done by the external forces (of the equilibrium state) acting through the displacements $\mathbf{u}_{\mathbf{i}}$ form the unstressed state to the state of equilibrium.

Proof: We are required to prove that

$$
\begin{equation*}
\int_{\tau} \mathrm{F}_{\mathrm{i}} \mathrm{u}_{\mathrm{i}} \mathrm{~d} \tau+\int_{\Sigma} \stackrel{v}{T_{i}} \mathrm{u}_{\mathrm{i}} \mathrm{~d} \sigma=2 \int_{\tau} \mathrm{W} \mathrm{~d} \tau \tag{1}
\end{equation*}
$$

where $\Sigma$ denotes the original surface of the unstressed region $\tau$ of the body and W is the energy density function representing the strain every per unit volume. Now

$$
\begin{aligned}
\int_{\Sigma} \stackrel{v}{T_{i}} \mathrm{u}_{\mathrm{i}} \mathrm{~d} \sigma & =\int_{\Sigma} \tau_{\mathrm{ij}} \mathrm{u}_{\mathrm{i}} v_{j} \mathrm{~d} \sigma \\
& =\int_{\tau}\left(\tau_{\mathrm{ij}} \mathrm{u}_{\mathrm{i}}\right)_{, \mathrm{j}} \mathrm{~d} \tau \\
& =\int_{\tau}\left\{\tau_{\mathrm{ij}, \mathrm{j}} \mathrm{u}_{\mathrm{i}}+\tau_{\mathrm{ij}} \mathrm{u}_{\mathrm{i}, \mathrm{j}}\right\} \mathrm{d} \tau
\end{aligned}
$$

$$
\begin{align*}
& =\int_{\tau} \tau_{\mathrm{ij}, \mathrm{j}} \mathrm{u}_{\mathrm{i}} \mathrm{~d} \tau+\int_{\tau} \tau_{\mathrm{ij}}\left\{\frac{u_{i, j}+u_{j, i}}{2}+\frac{u_{i, j}-u_{j, i}}{2}\right\} \mathrm{d} \tau \\
& =\int_{\tau} \tau_{\mathrm{ij}, \mathrm{j}} \mathrm{u}_{\mathrm{i}} \mathrm{~d} \tau+\int_{\tau} \tau_{\mathrm{ij}}\left(\mathrm{e}_{\mathrm{ij}}+\mathrm{w}_{\mathrm{ij}}\right) \mathrm{d} \tau \\
& =\int_{\tau}\left(\tau_{\mathrm{ij}, \mathrm{j}} \mathrm{u}_{\mathrm{i}}+\tau_{\mathrm{ij}} \mathrm{e}_{\mathrm{ij}}\right) \mathrm{d} \tau \tag{2}
\end{align*}
$$

since

$$
\mathrm{w}_{\mathrm{ij}}=-\mathrm{w}_{\mathrm{ji}} \text { and } \tau_{\mathrm{ij}}=\tau_{\mathrm{ji}} .
$$

Again from (2)

$$
\begin{equation*}
\int_{\Sigma} \stackrel{v}{T_{i}} \mathrm{u}_{\mathrm{i}} \mathrm{~d} \sigma=\int_{\tau}\left(-\mathrm{F}_{\mathrm{i}} \mathrm{u}_{\mathrm{i}}+2 \mathrm{~W}\right) \mathrm{d} \tau, \tag{3}
\end{equation*}
$$

since

$$
\tau_{\mathrm{ij}, \mathrm{j}}+\mathrm{F}_{\mathrm{i}}=0,
$$

being the equilibrium equations and

$$
\mathrm{W}=\frac{1}{2} \tau_{\mathrm{ij}} \mathrm{e}_{\mathrm{ij}} .
$$

From (3), we can write

$$
\begin{equation*}
\int_{\tau} \mathrm{F}_{\mathrm{i}} \mathrm{u}_{\mathrm{i}} \mathrm{~d} \tau+\int_{\Sigma} \stackrel{v}{i}^{v} \mathrm{u}_{\mathrm{i}} \mathrm{~d} \sigma=2 \int_{\tau} \mathrm{W} \mathrm{~d} \tau, \tag{4}
\end{equation*}
$$

proving the theorem.

### 5.4 RECIPROCAL THEOREM OF BETTI AND RAYLEIGH

Statement : If an elastic body is subjected to two systems of body and surface forces producing two equilibrium states, show that the work done by the system of forces in acting through the displacements of the second system is equal to the work done by the second system of forces in acting through the displacements of the first system.

Proof: Let the first system of body and surface forces $\left\{\mathrm{F}_{\mathrm{i}}, \stackrel{v}{T_{i}}\right\}$ produces the displacement $\mathrm{u}_{\mathrm{i}}$ and the second system $\left\{\mathrm{F}_{\mathrm{i}}{ }^{\prime}, \stackrel{v}{T_{i}}{ }^{\prime}\right\}$ produces displacements $\mathrm{u}_{\mathrm{i}}{ }^{\prime}$. Let
$\mathrm{W}_{1}=$ work done by the first system of forces in acting through the displacement of the second system.

Then

$$
\begin{align*}
\mathrm{W}_{1} & =\int_{V} F_{i} \mathrm{u}_{\mathrm{i}}^{\prime} \mathrm{dv}+\int_{S} T_{i} \mathrm{u}_{\mathrm{i}}^{\prime} \mathrm{ds} \\
& =\int_{V} F_{i} \mathrm{u}_{\mathrm{i}}^{\prime} \mathrm{dv}+\int_{S} \tau_{\mathrm{ij}} v_{\mathrm{j}} \mathrm{u}_{\mathrm{i}}^{\prime} \mathrm{ds} \\
& =\int_{V} F_{i} \mathrm{u}_{\mathrm{i}}^{\prime} \mathrm{dv}+\int_{V}\left(\tau_{\mathrm{ij}} \mathrm{u}_{\mathrm{i}}^{\prime}\right), \mathrm{j} \mathrm{dv} \\
& =\int_{V} F_{i} \mathrm{u}_{\mathrm{i}}^{\prime} \mathrm{dv}+\int_{V} \tau_{i j, j} \mathrm{u}_{\mathrm{i}}^{\prime} \mathrm{dv}+\int_{V} \tau_{i j} \mathrm{u}_{\mathrm{i}, \mathrm{j}}^{\prime} \mathrm{dv} \\
& =\int_{V}\left(\tau_{i j, j}+F_{i}\right) u_{i}^{\prime} \mathrm{dv}+\int_{V} \tau_{i j} \mathrm{e}_{\mathrm{ij}} \mathrm{dv}, \\
& =\int_{V} \tau_{i j} \mathrm{e}_{\mathrm{ij}}^{\prime \mathrm{d}} \mathrm{dv}, \tag{1}
\end{align*}
$$

using equations of equilibrium

$$
\tau_{\mathrm{ij}, \mathrm{j}}+\mathrm{F}_{\mathrm{i}}=0
$$

Hence

$$
\begin{align*}
\mathrm{W}_{1} & =\int_{V}\left[\lambda \delta_{\mathrm{ij}} \mathrm{e}_{\mathrm{kk}}+2 \mu \mathrm{e}_{\mathrm{ij}}\right] \mathrm{e}^{\prime}{ }_{\mathrm{ij}} \mathrm{dv} \\
& =\int_{V}\left[\lambda \mathrm{e}_{\mathrm{kk}} \mathrm{e}^{\prime}{ }_{\mathrm{kk}}+2 \mu \mathrm{e}_{\mathrm{ij}} \mathrm{e}^{\prime}{ }_{\mathrm{ij}}\right] \mathrm{dv} . \tag{2}
\end{align*}
$$

This expression is symmetric in primed and unprimed quantities.
We conclude that $W_{1}=W_{2}$ where $W_{2}$ is the workdone by the forces of the second system in acting through the displacements $u_{i}$ of the first system.

This completes the proof of the theorem.
Corollary : Let ${ }^{(1)} \tau_{i j}$ be the stresses corresponding to the strains $\stackrel{(1)}{(1)}_{i j}$ and ${ }^{(2)} \tau_{i j}$ be the stresses corresponding to the strains $\stackrel{(2)}{e_{i j}}$, in an elastic body. Prove that

$$
\begin{aligned}
& \text { (1) (2) (2) (1) } \\
& \tau_{i j} e_{i j}=\tau_{i j} e_{i j} \text {. }
\end{aligned}
$$

Remark 1: Reciprocal theorem relates the equilibrium states of an elastic solid under the action of different applied loads.

Remark 2: An alternative form of the reciprocal theorem is

$$
\int_{\Sigma} \mathrm{T}_{\mathrm{i}} \mathrm{u}_{\mathrm{i}}^{\prime} \mathrm{d} \sigma+\int_{\tau} \mathrm{F}_{\mathrm{i}} \mathrm{u}_{\mathrm{i}}^{\prime} \mathrm{d} \tau=\int_{\tau} \tau_{\mathrm{ij}} \mathrm{e}^{\prime}{ }_{\mathrm{ij}} \mathrm{~d} \tau .
$$

### 5.5. THEOREM OF MINIMUM POTENTIAL ENERGY

Now, we introduce an important functional, called the potential energy of deformation , and prove that this functional attains an absolute minimum when the displacements of the elastic body are those of the equilibrium configuration.

Statement : Of all displacements satisfying the given boundary conditions, those which satisfy the equilibrium equations make the potential energy an absolute minimum.

Proof : Let a body $\tau$ be in equilibrium under the action of specified body and surface forces. Suppose that the surface forces $T_{i}$ are prescribed only over a portion $\Sigma_{T}$ of the surface $\Sigma$, and over the remaining surface $\Sigma_{u}$ the displacements are known.

We denote the displacements of the equilibrium state by $u_{i}$. We consider a class of arbitrary displacements $u_{i}+\delta u_{i}$, consistent with constraints imposed on the elastic body. This means that

$$
\begin{equation*}
\delta u_{i}=0, \quad \text { on } \Sigma_{u} \tag{1}
\end{equation*}
$$

but $\delta u_{i}$ are arbitrary over the part $\Sigma_{\mathrm{T}}$, except for the condition that they belong to class $\mathrm{C}^{3}$ and are of the order of magnitude of displacements admissible in linear elasticity.

Displacements $\delta \mathrm{u}_{\mathrm{i}}$ are called virtual displacements .
We know that the strain energy U is given by the formula

$$
\begin{equation*}
\mathrm{U}=\int_{\tau} \mathrm{W} \mathrm{~d} \tau, \tag{2}
\end{equation*}
$$

where the strain energy function W is given by the formula

$$
\begin{equation*}
\mathrm{W}=\frac{1}{2} \lambda \mathrm{e}_{\mathrm{kk}} \mathrm{e}_{\mathrm{kk}}+\mu \mathrm{e}_{\mathrm{ij}} \mathrm{e}_{\mathrm{ij}} \tag{3}
\end{equation*}
$$

$\lambda$ and $\mu$ being Lame constants, and $\mathrm{e}_{\mathrm{ij}}$ strain tensor.
The strain energy $U$ is equal to the work done by the external forces on the elastic body in the process of bringing the body from the natural state to the equilibrium state characterized by the displacements $u_{i}$.

The virtual work $\delta \mathbf{U}$ performed by the external force $F_{i}$ and $T_{i}$ during the virtual displacements $\delta u_{i}$ is defined by the equation

$$
\begin{equation*}
\delta \mathrm{U}=\int_{\tau} \mathrm{F}_{\mathrm{i}} \delta \mathrm{u}_{\mathrm{i}} \mathrm{~d} \tau+\int_{\Sigma} \mathrm{T}_{\mathrm{i}} \delta \mathrm{u}_{\mathrm{i}} \mathrm{~d} \sigma . \tag{4}
\end{equation*}
$$

Since the volume $\tau$ is fixed and the forces $\mathrm{F}_{\mathrm{i}}$ and $\mathrm{T}_{\mathrm{i}}$ do not vary when the arbitrary variations $\delta u_{i}$ are considered, equation (4) can be written in the form

$$
\begin{equation*}
\delta \mathrm{U}=\delta\left(\int_{\tau} F_{i} u_{i} d \tau+\int_{\Sigma} T_{i} u_{i} d \sigma\right) \tag{5}
\end{equation*}
$$

From equation (2), we have

$$
\begin{equation*}
\delta \mathrm{U}=\delta\left(\int_{\tau} W d \tau\right) \tag{6}
\end{equation*}
$$

Equations (5) and (6) provide

$$
\begin{equation*}
\delta\left(\int_{\tau} W d \tau-\int_{\tau} F_{i} u_{i} d \tau-\int_{\Sigma} T_{i} u_{i} d \sigma\right)=0 \tag{7}
\end{equation*}
$$

## The potential energy $V$ is defined by the formula

$$
\begin{equation*}
\mathrm{V}=\int_{\tau} W d \tau-\int_{\tau} F_{i} u_{i} d \tau-\int_{\Sigma} T_{i} u_{i} d \sigma \tag{8}
\end{equation*}
$$

In view of equation (8), relation (7) reads

$$
\begin{equation*}
\delta \mathrm{V}=0 \tag{9}
\end{equation*}
$$

This formula shows that the potential energy functional V has a stationary value in a class of admissible variations $\delta u_{i}$ of the displacements $u_{i}$ of the equilibrium state.

We shall finally show that the functional V assumes a minimum value when the displacements $u_{i}$ are those of the equilibrium state.

To show this, we demonstrate that the increment $\Delta \mathrm{V}$ produced in V by replacing the equilibrium displacements $u_{i}$ by $u_{i}+\delta u_{i}$ is positive for all non vanishing variations $\delta \mathrm{u}_{\mathrm{i}}$.

First, we calculate the increment $\Delta \mathbf{W}$ in $\mathbf{W}$. From (3)

$$
\begin{equation*}
\Delta \mathrm{W}=\left.\left(\frac{\lambda}{2} v^{2}+\mu e_{i j} e_{i j}\right)\right|_{u_{i}+\delta u_{i}}-\left.\left(\frac{\lambda}{2} v^{2}+\mu e_{i j} e_{i j}\right)\right|_{u_{i}} \tag{10}
\end{equation*}
$$

Now

$$
\begin{align*}
\left.\mathrm{e}_{\mathrm{ij}}\right|_{u_{i}+\delta u_{i}} & =\left.\frac{1}{2}\left(u_{i, j}+u_{j, i}\right)\right|_{u_{i}+\delta u_{i}} \\
& =\frac{1}{2}\left(u_{i, j}+u_{j, i}\right)+\frac{1}{2}\left[\left(\delta u_{i}\right)_{, j}+\left(\delta u_{j}\right)_{i,}\right] \\
& =e_{i j}+\frac{1}{2}\left(\delta u_{i}\right)_{j}+\frac{1}{2}\left(\delta u_{j}\right)_{, i} \tag{11}
\end{align*}
$$

and

$$
\begin{align*}
\left.v\right|_{u_{i}+\delta u_{i}} & =\left.e_{k k}\right|_{u_{i}+\delta t_{i}} \\
& =\mathrm{e}_{\mathrm{ii}}+\left(\delta \mathrm{u}_{\mathrm{i}}\right)_{, \mathrm{i}} \\
& =v+\left(\delta \mathrm{u}_{\mathrm{i}}\right)_{, \mathrm{i}} \tag{12}
\end{align*}
$$

Therefore, equations (10) to (12) yield

$$
\begin{gather*}
\Delta \mathrm{W}=\left(\frac{\lambda}{2}\right)\left[v+\left(\delta \mathrm{u}_{\mathrm{i}}\right)_{\mathrm{i}}\right]\left[v+\left(\delta \mathrm{u}_{\mathrm{i}}\right)_{, \mathrm{i}}\right]+\mu\left[\mathrm{e}_{\mathrm{ij}}+\frac{1}{2}\left(\delta \mathrm{u}_{\mathrm{i}}\right)_{\mathrm{i}}+\frac{1}{2}\left(\delta \mathrm{u}_{\mathrm{j}}\right)_{\mathrm{i}}\right] \\
\times\left[\mathrm{e}_{\mathrm{ij}}+\frac{1}{2}\left(\delta \mathrm{u}_{\mathrm{i}}\right)_{\mathrm{j}}+\frac{1}{2}\left(\delta \mathrm{u}_{\mathrm{i}}\right)_{\mathrm{i}}\right]-\left(\frac{\lambda}{2}\right) v^{2}-\mu \mathrm{e}_{\mathrm{ij}} \mathrm{e}_{\mathrm{ij}} \\
=\lambda v\left(\delta \mathrm{u}_{\mathrm{i}}\right)_{, \mathrm{i}}+2 \mu \mathrm{e}_{\mathrm{ij}}\left(\delta \mathrm{u}_{\mathrm{i}}\right)_{, \mathrm{j}}+\mathrm{P}, \tag{13}
\end{gather*}
$$

where

$$
\begin{equation*}
\mathrm{P}=\left(\frac{\lambda}{2}\right)\left[\left(\delta \mathrm{u}_{\mathrm{i}}\right), \mathrm{i}\right]^{2}+\left(\frac{\mu}{4}\right)\left[\left(\delta \mathrm{u}_{\mathrm{i}}\right), \mathrm{j}+\left(\delta \mathrm{u}_{\mathrm{j}}\right)_{\mathrm{i}}\right]^{2} \geq 0 \tag{14}
\end{equation*}
$$

Equation (13) can be rewritten in the form

$$
\begin{align*}
\Delta \mathrm{W} & =\lambda v \delta_{\mathrm{ij}}\left(\delta \mathrm{u}_{\mathrm{i}}\right)_{\mathrm{j}}+2 \mu \mathrm{e}_{\mathrm{ij}}\left(\delta \mathrm{u}_{\mathrm{i}}\right)_{, \mathrm{j}}+\mathrm{P} \\
& =\left(\lambda v \delta_{\mathrm{ij}}+2 \mu \mathrm{e}_{\mathrm{ij}}\right)\left[\left(\delta \mathrm{u}_{\mathrm{i}}\right)_{, \mathrm{j}}\right]+\mathrm{P} \\
& =\tau_{\mathrm{ij}}\left[\left(\delta \mathrm{u}_{\mathrm{i}}\right)_{, \mathrm{j}}\right]+\mathrm{P} . \tag{15}
\end{align*}
$$

The increment $\Delta \mathrm{U}$ in strain energy is, therefore ,

$$
\begin{align*}
\Delta \mathrm{U} & =\int_{\tau} \Delta \mathrm{W} \mathrm{~d} \tau \\
& =\int_{\tau} \tau_{\mathrm{ij}}\left(\delta \mathrm{u}_{\mathrm{i}}\right), \mathrm{j} \mathrm{~d} \tau+\int_{\tau} \mathrm{Pd} \tau \\
& =\int_{\tau}\left[\left(\tau_{\mathrm{ij}} \delta \mathrm{u}_{\mathrm{i}}\right), \mathrm{j}-\tau_{\mathrm{ij}, \mathrm{j}} \delta \mathrm{u}_{\mathrm{i}}\right] \mathrm{d} \tau+\mathrm{Q} \\
& =\int_{\tau}\left(\tau_{\mathrm{ij}} \delta \mathrm{u}_{\mathrm{i}}\right), \mathrm{j} \mathrm{~d} \tau-\int_{\tau} \tau_{\mathrm{ij}, \mathrm{j}} \delta \mathrm{u}_{\mathrm{i}} \mathrm{~d} \tau+\mathrm{Q} \\
& =\int_{\Sigma} \tau_{\mathrm{ij}} \mathrm{v}_{\mathrm{j}} \delta \mathrm{u}_{\mathrm{i}} \mathrm{~d} \sigma-\int_{\tau} \tau_{\mathrm{ij}, \mathrm{j}} \delta \mathrm{u}_{\mathrm{i}} \mathrm{~d} \tau+\mathrm{Q} . \tag{16}
\end{align*}
$$

In equation (16), we have used divergence theorem and

$$
\begin{equation*}
\mathrm{Q}=\int_{\tau} \mathrm{Pd} \tau \geq 0 \tag{17}
\end{equation*}
$$

Since $P \geq 0$ by virtue of (14).
If the body is in equilibrium, then we have

$$
\begin{array}{ll}
\tau_{\mathrm{ij}, \mathrm{j}}=-\mathrm{F}_{\mathrm{i}}, & \text { in } \tau \\
\tau_{\mathrm{ij}} v_{\mathrm{j}}=\stackrel{v}{T_{i}} \quad, \quad \text { on } \Sigma \tag{19}
\end{array}
$$

and , therefore , equation (16) becomes

$$
\begin{equation*}
\Delta \mathrm{U}=\int_{\Sigma}{ }_{\mathrm{E}} T_{i} \delta \mathrm{u}_{\mathrm{i}} \mathrm{~d} \sigma+\int_{\tau} \mathrm{F}_{\mathrm{i}} \delta \mathrm{u}_{\mathrm{i}} \mathrm{~d} \tau+\mathrm{Q} . \tag{20}
\end{equation*}
$$

Using the definition (8) for potential energy. we get

$$
\begin{equation*}
\Delta \mathrm{V}=\Delta \mathrm{U}-\int_{\tau} \mathrm{F}_{\mathrm{i}} \delta \mathrm{u}_{\mathrm{i}} \mathrm{~d} \tau-\int_{\Sigma} \stackrel{v}{T_{i}} \delta \mathrm{u}_{\mathrm{i}} \mathrm{~d} \sigma . \tag{21}
\end{equation*}
$$

Substituting (20) in (21), we obtain

$$
\begin{gather*}
\Delta \mathrm{V}=\left[\int_{\Sigma} T_{i}^{v} \delta u_{i} d \sigma+\int_{\tau} F_{i} \delta u_{i} d \tau+Q\right]-\int_{\tau} \mathrm{F}_{\mathrm{i}} \delta \mathrm{u}_{\mathrm{i}} \mathrm{~d} \tau-\int_{\Sigma} \stackrel{v}{T_{i}} \delta \mathrm{u}_{\mathrm{i}} \mathrm{~d} \sigma \\
=\mathrm{Q} \tag{22}
\end{gather*}
$$

Since

$$
\mathrm{Q} \geq 0,
$$

( $\mathrm{Q}=0$ in the case of equilibrium only as $\mathrm{P}=0$ in this case), we find

$$
\begin{equation*}
\Delta \mathrm{V} \geq 0 \tag{23}
\end{equation*}
$$

This completes the proof of the theorem.
Converse : Assume that there is a set of admissible functions $u_{i}+\delta u_{i}$ which satisfy the prescribed boundary conditions and such that

$$
\begin{equation*}
\Delta \mathrm{V}=\left[\Delta \mathrm{U}-\int_{\Sigma} \stackrel{v}{T_{i}} \delta \mathrm{u}_{\mathrm{i}} \mathrm{~d} \sigma-\int_{\tau} \mathrm{F}_{\mathrm{i}} \delta \mathrm{u}_{\mathrm{i}} \mathrm{~d} \tau\right] \geq 0 \tag{24}
\end{equation*}
$$

on this set of functions.

## From equation (16), we write

$$
\Delta \mathrm{U}=\int_{\Sigma} \tau_{\mathrm{ij}} v_{\mathrm{j}} \delta \mathrm{u}_{\mathrm{i}} \mathrm{~d} \sigma-\int_{\tau} \tau_{\mathrm{ij}, \mathrm{j}} \delta \mathrm{u}_{\mathrm{i}} \mathrm{~d} \tau+\mathrm{Q},
$$

where Q is given in (17). Inserting this value of $\Delta \mathrm{U}$ in (24), we obtain

$$
\begin{equation*}
\left.\left[-\int_{\tau}\left(\tau_{\mathrm{ij}, \mathrm{j}}+\mathrm{F}_{\mathrm{i}}\right) \delta \mathrm{u}_{\mathrm{i}} \mathrm{~d} \tau+\int_{\Sigma} \tau_{\mathrm{ij}} v_{\mathrm{j}}-\stackrel{T}{i}^{v}\right) \delta \mathrm{u}_{\mathrm{i}} \mathrm{~d} \sigma+\mathrm{Q}\right] \geq 0 . \tag{25}
\end{equation*}
$$

On the part $\Sigma_{\mathrm{T}}$ of $\Sigma$, where $\mathrm{T}_{\mathrm{i}}$ are assigned,

$$
\begin{equation*}
\tau_{\mathrm{ij}} v_{\mathrm{j}}-\mathrm{T}_{\mathrm{i}}=0, \tag{26}
\end{equation*}
$$

and over the remaining part $\Sigma_{u}$ of $\Sigma$,

$$
\begin{equation*}
\delta u_{i}=0 . \tag{27}
\end{equation*}
$$

Therefore,

$$
\begin{equation*}
\int_{\Sigma}\left(\tau_{\mathrm{ij}} v_{\mathrm{j}}-\stackrel{v}{T}_{i}\right) \delta \mathrm{u}_{\mathrm{i}} \mathrm{~d} \sigma=0 \tag{28}
\end{equation*}
$$

Hence, equation (25) reduces to

$$
\begin{equation*}
\left[-\int_{\tau}\left(\tau_{\mathrm{ij}, \mathrm{j}}+\mathrm{F}_{\mathrm{i}}\right) \delta \mathrm{u}_{\mathrm{i}} \mathrm{~d} \tau+\mathrm{Q}\right] \geq 0 \tag{29}
\end{equation*}
$$

Since Q is essentially positive and the displacements $\delta \mathbf{u}_{\mathbf{i}}$ are arbitrary, the inequality (29) implies that

$$
\begin{equation*}
\tau_{i \mathrm{ij}, \mathrm{j}}+\mathrm{F}_{\mathrm{i}}=0 \tag{30}
\end{equation*}
$$

for every point interior to $\tau$.
Thus, the equations of equilibrium are satisfied for every interior point in $\tau$. This proves the converse part.

### 5.6. THEOREM OF MINIMUM COMPLEMENTARY ENERGY

Definition : The complementary energy $\mathrm{V}^{*}$ is defined by the formula

$$
\mathrm{V}^{*}=\mathrm{U}-\int_{\Sigma u} \mathrm{~T}_{\mathrm{i}} \mathrm{u}_{\mathrm{i}} \mathrm{~d} \sigma=\int_{\tau} \mathrm{W} \mathrm{~d} \tau-\int_{\Sigma u} \mathrm{~T}_{\mathrm{i}} \mathrm{u}_{\mathrm{i}} \mathrm{~d} \sigma,
$$

where U is the strain energy and W is the strain energy function.
Statement : The complementary energy $\mathrm{V}^{*}$ has an absolute minimum when the stress tensor $\tau_{\mathrm{ij}}$ is that of the equilibrium state and the varied states of stress fulfill the following conditions :
(i) $\left(\delta \tau_{\mathrm{ij}}\right),{ }_{\mathrm{j}}=0$ in $\tau$,
(ii) $\quad\left(\delta \tau_{\mathrm{ij}}\right) v_{\mathrm{j}}=0$ on $\Sigma_{\mathrm{T}}$,
(iii) $\delta \tau_{\mathrm{ij}}$ are arbitrary on $\Sigma_{\mathrm{u}}$.

Proof : Let a body $\tau$ be in equilibrium under the action of body forces $F_{i}$ and surface forces $\mathrm{T}_{\mathrm{i}}$ assigned over a part $\Sigma_{\mathrm{T}}$ of the surface $\Sigma$. On the remaining part $\Sigma_{\mathrm{u}}$ of $\Sigma$, the displacements $\mathrm{u}_{\mathrm{i}}$ are assumed to be known.

If the $\tau_{\mathrm{ij}}$ are the stress components of the equilibrium state, then we have

$$
\left.\begin{array}{cc}
\tau_{\mathrm{ij}, \mathrm{j}}+\mathrm{F}_{\mathrm{i}}=0 & \text { in } \tau,  \tag{1}\\
\tau_{\mathrm{ij}} v_{\mathrm{j}}=\mathrm{T}_{\mathrm{i}} & \text { on } \Sigma_{\mathrm{T}}, \\
\mathrm{u}_{\mathrm{i}}=\mathrm{f}_{\mathrm{i}} & \text { on } \Sigma_{\mathrm{u}}
\end{array}\right\}
$$

We introduce a set of functions $\tau_{i j}^{1}$ of class $\mathrm{C}^{2}$ the body $\tau$, which we shall also write as

$$
\begin{equation*}
{\stackrel{1}{\tau_{i j}}=\tau_{\mathrm{ij}}+\delta \tau_{\mathrm{ij}}, ., ~}_{\text {, }} \tag{4}
\end{equation*}
$$

satisfying the conditions :
(i) $\tau_{i j, j}^{1}+\mathrm{F}_{\mathrm{i}}=0$, in $\tau$
(ii) $\quad \tau_{i j} v_{\mathrm{j}}=\mathrm{T}_{\mathrm{i}}, \quad$ on $\Sigma_{\mathrm{T}}$
(iii) $\quad \stackrel{1}{\tau}_{i j}$ are arbitrary on the surface $\Sigma_{\mathrm{u}}$.

## From equations (4) and (5) at each point of $\tau$; we write

$$
\begin{align*}
& \left(\tau_{\mathrm{ij}}+\delta \tau_{\mathrm{ij}}\right)_{, \mathrm{j}}+\mathrm{F}_{\mathrm{i}}=0 \\
& \left(\tau_{\mathrm{ij}, \mathrm{j}}+\mathrm{F}_{\mathrm{i}}\right)\left(\delta \tau_{\mathrm{ij}}\right)_{, \mathrm{j}}=0 \\
& \left(\delta \tau_{\mathrm{ij}}\right)_{\mathrm{j}}=0 \quad \text { in } \tau . \tag{7}
\end{align*}
$$

Also from equations (4) and (6), we have

$$
\begin{array}{ll}
\left(\tau_{\mathrm{ij}}+\delta \tau_{\mathrm{ij}}\right) v_{\mathrm{j}}=\mathrm{T}_{\mathrm{i}} & \text { on } \Sigma_{\mathrm{T}} \\
\tau_{\mathrm{ij}} v_{\mathrm{j}}+\left(\delta \tau_{\mathrm{ij}}\right) v_{\mathrm{j}}=\mathrm{T}_{\mathrm{i}} & \text { on } \Sigma_{\mathrm{T}}, \\
\left(\delta \tau_{\mathrm{ij}}\right) v_{\mathrm{j}}=0 & \text { on } \Sigma_{\mathrm{T}} . \tag{8}
\end{array}
$$

Since $\stackrel{1}{\tau}_{i j}$ are arbitrary on $\Sigma_{\mathrm{u}}$, so the variations $\delta \tau_{\mathrm{ij}}$ are arbitrary on $\Sigma_{\mathrm{u}}$.

As the stresses $\tau_{\mathrm{ij}}$ are associated with the equilibrium state of the body, so $\tau_{\mathrm{ij}}$ satisfy the Biltrami - Michell compatibility equations. Let W denote the strain - energy density function. It is given by the formula

$$
\begin{equation*}
\mathrm{W}=\left(\frac{1+\sigma}{2 E}\right) \tau_{\mathrm{ij}} \tau_{\mathrm{ij}}-\left(\frac{\sigma}{2 E}\right) \tau_{\mathrm{ii}} \tau_{\mathrm{ii}} \tag{9}
\end{equation*}
$$

where $\sigma=$ Poisson ration and $\mathrm{E}=$ Young's modulus.
The increment $\Delta \mathrm{U}$ in the strain energy U is given by the formula

$$
\begin{equation*}
\Delta \mathrm{U}=\int_{\tau} \Delta \mathrm{W} \mathrm{~d} \tau \tag{10}
\end{equation*}
$$

where the increment $\Delta \mathrm{W}$ in W is produced by replacing $\tau_{\mathrm{ij}}$ in (9) by $\stackrel{1}{\tau}_{i j}=\tau_{\mathrm{ij}}+$ $\delta \tau_{\mathrm{ij}}$. That is ,

$$
\begin{aligned}
\mathrm{W}+\Delta \mathrm{W} & =\left(\frac{1+\sigma}{2 E}\right)\left(\tau_{\mathrm{ij}}+\delta \tau_{\mathrm{ij}}\right)\left(\tau_{\mathrm{ij}}+\delta \tau_{\mathrm{ij}}\right)-\left(\frac{\sigma}{2 E}\right)\left(\tau_{\mathrm{ii}}+\delta \tau_{\mathrm{ii}}\right)^{2} \\
& =\mathrm{W}+\left(\frac{1+\sigma}{2 E}\right)\left[2 \tau_{\mathrm{ij}}\left(\delta \tau_{\mathrm{ij}}\right)+\left(\delta \tau_{\mathrm{ij}}\right)\right]-\left(\frac{\sigma}{2 E}\right)\left[2 \tau_{\mathrm{ii}}\left(\delta \tau_{\mathrm{ii}}\right)+\left(\delta \tau_{\mathrm{ii}}\right)^{2}\right]
\end{aligned}
$$

Hence,
$\Delta \mathrm{W}=\left(\frac{1+\sigma}{E}\right) \tau_{\mathrm{ij}}\left(\delta \tau_{\mathrm{ij}}\right)-\frac{\sigma}{E} \tau_{\mathrm{ii}}\left(\delta \tau_{\mathrm{iij}}\right)+\mathrm{W}\left(\delta \tau_{\mathrm{ij}}\right)$,
where
$\mathrm{W}\left(\delta \tau_{\mathrm{ij}}\right)=\left[\left(\frac{1+\sigma}{2 E}\right)\left(\delta \tau_{i j}\right)\left(\delta \tau_{i j}\right)-\left(\frac{\sigma}{2 E}\right)\left(\delta \tau_{i i}\right)^{2}\right] \geq 0$,
since the strain energy function W is a positive definite quadric form in its variables.

From Hooke's law for isotropic solids, we have

$$
\begin{equation*}
\mathrm{e}_{\mathrm{ij}}=\left(\frac{1+\sigma}{E}\right) \tau_{\mathrm{ij}}-\left(\frac{\sigma}{E}\right) \tau_{\mathrm{kk}} \delta_{\mathrm{ij}} \tag{13}
\end{equation*}
$$

From equations (11) to (12), we write

$$
\begin{aligned}
\Delta \mathrm{W} & =\left[\left(\frac{1+\sigma}{2 E}\right) \tau_{i j}-\left(\frac{\sigma}{E}\right) \tau_{k k} \delta_{i j}\right]\left(\delta \tau_{\mathrm{ij}}\right)+\mathrm{W}\left(\delta \tau_{\mathrm{ij}}\right) \\
& =\left(\mathrm{e}_{\mathrm{ij}}\right)\left(\delta \tau_{\mathrm{ij}}\right)+\mathrm{W}\left(\delta \tau_{\mathrm{ij}}\right) \\
& =\left(\frac{u_{i, j}+u_{j, i}}{2}\right)\left(\delta \tau_{\mathrm{ij}}\right)+\mathrm{W}\left(\delta \tau_{\mathrm{ij}}\right)
\end{aligned}
$$

$$
\begin{align*}
& =\left(\mathrm{u}_{\mathrm{i}, \mathrm{j}}\right)\left(\delta \tau_{\mathrm{ij}}\right)+\mathrm{W}\left(\delta \tau_{\mathrm{ij}}\right) \\
& =\left[\left(\mathrm{u}_{\mathrm{i}} \delta \tau_{\mathrm{ij}}\right)_{\mathrm{j}}-\mathrm{u}_{\mathrm{i}}\left(\delta \tau_{\mathrm{ij}}\right)_{\mathrm{j}}\right]+\mathrm{W}\left(\delta \tau_{\mathrm{ij}}\right) . \tag{14}
\end{align*}
$$

Since the stress components $\tau_{\mathrm{ij}}$ were assumed to satisfy the Beltrami - Michell compatibility equations, therefore, the displacements $u_{i}$ appearing in (14) are those of the actual equilibrium state of the body.

Using (14) in (10), the increment $\Delta \mathrm{U}$ in the strain energy becomes

$$
\begin{align*}
\Delta \mathrm{U} & =\int_{\tau}\left[\left(\mathrm{u}_{\mathrm{i}} \delta \tau_{\mathrm{ij}}\right)_{, \mathrm{j}}-\mathrm{u}_{\mathrm{i}}\left(\delta \tau_{\mathrm{ij}}\right)_{, \mathrm{j}}+\mathrm{W}\left(\delta \tau_{\mathrm{ij}}\right)\right] \mathrm{d} \tau \\
& =\int_{\tau}\left(\mathrm{u}_{\mathrm{i}} \delta \tau_{\mathrm{ij}}\right)_{\mathrm{j}} \mathrm{~d} \tau-\int_{\tau} \mathrm{u}_{\mathrm{i}}\left(\delta \tau_{\mathrm{ij}}\right)_{, \mathrm{j}} \mathrm{~d} \tau+\int_{\tau} \mathrm{W}\left(\delta \tau_{\mathrm{ij}}\right) \mathrm{d} \tau \\
& =\int_{\Sigma_{u}}\left(u_{i} \delta \tau_{\mathrm{ij}}\right) v_{\mathrm{j}} \mathrm{~d} \sigma+\mathrm{P}, \tag{15}
\end{align*}
$$

using the Gauss divergence theorem and equations (7) and (8). In equation (15)

$$
\begin{equation*}
\mathrm{P}=\int_{\tau} \mathrm{W}\left(\delta \tau_{\mathrm{ij}}\right) \mathrm{d} \tau \geq 0 \tag{16}
\end{equation*}
$$

As the variations $\delta \tau_{\mathrm{ij}}$ are arbitrary on the surface $\Sigma_{\mathrm{u}}$, we write

$$
\begin{equation*}
\left(\delta \tau_{\mathrm{ij}}\right) v_{\mathrm{j}}=\Delta \mathrm{T}_{\mathrm{i}}, \quad \text { on } \Sigma_{\mathrm{u}} \tag{17}
\end{equation*}
$$

then, equation (15) reads as

$$
\begin{equation*}
\Delta \mathrm{U}=\int_{\Sigma_{u}} u_{i} \Delta \mathrm{~T}_{\mathrm{i}} \mathrm{~d} \sigma+\mathrm{P} \tag{18}
\end{equation*}
$$

Since the displacements $u_{i}$ are assigned on the surface $\Sigma_{u}$, we can write (18) as

$$
\Delta\left(\mathrm{U}-\int_{\Sigma_{u}} u_{i} \mathrm{~T}_{\mathrm{i}}\right)=\mathrm{P} \geq 0
$$

or

$$
\begin{equation*}
\Delta \mathrm{V}^{*} \geq 0 \tag{19}
\end{equation*}
$$

That is , the increment $\Delta \mathrm{V}^{*}$ in the complementary energy $\mathrm{V}^{*}$ (for the equilibrium state) is essentially positive.

Hence, the complementary energy functional $\mathrm{V}^{*}$ has an absolute minimum in the case of an equilibrium state of the body.

This completes the proof.

### 5.7. THEOREM OF MINIMUM STRAIN

## ENERGY

Statement : The strain energy $U$ of an elastic body in equilibrium under the action of prescribed surface forces is an absolute minimum on the set of all values of the functional $U$ determined by the solutions of the system

$$
\tau_{\mathrm{i}, \mathrm{j},}+\mathrm{F}_{\mathrm{i}}=0 \quad \text { in } \tau \quad, \quad \tau_{\mathrm{ij}} v_{\mathrm{j}}=\mathrm{T}_{\mathrm{i}} \quad \text { on } \Sigma .
$$

Proof : Continuing from the previous theorem on complementary energy, we write

$$
\left(\delta \tau_{\mathrm{ij}}\right) v_{\mathrm{j}}=0 \quad \text { on } \Sigma=\Sigma_{\mathrm{T}} \mathrm{U} \Sigma_{\mathrm{u}},
$$

and equation (15) reduces to

$$
\Delta \mathrm{U}=\mathrm{P} \geq 0,
$$

showing that the increment $\Delta U$ in the strain energy $U$ of a body in equilibrium state is positive. Therefore, U is an absolute minimum.

Hence the result.

## Chapter-6

## Two-Dimensional Problems

### 6.1 INTRODUCTION

Many physical problems regarding the deformation of elastic solids are reducible to two-dimensional elastostatic problems. This reduction facilitates an easy solution.

### 6.2 PLANE STRAIN DEFORMATION

An elastic body is said to be in the state of plane strain deformation, parallel to the $\mathrm{x}_{1} \mathrm{x}_{2}$-plane, if the displacement component $\mathrm{u}_{3}$ vanishes identically and the other two displacement components $u_{1}$ and $u_{2}$ are function of $x_{1}$ and $x_{2}$ coordinates only and independent of $x_{3}$ coordinate.

Thus, the state of plane strain deformation (parallel to $\mathrm{x}_{1} \mathrm{x}_{2}$-plane) is characterised by the displacement components of the following type

$$
\begin{equation*}
\mathrm{u}_{1}=\mathrm{u}_{1}\left(\mathrm{x}_{1}, \mathrm{x}_{2}\right), \mathrm{u}_{2}=\mathrm{u}_{2}\left(\mathrm{x}_{1}, \mathrm{x}_{2}\right), \mathrm{u}_{3}=0 . \tag{1}
\end{equation*}
$$

The plane strain deformation is a two-dimensional approximation.
A plane strain state is used for a body in which one dimension is much larger than the other two.

For example, a long pressurized pipe or a dam between two massive end walls is a suitable case of plane strain deformation.

The maintenance of a state of plane strain requires the application of tension or pressure over the terminal sections, adjusted so as to keep constant the lengths of all the longitudinal filaments.
The states of plane strain deformation can be maintained in bodies of cylindrical form by suitable forces. We take the generators of the cylindrical bounding surface to be parallel to the $\mathrm{x}_{3}$-axis. We further suppose that the terminal sections are at right angles to this axis. The body force, if any, must be at right angles to the $\mathrm{x}_{3}$-axis and independent of it.

The strain components, $\mathrm{e}_{\mathrm{ij}}$ are given by the following strain-displacement relation

$$
\begin{equation*}
\mathrm{e}_{\mathrm{ij}}=\frac{1}{2}\left(\mathrm{u}_{\mathrm{i}, \mathrm{j}}+\mathrm{u}_{\mathrm{j}, \mathrm{i}}\right) . \tag{2a}
\end{equation*}
$$

We find, for plane strain deformation parallel to $\mathrm{x}_{1} \mathrm{x}_{2}$-plane,

$$
\begin{equation*}
e_{13}=e_{23}=e_{33}=0, \tag{3}
\end{equation*}
$$

and

$$
\begin{align*}
& \mathrm{e}_{11}=\frac{\partial \mathrm{u}_{1}}{\partial \mathrm{x}_{1}}, \mathrm{e}_{22}=\frac{\partial \mathrm{u}_{2}}{\partial \mathrm{x}_{2}}, \\
& \mathrm{e}_{12}=\frac{1}{2}\left(\frac{\partial \mathrm{u}_{1}}{\partial \mathrm{x}_{2}}+\frac{\partial \mathrm{u}_{2}}{\partial \mathrm{x}_{1}}\right), \tag{4}
\end{align*}
$$

which are independent of $x_{3}$.
It shows that all non-zero strains are on the $\mathrm{x}_{1} \mathrm{x}_{2}$-plane and $\mathrm{x}_{3}$-axis is strain-free/extension-free.

The strain matrix is, thus

$$
\left(\mathrm{e}_{\mathrm{ij}}\right)=\left(\begin{array}{lll}
\frac{\partial \mathrm{u}_{1}}{\partial \mathrm{x}_{1}} & \frac{1}{2}\left(\frac{\partial \mathrm{u}_{1}}{\partial \mathrm{x}_{2}}+\frac{\partial \mathrm{u}_{2}}{\partial \mathrm{x}_{1}}\right) & 0 \\
\frac{1}{2}\left(\frac{\partial \mathrm{u}_{1}}{\partial \mathrm{x}_{2}}+\frac{\partial \mathrm{u}_{2}}{\partial \mathrm{x}_{1}}\right) & \frac{\partial \mathrm{u}_{2}}{\partial \mathrm{x}_{2}} & 0 \\
0 & 0 & 1
\end{array}\right)
$$

The stress components $\tau_{i \mathrm{ij}}$ in terms of strains components $\mathrm{e}_{\mathrm{ij}}$ are governed by generalized Hooke's law for isotropic elastic solids

$$
\begin{equation*}
\tau_{\mathrm{ij}}=\lambda \delta_{\mathrm{ij}} \mathrm{e}_{\mathrm{kk}}+2 \mu \mathrm{e}_{\mathrm{ij}} . \tag{5a}
\end{equation*}
$$

We find, for plane strain deformation parallel to $\mathrm{x}_{1} \mathrm{x}_{2}$-plane,

$$
\begin{align*}
& \tau_{11}=\lambda\left(\mathrm{e}_{11}+\mathrm{e}_{22}\right)+2 \mu \mathrm{e}_{11}=(\lambda+2 \mu) \mathrm{e}_{11}+\lambda \mathrm{e}_{22},  \tag{5b}\\
& \tau_{22}=\lambda\left(\mathrm{e}_{11}+\mathrm{e}_{22}\right)+2 \mu \mathrm{e}_{22}=\lambda \mathrm{e}_{11}+(\lambda+2 \mu) \mathrm{e}_{22},  \tag{5c}\\
& \tau_{12}=2 \mu \mathrm{e}_{12},  \tag{5~d}\\
& \tau_{13}=\tau_{23}=0,  \tag{6a}\\
& \tau_{33}=\lambda\left(\mathrm{e}_{11}+\mathrm{e}_{22}\right)=\sigma\left(\tau_{11}+\tau_{22}\right)
\end{align*}
$$

$$
\begin{equation*}
=\frac{\lambda}{2(\lambda+\mu)}\left(\tau_{11}+\tau_{22}\right) \tag{6b}
\end{equation*}
$$

where

$$
\sigma=\frac{\lambda}{2(\lambda+\mu)},
$$

is the Poisson's ratio.
These relations shows all stress components are also independent of $\mathrm{x}_{3}$ coordinate.

Since stress

$$
\tau_{13}=\tau_{23}=0, \text { but } \tau_{33} \neq 0,
$$

so the strain-free axis ( $\mathrm{x}_{3}$-axis) is not stress-free, in general.
The Cauchy's equilibrium equations for an elastic solid are

$$
\begin{equation*}
\tau_{\mathrm{ij}, \mathrm{j}}+\mathrm{f}_{\mathrm{i}}=0, \tag{7}
\end{equation*}
$$

where $\underline{f}=f_{i}$ is the body force per unit volume.
In the case of plane strain deformation parallel to $\mathrm{x}_{1} \mathrm{x}_{2}$-plane, these equations reduce to, using equations (6),

$$
\begin{array}{r}
\tau_{11,1}+\tau_{12,2}+f_{1}=0 \\
\tau_{12,1}+\tau_{22,2}+f_{2}=0 \\
f_{3}=0 \tag{8}
\end{array}
$$

It shows that, for a plane strain deformation parallel to $\mathrm{x}_{1} \mathrm{x}_{2}$-plane, the body force is also independent of $x_{3}$ coordinate and the body force must be perpendicular to $\mathrm{x}_{3}$-direction.
In general, there are 6 Saint-Venant compatibility conditions for infinitesimal strain components. In the state of plane strain deformation five out of these 6 conditions are identically satisfied and the only compatibility condition to be considered further, for plane strain deformation parallel to $\mathrm{x}_{1} \mathrm{X}_{2}$-plane, is

$$
\begin{equation*}
\mathrm{e}_{11,22}+\mathrm{e}_{22,11}=2 \mathrm{e}_{12,12} \tag{9}
\end{equation*}
$$

Remark :- To distinguish plane strain case from the general case, we shall use subscripts $\alpha, \beta$ instead of $i, j$. We shall also assume that $\alpha, \beta$ vary form 1 to 2 .
From equations (5a) and (8a, b); we write as follows:

$$
\begin{align*}
& \tau_{\alpha \beta, \beta}+\mathrm{f}_{\alpha}=0, \quad \text { for } \alpha=1,2, \\
& \lambda \delta_{\alpha \beta}\left(\mathrm{e}_{11}+\mathrm{e}_{22}\right)_{, \beta}+\mu\left(\mathrm{u}_{\alpha, \beta \beta}+\mathrm{u}_{\beta, \alpha \beta}\right)+\mathrm{f}_{\alpha}=0, \\
& \lambda\left(\mathrm{e}_{11}+\mathrm{e}_{22}\right)_{, \alpha}+\mu\left[\nabla^{2} \mathrm{u}_{\alpha}+\left(\mathrm{e}_{11}+\mathrm{e}_{22}\right)_{, \alpha}\right]+\mathrm{f}_{\alpha}=0 \\
& (\lambda+\mu) \frac{\partial}{\partial \mathrm{x}_{\alpha}}\left(\mathrm{e}_{11}+\mathrm{e}_{22}\right)+\mu \nabla^{2} \mathrm{u}_{\alpha}+\mathrm{f}_{\alpha}=0 \tag{10}
\end{align*}
$$

for $\alpha=1,2$ and

$$
\begin{equation*}
\nabla^{2}=\frac{\partial^{2}}{\partial x_{1}^{2}}+\frac{\partial^{2}}{\partial x_{2}^{2}} \tag{10a}
\end{equation*}
$$

These are the equations of equilibrium (or Navier equations) for plane strain deformation, parallel to $\mathrm{x}_{1} \mathrm{x}_{2}$-plne.

## Beltrami - Michell Conditions of Compatibility for Plane Strain Deformation Parallel to $\mathbf{x}_{1} \mathbf{x}_{2}$-Plane.

Solving equations ( $5 \mathrm{~b}, \mathrm{c}, \mathrm{d}$ ) for strain components in terms of stresses, we write

$$
\begin{align*}
& \mathrm{e}_{11}=\frac{(\lambda+2 \mu) \tau_{11}-\lambda \tau_{22}}{4 \mu(\lambda+\mu)}, \\
& \mathrm{e}_{22}=\frac{(\lambda+2 \mu) \tau_{22}-\lambda \tau_{11}}{4 \mu(\lambda+\mu)}, \\
& 2 \mathrm{e}_{12}=\frac{1}{\mu} \tau_{12} . \tag{11}
\end{align*}
$$

Substituting the values of these strain components into Saint-Venant compatibility condition (9), we obtain

$$
\begin{gather*}
\frac{1}{4 \mu(\lambda+\mu)}\left[(\lambda+2 \mu) \tau_{11,22}-\lambda \tau_{22,22}+(\lambda+2 \mu) \tau_{22,11}-\lambda \tau_{11,11}\right]=\frac{1}{\mu} \tau_{12,12} \\
\quad(\lambda+2 \mu)\left(\tau_{11,22}+\tau_{22,11}\right)-\lambda\left(\tau_{11,11}+\tau_{22,22}\right)=4(\lambda+\mu) \tau_{12,12} \tag{12}
\end{gather*}
$$

Differentiating the equilibrium equations

$$
\tau_{\alpha \beta, \beta}+f_{\alpha}=0,
$$

w.r.t. $\mathrm{x}_{\alpha}$ and adding under summation convention, we find

$$
\begin{align*}
& \tau_{\alpha \beta, \alpha \beta}+\mathrm{f}_{\alpha, \alpha}=0, \\
& \tau_{11,11}+\tau_{22,22}+2 \tau_{12,12}+\mathrm{f}_{\alpha, \alpha}=0 . \tag{13}
\end{align*}
$$

Eliminating $\tau_{12,12}$ from (12) and (13), we get

$$
\begin{align*}
(\lambda+2 \mu)\left(\tau_{11,22}+\right. & \left.\tau_{22,11}\right)-\lambda\left(\tau_{11,11}+\tau_{22,22}\right)+2(\lambda+\mu)\left[\tau_{11,11}+\tau_{22,22}+f_{\alpha, \alpha}\right]=0 \\
& (\lambda+2 \mu)\left[\tau_{11,22}+\tau_{22,11}+\tau_{11,11} \tau_{22,22}\right]+2(\lambda+\mu) \mathrm{f}_{\alpha, \alpha}=0 \\
& \left(\frac{\partial^{2}}{\partial \mathrm{x}_{1}{ }^{2}}+\frac{\partial^{2}}{\partial \mathrm{x}_{2}{ }^{2}}\right)\left(\tau_{11}+\tau_{22}\right)+\frac{2(\lambda+\mu)}{\lambda+2 \mu}\left(\mathrm{f}_{1,1}+\mathrm{f}_{2,2}\right)=0 \\
& \left(\frac{\partial^{2}}{\partial \mathrm{x}_{1}{ }^{2}}+\frac{\partial^{2}}{\partial \mathrm{x}_{2}{ }^{2}}\right)\left(\tau_{11}+\tau_{22}\right)+\frac{1}{1-\sigma} \operatorname{div} \overline{\mathrm{f}}=0 . \tag{14}
\end{align*}
$$

When the body force is constant or absent, then the Beltrami-Michell compatibility condition (14) for plane strain deformation (parallel to $\mathrm{x}_{1} \mathrm{X}_{2^{-}}$ plane) reduces to

$$
\begin{equation*}
\left(\frac{\partial^{2}}{\partial \mathrm{x}_{1}{ }^{2}}+\frac{\partial^{2}}{\partial \mathrm{x}_{2}{ }^{2}}\right)\left(\tau_{11}+\tau_{22}\right)=0 \tag{15}
\end{equation*}
$$

Equation (15) shows that the stress $\tau_{11}+\tau_{22}$ is harmonic, when the body force is either absent or constant, and consequently $\left(e_{11}+e_{22}\right)$ is harmonic.
Note :- The generalized Hooke's law

$$
\tau_{\mathrm{ij}}=\frac{\mathrm{E}}{1+\sigma}\left[\mathrm{e}_{\mathrm{ij}}+\frac{\sigma}{1-2 \sigma} \delta_{\mathrm{ij}} \mathrm{e}_{\mathrm{kk}}\right],
$$

may also be used to calculate the stress components for plane strain deformation parallel to $\mathrm{x}_{1} \mathrm{x}_{2}$-plane is term of elastic modulli E , $\sigma$ instead of $\lambda$, $\mu$.

## Examples of Plane strain deformations

(A) The problem of stresses in an elastic semi-infinite medium subjected to a vertical line-load is a plane strain problem.


Here, the line-load extends to infinity on both sides of the origin. The displacement components are of the type

$$
\mathrm{u}_{1}=0, \mathrm{u}_{2}=\mathrm{u}_{2}\left(\mathrm{x}_{2}, \mathrm{x}_{3}\right), \mathrm{u}_{3}=\mathrm{u}_{3}\left(\mathrm{x}_{2}, \mathrm{x}_{3}\right),
$$

(B) The problem of determination of stresses resulting from a tangential lineload at the surface of a semi-infinite medium is a plane strain problem.

(C) The stresses and displacements in a semi-infinite elastic medium subjected to inclined loads can be obtained by superposition of the vertical and horizontal cases. If the components of the line-load are $\mathrm{q} \cos \alpha$ and $\mathrm{q} \sin \alpha$, the stresses can be determined.

(D) The problem of deformation of an infinite cylinder by a force in the $\mathrm{x}_{1} \mathrm{x}_{2}-$ plane is a plane strain problem.


In Cartesian coordinates

$$
\mathrm{u}_{1}=\mathrm{u}_{1}\left(\mathrm{x}_{1}, \mathrm{x}_{2}\right), \mathrm{u}_{2}=\mathrm{u}_{2}\left(\mathrm{x}_{1}, \mathrm{x}_{2}\right), \mathrm{u}_{3}=0
$$

In cylindrical coordinates

$$
\mathrm{u}_{\mathrm{r}}=\mathrm{u}(\mathrm{r}, \theta), \mathrm{u}_{\theta}=\mathrm{v}(\mathrm{r}, \theta), \mathrm{u}_{\mathrm{z}}=0,
$$

## Principal Strains And Directions For Plane Strain Deformation

A deformation for which the strain components $\mathrm{e}_{11}, \mathrm{e}_{22}$ and $\mathrm{e}_{12}$ are independent of $x_{3}$ and $e_{13}=e_{23}=e_{33} \equiv 0$ is called a plane strain deformation parallel to the $\mathrm{x}_{1} \mathrm{x}_{2}$-plane.

For such a deformation, the principal strain in the direction of $\mathrm{x}_{3}$-axis is zero and the strain quadric of Cauchy

$$
\begin{equation*}
\mathrm{e}_{\mathrm{ij}} \mathrm{x}_{\mathrm{i}} \mathrm{x}_{\mathrm{j}}= \pm \mathrm{k}^{2}, \tag{1}
\end{equation*}
$$

becomes

$$
\begin{equation*}
e_{11} x_{1}^{2}+2 e_{12} x_{1} x_{2}+e_{22} x_{2}^{2}= \pm k^{2}, \tag{2}
\end{equation*}
$$

which represents a cylinder in three-dimensions. Let the axes be rotated about $\mathrm{x}_{3}$-axis through an angle $\theta$ to get new axes $\mathrm{Ox}_{1}{ }^{\prime} \mathrm{x}_{2}{ }^{\prime} \mathrm{x}_{3}{ }^{\prime}$.


Let

$$
\begin{equation*}
\mathrm{a}_{\mathrm{ij}}=\cos \left(\mathrm{x}_{\mathrm{i}}^{\prime}, \mathrm{x}_{\mathrm{j}}\right) . \tag{3}
\end{equation*}
$$

Then

|  | $\mathrm{x}_{1}$ | $\mathrm{x}_{2}$ | $\mathrm{x}_{3}$ |
| :--- | :--- | :--- | :--- |
| $\mathrm{x}_{1}{ }^{\prime}$ | $\cos \theta$ | $\sin \theta$ | 0 |
| $\mathrm{x}_{2}{ }^{\prime}$ | $-\sin \theta$ | $\cos \theta$ | 0 |
| $\mathrm{x}_{3}{ }^{\prime}$ | 0 | 0 | 1 |

The strains $\mathrm{e}_{\mathrm{pq}}^{\prime}$ relative to primed system are given by the law

$$
\begin{equation*}
\mathrm{e}_{\mathrm{pq}}^{\prime}=\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qj}} \mathrm{e}_{\mathrm{ij}} \tag{5}
\end{equation*}
$$

for $(\mathrm{ij})=(11),(22),(12),(21)$. We find

$$
\begin{align*}
e_{11}^{\prime} & =a_{1 i} a_{1 j} e_{i j}=a_{11}^{2} e_{11}+a_{12}^{2} e_{22}+a_{11} a_{12} e_{12}+a_{12} a_{11} e_{12} \\
& =\cos ^{2} \theta \cdot e_{11}+\sin ^{2} \theta e_{22}+2 \sin \theta \cos \theta e_{12} \\
& =e_{11}\left(\frac{1+\cos 2 \theta}{2}\right)+e_{22}\left(\frac{1-\cos 2 \theta}{2}\right)+e_{12} \sin 2 \theta \\
& =\frac{1}{2}\left(e_{11}+e_{22}\right)+\frac{1}{2}\left(e_{11}-e_{22}\right) \cos 2 \theta+e_{12} \sin 2 \theta, \ldots(6 \tag{6a}
\end{align*}
$$

Similarly

$$
\begin{align*}
& \mathrm{e}_{22}^{\prime}=\frac{1}{2}\left(\mathrm{e}_{11}+\mathrm{e}_{22}\right)-\frac{1}{2}\left(\mathrm{e}_{11}-\mathrm{e}_{22}\right) \cos 2 \theta-\mathrm{e}_{12} \sin 2 \theta,  \tag{6b}\\
& \mathrm{e}_{12}^{\prime}=-\frac{1}{2}\left(\mathrm{e}_{11}-\mathrm{e}_{22}\right) \sin 2 \theta+\mathrm{e}_{12} \cos 2 \theta  \tag{6c}\\
& \mathrm{e}_{31}^{\prime}=\mathrm{e}_{32}^{\prime}=\mathrm{e}_{33}^{\prime}=0 . \tag{6d}
\end{align*}
$$

The principal directions in the $\mathrm{x}_{1} \mathrm{x}_{2}$-plane are given by

$$
\mathrm{e}_{12}^{\prime}=0
$$

This gives

$$
\begin{equation*}
\frac{\sin 2 \theta}{e_{12}}=\frac{\cos 2 \theta}{\frac{1}{2}\left(e_{11}-e_{22}\right)}=\frac{1}{\sqrt{e_{12}^{2}+\frac{1}{4}\left(e_{11}-e_{22}\right)^{2}}} \tag{7a}
\end{equation*}
$$

and

$$
\begin{equation*}
\tan 2 \theta=\frac{\mathrm{e}_{12}}{\frac{1}{2}\left(\mathrm{e}_{11}-\mathrm{e}_{22}\right)}=\frac{2 \mathrm{e}_{12}}{\mathrm{e}_{11}-\mathrm{e}_{22}} . \tag{7b}
\end{equation*}
$$

Let $\phi$ be the angle which the principal directions $\mathrm{O}_{1}$ and $\mathrm{O}_{2}$ make with the old axes in the $\mathrm{x}_{1} \mathrm{x}_{2}$-plane. Then

$$
\begin{equation*}
\tan 2 \phi=\frac{2 \mathrm{e}_{12}}{\mathrm{e}_{11}-\mathrm{e}_{22}} \tag{8}
\end{equation*}
$$

The principal strains $e_{1}$ and $e_{2}$ given by equations (6a,b) and (7a). We find ( $e_{1}$
$\left.=\mathrm{e}_{11}^{1}, \mathrm{e}_{2}=\mathrm{e}_{22}^{1}\right)$

$$
\begin{equation*}
e_{1}, e_{2}=\frac{1}{2}\left(e_{11}+e_{22}\right) \pm \sqrt{\frac{1}{4}\left(e_{11}-e_{22}\right)^{2}+e_{12}^{2}}, \tag{9}
\end{equation*}
$$

the shearing strain $\mathrm{e}^{\prime}{ }_{12}$ will be maximum when

$$
\begin{align*}
& \frac{\mathrm{d}}{\mathrm{~d} \theta} \mathrm{e}_{12}^{\prime}=0 \\
& -\left(\mathrm{e}_{11}-\mathrm{e}_{22}\right) \cos 2 \theta-2 \mathrm{e}_{12} \sin 2 \theta=0 \\
& \frac{\cos 2 \theta}{\mathrm{e}_{12}}=\frac{\sin 2 \theta}{-\frac{1}{2}\left(\mathrm{e}_{11}-\mathrm{e}_{22}\right)}=\frac{1}{\sqrt{\mathrm{e}_{12}^{2}+\frac{1}{4}\left(\mathrm{e}_{11}-\mathrm{e}_{22}\right)^{2}}} \tag{10a}
\end{align*}
$$

This gives the direction in which the shearing strain $\mathrm{e}^{\prime}{ }_{12}$ is maximum and maximum value of $\mathrm{e}^{\prime}{ }_{12}$ is given by equations (6c) and (10a). We find

$$
\begin{equation*}
\mathbf{d}_{12}^{\prime}-\quad=\sqrt{\mathrm{e}_{12}^{2}+\frac{1}{4}\left(\mathrm{e}_{11}-\mathrm{e}_{22}\right)^{2}} \tag{10b}
\end{equation*}
$$

From equations (9) and (10b), we obtain

$$
\begin{equation*}
\frac{\mathrm{e}_{1}-\mathrm{e}_{2}}{2}=\mathbf{d}_{12}^{\prime} \text { max. } \tag{11}
\end{equation*}
$$

This shows that maximum value of shearing strain is half of the difference of two principal strains in the $\mathrm{x}_{1} \mathrm{x}_{2}$ plane.

### 6.3 ANTIPLANE STRAIN DEFORMATION PARALLEL TO $\mathbf{x}_{1} \mathbf{x}_{2}$-PLANE

This deformation is characterised by

$$
\mathrm{u}_{1}=\mathrm{u}_{2} \equiv 0, \mathrm{u}_{3}=\mathrm{u}_{3}\left(\mathrm{x}_{1}, \mathrm{x}_{2}\right) .
$$

The strains are

$$
\begin{align*}
& \mathrm{e}_{11}=0, \mathrm{e}_{22}=\mathrm{e}_{12}=\mathrm{e}_{33}=0,  \tag{1}\\
& \mathrm{e}_{13}=\frac{1}{2} \frac{\partial \mathrm{u}_{3}}{\partial \mathrm{x}_{1}}, \mathrm{e}_{23}=\frac{1}{2} \frac{\partial \mathrm{u}_{3}}{\partial \mathrm{x}_{2}} . \tag{2}
\end{align*}
$$

Thus, only shear strains in the $\mathrm{x}_{3}$-direction are non-zero we can now find stresses from the Hooke's law

$$
\tau_{\mathrm{ij}}=\lambda v \cdot \delta_{\mathrm{ij}}+2 \mu \mathrm{e}_{\mathrm{ij}}
$$

giving $\quad \tau_{11}=\tau_{22}=\tau_{33}=0, \tau_{12}=0$,
and non-zero shear stresses are

$$
\begin{equation*}
\tau_{13}=\mu \frac{\partial \mathrm{u}_{3}}{\partial \mathrm{x}_{1}}, \tau_{23}=\mu \frac{\partial \mathrm{u}_{3}}{\partial \mathrm{x}_{2}} \tag{4}
\end{equation*}
$$

The equations of equilibrium are

$$
\tau_{\mathrm{ij}, \mathrm{j}}+\mathrm{F}_{\mathrm{i}}=0
$$

Using the above values of stresses, we see that for $i=1,2$, we must have

$$
\begin{equation*}
\mathrm{F}_{1}=\mathrm{F}_{2}=0, \tag{5}
\end{equation*}
$$

and for $\mathrm{i}=3$,

$$
\begin{align*}
& \tau_{31,1}+\tau_{32,2}+\tau_{33,3}+\mathrm{F}_{3}=0 \\
& \frac{\partial \tau_{31}}{\partial \mathrm{x}_{1}}+\frac{\partial \tau_{32}}{\partial \mathrm{x}_{2}}+\mathrm{F}_{3}=0 \tag{6}
\end{align*}
$$

In term of $u_{3}$, this may be written as

$$
\begin{aligned}
\mu\left(\frac{\partial^{2} \mathbf{u}_{3}}{\partial \mathbf{x}_{1}^{2}}+\frac{\partial^{2} \mathbf{u}_{3}}{\partial \mathbf{x}_{2}^{2}}\right)+\mathrm{F}_{3} & =0 . \\
\mu \nabla^{2} \mathbf{u}_{3}+\mathrm{F}_{3} & =0
\end{aligned}
$$

## Example of Anti-plane Deformation

Suppose that a force is applied along the line which is parallel to $\mathrm{x}_{1}$-axis and is situated at a depth $h$ below the free-surface of an elastic isotropic half-space.


The resulting deformation is that of anti-plane strain deformation with

$$
\mathrm{u}_{1}=\mathrm{u}_{1}\left(\mathrm{x}_{2}, \mathrm{x}_{3}\right), \mathrm{u}_{2}=\mathrm{u}_{3}=0
$$

Remark :- Two-dimensional problems in acoustics are antiplane strain problems.

### 6.4 PLANE STRESS

An elastic body is said to be in a state of plane stress parallel to the $\mathrm{x}_{1} \mathrm{x}_{2}$-plane if

$$
\begin{equation*}
\tau_{31}=\tau_{32}=\tau_{33}=0 \tag{1}
\end{equation*}
$$

and the remaining stress components $\tau_{11}, \tau_{22}, \tau_{12}$ are independent of $\mathrm{x}_{3}$.

The equilibrium equations

$$
\tau_{\mathrm{i}, \mathrm{j},}+\mathrm{f}_{\mathrm{i}}=0,
$$

for the case of plane stress reduce to

$$
\begin{array}{r}
\tau_{11,1}+\tau_{12,2} \mathrm{f}_{1}=0, \\
\tau_{12,1}+\tau_{22,2}+\mathrm{f}_{2}=0, \\
\mathrm{f}_{3}=0, \tag{2c}
\end{array}
$$

which are the same as for the case of plane strain deformation parallel to $\mathrm{x}_{1} \mathrm{x}_{2}-$ plane. In the state of plane stress, the body force $\bar{f}=\left(f_{1}, f_{2}, 0\right)$ must be independent of $\mathbf{x}_{3}$ as various stress components in Cauchy's equilibrium equations in (2) are independent of $x_{3}$.
The strain components $\mathrm{e}_{\mathrm{ij}}$ and stresses components $\tau_{\mathrm{ij}}$ are connected by the Hooke's law

$$
\begin{equation*}
\tau_{\mathrm{ij}}=\lambda \delta_{\mathrm{ij}} \mathrm{e}_{\mathrm{kk}}+2 \mu \mathrm{e}_{\mathrm{ij}} \tag{3}
\end{equation*}
$$

This gives

$$
\mathrm{e}_{12}=\frac{1}{2 \mu} \tau_{12}, \mathrm{e}_{13}=0, \mathrm{e}_{23}=0
$$

and

$$
\begin{align*}
& \tau_{33}=\lambda\left(\mathrm{e}_{11}+\mathrm{e}_{22}+\mathrm{e}_{33}\right)+2 \mu \mathrm{e}_{33} \\
& \mathrm{e}_{33}=-\frac{\lambda\left(\mathrm{e}_{11}+\mathrm{e}_{22}\right)}{\lambda+2 \mu} \neq 0 . \tag{5a}
\end{align*}
$$

Hence

$$
e_{k k}=e_{11}+e_{22}+e_{33}=\left(e_{11}+e_{22}\right)-\frac{\lambda\left(e_{11}+e_{22}\right)}{\lambda+2 \mu}
$$

$$
\begin{equation*}
=\frac{2 \mu}{\lambda+2 \mu}\left(\mathrm{e}_{11}+\mathrm{e}_{22}\right) . \tag{5b}
\end{equation*}
$$

This shows that strain component $\mathrm{e}_{13}$ and $\mathrm{e}_{23}$ are zero but $\mathrm{e}_{33}$ is not zero.
Hence, a state of plane stress does not imply a corresponding state of plane strain.
In view of Hooke's law (3), the strain components also do not depend upon $x_{3}$.
Let

$$
\begin{equation*}
\bar{\lambda}=\frac{2 \lambda \mu}{\lambda+2 \mu} \tag{6}
\end{equation*}
$$

From (3), we write
$\tau_{11}=\lambda \frac{2 \mu}{\lambda+2 \mu}\left(\mathrm{e}_{11}+\mathrm{e}_{22}\right)+2 \mu \mathrm{e}_{11}=(\bar{\lambda}+2 \mu) \mathrm{e}_{11}+\bar{\lambda} \mathrm{e}_{22}$

$$
\begin{equation*}
\tau_{22}=\bar{\lambda}\left(\mathrm{e}_{11}+\mathrm{e}_{22}\right)+2 \mu \mathrm{e}_{22}=\bar{\lambda} \mathrm{e}_{11}+(\bar{\lambda}+2 \mu) \mathrm{e}_{22} \tag{7b}
\end{equation*}
$$

comparing equations (7a, b) with the corresponding relations for plane strain deformation parallel to $\mathrm{x}_{1} \mathrm{x}_{2}$-plane, it is evident that solutions of plane stress problems can be obtained from the solutions of corresponding plane strain problems on replacing the true value of $\lambda$ by the apparent value $\bar{\lambda}=\frac{2 \lambda \mu}{\lambda+2 \mu}$.

## Strain Components in terms of Stress Components

Solving equations ( $7 \mathrm{a}, \mathrm{b}$ ) for $\mathrm{e}_{11}$ and $\mathrm{e}_{22}$, we find

$$
\begin{align*}
& \mathrm{e}_{11}=\frac{2(\lambda+\mu) \tau_{11}-\lambda \tau_{22}}{2 \mu(3 \lambda+2 \mu)}, \\
& \mathrm{e}_{22}=\frac{2(\lambda+\mu) \tau_{22}-\lambda \tau_{11}}{2 \mu(3 \lambda+2 \mu)} . \tag{8}
\end{align*}
$$

Substituting these values of $\mathrm{e}_{11}$ and $\mathrm{e}_{22}$ into equation (5a), we find

$$
\begin{equation*}
\mathrm{e}_{33}=\frac{-\lambda\left(\tau_{11}+\tau_{22}\right)}{2 \mu(3 \lambda+2 \mu)} \tag{9}
\end{equation*}
$$

Other strain components have been obtained in equation (4) already. All strain components are independent of $\mathrm{x}_{3}$ by Hooke's law.

In this plane stress problem, two compatibility equations are identically satisfied and the remaining four are

$$
\begin{align*}
& e_{11,22}+e_{22,11}=2 e_{12,12}  \tag{10a}\\
& e_{33,11}=e_{33,22}=e_{33,12}=0 . \tag{10b}
\end{align*}
$$

Since $e_{33}$ is independent of $x_{3}$ and satisfies all conditions in (10b), so $e_{33}$ must be of the type

$$
\begin{equation*}
e_{33}=c_{1}+c_{2} x_{2}+c_{3} x_{1}, \tag{11}
\end{equation*}
$$

where $\mathrm{c}_{1}, \mathrm{c}_{2}, \mathrm{c}_{3}$ are constants.
In most problems, equation (10a) is taken into consideration and requirement of equations in (10b) is ignored. This is possible, although approximately, when the dimension of the elastic body in the $\mathrm{x}_{3}$-direction is small.
In the plane stress state, strain components $e_{11}, e_{22}, e_{33}$ are independent of $x_{3}$ but the displacements may depend upon $\mathrm{x}_{3}$.

Hence, plane stress problems are not truly two-dimensional.

## Compatibility Equation in terms of Stresses

From equations (4), (8) and (10a), we write

$$
\begin{gather*}
\frac{1}{2 \mu(3 \lambda+2 \mu)}\left[2(\lambda+\mu) \tau_{11,22}-\lambda \tau_{22,22}+2(\lambda+\mu) \tau_{22,11}-\lambda \tau_{11,11}\right]=\frac{2}{2 \mu} \tau_{12,12} \\
2(\lambda+\mu)\left(\tau_{11,22}+\tau_{22,11}\right)-\lambda\left(\tau_{22,22} \quad \tau_{11,11}\right)=2(3 \lambda+2 \mu) \tau_{12,12} . \tag{12}
\end{gather*}
$$

From equation in (2), we write

$$
\begin{align*}
& \tau_{11,11}+\tau_{12,12}+\mathrm{f}_{1,1}=0, \tau_{12,12}+\tau_{22,22}+\mathrm{f}_{2,2}=0 \\
& \left(\tau_{11,11}+\tau_{22,22}\right)+\left(\mathrm{f}_{1,1}+\mathrm{f}_{2,2}\right)=-2 \tau_{12,12} . \tag{13}
\end{align*}
$$

From (12) and (13), we have

$$
\begin{align*}
& 2(\lambda+\mu)\left(\tau_{11,22}+\tau_{22,11}\right)-\lambda\left(\tau_{22,22}+\tau_{11,11}\right)=-(3 \lambda+2 \mu)\left[\tau_{11,11}+\tau_{22,22}+\mathrm{f}_{1,1}+\mathrm{f}_{2,2}\right] \\
& 2(\lambda+\mu)\left(\tau_{11,22}+\tau_{22,11}+(2 \lambda+2 \mu)\left(\tau_{22,22}+\tau_{11,11}\right)+(3 \lambda+2 \mu)\left(\mathrm{f}_{1,1}+\mathrm{f}_{2,2}\right)=0\right. \\
& 2(\lambda+\mu)\left[\left(\tau_{11,22}+\tau_{22,22}\right)+\left(\tau_{11,11}+\tau_{22,11}\right)\right]+(3 \lambda+2 \mu)\left(\mathrm{f}_{1,1}+\mathrm{f}_{2,2}\right)=0 \\
& \left(\frac{\partial^{2}}{\partial \mathrm{x}_{1}{ }^{2}}+\frac{\partial^{2}}{\partial \mathrm{x}_{2}^{2}}\right)\left(\tau_{11}+\tau_{22}\right)+\frac{3 \lambda+2 \mu}{2(\lambda+\mu)}\left(\mathrm{f}_{1,1}+\mathrm{f}_{2,2}\right)=0, \quad \ldots(14) \tag{14}
\end{align*}
$$

which is the same as obtained from the corresponding equation for plane strain deformation parallel to $\mathrm{x}_{1} \mathrm{x}_{2}$-plane on replacing $\lambda$ by $\frac{2 \lambda \mu}{\lambda+2 \mu}$.

Remark :- Since $\tau_{31}=\tau_{32}=0$, so $x_{3}$-axis a principal axis of stress and the corresponding principal stress $\tau_{3}$ is zero because

$$
\tau_{31}=\tau_{32}=\tau_{33}=0
$$

In the state of plane stress, one principal stress is zero or when one of the principal stress is zero, the state of stress is known as plane stress state.
Note :- A state of plane stress is obviously a possibility for bodies with one dimension much smaller than the other two. This type of state appears in the study of the deformation of a thin sheet plate when the plate is loaded by force applied at the boundary.
When the lengths of the generators in a cylindrical body are small in comparison with the linear dimensions of the cross-section, the body becomes a plate and the terminal sections are its faces.
The maintenance in a plate of a state of plane stress does not require the application of traction to the faces of the plate, but it required the body forces and tractions at the edge (or curved boundary) to be distributed in certain special ways.
In such a state, the stress components in the direction of the thickness of the plate are zero on both faces of the plate.
Question :- Discuss the principal stresses and principal directions of stress in a state of plane stress.
Answer :- Let an elastic body be in the state of plane stress parallel to the $\mathrm{x}_{1} \mathrm{x}_{2^{-}}$ plane. Then the stress components $\tau_{31}, \tau_{32}, \tau_{33}$ vanishes, i.e.,

$$
\tau_{31}=\tau_{32}=\tau_{33}=0
$$

The equation of stress quadric in the state of
plane stress becomes

$$
\tau_{11} \mathrm{x}^{2}+\tau_{22} \mathrm{y}^{2}+2 \tau_{12} \mathrm{xy}= \pm \mathrm{k}^{2}
$$

Let us rotate the $0 \mathrm{x}_{1} \mathrm{x}_{2} \mathrm{X}_{3}$ system about $0 \mathrm{x}_{3}$-axis by an amount $\theta$.


Then

$$
\begin{aligned}
& a_{11}=\cos \left(x_{1}{ }^{\prime}, x_{1}\right)=\cos \theta, \\
& a_{12}=\cos \left(x_{1^{\prime}}, x_{2}\right)=+\sin \theta, \\
& a_{13}=0 \\
& a_{21}=\cos \left(x_{2}^{\prime}, x_{1}\right)=-\sin \theta, \\
& a_{22}=\cos \left(x_{2^{\prime}}, x_{2}\right)=\cos \theta, \\
& a_{23}=\cos \left(x_{2}^{\prime}, x_{3}\right)=0 \\
& a_{31}=0, a_{32}=0, a_{33}=1
\end{aligned}
$$

Let $\tau_{\mathrm{pq}}$ be the stresses relative to new system. Then

$$
\tau_{\mathrm{pq}}^{\prime}=\mathrm{a}_{\mathrm{pi}} \mathrm{a}_{\mathrm{qj}} \tau_{\mathrm{ij}} .
$$

This gives

$$
\begin{aligned}
\tau_{11}^{\prime} & =a_{1 i} a_{1 j} \tau_{\mathrm{ij}} \\
& =\tau_{11} \cos ^{2} \theta+\tau_{22} \sin ^{2} \theta+2 \tau_{12} \cos \theta \sin \theta \\
& =\frac{1}{2}\left(\tau_{11}+\tau_{22}\right)+\frac{1}{2}\left(\tau_{11}-\tau_{22}\right) \cos 2 \theta+\tau_{12} \sin 2 \theta
\end{aligned}
$$

Similarly,

$$
\begin{aligned}
& \tau_{22}^{\prime}=\frac{1}{2}\left(\tau_{11}+\tau_{22}\right)-\frac{1}{2}\left(\tau_{11}-\tau_{22}\right) \cos 2 \theta-\tau_{12} \sin 2 \theta, \\
& \tau_{12}^{\prime}=-\frac{1}{2}\left(\tau_{11}-\tau_{22}\right) \sin 2 \theta+\tau_{12} \cos 2 \theta, \\
& \tau_{31}^{\prime}=\tau_{32}^{\prime}=\tau_{33}^{\prime}=0 .
\end{aligned}
$$

To obtain the other two principal directions of stress, we put

$$
\tau_{12}^{\prime}=0 .
$$

This gives

$$
\begin{aligned}
& \frac{1}{2}\left(\tau_{11}-\tau_{22}\right) \sin 2 \theta=\tau_{12} \cos 2 \theta \\
& \frac{\cos 2 \theta}{\frac{1}{2}\left(\tau_{11}-\tau_{22}\right)}=\frac{\sin 2 \theta}{\tau_{12}}=\frac{1}{\sqrt{\frac{1}{4}\left(\tau_{11}-\tau_{22}\right)^{2}+\tau_{12}^{2}}} \\
& \tan 2 \theta=\frac{2 \tau_{12}}{\tau_{11}-\tau_{22}} .
\end{aligned}
$$

This determines $\theta$ and hence the directions of two principal stresses $0 \mathrm{x}_{1}^{\prime}$ and $0 \mathrm{x}_{2}^{\prime}$.

Let $\tau_{1}$ and $\tau_{2}$ be the principal stresses in the directions $0 x_{1}^{\prime}$ and $0 x_{2}^{\prime}$ respectively. Then

$$
\begin{aligned}
\tau_{1}= & \frac{1}{2}\left(\tau_{11}+\tau_{22}\right)+\frac{1}{4} \frac{\left(\tau_{11}-\tau_{22}\right)^{2}}{\sqrt{\frac{1}{4}\left(\tau_{11}-\tau_{22}\right)^{2}+\tau_{12}^{2}}}+\frac{\tau_{22}^{2}}{\sqrt{\frac{1}{4}\left(\tau_{11}-\tau_{22}\right)^{2}+\tau_{12}^{2}}} \\
& =\frac{1}{2}\left(\tau_{11}-\tau_{22}\right)+\sqrt{\frac{1}{4}\left(\tau_{11}-\tau_{22}\right)^{2}+\tau_{12}^{2}},
\end{aligned}
$$

and

$$
\tau_{2}=\frac{1}{2}\left(\tau_{11}-\tau_{22}\right)-\sqrt{\frac{1}{4}\left(\tau_{11}-\tau_{22}\right)^{2}+\tau_{12}^{2}}
$$

The principal stress in the direction $0 \mathrm{x}_{3}$ or $\mathrm{ox}_{3}{ }^{\prime}$ is

$$
\tau_{3}=0
$$

The stress quadric with respect to principal axes becomes

$$
\tau_{1} \mathrm{x}_{1}^{, 2}+\tau_{2} \mathrm{x}_{2}^{, 2}= \pm \mathrm{k}^{2}
$$

which is a cylinder whose base is a conic (where may be called stress conic); its plane contains the directions of the two principal stresses which do not vanish.

Note : (The stress of plane stress is also defined as the one in which one principal stress is zero).

### 6.5 GENERALIZED PLANE STRESS

Consider a thin flat plate of thickness 2 h . We take the middle plane of the plate as $x_{3}=0$ plane so that the two faces of the plate are $x_{3}=h$ and $x_{3}=-h$. We make the following assumptions :

(a) The faces of plate are free from applied loads.
(b) The surface forces acting on the edge (curved surface) of the plate lie in planes, parallel to the middle plane ( $\mathrm{x}_{3}=0$ ), i.e., parallel to $\mathrm{x}_{1} \mathrm{x}_{2}$-plane and are symmetrically distributed w.r.t the middle plane $x_{3}=0$.
(c) $f_{3}=0$ and components $f_{1}$ and $f_{2}$ of the body force are symmetrically distributed w.r.t the middle plane.

Under these assumptions, the points of the middle plane will not undergo any deformation in the $\mathrm{x}_{3}$-direction. Let

$$
\begin{equation*}
\overline{\mathrm{u}}_{3}\left(\mathrm{x}_{1}, \mathrm{x}_{2}\right)=\frac{1}{2 \mathrm{~h}} \int_{-\mathrm{h}}^{\mathrm{h}} \mathrm{u}_{3}\left(\mathrm{x}_{1}, \mathrm{x}_{2}, \mathrm{x}_{3}\right) \mathrm{dx}_{3} \tag{1}
\end{equation*}
$$

denote the mean value of $u_{3}$ over the thickness of the plate. Then $\bar{u}_{3}\left(x_{1}, x_{2}\right)$ is independent of $x_{3}$. The symmetrical distribution of external forces w.r.t. the middle plane implies that

$$
\begin{equation*}
\overline{\mathrm{u}}_{3}\left(\mathrm{x}_{1}, \mathrm{x}_{2}\right)=0 . \tag{2}
\end{equation*}
$$

Since the faces $\mathrm{x}_{3}= \pm \mathrm{h}$ of the plate are free from applied loads (as assumed in
(a), so

$$
\begin{equation*}
\tau_{31}\left(\mathrm{x}_{1}, \mathrm{x}_{2}, \pm \mathrm{h}\right)=\tau_{32}\left(\mathrm{x}_{1}, \mathrm{x}_{2}, \pm \mathrm{h}\right)=\tau_{33}\left(\mathrm{x}_{1}, \mathrm{x}_{2}, \pm \mathrm{h}\right)=0 \tag{3a}
\end{equation*}
$$

for all admissible values of $x_{1}$ and $x_{2}$. Hence

$$
\begin{equation*}
\tau_{31,1}=\tau_{32,2}=0, \quad \text { at } \mathrm{x}= \pm \mathrm{h} . \tag{3b}
\end{equation*}
$$

Third equilibrium equation (with $f_{3}=0$ ) is

$$
\begin{equation*}
\tau_{31,1}+\tau_{32,2}+\tau_{33,3}=0 \tag{4}
\end{equation*}
$$

Using (2b), equation (4) reduce to

$$
\begin{equation*}
\tau_{33,3}\left(\mathrm{x}_{1}, \mathrm{x}_{2}, \pm \mathrm{h}\right)=0 . \tag{5}
\end{equation*}
$$

We note that $\tau_{33}$ and its derivative w.r.t. $x_{3}$ vanish on the faces of the plate. Since the thickness of plate is assumed to be very small, the stress component $\tau_{33}$ is small throughout of plate. Therefore we make assumption that

$$
\begin{equation*}
\tau_{33}=0 \tag{6}
\end{equation*}
$$

throughout the plate.
Now, we make the following definition.
Definition :- The stressed state of a thin plate for which $\tau_{33}=0$ everywhere and $\tau_{31}, \tau_{32}$ vanish on the two faces of the plate is known as generalized plane stress.
The remaining equilibrium equations for an elastic body are

$$
\tau_{\alpha 1,1}+\tau_{\alpha 2,2}+\tau_{\alpha 3,3}+f_{\alpha}=0 \quad \text { for } \alpha=1,2 .
$$

Integrating w.r.t $\mathrm{x}_{3}$ between the limits -h and +h , we obtain

$$
\frac{1}{2 h} \int_{-h}^{h}\left[\tau_{\alpha 1,1}+\tau_{\alpha 2,2}+\tau_{\alpha 3,3}+f_{\alpha}\right] d x_{3}=0
$$

$$
\begin{equation*}
\bar{\tau}_{\alpha 1,1}+\bar{\tau}_{\alpha 2,2}+\overline{\mathrm{f}}_{\alpha}=0 \tag{7}
\end{equation*}
$$

for $\alpha=1,2$, because,

$$
\begin{align*}
\int_{-\mathrm{h}}^{\mathrm{h}} \tau_{\alpha 3,3} \mathrm{dx}_{3} & =\tau_{\alpha 3}\left(\mathrm{x}_{1}, \mathrm{x}_{2}, \mathrm{~h}\right)-\tau_{\alpha 3}\left(\mathrm{x}_{1}, \mathrm{x}_{2},-\mathrm{h}\right) \\
& =0-0=0 \tag{7a}
\end{align*}
$$

as $\tau_{31}$ and $\tau_{32}$ vanish on the two faces of the plate.
Equations in (7) are the equilibrium equations for the mean values of the stresses and forces. Here $\bar{\tau}_{\alpha 1}$ etc. represents mean values.
When a plate is thin, the determination of the mean values of the components of displacement, strain and stress, taken over the thickness of the plate, may lead to knowledge nearly as useful as that of the actual values at each point. The actual values of the stresses, strains and displacements produced in the plate are determined in the case of plane stress state.
We note that the mean values of the displacements and stresses (which are independent of $\mathrm{x}_{3}$ ) for the generalized plane stress problem satisfy the same set of equations that govern the plane strain problem, the only difference being is that we have to replace $\lambda$ by $\frac{2 \lambda \mu}{\lambda+2 \mu}$.

## The state of generalized stress is purely two-dimensional and similar to the plane strain deformation, parallel to $\mathbf{x}_{1} \mathbf{x}_{2}$-plane.

We introduce the average field quantities $\overline{\mathrm{u}_{\mathrm{i}}}, \overline{\mathrm{e}_{\mathrm{ij}}}, \overline{\tau_{\mathrm{ij}}}$ as defined in equation (1) for $u_{3}$. Then

$$
\begin{equation*}
\overline{\mathrm{u}}_{1}=\overline{\mathrm{u}}_{1}\left(\mathrm{x}_{1}, \mathrm{x}_{2}\right), \quad \overline{\mathrm{u}}_{2}=\overline{\mathrm{u}}_{2}\left(\mathrm{x}_{1}, \mathrm{x}_{2}\right), \overline{\mathrm{u}}_{3}=0 \tag{8}
\end{equation*}
$$

Since $\tau_{33}=0$, so

$$
\begin{align*}
& \lambda\left(\mathrm{e}_{11}+\mathrm{e}_{22}+\mathrm{e}_{33}\right)+2 \mu \mathrm{e}_{33}=0 \\
& \mathrm{e}_{33}=\frac{-\lambda}{\lambda+2 \mu}\left(\mathrm{e}_{11}+\mathrm{e}_{22}\right) \tag{8a}
\end{align*}
$$

$$
\begin{equation*}
e_{11}+e_{22}+e_{33}=\left(1-\frac{\lambda}{\lambda+2 \mu}\right)\left(e_{11}+e_{22}\right)=\frac{2 \mu}{\lambda+2 \mu}\left(e_{11}+e_{22}\right) . \tag{8b}
\end{equation*}
$$

The generalized Hooke's law gives

$$
\begin{aligned}
& \tau_{\alpha \beta}=\lambda \delta_{\alpha \beta}\left(\mathrm{e}_{11}+\mathrm{e}_{22}+\mathrm{e}_{33}\right)+2 \mu \mathrm{e}_{\alpha \beta} \text { for } \alpha, \beta=1,2 \\
& \tau_{\alpha \beta}=\frac{2 \lambda \mu}{\lambda+2 \mu} \cdot \delta_{\alpha \beta} \cdot\left(\mathrm{e}_{11}+\mathrm{e}_{22}\right)+2 \mu \mathrm{e}_{\alpha \beta} .
\end{aligned}
$$

Integrating over $x_{3}$ and taking mean value over the thickness, we find

$$
\begin{equation*}
\bar{\tau}_{\alpha \beta}=\frac{2 \lambda \mu}{\lambda+2 \mu} \delta_{\alpha \beta}\left(\overline{\mathrm{e}}_{11}+\overline{\mathrm{e}}_{22}\right)+2 \mu \overline{\mathrm{e}}_{\alpha \beta}, \tag{9}
\end{equation*}
$$

for $\alpha, \beta=1,2$.
The five equations consisting of equations in (6) and (9) serve to determine the five unknown mean values $\bar{u}_{1}, \overline{\mathbf{u}}_{2}, \bar{\tau}_{11}, \bar{\tau}_{22}, \bar{\tau}_{12}$.
The substitution from (9) into (6) yields two equations of the Navier type

$$
\begin{equation*}
(\bar{\lambda}+\mu) \frac{\partial}{\partial \mathrm{x}_{\alpha}}\left(\overline{\mathrm{e}}_{11}+\overline{\mathrm{e}}_{22}\right)+\mu \nabla^{2} \overline{\mathrm{u}}_{\alpha}+\overline{\mathrm{F}}_{\alpha}\left(\mathrm{x}_{1}, \mathrm{x}_{2}\right)=0, \tag{10}
\end{equation*}
$$

from which the average displacements $\overline{\mathrm{u}}_{\alpha}$ can be determined when the values of the $\bar{u}_{\alpha}$ are specified on the contour. Here $\bar{\lambda}=2 \lambda \mu(\lambda+2 \mu)$.
Example 1 :- Thick-walled Tube Under External and Internal Pressures
We consider a cross-section of a thick-walled cylindrical tube whose inner radius is a and external radius $b$. We shall determine the deformation of the tube due to uniform internal pressure $\mathrm{p}_{1}$ and external pressure $\mathrm{p}_{2}$ acting on it.


We shall use the cylindrical coordinates ( $\mathrm{r}, \theta, \mathrm{z}$ ) to solve the problem and axis of cylinder tube is taken as z -axis.

This problem is a plane strain problem and due to symmetry, the cylindrical components of displacement are of the type

$$
\begin{equation*}
\mathrm{u}_{\mathrm{r}}=\mathrm{u}(\mathrm{r}), \mathrm{u}_{\theta}=0, \mathrm{u}_{\mathrm{z}}=0 \tag{1}
\end{equation*}
$$

We know that for an isotropic elastic medium, the Stoke's Navier equation of equilibrium for zero body force is

$$
\begin{equation*}
(\lambda+2 \mu) \text { grad div } \bar{u}-\mu \text { curl curl } \overline{\mathrm{u}}=\overline{0} . \tag{2}
\end{equation*}
$$

We find

$$
\operatorname{curl} \overline{\mathrm{u}}=\frac{1}{\mathrm{r}}\left|\begin{array}{ccc}
\hat{\mathrm{e}}_{\mathrm{r}} & \mathrm{re}_{\theta} & \hat{\mathrm{e}}_{\mathrm{z}}  \tag{3}\\
\partial / \partial \mathrm{r} & \partial / \partial \theta & \partial / \partial \mathrm{z} \\
\mathrm{u}_{\mathrm{r}} & 0 & 0
\end{array}\right|=\overline{0}
$$

and

$$
\operatorname{div} \overline{\mathrm{u}}=\frac{1}{\mathrm{r}}\left[\frac{\mathrm{~d}}{\mathrm{dr}}(\mathrm{ru})\right]=\frac{\mathrm{du}}{\mathrm{dr}}+\frac{\mathrm{u}}{\mathrm{r}},
$$

as u is a function fr only. We have used $\frac{\mathrm{d}}{\mathrm{dr}}$ instead of $\frac{\partial}{\partial \mathrm{r}}$.
From equations (2)-(4), we find

$$
\begin{align*}
\operatorname{grad} \operatorname{div} \overline{\mathrm{u}} & =\overline{0} \\
\operatorname{div} \overline{\mathrm{u}} & =\text { constt. } 2 \mathrm{~A}(\text { say }) \\
\frac{1}{\mathrm{r}} \frac{\mathrm{~d}}{\mathrm{dr}}(\mathrm{r} \mathrm{u}) & =2 \mathrm{~A} \\
\frac{\mathrm{~d}}{\mathrm{dr}}(\mathrm{ru}) & =2 \mathrm{Ar} \\
\mathrm{ru} & =\mathrm{Ar}^{2}+\mathrm{B}, \mathrm{~B}=\mathrm{constt} . \\
\mathrm{U} & =\mathrm{Ar}+\mathrm{B} / \mathrm{r} \tag{5}
\end{align*}
$$

where $A$ and $B$ are constants to be determined from the boundary conditions.
The strains in cylindrical coordinates ( $\mathrm{r}, \theta, \mathrm{z}$ ), using (1) and (5), are found to be

$$
\begin{align*}
& e_{\mathrm{rr}}=\mathrm{u}, \mathrm{r}=\mathrm{A}-\mathrm{B} / \mathrm{r}^{2} \\
& \mathrm{e}_{\theta \theta}=\frac{\mathrm{u}}{\mathrm{r}}=\mathrm{A}+\mathrm{B} / \mathrm{r}^{2}  \tag{6}\\
& \mathrm{e}_{\mathrm{r} \theta}=\mathrm{e}_{\theta \mathrm{z}}=\mathrm{e}_{\mathrm{rz}}=\mathrm{e}_{\mathrm{zz}}=0 \\
& \operatorname{div} \overline{\mathrm{u}}=\mathrm{e}_{\mathrm{rr}}+\mathrm{e}_{\theta \theta}=2 \mathrm{~A} .
\end{align*}
$$

The generalized Hooke's law for an isotropic material gives the expressions for stresses. We find

$$
\begin{align*}
& \tau_{\mathrm{rr}}=\lambda \operatorname{div} \overline{\mathrm{u}}+2 \mu \mathrm{e}_{\mathrm{rr}}=2 \mathrm{~A} \lambda+2 \mu\left(\mathrm{~A}-\mathrm{B} / \mathrm{r}^{2}\right)=2(\lambda+\mu) \mathrm{A}-2 \mu \frac{\mathrm{~B}}{\mathrm{r}^{2}}, \\
& \tau_{\theta \theta}=\lambda \operatorname{div} \overline{\mathrm{u}}+2 \mu \mathrm{e}_{\theta \theta}=2(\lambda+\mu) \mathrm{A}+2 \mu \mathrm{~B} / \mathrm{r}^{2}, \\
& \tau_{\mathrm{zz}}=\lambda \operatorname{div} \overline{\mathrm{u}}+2 \mu \mathrm{e}_{\mathrm{zz}}=2 \mathrm{~A} \lambda=\sigma\left(\tau_{\mathrm{rr}}+\tau_{\theta \theta}\right),  \tag{7}\\
& \tau_{\mathrm{r} \theta}=\tau_{\theta \mathrm{z}}=\tau_{\mathrm{rz}}=0 .
\end{align*}
$$

Here,

$$
\sigma=\frac{\lambda}{2(\lambda+\mu)} .
$$

Boundary conditions : The boundary conditions on the curved surface of the tube (or cross-section) are

$$
\left.\begin{array}{ll}
\tau_{\mathrm{rr}}=-\mathrm{p}_{1} & \text { at } \mathrm{r}=\mathrm{a} \\
\tau_{\mathrm{rr}}=-\mathrm{p}_{2} & \text { at } \mathrm{r}=\mathrm{b} \tag{8}
\end{array}\right\}
$$

From equations (7) and (8), we write

$$
\begin{aligned}
& -\mathrm{p}_{1}=2(\lambda+\mu) \mathrm{A}-2 \mu \mathrm{~B} / \mathrm{a}^{2}, \\
& -\mathrm{p}_{2}=2(\lambda+\mu) \mathrm{A}-2 \mu \mathrm{~B} / \mathrm{b}^{2} .
\end{aligned}
$$

Solving these equations for A and B , we find

$$
\mathrm{A}=\frac{\mathrm{p}_{1} \mathrm{a}^{2}-\mathrm{p}_{2} \mathrm{~b}^{2}}{2(\lambda+\mu)\left(\mathrm{b}^{2}-\mathrm{a}^{2}\right)},
$$

$$
\begin{equation*}
\mathrm{B}=-\frac{\mathrm{a}^{2} \mathrm{~b}^{2}\left(\mathrm{p}_{2}-\mathrm{p}_{1}\right)}{2 \mu\left(\mathrm{~b}^{2}-\mathrm{a}^{2}\right)} \tag{9}
\end{equation*}
$$

Putting these values of A and B in (5) and (7), we obtain the expressions for the displacement and stress which are

$$
\begin{equation*}
\mathrm{u}_{\mathrm{r}}=\mathrm{u}(\mathrm{r})=\frac{\mathrm{p}_{1} \mathrm{a}^{2}-\mathrm{p}_{2} \mathrm{~b}^{2}}{2(\lambda+\mu)\left(\mathrm{b}^{2}-\mathrm{a}^{2}\right)} \cdot \mathrm{r}-\frac{\left(\mathrm{p}_{2}-\mathrm{p}_{1}\right) \mathrm{a}^{2} \mathrm{~b}^{2}}{2 \mu\left(\mathrm{~b}^{2}-\mathrm{a}^{2}\right)} \cdot \frac{1}{\mathrm{r}}, \tag{10}
\end{equation*}
$$

This gives the displacement

$$
\overline{\mathrm{u}}=\mathrm{u} \hat{\mathrm{e}}_{\mathrm{r}}
$$

that occurs at a point distance $r$ from the axis of the tube $(a<r<b)$.
The stresses are

$$
\begin{align*}
& \tau_{\mathrm{rr}}=\frac{\mathrm{p}_{1} \mathrm{a}^{2}-\mathrm{p}_{2} \mathrm{~b}^{2}}{\mathrm{~b}^{2}-\mathrm{a}^{2}}+\frac{\mathrm{a}^{2} \mathrm{~b}^{2}\left(\mathrm{p}_{1}-\mathrm{p}_{2}\right)}{\left(\mathrm{b}^{2}-\mathrm{a}^{2}\right)} \cdot \frac{1}{\mathrm{r}^{2}}  \tag{11a}\\
& \tau_{\theta \theta}=\frac{\mathrm{p}_{1} \mathrm{a}^{2}-\mathrm{p}_{2} \mathrm{~b}^{2}}{\mathrm{~b}^{2}-\mathrm{a}^{2}}-\frac{\mathrm{a}^{2} \mathrm{~b}^{2}\left(\mathrm{p}_{2}-\mathrm{p}_{1}\right)}{\left(\mathrm{b}^{2}-\mathrm{a}^{2}\right)} \cdot \frac{1}{\mathrm{r}^{2}},  \tag{11b}\\
& \tau_{\mathrm{zz}}=\sigma\left(\tau_{\mathrm{rr}}+\tau_{\theta \theta}\right)=\frac{\lambda}{\lambda+\mu}\left(\frac{\mathrm{p}_{1} \mathrm{a}^{2}-\mathrm{p}_{2} \mathrm{~b}^{2}}{\mathrm{~b}^{2}-\mathrm{a}^{2}}\right) \tag{11c}
\end{align*}
$$

Equation (11c) show that $\tau_{z z}$ is constant. Hence, there is a uniform extension contraction in the direction of the axis of the tube. Moreover, cross-sections perpendicular to this axis remain plane after deformation.

## Rotating Shaft

Suppose that a solid long right circular cylinder (without an axle-hole) of radius a is rotating about its axis with uniform (constant) angular velocity $\omega$.
We assume that the cylinder is not free to deform longitudinally.
We shall be using the cylindrical co-ordinate system (r, $\theta, \mathrm{z}$ ) to determine the displacement and stresses at any point of the cylinder.
We consider a cross-section of long right circular cylinder of radius a. This cross-section is a circle with radius $a$. Consider a point $\mathrm{P}(\mathrm{r}, \theta)$ at a distance r from the origin.


Due to symmetry, the displacement components are dependent on $r$ only and

$$
\begin{equation*}
\mathrm{u}_{\mathrm{r}}=\mathrm{u}(\mathrm{r}), \mathrm{u}_{\theta}=\mathrm{u}_{\mathrm{z}}=0 . \tag{1}
\end{equation*}
$$

This problem is a plane strain problem.
We know that radial and transverse components of acceleration are

$$
\ddot{\mathrm{r}}-\mathrm{r} \dot{\theta}^{2}, 2 \dot{\mathrm{r}} \dot{\theta}+\mathrm{r} \ddot{\theta} .
$$

Since

$$
\dot{\mathrm{r}}=0, \dot{\theta}=\omega, \ddot{\mathrm{r}}=\ddot{\theta}=0,
$$

so the components of acceleration are

$$
-\mathrm{r} \omega^{2}, 0 .
$$

Hence equation of motion is

$$
\begin{equation*}
(\lambda+2 \mu) \operatorname{grad} \operatorname{div} \overline{\mathrm{u}}-\mu \text { curl curl } \overline{\mathrm{u}}=-\rho \mathrm{r} \omega^{2} \hat{\mathrm{e}}_{\mathrm{r}}, \tag{2}
\end{equation*}
$$

where $\rho$ is the density of the shaft. In view of (1),

$$
\begin{equation*}
\text { curl curl } \overline{\mathrm{u}}=\overline{0}, \operatorname{div} \overline{\mathrm{u}}=\frac{\mathrm{du}}{\mathrm{dr}}+\mathrm{u} / \mathrm{r} . \tag{3}
\end{equation*}
$$

From (2) and (3), we find

$$
(\lambda+2 \mu) \operatorname{grad}\left(\frac{\mathrm{du}}{\mathrm{dr}}+\mathrm{u} / \mathrm{r}\right)+\rho \mathrm{r} \omega^{2} \hat{\mathrm{e}}_{\mathrm{r}}=0
$$

$$
\frac{\mathrm{d}}{\mathrm{dr}}\left(\frac{\mathrm{du}}{\mathrm{dr}}+\mathrm{u} / \mathrm{r}\right)+\frac{\rho \omega^{2}}{\lambda+2 \mu} \mathrm{r}=0
$$

Integrating, we obtain

$$
\begin{align*}
& \frac{\mathrm{du}}{\mathrm{dr}}+\mathrm{u} / \mathrm{r}+\frac{\rho \omega^{2}}{\lambda+2 \mu} \cdot \frac{\mathrm{r}^{2}}{2}=2 \mathrm{~A} \\
& \frac{1}{\mathrm{r}} \frac{\mathrm{~d}}{\mathrm{dr}}(\mathrm{u} . \mathrm{r})+\frac{\rho \omega^{2}}{\lambda+2 \mu} \cdot \frac{\mathrm{r}^{2}}{2}=2 \mathrm{~A} \\
& \frac{\mathrm{~d}}{\mathrm{dr}}(\mathrm{u} . \mathrm{r})+\frac{\rho \omega^{2}}{\lambda+2 \mu} \cdot \frac{\mathrm{r}^{3}}{2}=2 \mathrm{Ar} \\
& \mathrm{u} \cdot \mathrm{r}+\frac{\rho \omega^{2}}{\lambda+2 \mu} \cdot \frac{\mathrm{r}^{4}}{8}=\mathrm{Ar}+\mathrm{B} \\
& \mathrm{u}(\mathrm{r})=\mathrm{Ar}+\frac{\mathrm{B}}{\mathrm{r}}-\frac{\rho \omega^{2}}{\lambda+2 \mu} \cdot \frac{\mathrm{r}^{3}}{8}, \tag{4}
\end{align*}
$$

where A and B are constants to be determined from boundary conditions.
Since cylinder is a solid cylinder, we must take

$$
\mathrm{B}=0,
$$

since, otherwise, $|\mathrm{u}| \rightarrow \infty$ as $\mathrm{r} \rightarrow 0$. So (4) reduces to

$$
\begin{equation*}
\mathrm{u}(\mathrm{r})=\operatorname{Ar}-\frac{\rho \omega^{2}}{\lambda+2 \mu}\left(\frac{\mathrm{r}^{3}}{8}\right) \tag{5}
\end{equation*}
$$

We know that the generalized Hooke's law in term of cylindrical coordinates gives

$$
\begin{aligned}
\tau_{\mathrm{rr}} & =\lambda \operatorname{div} \overline{\mathrm{u}}+2 \mu \mathrm{e}_{\mathrm{rr}}=\lambda \operatorname{div} \overline{\mathrm{u}}+2 \mu \frac{\mathrm{du}}{\mathrm{dr}} \\
& =\lambda\left[\frac{\mathrm{du}}{\mathrm{dr}}+\frac{\mathrm{u}}{\mathrm{r}}\right]+2 \mu \frac{\mathrm{du}}{\mathrm{dr}}
\end{aligned}
$$

$$
\begin{align*}
& =\lambda\left[2 \mathrm{~A}-\frac{\rho \omega^{2}}{\lambda+2 \mu} \cdot \frac{\mathrm{r}^{2}}{2}\right]+2 \mu\left[\mathrm{~A}-\frac{\rho \omega^{2}}{\lambda+2 \mu} \cdot \frac{3 \mathrm{r}^{2}}{8}\right] \\
& =2(\lambda+\mu) \mathrm{A}-\frac{\rho \omega^{2}}{\lambda+2 \mu}(2 \lambda+3 \mu) \frac{\mathrm{r}^{2}}{4} . \tag{6}
\end{align*}
$$

The surface $\mathrm{r}=\mathrm{a}$ of the shaft is traction free, so the boundary condition is

$$
\begin{equation*}
\tau_{\mathrm{rr}}=0 \text { at } \mathrm{r}=\mathrm{a} \tag{7}
\end{equation*}
$$

From equations (7) and (8); we find

$$
\begin{equation*}
A=\frac{\rho \omega^{2}(2 \lambda+3 \mu)}{(\lambda+\mu)(\lambda+2 \mu)} \cdot \frac{a^{2}}{8} \tag{8}
\end{equation*}
$$

Hence, at any point $\mathrm{P}(\mathrm{r}, \theta)$ of the shaft, the displacement and stress $\tau_{\mathrm{rr}}$ due to rotation with angular velocity $\omega$ is

$$
\begin{align*}
& \overline{\mathrm{u}}=\mathrm{u}(\mathrm{r}) \cdot \hat{\mathrm{e}}_{\mathrm{r}}=\frac{\rho \omega^{2} \mathrm{r}}{8(\lambda+2 \mu)}\left[\frac{2 \lambda+3 \mu}{\lambda+\mu} \mathrm{a}^{2}-\mathrm{r}^{2}\right] \hat{\mathrm{e}}_{\mathrm{r}},  \tag{9}\\
& \tau_{\mathrm{rr}}=\frac{\rho \omega^{2}}{4}\left(\frac{2 \lambda+3 \mu}{\lambda+2 \mu}\right)\left[\mathrm{a}^{2}-\mathrm{r}^{2}\right] . \tag{10}
\end{align*}
$$

The other non-zero stresses are

$$
\begin{gather*}
\tau_{\theta \theta}=\lambda \operatorname{div} \overline{\mathrm{u}}+2 \mu \mathrm{e}_{\theta \theta}=\lambda \operatorname{div} \overline{\mathrm{u}}+2 \mu \cdot\left(\frac{\mathrm{u}}{\mathrm{r}}\right) \\
=\frac{\rho \omega^{2}}{4(\lambda+2 \mu)}\left[(2 \lambda+3 \mu) \mathrm{a}^{2}-(2 \lambda+\mu) \mathrm{r}^{2}\right],  \tag{11a}\\
\tau_{\mathrm{zz}}=\sigma\left(\tau_{\mathrm{rr}}+\tau_{\theta \theta}\right)=\lambda \operatorname{div} \overline{\mathrm{u}}=\lambda\left[2 \mathrm{~A}-\frac{\rho \omega^{2}}{\lambda+2 \mu} \cdot \frac{\mathrm{r}^{2}}{2}\right]
\end{gather*}
$$

$$
\begin{equation*}
=\frac{\lambda \rho \omega^{2}}{2(\lambda+2 \mu)}\left[\frac{2 \lambda+3 \mu}{2(\lambda+\mu)} \mathrm{a}^{2}-\mathrm{r}^{2}\right] \tag{11b}
\end{equation*}
$$

### 6.6 AIRY'S STRES FUNCTION FOR PLANE STRAIN PROBLEMS

In plane elastostatic problems, it is convenient to use the standard notation $\mathrm{x}, \mathrm{y}$, z instead of $\mathrm{x}_{1}, \mathrm{x}_{2}, \mathrm{x}_{3}$ for Cartesian system.
The plane strain Cauchy's equilibrium equations on the xy-plane are

$$
\begin{align*}
& \tau_{\mathrm{xx}, \mathrm{x}}+\tau_{\mathrm{xy}, \mathrm{y}}+\mathrm{f}_{\mathrm{x}}=0  \tag{1}\\
& \tau_{\mathrm{xy}, \mathrm{x}}+\tau_{\mathrm{yy}, \mathrm{y}}+\mathrm{f}_{\mathrm{y}}=0 \tag{2}
\end{align*}
$$

Assume that the external body force is conservative, so that

$$
\underline{f}=-\underline{\nabla} \mathrm{V},
$$

where V is the force potential. This gives

$$
\mathrm{f}_{\mathrm{x}}=-\mathrm{V}, \mathrm{x} \text { and } \mathrm{f}_{\mathrm{y}}=-\mathrm{V}, \mathrm{y} .
$$

Using this, equations (1) and (2) can be put in the form

$$
\begin{align*}
& \left(\tau_{\mathrm{xx}}-\mathrm{V}\right)_{, \mathrm{x}}+\tau_{\mathrm{xy}, \mathrm{y}}=0  \tag{3}\\
& \tau_{\mathrm{xy}, \mathrm{x}}+\left(\tau_{\mathrm{yy}}-\mathrm{V}\right)_{, \mathrm{y}}=0 \tag{4}
\end{align*}
$$

Equations (3) and (4) can be satisfied identically through the introduction of a stress function $\quad \Phi=\Phi(\mathrm{x}, \mathrm{y})$ such that

$$
\begin{align*}
& \tau_{\mathrm{xx}}=\frac{\partial^{2} \phi}{\partial \mathrm{y}^{2}}+\mathrm{V} \\
& \tau_{\mathrm{yy}}=\frac{\partial^{2} \phi}{\partial \mathrm{x}^{2}}+\mathrm{V} \\
& \tau_{\mathrm{xy}}=\frac{-\partial^{2} \phi}{\partial \mathrm{x} \partial \mathrm{y}} \tag{5}
\end{align*}
$$

The function $\Phi$ is know as Airy's stress function, after the name of a British astronomer G. B. Airy. The function $\phi$ is called a stress function as $\phi$ generates stresses.

The Beltrani-Michell compatibility equation for plane strain deformation (in term of stresses) is

$$
\nabla^{2}\left(\tau_{\mathrm{xx}}+\tau_{\mathrm{yy}}\right)+\frac{2(\lambda+\mu)}{\lambda+2 \mu}\{\operatorname{div} \overrightarrow{\mathrm{f}}\}=0,
$$

which now in the case of conservative body force and stress function $\Phi$ becomes, using (5),

$$
\begin{align*}
& \nabla^{2}\left\{\nabla^{2} \phi+2 \mathrm{~V}\right\}+\frac{2(\lambda+\mu)}{\lambda+2 \mu}\left\{-\nabla^{2} \mathrm{~V}\right\}=0 \\
& \nabla^{2} \nabla^{2} \phi+\left(\frac{2 \mu}{\lambda+2 \mu}\right) \nabla^{2} \mathrm{~V}=0 \tag{6}
\end{align*}
$$

where $\quad \nabla^{2}=\frac{\partial^{2}}{\partial \mathrm{x}^{2}}+\frac{\partial^{2}}{\partial \mathrm{y}^{2}}$.

Equation (6) shows that the stress-function $\Phi$ is a biharmonic function whenever V is harmonic.

If the body force is absent/vanish, then for plane strain problem, the stress function $\Phi$ satisfies the biharmonic equation

$$
\begin{equation*}
\nabla^{2} \nabla^{2} \phi=0 \tag{7}
\end{equation*}
$$

i.e. $\quad \frac{\partial^{4} \phi}{\partial x^{4}}+2 \frac{\partial^{4} \phi}{\partial x^{2} \partial y^{2}}+\frac{\partial^{4} \phi}{\partial y^{4}}=0$.

The formula's for displacements in terms of stress function $\phi$ can be obtained by integrating the stress-strain relations for plane strain, which are

$$
\begin{align*}
& \tau_{11}=\phi_{, 22}=(\lambda+2 \mu) \mathrm{e}_{11}+\lambda \mathrm{e}_{22},  \tag{8}\\
& \tau_{22}=\phi_{, 11}=\lambda \mathrm{e}_{11}+(\lambda+2 \mu) \mathrm{e}_{22}, \tag{9}
\end{align*}
$$

$$
\begin{equation*}
\tau_{12}=-\phi_{, 12}=2 \mu \mathrm{e}_{12}=\mu\left(\mathrm{u}_{1,2}+\mathrm{u}_{2,1}\right) . \tag{10}
\end{equation*}
$$

Solving equations (8) and (9), for $\mathrm{e}_{11}$ and $\mathrm{e}_{22}$, we get

$$
\begin{align*}
& 2 \mu \mathrm{u}_{1,1}=2 \mu \mathrm{e}_{11}=-\phi_{, 11}+\frac{\lambda+2 \mu}{2(\lambda+\mu)} \nabla^{2} \phi,  \tag{11}\\
& 2 \mu \mathrm{u}_{2,2}=2 \mu \mathrm{e}_{22}=-\phi_{, 22}+\frac{\lambda+2 \mu}{2(\lambda+\mu)} \nabla^{2} \phi . \tag{12}
\end{align*}
$$

The integration of equations (11) and (12) yields

$$
\begin{align*}
& 2 \mu u_{1}=-\phi_{, 1}+\frac{\lambda+2 \mu}{2(\lambda+\mu)} \int \nabla^{2} \phi d x+f(y),  \tag{13}\\
& 2 \mu u_{2}=-\phi_{, 2}+\frac{\lambda+2 \mu}{2(\lambda+\mu)} \int \nabla^{2} \phi d y+g(x), \tag{14}
\end{align*}
$$

where $f$ and $g$ are arbitrary functions. Substituting the values of $u_{1}$ and $u_{2}$ from (13) and (14) into equation (10),, we find that

$$
\begin{aligned}
& f^{\prime}(y)+g^{\prime}(x)=0 \\
& f^{\prime}(y)=-g^{\prime}(x)=\text { constt }=\alpha(\text { say }) \\
& f(y)=\alpha y+\beta \\
& g(x)=-\alpha x+\gamma
\end{aligned}
$$

where $\alpha, \beta, \gamma$ are constants.
The form of $f$ and $g$ indicates that they represent a rigid body displacement and can thus be ignored in the analysis of deformation.
Thus, whenever $\phi$ becomes knows; the displacements, strains and stresses can be obtained for the plane strain problem.

## Biharmonic Boundary-Value Problems

Let the body occupies the region bounded by the curve C. Then the boundary conditions are of the form

$$
\begin{equation*}
\tau_{\alpha \beta} v_{\beta}=\mathrm{T}_{\alpha}(\mathrm{s}), \tag{1}
\end{equation*}
$$

where $\mathrm{T}_{\alpha}(\mathrm{s})$ are known functions of the are parameter s on C , the are length s being measured along C from a fixed point, say A .


The direction cosines of a tangent are $\left(\frac{\mathrm{dx}_{1}}{\mathrm{ds}}, \frac{\mathrm{dx}_{2}}{\mathrm{ds}}\right)$, so the d . c '.s of a normal to the curve is

$$
\left(\frac{\mathrm{dx}_{2}}{\mathrm{ds}},-\frac{\mathrm{dx}_{1}}{\mathrm{ds}}\right) .
$$

Choosing

$$
v_{1}=\frac{\mathrm{dx}_{2}}{\mathrm{ds}}, \quad v_{2}=-\frac{\mathrm{dx}_{1}}{\mathrm{ds}},
$$

equation (1) can be written as

$$
\left.\begin{array}{l}
\tau_{11} \frac{\mathrm{dx}_{2}}{\mathrm{ds}}-\tau_{12} \frac{\mathrm{dx}_{1}}{\mathrm{ds}}=\mathrm{T}_{1}(\mathrm{~s}) \\
\tau_{21} \frac{\mathrm{dx}_{2}}{\mathrm{ds}}-\tau_{22} \frac{\mathrm{dx}_{1}}{\mathrm{ds}}=\mathrm{T}_{2}(\mathrm{~s}) \tag{2}
\end{array}\right\}
$$

The Airy stress function $\phi$ generates the stresses as given below

$$
\begin{equation*}
\tau_{11}=\frac{\partial^{2} \phi}{\partial \mathrm{x}_{2}^{2}}, \tau_{22}=\frac{\partial^{2} \phi}{\partial \mathrm{x}_{1}^{2}}, \tau_{12}=-\frac{\partial^{2} \phi}{\partial \mathrm{x}_{1} \partial \mathrm{x}_{2}} . \tag{3}
\end{equation*}
$$

From (2) and (3), we write

$$
\begin{aligned}
& \phi_{, 22} \frac{\mathrm{dx}_{2}}{\mathrm{ds}}+\phi_{, 12} \frac{\mathrm{dx}_{1}}{\mathrm{ds}}=\mathrm{T}_{1}(\mathrm{~s}) \\
& \Rightarrow \frac{\partial}{\partial \mathrm{x}_{2}}\left\{\frac{\partial \phi}{\partial \mathrm{x}_{2}}\right\} \frac{\mathrm{dx}_{2}}{\mathrm{ds}}+\frac{\partial}{\partial \mathrm{x}_{1}}\left\{\frac{\partial \phi}{\partial \mathrm{x}_{2}}\right\} \frac{\mathrm{dx}_{1}}{\mathrm{ds}}=\mathrm{T}_{1}(\mathrm{~s}),
\end{aligned}
$$

and

$$
\begin{aligned}
& -\phi, 12 \frac{\mathrm{dx}_{2}}{\mathrm{ds}}-\phi_{11} \frac{\mathrm{dx}_{1}}{\mathrm{ds}}=\mathrm{T}_{2}(\mathrm{~s}) \\
& \Rightarrow \quad \frac{\partial}{\partial \mathrm{x}_{2}}\left\{\frac{\partial \phi}{\partial \mathrm{x}_{1}}\right\} \frac{\mathrm{dx}_{2}}{\mathrm{ds}}+\frac{\partial}{\partial \mathrm{x}_{1}}\left\{\frac{\partial \phi}{\partial \mathrm{x}_{1}}\right\} \frac{\mathrm{dx}_{1}}{\mathrm{ds}}=-\mathrm{T}_{2}(\mathrm{~s}) .
\end{aligned}
$$

Using chain rule, we write

$$
\begin{gathered}
\frac{\mathrm{d}}{\mathrm{ds}}\left(\frac{\partial \phi}{\partial \mathrm{x}_{2}}\right)=\mathrm{T}_{1}(\mathrm{~s}) \\
\frac{\mathrm{d}}{\mathrm{ds}}\left(\frac{\partial \phi}{\partial \mathrm{x}_{1}}\right)=-\mathrm{T}_{2}(\mathrm{~s}) .
\end{gathered}
$$

Integrating these equations along C, we get

$$
\begin{align*}
& \frac{\partial \phi}{\partial \mathrm{x}_{1}}=-\int \mathrm{T}_{2}(\mathrm{~s}) \mathrm{ds}=\mathrm{f}_{1}(\mathrm{~s}), \\
& \ldots(4 \mathrm{a}) \\
& \frac{\partial \phi}{\partial \mathrm{x}_{2}}=\int \mathrm{T}_{1}(\mathrm{~s}) \mathrm{ds}=\mathrm{f}_{2}(\mathrm{~s}) \tag{4b}
\end{align*}
$$

Hence, the stress boundary problem of elasticity is related to the boundary value problem of the type

$$
\nabla^{2} \nabla^{2} \phi=0 \quad \text { in } \mathrm{R},
$$

$$
\begin{equation*}
\phi_{, \alpha}=\mathrm{f}_{\alpha}(\mathrm{s}) \quad \text { on } \mathrm{C}, \tag{5}
\end{equation*}
$$

where $f_{\alpha}(s)$ are known functions.
The boundary value problem (5) is know as the fundamental biharmonic boundary-value problem.
This boundary-value problem can be phrased in the following form.

Normal derivative of $\phi=\frac{\partial \phi}{\partial v}=\underline{\nabla} \phi \cdot \hat{v}$

$$
=\phi_{, \alpha} u_{\alpha}
$$

$$
=\mathrm{f}_{1}(\mathrm{~s}) \frac{\mathrm{dx}_{2}}{\mathrm{ds}}-\mathrm{f}_{2}(\mathrm{~s}) \frac{\mathrm{dx}_{1}}{\mathrm{ds}}
$$

$$
\begin{equation*}
\equiv \mathrm{g}(\mathrm{~s}), \text { say } \tag{6}
\end{equation*}
$$

on C,

Since

$$
\begin{align*}
\mathrm{d} \phi & =\phi, \alpha \mathrm{dx} \mathrm{x}_{\alpha} \\
& =\mathrm{f}_{\alpha} \mathrm{dx}_{\alpha} \\
\phi & =\int \mathrm{f}_{\alpha} \mathrm{dx}_{\alpha} \\
& =\int\left(\mathrm{f}_{\alpha} \frac{\mathrm{dx}_{\alpha}}{\mathrm{ds}}\right) \mathrm{ds} \\
& =\mathrm{f}(\mathrm{~s}) \text { say } \tag{7}
\end{align*}
$$

on C.
Thus, the knowledge of the $\phi_{, \alpha}(\mathrm{s})$ on C leads to compute the value of $\phi(\mathrm{s})$ and its normal derivative $\frac{\partial \phi}{\partial v}$ on C .
Conversely, if $\phi$ and $\frac{\partial \phi}{\partial v}$ are know on C, we can compute $\phi_{, \alpha}(\mathrm{s})$.
Consequently, the boundary value problem in (5) can be written in an equivalent form

$$
\left.\begin{array}{l}
\nabla^{2} \nabla^{2} \phi=0 \text { in } \mathrm{R}  \tag{8}\\
\phi=\mathrm{f}(\mathrm{~s}) \text { and } \frac{\mathrm{d} \phi}{\mathrm{~d} v}=\mathrm{g}(\mathrm{~s}) \text { on } \mathrm{C} .
\end{array}\right\}
$$

The boundary value problem in (8) is more convenient in some problems.

### 6.7 STRESS FUNCTION FOR PLANE STRESS PROBLEM

For plane stress case, the Airy stress function $\phi$ is defined in the same way as for the plane strain problem. The generalized Hooke's law gives

$$
\begin{aligned}
& \mathrm{e}_{11}=\frac{1}{\mathrm{E}}\left(\tau_{11}-\sigma \tau_{22}\right), \\
& \mathrm{e}_{22}=\frac{1}{\mathrm{E}}\left(\tau_{22}-\sigma \tau_{11}\right), \\
& \mathrm{e}_{12}=\frac{1+\sigma}{\mathrm{E}} \tau_{12} . \\
& \mathrm{e}_{11}=\frac{1}{\mathrm{E}}\left[\phi_{, 22}-\sigma \phi_{, 11}+(1-\sigma) \mathrm{V}\right], \\
& \mathrm{e}_{22}=\frac{1}{\mathrm{E}}\left[\phi_{, 11}-\sigma \phi_{, 22}+(1-\sigma) \mathrm{V}\right], \\
& \mathrm{E}_{12}=-\frac{1+\sigma}{2} \phi_{, 12} .
\end{aligned}
$$

Substitution into the Cauchy's compatibility equation, we find

$$
\nabla^{2} \nabla^{2} \phi+(1-\sigma) \nabla^{2} \mathrm{~V}=0
$$

For zero body forces, the stress function $\phi$ in plane stress problems, satisfies the biharmonic equation

$$
\nabla^{2} \nabla^{2} \phi=0
$$

Exercise :- Show that, when body force is absent, the stress function $\phi$ is a biharmonic function for both elasto-static problems-plane strain problem and plane stress state of the body

## Airy Stress Function in Polar Coordinate

(for both cases of plane stress and plane strain)
For zero body force, components $f_{r}, f_{\theta}$; the equilibrium equations in terms of 2-D polar coordinates $\quad(r, \theta)$

$$
\begin{aligned}
& \frac{1}{\mathrm{r}} \frac{\partial}{\partial \mathrm{r}}\left(\mathrm{r} \tau_{\mathrm{rr}}\right)+\frac{1}{\mathrm{r}} \frac{\partial}{\partial \theta} \tau_{\mathrm{r} \theta}-\frac{\tau_{\theta \theta}}{\mathrm{r}}=0, \\
& \frac{1}{\mathrm{r}^{2}} \frac{\partial}{\partial \mathrm{r}}\left(\mathrm{r}^{2} \tau_{\theta \mathrm{r}}\right)+\frac{1}{\mathrm{r}} \frac{\partial}{\partial \theta} \tau_{\theta \theta}=0,
\end{aligned}
$$

are identically satisfied if the stresses are derived from a function $\phi=\phi(\mathrm{r}, \theta)$

$$
\begin{aligned}
& \tau_{\mathrm{rr}}=\frac{1}{\mathrm{r}} \frac{\partial \phi}{\partial \mathrm{r}}+\frac{1}{\mathrm{r}^{2}} \frac{\partial^{2} \phi}{\partial \theta^{2}}, \\
& \tau_{\theta \theta}=\frac{\partial^{2} \phi}{\partial \mathrm{r}^{2}}, \\
& \tau_{\mathrm{r} \theta}=-\frac{\partial}{\partial \mathrm{r}}\left(\frac{1}{\mathrm{r}} \frac{\partial \phi}{\partial \theta}\right)
\end{aligned}
$$

The function $\phi(\mathrm{r}, \theta)$ is the Airy stress function.
In this case, the compatibility equation shows that $\phi(\mathrm{r}, \theta)$ satisfies the biharmonic equation

$$
\begin{align*}
& \nabla^{2} \nabla^{2} \phi=0 \\
& \left(\frac{\partial^{2}}{\partial \mathrm{r}^{2}}+\frac{1}{\mathrm{r}} \frac{\partial}{\partial \mathrm{r}}+\frac{1}{\mathrm{r}^{2}} \frac{\partial^{2}}{\partial \theta^{2}}\right)\left(\frac{\partial^{2} \phi}{\partial \mathrm{r}^{2}}+\frac{1}{\mathrm{r}} \frac{\partial \phi}{\partial \mathrm{r}}+\frac{1}{\mathrm{r}^{2}} \frac{\partial^{2} \phi}{\partial \theta^{2}}\right)=0 \tag{*}
\end{align*}
$$

which is know as the compatibility equation to be satisfied by the Airy stress function $\phi(\mathrm{r}, \theta)$.

If one can find a solution of $(*)$ that also satisfies the given boundary conditions, then the problem is solved, since by uniqueness theorem (called Kirchoff's uniqueness theorem), such a solution is unique.

### 6.8. DEFORMATION OF A SEMI-INFINITE ELASTIC ISOTROPIC SOLID WITH DISPLACEMENTS OR STRESSES PRESCRIBED ON THE PLANE BOUNDARY

We consider a semi-infinite elastic medium with $\mathrm{x}_{1}$-axis pointing into the medium so that the medium occupies the region $\mathrm{x}_{1}>0$ and $\mathrm{x}_{1}=0$ is the plane boundary.

$$
\mathrm{x}_{1}>0
$$

$\quad \lambda, \mu$
$\mathrm{x}_{1}$
( $\mathrm{x}_{3}=0$ cross-section of the medium)
We consider the plane strain deformation parallel to $\mathrm{x}_{1} \mathrm{x}_{2}$-plane. Then, the displacements are of the type

$$
\begin{equation*}
\mathrm{u}_{1}=\mathrm{u}_{1}\left(\mathrm{x}_{1}, \mathrm{x}_{2}\right), \mathrm{u}_{2}=\mathrm{u}_{2}\left(\mathrm{x}_{1}, \mathrm{x}_{2}\right), \mathrm{u}_{3}=0 \tag{1}
\end{equation*}
$$

The stresses are generated by the Airy stress function $\phi=\phi\left(\mathrm{x}_{1}, \mathrm{x}_{2}\right)$ such that

$$
\begin{equation*}
\tau_{11}=\frac{\partial^{2} \phi}{\partial \mathrm{x}_{2}^{2}}, \quad \tau_{22}=\frac{\partial^{2} \phi}{\partial \mathrm{x}_{1}^{2}}, \quad \tau_{12}=-\frac{\partial^{2} \phi}{\partial \mathrm{x}_{1} \partial \mathrm{x}_{2}}, \tag{2}
\end{equation*}
$$

where $\phi$ satisfies the biharmonic equation

$$
\begin{equation*}
\frac{\partial^{4} \phi}{\partial \mathrm{x}_{1}^{4}}+2 \frac{\partial^{4} \phi}{\partial \mathrm{x}_{1}^{2} \partial \mathrm{x}_{2}^{2}}+\frac{\partial^{4} \phi}{\partial \mathrm{x}_{2}^{4}}=0 \tag{3}
\end{equation*}
$$

For convenience, we write

$$
\begin{equation*}
\left(\mathrm{x}_{1}, \mathrm{x}_{2}, \mathrm{x}_{3}\right) \equiv(\mathrm{x}, \mathrm{y}, \mathrm{z}),\left(\mathrm{u}_{1}, \mathrm{u}_{2}, \mathrm{u}_{3}\right)=(\mathrm{u}, \mathrm{v}, \mathrm{w}) . \tag{4}
\end{equation*}
$$

We use the Fourier transform method to solve the biharmonic equation (4). We use $\bar{f}(x, k) \quad$ to denote the Fourier transform of $f(x, y)$. Then

$$
\begin{equation*}
\overline{\mathrm{f}}(\mathrm{x}, \mathrm{k})=\mathrm{F}[\mathrm{f}]=\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{\infty} \mathrm{f}(\mathrm{x}, \mathrm{y}) \exp \{-\mathrm{iky}\} \mathrm{dy} \tag{5}
\end{equation*}
$$

If $f(x, y)$ satisfies the Dirichlet conditions then at the points where it is continuous, we have the inverse transformation

$$
\begin{equation*}
\mathrm{f}(\mathrm{x}, \mathrm{y})=\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{\infty} \overline{\mathrm{f}}(\mathrm{x}, \mathrm{k}) \exp (\mathrm{iky}) \mathrm{dk} \tag{6}
\end{equation*}
$$

for $-\infty<y<\infty$. Also if

$$
\frac{\partial^{\mathrm{n}} \mathrm{f}}{\partial \mathrm{y}^{\mathrm{n}}} \rightarrow 0, \quad \text { when } \mathrm{y} \rightarrow \pm \infty
$$

for the case $\mathrm{n}=0,1,2, \ldots, \mathrm{r}-1$, then

$$
\begin{equation*}
\mathrm{F}\left[\frac{\partial^{\mathrm{r}} \mathrm{f}}{\partial \mathrm{y}^{\mathrm{r}}}\right]=(-\mathrm{ik})^{\mathrm{r}} \overline{\mathrm{f}}(\mathrm{x}, \mathrm{k}) . \tag{7}
\end{equation*}
$$

Taking the Fourier transform of equation (3) w.r.t. the variable y, we obtain

$$
\begin{equation*}
\frac{d^{4} \Phi}{d^{4}}-2 \mathrm{k}^{2} \frac{\mathrm{~d}^{4} \Phi}{\mathrm{dx}^{2}}+\mathrm{k}^{4} \Phi=0 \tag{8}
\end{equation*}
$$

where $\Phi(\mathrm{x}, \mathrm{k})$ is the Fourier transform of $\phi(\mathrm{x}, \mathrm{y})$. Equation (8) is an ODE of fourth order. Its solution is

$$
\begin{equation*}
\Phi(\mathrm{x}, \mathrm{k})=(\mathrm{A}+\mathrm{Bx}) \exp (-|\mathrm{k}| \mathrm{x})+(\mathrm{C}+\mathrm{Dx}) \exp (|\mathrm{k}| \mathrm{x}), \tag{9}
\end{equation*}
$$

where $\mathrm{A}, \mathrm{B}, \mathrm{C}, \mathrm{D}$ are constants and they may depend upon k also.
Since we require that (9) is bounded as $\mathrm{x} \rightarrow \infty$, we conclude that

$$
\begin{equation*}
\mathrm{C}=\mathrm{D}=0, \tag{10}
\end{equation*}
$$

so that (9) becomes

$$
\begin{equation*}
\Phi(\mathrm{x}, \mathrm{k})=(\mathrm{A}+\mathrm{Bx}) \exp (-|\mathrm{k}| \mathrm{x}) . \tag{11}
\end{equation*}
$$

Inverting (11) with the help of (6), we write

$$
\begin{equation*}
\phi(\mathrm{x}, \mathrm{y})=\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{\infty}[(\mathrm{A}+\mathrm{Bx}) \exp (-|\mathrm{k}| \mathrm{x})] \exp (\mathrm{iky}) \mathrm{dk} \tag{12}
\end{equation*}
$$

From equations (2) and (12), the stresses are found to be
$\tau_{11}=\frac{-1}{\sqrt{2 \pi}} \int_{-\infty}^{\infty} \mathrm{k}^{2}[(\mathrm{~A}+\mathrm{Bx}) \exp (-|\mathrm{k}| \mathrm{x})] \exp (\mathrm{iky}) \mathrm{dk}$,
$\tau_{12}=\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{\infty}|\mathrm{k}|[|\mathrm{k}| \mathrm{A}+\mathrm{B}(|\mathrm{k}| \mathrm{x}-2)] \exp (-|\mathrm{k}| \mathrm{x}) \exp (\mathrm{iky}) \mathrm{dk}$,
$\tau_{12}=-\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{\infty} \mathrm{ik}[|\mathrm{k}| \mathrm{A}-\mathrm{B}(1-|\mathrm{k}| \mathrm{x})] \exp (-|\mathrm{k}| \mathrm{x}) \exp (\mathrm{iky}) \mathrm{dk}$.
We know that the displacements $\mathrm{u}(\mathrm{x}, \mathrm{y})$ and $\mathrm{v}(\mathrm{x}, \mathrm{y})$ are given by

$$
\begin{align*}
& 2 \mu \mathrm{u}=-\frac{\partial \phi}{\partial \mathrm{x}}+\frac{1}{2 \alpha} \int \nabla^{2} \phi \mathrm{dx}  \tag{16}\\
& 2 \mu \mathrm{v}=-\frac{\partial \phi}{\partial \mathrm{x}}+\frac{1}{2 \alpha} \int \nabla^{2} \phi \mathrm{dy} \tag{17}
\end{align*}
$$

after neglecting the rigid body displacements, and

$$
\begin{equation*}
\alpha=\frac{\lambda+\mu}{\lambda+2 \mu} \tag{18}
\end{equation*}
$$

we find

$$
\begin{align*}
\nabla^{2} \phi & =\tau_{11}+\tau_{22} \\
& =\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{\infty}(-2 B)|k| \exp (-|k| x) \exp (\mathrm{iky}) \mathrm{dk} \tag{19}
\end{align*}
$$

Hence, we find
$2 \mu u(x, y)=\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{\infty}\left[+|k| A+B\left(-1+\frac{1}{\alpha}+|k| x\right) \exp (-|k| x) \exp (i k y) d k\right.$,
$2 \mu v(\mathrm{x}, \mathrm{y})=\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{\infty}\left[(-\mathrm{ik}) \mathrm{A}+\mathrm{B}\left(-\mathrm{ikx}-\frac{\mathrm{i}|\mathrm{k}|}{\mathrm{k} \alpha}\right)\right] \exp (-|\mathrm{k}| \mathrm{x}) \exp (\mathrm{iky}) \mathrm{dk}$.

## Case I. When stresses are prescribed on $\mathbf{x}=0$.

Let the boundary conditions be

$$
\begin{align*}
& \tau_{11}(0, y)=\mathrm{h}(\mathrm{y})  \tag{22}\\
& \tau_{12}(0, \mathrm{y})=\mathrm{g}(\mathrm{y}) \tag{23}
\end{align*}
$$

where $\mathrm{h}(\mathrm{y})$ and $\mathrm{g}(\mathrm{y})$ are known functions of y . Then

$$
\begin{align*}
& \bar{\tau}_{11}=\overline{\mathrm{g}}(\mathrm{k})=\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{\infty} \mathrm{g}(\mathrm{y}) \exp (-\mathrm{iky}) \mathrm{dy}  \tag{24}\\
& \bar{\tau}_{12}=\overline{\mathrm{h}}(\mathrm{k})=\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{\infty} \mathrm{h}(\mathrm{y}) \exp (-\mathrm{iky}) \mathrm{dy} \tag{25}
\end{align*}
$$

so that

$$
\begin{align*}
& \tau_{11}(0, y)=h(y)=\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{\infty} \bar{h}(k) \exp (i k y) d k  \tag{26}\\
& \tau_{12}(0, y)=g(y)=\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{\infty} \overline{\mathrm{g}}(\mathrm{k}) \exp (\mathrm{iky}) \mathrm{dk} \tag{27}
\end{align*}
$$

Putting $\mathrm{x}=0$ in equations (13) and (15), and comparing with respective equations (26) and (27), we obtain

$$
\begin{align*}
& -\mathrm{k}^{2}[\mathrm{~A}]=\overline{\mathrm{h}}(\mathrm{k})  \tag{28}\\
& \operatorname{ik}[|\mathrm{k}| \mathrm{A}-\mathrm{B}]=\overline{\mathrm{g}}(\mathrm{k}) \tag{29}
\end{align*}
$$

Solving these equations for $A$ and $B$, we obtain

$$
\begin{align*}
& A=-\frac{[\overline{\mathrm{h}}(\mathrm{k})]}{\mathrm{k}^{2}}  \tag{30}\\
& B=-\left(|\mathrm{k}| \frac{\overline{\mathrm{h}}(\mathrm{k})}{\mathrm{k}^{2}}+\frac{\overline{\mathrm{g}}(\mathrm{k})}{\mathrm{ik}}\right) . \tag{31}
\end{align*}
$$

Putting the values of coefficients A and B from equations (30) and (31) into (13) - (15) and (20) and (21), we obtain the integral expressions for displacements and stresses.

$$
\begin{gather*}
2 \mu \mathrm{u}(\mathrm{x}, \mathrm{y})=\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{\infty}\left[\left(\frac{1}{\mathrm{ik}}-\frac{1}{\mathrm{i} \alpha \mathrm{k}}-\frac{|\mathrm{k}| \mathrm{x}}{\mathrm{ik}}\right) \overline{\mathrm{g}}(\mathrm{k})+\left(-\frac{1}{\alpha|\mathrm{k}|} \mathrm{x}\right) \overline{\mathrm{h}}(\mathrm{k})\right], \\
\quad \exp (-|\mathrm{k}| \mathrm{x}) \exp (\mathrm{iky}) \mathrm{dk} \tag{32}
\end{gather*}
$$

$2 \mu \mathrm{~V}(\mathrm{x}, \mathrm{y})=\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{\infty}\left[\left(\frac{-1}{\alpha|\mathrm{k}|}+\mathrm{x}\right) \overline{\mathrm{g}}(\mathrm{k})+\left(\frac{\mathrm{i}}{\mathrm{k}}+\frac{\mathrm{ikx}}{|\mathrm{k}|}-\frac{\mathrm{i}}{\mathrm{k} \alpha}\right) \overline{\mathrm{h}}(\mathrm{k})\right]$.

$$
\begin{equation*}
\exp (-|k| x) \exp (i k y) d k \tag{33}
\end{equation*}
$$

$\tau_{11}(\mathrm{x}, \mathrm{y})=\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{\infty}[(-\mathrm{ikx}) \overline{\mathrm{g}}(\mathrm{k})+(1+\mathrm{x}|\mathrm{k}|) \overline{\mathrm{h}}(\mathrm{k})] \exp (-|\mathrm{k}| \mathrm{x}) \exp (\mathrm{iky}) \mathrm{dk}$,
$\tau_{22}(\mathrm{x}, \mathrm{y})=\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{\infty}\left[\left(\mathrm{ixk}-\frac{2|\mathrm{k}|}{\mathrm{k}}\right) \overline{\mathrm{g}}(\mathrm{k})+(1-\mathrm{x}|\mathrm{k}|) \overline{\mathrm{h}}(\mathrm{k})\right]$,
$\exp (-|k| x) \exp (i k y) d k$

$$
\begin{equation*}
\tau_{12}(\mathrm{x}, \mathrm{y})=\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{\infty}-\mathrm{x}|\mathrm{k}| \overline{\overline{\mathrm{g}}}(\mathrm{k})+(-\mathrm{ixk}) \overline{\mathrm{h}}(\mathrm{k}) \underline{-}-\overline{\exp }(-|\mathrm{k}| \mathrm{x}) \exp (\mathrm{iky}) \mathrm{dk} \tag{36}
\end{equation*}
$$

Now we consider two particular situations in which specific surface stresses are known.

## Particular Cases

(a) Normal line-load :- In this particular case, a normal line-load, P , per unit length acts on the z -axis, then

$$
\begin{align*}
& \mathrm{h}(\mathrm{y})=-\mathrm{P} \delta(\mathrm{y}),  \tag{37}\\
& \mathrm{g}(\mathrm{y})=0 . \tag{38}
\end{align*}
$$

Consequently,

$$
\begin{align*}
& \overline{\mathrm{h}}(\mathrm{k})=\frac{-\mathrm{P}}{\sqrt{2 \pi}} .  \tag{39}\\
& \overline{\mathrm{g}}(\mathrm{k})=0 . \tag{40}
\end{align*}
$$

Putting the values of $\overline{\mathrm{h}}(\mathrm{k})$ and $\overline{\mathrm{g}}(\mathrm{k})$ from equations (39) and (40) into equations (34) to (36), we find the following integral expressions for stresses at any point of isotropic elastic half-space due to a normal line-load acting on the z-axis.

$$
\begin{align*}
& \tau_{11}(\mathrm{x}, \mathrm{y})=\frac{-\mathrm{P}}{2 \pi} \int_{-\infty}^{\infty}(1+\mathrm{x}|\mathrm{k}|) \exp (-|\mathrm{k}| \mathrm{x}) \exp (\mathrm{iky}) \mathrm{dk}  \tag{41}\\
& \tau_{22}(\mathrm{x}, \mathrm{y})=\frac{-\mathrm{P}}{2 \pi} \int_{-\infty}^{\infty}(1-\mathrm{x}|\mathrm{k}|) \exp (-|\mathrm{k}| \mathrm{x}) \exp (\mathrm{iky}) \mathrm{dk}  \tag{42}\\
& \tau_{12}(\mathrm{x}, \mathrm{y})=\frac{-\mathrm{P}}{2 \pi} \int_{-\infty}^{\infty}(-\mathrm{ixk}) \exp (-|\mathrm{k}| \mathrm{x}) \exp (\mathrm{iky}) \mathrm{dk} \tag{43}
\end{align*}
$$

We shall evaluate the integrals (41) to (43). We find

$$
\begin{align*}
\tau_{11}(x, y) & =\frac{-P}{2 \pi}\left[\frac{2 x}{\left(y^{2}+x^{2}\right)}+\frac{2 x\left(x^{2}-y^{2}\right)}{\left(x^{2}+y^{2}\right)^{2}}\right] \\
& =\frac{-2 P}{\pi}\left[\frac{x^{3}}{\left(x^{2}+y^{2}\right)^{2}}\right]  \tag{44}\\
\tau_{22}(x, y) & =\frac{-P}{2 \pi}\left[\frac{2 x}{y^{2}+x^{2}}-\frac{2 x\left(x^{2}-y^{2}\right)}{\left(x^{2}+y^{2}\right)^{2}}\right] \\
& =\frac{-2 P}{\pi}\left[\frac{x y^{2}}{\left(x^{2}+y^{2}\right)^{2}}\right]  \tag{45}\\
\tau_{12}(x, y) & =\frac{-P}{2 \pi}\left[-i x\left\{\frac{4 i x y}{\left.\left(x^{2}+y^{2}\right)^{2}\right\}}\right]\right] \\
& =\frac{-2 P}{\pi}\left[\frac{x^{2} y}{\left(x^{2}+y^{2}\right)^{2}}\right] \tag{46}
\end{align*}
$$

using the following standard integrals.
(1)

$$
\int_{-\infty}^{\infty} \mathrm{e}^{+\mathrm{ik}(\mathrm{y})} \mathrm{dk}=2 \pi \delta(-\mathrm{y})
$$

$$
\begin{equation*}
\int_{-\infty}^{\infty}(|k|)^{-1} \mathrm{e}^{-|k| x} \mathrm{e}^{\mathrm{iky}} \mathrm{dk}=-\log \left(\mathrm{y}^{2}+\mathrm{x}^{2}\right) \tag{2}
\end{equation*}
$$

(3)

$$
\int_{-\infty}^{\infty} \mathrm{k}^{-1} \mathrm{e}^{-|k| x} \mathrm{e}^{i k y} d k=2 i \tan ^{-1}\left(\frac{\mathrm{y}}{\mathrm{x}}\right)
$$

(4)

$$
\int_{-\infty}^{\infty} e^{-|k| x} e^{i k y} d k=\frac{2 x}{y^{2}+x^{2}}
$$

$$
\begin{equation*}
\int_{-\infty}^{\infty} \frac{\mathrm{k}}{|\mathrm{k}|} \mathrm{e}^{-|\mathrm{k}| \mathrm{x}} \mathrm{e}^{i \mathrm{ky}} \mathrm{dk}=\frac{-2 \mathrm{i}(-\mathrm{y})}{\mathrm{y}^{2}+\mathrm{x}^{2}}=\frac{2 \mathrm{i} y}{\mathrm{x}^{2}+\mathrm{y}^{2}} \tag{5}
\end{equation*}
$$

$$
\begin{equation*}
\int_{-\infty}^{\infty} \mathrm{ke}^{-|k| x} \mathrm{e}^{\mathrm{iky}} \mathrm{dk}=\frac{4 \mathrm{iyx}}{\left(\mathrm{y}^{2}+\mathrm{x}^{2}\right)^{2}} \tag{6}
\end{equation*}
$$

$$
\begin{equation*}
\int_{-\infty}^{\infty}|k| e^{-|k| x} e^{i k y} d k=\frac{2\left(x^{2}-y^{2}\right)}{\left(x^{2}+y^{2}\right)^{2}} \tag{7}
\end{equation*}
$$

$$
\begin{equation*}
\int_{-\infty}^{\infty} k^{2} e^{-|k| x} e^{i k y} d k=\frac{4 x\left(x^{2}-3 y^{2}\right)}{\left(x^{2}+y^{2}\right)^{3}} \tag{8}
\end{equation*}
$$

$$
\begin{equation*}
\int_{-\infty}^{\infty} k|k| e^{-|k| x \mid} e^{i k y} d k=\frac{4 i y\left(3 x^{2}-y^{2}\right)}{\left(x^{2}+y^{2}\right)^{3}} \tag{9}
\end{equation*}
$$

The corresponding displacements can be found from equations (32), (33), (39) and (40). We find

$$
\begin{align*}
2 \mu \mathrm{u} & =\frac{-\mathrm{P}}{2 \pi} \int_{-\infty}^{\infty}\left(-\mathrm{x}-\frac{1}{\alpha|\mathrm{k}|}\right) \exp (-|\mathrm{k}| \mathrm{x}) \exp (\mathrm{iky}) \mathrm{dk} \\
& =\frac{\mathrm{P}}{\pi}\left[\frac{\mathrm{x}^{2}}{\mathrm{y}^{2}+\mathrm{x}^{2}}-\frac{\log \left(\mathrm{x}^{2}+\mathrm{y}^{2}\right)}{2 \alpha}\right]  \tag{47}\\
2 \mu \mathrm{v} & =\frac{-\mathrm{P}}{2 \pi} \int_{-\infty}^{\infty}\left(\frac{i}{\mathrm{k}}+\frac{\mathrm{ixk}}{|\mathrm{k}|}+\frac{1}{\mathrm{i} \alpha \mathrm{k}}\right) \exp (-|\mathrm{k}| \mathrm{x}) \exp (\mathrm{iky}) \mathrm{dk} \\
& =\frac{\mathrm{P}}{\pi}\left[\frac{\mathrm{xy}}{\mathrm{x}^{2}+\mathrm{y}^{2}}+\left(\frac{\alpha-1}{\alpha}\right) \tan ^{-1}\left(\frac{\mathrm{y}}{\mathrm{x}}\right)\right] . \tag{48}
\end{align*}
$$

(b) Normal pressure : Suppose that a uniform normal pressure $p_{0}$ acts over the strip $-\mathrm{a} \leq \mathrm{y} \leq \mathrm{a}$ on the surface $\mathrm{x}=0$ in the positive x -direction. The corresponding boundary conditions give

$$
\begin{align*}
& \mathrm{h}(\mathrm{y})=\left\{\begin{array}{cc}
-\mathrm{p}_{0} & |\mathrm{y}| \leq \mathrm{a} \\
0 & |\mathrm{y}|>\mathrm{a}
\end{array},\right.  \tag{49}\\
& \mathrm{g}(\mathrm{y})=0 . \tag{50}
\end{align*}
$$

We find

$$
\begin{equation*}
\overline{\mathrm{h}}(\mathrm{k})=-2 \mathrm{p}_{0}\left(\frac{\sin \mathrm{ka}}{\mathrm{k}}\right), \quad \overline{\mathrm{g}}(\mathrm{k})=0 . \tag{51}
\end{equation*}
$$

Proceeding as in the previous case, we find the following integral expressions for the stresses and displacements at any point of an isotropic elastic half-space due to normal pressure.

$$
\begin{align*}
& \tau_{11}(x, y)=-\frac{2 p_{0}}{\sqrt{2 \pi}} \int_{-\infty}^{\infty}(1+x|k|)\left(\frac{\sin \mathrm{ka}}{\mathrm{k}}\right) \exp (-|\mathrm{k}| \mathrm{x}) \exp (\mathrm{iky}) \mathrm{dk}  \tag{52}\\
& \tau_{22}(\mathrm{x}, \mathrm{y})=\frac{-2 \mathrm{p}_{0}}{\sqrt{2 \pi}} \int_{-\infty}^{\infty}(1-\mathrm{x}|\mathrm{k}|)\left(\frac{\sin \mathrm{ka}}{\mathrm{k}}\right) \exp (-|\mathrm{k}| \mathrm{x}) \exp (\mathrm{iky})  \tag{53}\\
& \tau_{12}(\mathrm{x}, \mathrm{y})=\frac{-2 \mathrm{p}_{0}}{\sqrt{2 \pi}} \int_{-\infty}^{\infty}(-\mathrm{ixk})\left(\frac{\sin \mathrm{ka}}{\mathrm{k}}\right) \exp (-|\mathrm{k}| \mathrm{x}) \exp (\mathrm{iky}) \mathrm{dk} \tag{54}
\end{align*}
$$

$2 \mu \mathrm{u}(\mathrm{x}, \mathrm{y})=\frac{-2 \mathrm{p}_{0}}{\sqrt{2 \pi}} \int_{-\infty}^{\infty}\left(-\mathrm{x}-\frac{1}{\alpha|\mathrm{k}|}\right)\left(\frac{\sin \mathrm{ka}}{\mathrm{k}}\right) \exp (-|\mathrm{k}| \mathrm{x}) \exp (\mathrm{iky}) \mathrm{dk}$,
$2 \mu v(x, y)=\frac{-2 p_{0}}{\sqrt{2 \pi}} \int_{-\infty}^{\infty} \mathrm{i}\left(\frac{1}{\mathrm{k}}+\frac{\mathrm{xk}}{|\mathrm{k}|}-\frac{1}{\alpha \mathrm{k}}\right)\left(\frac{\sin \mathrm{ka}}{\mathrm{k}}\right) \exp (-|\mathrm{k}| \mathrm{x}) \exp (\mathrm{iky}) \mathrm{dk}$.

The earlier case of the normal line-load, P per unit length, becomes the particular case of the above uniform normal strip-loading case, by taking

$$
\begin{equation*}
\mathrm{P}_{0}=\frac{\mathrm{P}}{2 \mathrm{a}} \tag{57}
\end{equation*}
$$

and using the relation

$$
\begin{equation*}
\lim _{\mathrm{a} \rightarrow 0}\left(\frac{\sin \mathrm{ka}}{\mathrm{ka}}\right)=1 . \tag{58}
\end{equation*}
$$

Case 2. When surface displacements are prescribed on the boundary $\mathrm{x}=0$.
Let the boundary conditions be

$$
\begin{align*}
& \mathrm{u}(0, \mathrm{y})=\mathrm{h}_{1}(\mathrm{y})  \tag{59}\\
& \mathrm{v}(0, \mathrm{y})=\mathrm{g}_{1}(\mathrm{y}) \tag{60}
\end{align*}
$$

Then, we write, as before,

$$
\begin{align*}
& u(0, y)=\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{\infty} \overline{\mathrm{h}}_{1}(\mathrm{k}) \exp (\mathrm{iky}) \mathrm{dk}  \tag{61}\\
& \mathrm{v}(0, \mathrm{y})=\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{\infty} \overline{\mathrm{g}}_{1}(\mathrm{k}) \exp (\mathrm{iky}) \mathrm{dk} \tag{62}
\end{align*}
$$

Proceeding as earlier, we shall get the result.

### 6.9 GENERAL SOLUTION OF THE BIHARMONIC EQUATION

The biharmonic equation in two-dimension is

$$
\begin{equation*}
\nabla^{2} \nabla^{2} \phi\left(\mathrm{x}_{1}, \mathrm{x}_{2}\right)=0 \quad \text { in } \mathrm{R} . \tag{1}
\end{equation*}
$$

Here R is a region of $\mathrm{x}_{1} \mathrm{x}_{2}$-plane.
Let $\quad \nabla^{2} \phi=\mathrm{P}_{1}\left(\mathrm{x}_{1}, \mathrm{x}_{2}\right)$.
Then

$$
\begin{equation*}
\nabla^{2} \mathrm{P}_{1}=0, \tag{3}
\end{equation*}
$$

showing that $P_{1}\left(x_{1}, x_{2}\right)$ is a harmonic function. Let $P_{2}\left(x_{1}, x_{2}\right)$ be the conjugate harmonic function. Then, the function

$$
\begin{equation*}
F(z)=P_{1}+i P_{2}, \tag{4}
\end{equation*}
$$

is an analytic function of the complex variable

$$
\mathrm{z}=\mathrm{x}_{1}+\mathrm{ix}_{2}
$$

satisfying the Cauchy-Riemann equations

$$
\begin{align*}
& \mathrm{P}_{1,1}=\mathrm{P}_{2,2} \\
& \mathrm{P}_{2,1}=-\mathrm{P}_{1,2} \tag{5}
\end{align*}
$$

Let $\quad \mathrm{G}(\mathrm{z})=\frac{1}{4} \int \mathrm{~F}(\mathrm{z}) \mathrm{dz}=\mathrm{p}_{1}+\mathrm{ip}_{2}$.
Then $G(z)$ is also an analytic function such that

$$
\begin{equation*}
\mathrm{G}^{\prime}(\mathrm{z})=\frac{1}{4} \mathrm{~F}(\mathrm{z}), \tag{7}
\end{equation*}
$$

and by virtue of CR-equations for $\mathrm{G}(\mathrm{z})$,

$$
\begin{align*}
& \mathrm{p}_{1,1}=\mathrm{p}_{2,2}, \\
& \mathrm{p}_{1,2}=-\mathrm{p}_{2,1}, \tag{8}
\end{align*}
$$

we find $\quad \frac{\partial \mathrm{p}_{1}}{\partial \mathrm{x}_{1}}+\frac{\mathrm{i} \partial \mathrm{p}_{2}}{\partial \mathrm{x}_{1}}=\frac{1}{4}\left(\mathrm{P}_{1}+\mathrm{iP}_{2}\right)$.
This gives $\quad \mathrm{p}_{1,1}=\mathrm{p}_{2,2}=\frac{1}{4} \mathrm{P}_{1}$,

$$
\begin{equation*}
\mathrm{P}_{1,2}=-\mathrm{p}_{2,1}=\frac{-1}{4} \mathrm{P}_{2} . \tag{10}
\end{equation*}
$$

Now

$$
\begin{align*}
\nabla^{2}\left(\mathrm{p}_{1} \cdot \mathrm{x}_{1}\right) & =\mathrm{p}_{1} \nabla_{2} \mathrm{x}_{1}+\mathrm{x}_{1} \nabla^{2} \mathrm{p}_{1}+2 \nabla \mathrm{p}_{1} \cdot \nabla \mathrm{x}_{1} \\
& =0+0+2\left[\frac{\partial \mathrm{p}_{1}}{\partial \mathrm{x}_{1}} \frac{\partial \mathrm{x}_{1}}{\partial \mathrm{x}_{1}}+\frac{\partial \mathrm{p}_{1}}{\partial \mathrm{x}_{2}} \cdot \frac{\partial \mathrm{x}_{1}}{\partial \mathrm{x}_{2}}\right] \\
& =2\left[\mathrm{p}_{1,1}+0\right] \\
& =2 \mathrm{p}_{1,1} \tag{11}
\end{align*}
$$

Similarly, $\quad \nabla^{2}\left(\mathrm{p}_{2} \cdot \mathrm{x}_{2}\right)=2 \mathrm{p}_{2,2}$.

As $\mathrm{p}_{1}$ and $\mathrm{p}_{2}$ are harmonic functions therefore,

$$
\begin{aligned}
\nabla^{2}\left(\phi-\mathrm{p}_{1} \mathrm{x}_{1}-\mathrm{p}_{2} \mathrm{x}_{2}\right) & =\nabla^{2} \phi-2\left(\mathrm{p}_{1,1}+\mathrm{p}_{2,2}\right) \\
& =\mathrm{P}_{1}-2\left[\frac{1}{4} \mathrm{P}_{1}+\frac{1}{4} \mathrm{P}_{1}\right] .
\end{aligned}
$$

This implies

$$
\begin{equation*}
\nabla^{2}\left(\phi-\mathrm{p}_{1} \mathrm{x}_{1}-\mathrm{p}_{2} \mathrm{x}_{2}\right)=0 \text { in } \mathrm{R} \tag{13}
\end{equation*}
$$

Let

$$
\begin{equation*}
\phi-\mathrm{p}_{1} \mathrm{x}_{1}-\mathrm{p}_{2} \mathrm{x}_{2}=\mathrm{q}_{1}\left(\mathrm{x}_{1}, \mathrm{x}_{2}\right), \text { say . } \tag{14}
\end{equation*}
$$

Then $\mathrm{q}_{1}\left(\mathrm{x}_{1}, \mathrm{x}_{2}\right)$ is a harmonic function in R . Let $\mathrm{q}_{2}\left(\mathrm{x}_{1}, \mathrm{x}_{2}\right)$ be a conjugate function of $\mathrm{q}_{1}\left(\mathrm{x}_{1}, \mathrm{x}_{2}\right)$ and let

$$
\begin{equation*}
\mathrm{H}(\mathrm{z})=\mathrm{q}_{1}+\mathrm{iq}_{2} . \tag{15}
\end{equation*}
$$

Then $\mathrm{H}(\mathrm{z})$ is an analytic function of z .
From (14), we write

$$
\begin{align*}
& \phi\left(\mathrm{x}_{1}, \mathrm{x}_{2}\right)=\mathrm{p}_{1} \mathrm{x}_{1}+\mathrm{p}_{2} \mathrm{x}_{2}+\mathrm{q}_{1} \\
& \phi\left(\mathbf{x}_{1}, \mathbf{x}_{2}\right)=\operatorname{Re}[\overline{\mathbf{z}} \mathbf{G}(\mathbf{z})+\mathbf{H}(\mathbf{z})] . \tag{16}
\end{align*}
$$

Here, $\quad \bar{z}=x_{1}-i x_{2}$ and Re denotes the real part of the bracketed expression. The representation (16) of the biharmonic function $\phi\left(x_{1}, x_{2}\right)$ in terms of two analytic functions $\mathrm{G}(\mathrm{z})$ and $\mathrm{H}(\mathrm{z})$ was first obtained by GOURSAT.

Deduction : From (16), we write

$$
\begin{align*}
& \phi=\mathrm{p}_{1} \mathrm{x}_{1}+\mathrm{p}_{2} \mathrm{x}_{2}+\mathrm{q}_{1} \\
& \phi_{(\mathrm{z}, \mathrm{z})}=\frac{1}{2}[\overline{\mathrm{z}} \mathrm{G}(\mathrm{z})+\mathrm{z} \overline{\mathrm{G}(\mathrm{z})}+\mathrm{H}(\mathrm{z})+\overline{\mathrm{H}(\mathrm{z})}], \tag{17}
\end{align*}
$$

since $H(z)+\overline{H(z)}=2 q_{1}$ and $G(z)=p_{1}+i p_{2}, \overline{G(z)}=p_{1}-i p_{2}$.

## Some Terminology Involving Conjugate of Complex Functions

Suppose that t is a real variable and

$$
\begin{equation*}
\mathrm{f}(\mathrm{t})=\mathrm{U}(\mathrm{t})+\mathrm{iV}(\mathrm{t}) . \tag{1}
\end{equation*}
$$

is a complex valued function of the real variable $t$, where $U(t)$ and $V(t)$ are functions of $t$ with real coefficients. The conjugate $\bar{f}(t)$ of $f(t)$ is defined

$$
\begin{equation*}
\overline{\mathrm{f}}(\mathrm{t})=\mathrm{U}(\mathrm{t})-\mathrm{i} \mathrm{~V}(\mathrm{t}) \tag{2}
\end{equation*}
$$

If we replace the real variable $t$ by the complex variable $z(=x+i y)$, then $f(z)$ and $\overline{\mathrm{f}}(\mathrm{z})$ are defined to be

$$
\left.\begin{array}{c}
\mathrm{f}(\mathrm{z})=\mathrm{U}(\mathrm{z})+\mathrm{iV}(\mathrm{z}) \\
\overline{\mathrm{f}}(\mathrm{z})=\mathrm{U}(\mathrm{z})-\mathrm{iV}(\mathrm{z}) \tag{3}
\end{array}\right\}
$$

Similarly $\mathrm{f}(\overline{\mathrm{z}})$ and $\overline{\mathrm{f}}(\overline{\mathrm{z}})$, functions of the complex variable $\overline{\mathrm{z}}=\overline{\mathrm{x}+\mathrm{iy}}=\mathrm{x}$ -iy, are defined to be

$$
\left.\begin{array}{c}
\mathrm{f}(\overline{\mathrm{z}})=\mathrm{U}(\overline{\mathrm{z}})+\mathrm{iV}(\overline{\mathrm{z}})  \tag{4}\\
\overline{\mathrm{f}}(\overline{\mathrm{z}})=\mathrm{U}(\overline{\mathrm{z}})-\mathrm{iV}(\overline{\mathrm{z}})
\end{array}\right\}
$$

Now $\quad \overline{f(z)}=$ conjugate of $f(z)$

$$
\begin{aligned}
& =\text { conjugate of }\{\mathrm{U}(\mathrm{z})+\mathrm{iV}(\mathrm{z})\} \\
& =\mathrm{U}(\overline{\mathrm{z}})-\mathrm{iV}(\overline{\mathrm{z}})(\text { on changing } \mathrm{i} \rightarrow-\mathrm{i} \text { and } \mathrm{z} \text { to } \overline{\mathrm{z}}) \\
& =\overline{\mathrm{f}}(\overline{\mathrm{z}}) \text {, using }(4)
\end{aligned}
$$

(i) Thus $\overline{\mathrm{f}(\mathrm{z})}=\overline{\mathrm{f}}(\overline{\mathrm{z}})$.
(ii)

$$
\begin{align*}
\frac{\mathrm{d}}{\mathrm{dz}}[\overline{\mathrm{f}(\overline{\mathrm{z}})}] & =\frac{\mathrm{d}}{\mathrm{dz}}[\mathrm{U}(\mathrm{z})-\mathrm{iV}(\mathrm{z})], \text { using }(4 \mathrm{a}) \\
& =\mathrm{U}^{\prime}(\mathrm{z})-\mathrm{i} \mathrm{~V}^{\prime}(\mathrm{z}) \\
& =\mathrm{U}^{\prime}(\overline{\mathrm{z}})+\mathrm{i} \mathrm{~V}^{\prime}(\overline{\mathrm{z}}) \\
& =\overline{[\mathrm{U}(\overline{\mathrm{z}})+\mathrm{iV}(\overline{\mathrm{z}})]^{\prime}} \\
& =\overline{\mathrm{f}^{\prime}(\overline{\mathrm{z}})} . \tag{6}
\end{align*}
$$

Similarly

$$
\begin{equation*}
\left.\frac{\mathrm{d}}{\mathrm{dz}}[\overline{\mathrm{f}} \overline{\mathrm{z}})\right]=\mathrm{U}^{\prime}(\overline{\mathrm{z}})-\mathrm{i} \mathrm{~V}^{\prime}(\overline{\mathrm{z}})=\left[\overline{\mathrm{f}^{\prime}(\mathrm{z})}\right] . \tag{7}
\end{equation*}
$$

Stresses and Displacements in terms of Analytic Functions $\mathbf{G}(\mathbf{z})$ and $\mathbf{H}(\mathbf{z})$.
The stresses in terms of Airy stress function $\phi=\phi\left(\mathrm{x}_{1}, \mathrm{x}_{2}\right)$ are given by

$$
\begin{equation*}
\tau_{11}=\phi_{, 22}, \tau_{22}=\phi_{, 11}, \tau_{12}=-\phi_{, 12} . \tag{1}
\end{equation*}
$$

This gives

$$
\begin{align*}
\tau_{11}+\mathrm{i} \tau_{12} & =\phi_{, 22}-\mathrm{i} \phi_{, 12} \\
& =-\mathrm{i} \frac{\partial}{\partial \mathrm{x}_{2}}\left[\phi_{, 1}+\mathrm{i} \phi_{, 2}\right],  \tag{2}\\
\tau_{22}-\mathrm{i} \tau_{12} & =\phi_{, 11}+\mathrm{i} \phi_{, 12} \\
& =\frac{\partial}{\partial \mathrm{x}_{1}}\left[\phi_{, 1}+\mathrm{i} \phi_{, 2}\right] . \tag{3}
\end{align*}
$$

We have

$$
\begin{align*}
& \mathrm{z}=\mathrm{x}_{1}+\mathrm{i} \mathrm{x}_{2}, \quad \overline{\mathrm{z}}=\mathrm{x}_{1}-\mathrm{i} \mathrm{x}_{2} \\
& \mathrm{x}_{1}=\frac{1}{2}(\mathrm{z}+\overline{\mathrm{z}}), \mathrm{x}_{2}=\frac{1}{2 \mathrm{i}}(\mathrm{z}-\overline{\mathrm{z}}) . \tag{4}
\end{align*}
$$

By chain rule

$$
\begin{align*}
& \frac{\partial}{\partial \mathbf{x}_{1}}=\frac{\partial}{\partial \mathbf{z}}+\frac{\partial}{\partial \overline{\mathbf{z}}}, \\
& \frac{\partial}{\partial \mathbf{x}_{2}}=\mathrm{i}\left(\frac{\partial}{\partial \mathbf{z}}-\frac{\partial}{\partial \mathbf{z}}\right) . \tag{5}
\end{align*}
$$

Now

$$
\begin{align*}
\frac{\partial \phi}{\partial \mathbf{x}_{1}}+\mathrm{i} \frac{\partial \phi}{\partial \mathbf{x}_{2}} & =\left(\frac{\partial}{\partial \mathbf{x}_{1}}+\mathrm{i} \frac{\partial}{\partial \mathbf{x}_{2}}\right) \phi \\
& =\left[\frac{\partial}{\partial \mathbf{z}}+\frac{\partial}{\partial \mathbf{z}}-\left(\frac{\partial}{\partial \mathbf{z}}-\frac{\partial}{\partial \mathbf{z}}\right)\right] \phi \\
& =2 \frac{\partial \phi}{\partial \mathbf{z}} \tag{6}
\end{align*}
$$

Since

$$
\begin{equation*}
\phi=\frac{1}{2}[\overline{\mathrm{z}} \mathrm{G}(\mathrm{z})+\mathrm{z} \overline{\mathrm{G}(\mathrm{z})}+\mathrm{H}(\mathrm{z})+\overline{\mathrm{H}(\mathrm{z})}] \tag{7}
\end{equation*}
$$

where G and H are analytical functions, we find

$$
\begin{align*}
2 \frac{\partial \phi}{\partial \overline{\mathrm{z}}} & =\mathrm{G}(\mathrm{z})+\mathrm{z} \frac{\partial}{\partial \overline{\mathrm{z}}}\{\overline{\mathrm{G}(\mathrm{z})}\}+0+\frac{\partial}{\partial \overline{\mathrm{z}}}\{\overline{\mathrm{H}(\mathrm{z})}\} \\
& =\mathrm{G}(\mathrm{z})+\mathrm{z} \overline{\mathrm{G}^{\prime}(\mathrm{z})}+\overline{\mathrm{H}^{\prime}(\mathrm{z})} \\
& =\mathrm{G}(\mathrm{z})+\mathrm{z} \overline{\mathrm{G}^{\prime}(\mathrm{z})}+\overline{\mathrm{K}(\mathrm{z})} \tag{8}
\end{align*}
$$

where

$$
\begin{equation*}
\mathrm{K}(\mathrm{z})=\mathrm{H}^{\prime}(\mathrm{z}) . \tag{9}
\end{equation*}
$$

From equations (2), (3), (5), (6) and (8), we find

$$
\begin{align*}
\tau_{11}+\mathrm{i} \tau_{12} & =\left(\frac{\partial}{\partial \mathrm{z}}-\frac{\partial}{\partial \overline{\mathrm{z}}}\right)\left[\mathrm{G}(\mathrm{z})+\mathrm{z} \overline{\mathrm{G}^{\prime}(\mathrm{z})}+\overline{\mathrm{K}(\mathrm{z})}\right] \\
& =\left\{\mathrm{G}^{\prime}(\mathrm{z})+1 \overline{\mathrm{G}^{\prime}(\mathrm{z})}\right\}-\left\{\mathrm{z} \cdot \overline{\mathrm{G}^{\prime \prime}(\mathrm{z})}+\overline{\mathrm{K}^{\prime}(\mathrm{z})}\right\} \\
& =\mathrm{G}^{\prime}(\mathrm{z})+\overline{\mathrm{G}^{\prime}(\mathrm{z})}-\mathrm{z} \cdot \overline{\mathrm{G}^{\prime}(\mathrm{z})}-\overline{\mathrm{K}^{\prime}(\mathrm{z})} \tag{10}
\end{align*}
$$

because

$$
\begin{equation*}
\frac{\partial}{\partial \bar{Z}} \overline{\mathbf{G}^{\prime}(\mathrm{z})}=\overline{\frac{\partial}{\partial \mathbf{z}}\left(\mathrm{G}^{\prime}(\mathrm{z})\right.}=\overline{\mathrm{G}^{\prime}(\mathrm{z})} . \tag{11}
\end{equation*}
$$

Also

$$
\begin{align*}
\tau_{22}-\mathrm{i} \tau_{12} & =\left(\frac{\partial}{\partial \mathrm{z}}+\frac{\partial}{\partial \overline{\mathrm{z}}}\right)\left[\mathrm{G}(\mathrm{z})+\mathrm{z} \cdot \overline{\mathrm{G}^{\prime}(\mathrm{z})}+\overline{\mathrm{K}(\mathrm{z})}\right] \\
& =\mathrm{G}^{\prime}(\mathrm{z})+\overline{\mathrm{G}^{\prime}(\mathrm{z})}+\mathrm{z} \cdot \overline{\mathrm{G}^{\prime \prime}(\mathrm{z})}+\overline{\mathrm{K}^{\prime}(\mathrm{z})} \tag{12}
\end{align*}
$$

Adding and subtracting (10) from (10a), we obtain

$$
\tau_{11}+\tau_{22}=2\left[\mathrm{G}^{\prime}(\mathrm{z})+\overline{\mathrm{G}^{\prime}(\mathrm{z})}\right]
$$

$$
\begin{align*}
& =2.2 \operatorname{Re}\left\{\mathrm{G}^{\prime}(\mathrm{z})\right\} \\
& =4 \operatorname{Re}\left[\mathrm{G}^{\prime}(\mathrm{z})\right] \tag{13}
\end{align*}
$$

$$
\tau_{22}-\tau_{11}-2 \mathrm{i} \tau_{12}=2\left[\overline{\mathrm{zG} \mathrm{G}^{\prime}(\mathrm{z})}+\overline{\mathrm{K}^{\prime}(\mathrm{z})}\right]
$$

On taking conjugate both sides, we obtain

$$
\begin{equation*}
\tau_{22}-\tau_{11}+2 \mathrm{i} \tau_{12}=2\left[\mathrm{z} \mathrm{G}{ }^{\prime}(\mathrm{z})+\mathrm{K}^{\prime}(\mathrm{z})\right] . \tag{14}
\end{equation*}
$$

Equations (13) and (14) provide expressions for stresses for plane strain deformation in terms of analytic functions $\mathrm{G}(\mathrm{z})$ and $\mathrm{H}(\mathrm{z})$.

Expressions for Displacements :- We know that, for plane strain deformation parallel to $\mathrm{x}_{1} \mathrm{x}_{2}$-plane,

$$
\begin{align*}
& \tau_{11}=\phi_{, 22}=(\lambda+2 \mu) u_{1,1}+\lambda u_{2,2},  \tag{15}\\
& \tau_{22}=\phi_{, 11}=\lambda u_{1,1}+(\lambda+2 \mu) u_{2,2},  \tag{16}\\
& \tau_{12}=-\phi_{, 12}=\mu\left(u_{1,2}+u_{2,1}\right) . \tag{17}
\end{align*}
$$

Solving equations (15) and (16) for $u_{1,1}$ and $u_{2,2}$ in terms of $\phi_{, 11}$ and $\phi_{, 22}$, we find

$$
\begin{align*}
& 2 \mu \mathrm{u}_{1,1}=\phi_{, 11}+\frac{\lambda+2 \mu}{2(\lambda+\mu)} \nabla^{2} \phi  \tag{18}\\
& 2 \mu \mathrm{u}_{2,2}=-\phi_{, 22}+\frac{\lambda+2 \mu}{2(\lambda+\mu)} \nabla^{2} \phi \tag{19}
\end{align*}
$$

We know that

$$
\begin{equation*}
\nabla^{2} \phi=\mathrm{P}_{1}=4 \mathrm{p}_{1,1}=4 \mathrm{p}_{2,2} . \tag{20}
\end{equation*}
$$

Putting the values of $\nabla^{2} \phi$ from equations (20) into equations (18) and (17), we obtain

$$
\begin{align*}
& 2 \mu u_{1,1}=-\phi_{, 11}+\frac{2(\lambda+2 \mu)}{(\lambda+\mu)} p_{1,1},  \tag{21}\\
& 2 \mu u_{2,2}=-\phi_{, 22}+\frac{2(\lambda+2 \mu)}{\lambda+\mu} p_{2,2} . \tag{22}
\end{align*}
$$

The integration of these equations yields,

$$
\begin{align*}
& 2 \mu u_{1}=-\phi_{, 1}+\frac{2(\lambda+2 \mu)}{\lambda+\mu} p_{1}+f\left(x_{2}\right),  \tag{23}\\
& 2 \mu u_{2}=-\phi_{, 2}+\frac{2(\lambda+2 \mu)}{\lambda+\mu} p_{2}+g\left(x_{1}\right) . \tag{24}
\end{align*}
$$

where $f\left(x_{2}\right)$ and $g\left(x_{1}\right)$ are, as yet, arbitrary functions. Equation (17) serves to determine $f$ and $g$.

Since $p_{1,2}=-p_{2,1}$, we easily obtain from equations (17), (23) and (24) that

$$
\mathrm{f}^{\prime}\left(\mathrm{x}_{2}\right)+\mathrm{g}^{\prime}\left(\mathrm{x}_{1}\right)=0
$$

$$
f^{\prime}\left(x_{2}\right)=-g^{\prime}\left(x_{1}\right)=\text { constt }=\alpha(\text { say }) .
$$

Hence $\quad f\left(x_{2}\right)=\alpha x_{2}+\beta$,

$$
\begin{equation*}
\mathrm{g}\left(\mathrm{x}_{1}\right)=-\alpha \mathrm{x}_{1}+\gamma \tag{25}
\end{equation*}
$$

where $\alpha, \beta$ and $\gamma$ are constants.
From equations (23), (24) and (25), we note that $f$ and $g$ represent a rigid body displacements and therefore can be neglected in the analysis of deformation.
Setting $f=g=0$ in (23) and (24), we write

$$
\begin{align*}
& 2 \mu \mathrm{u}_{1}=-\phi_{, 1}+\frac{2(\lambda+2 \mu)}{\lambda+\mu} p_{1},  \tag{26}\\
& 2 \mu \mathrm{u}_{2}=-\phi_{, 2}+\frac{2(\lambda+2 \mu)}{\lambda+\mu} \mathrm{p}_{2} . \tag{27}
\end{align*}
$$

This implies, in compact form,

$$
\begin{align*}
2 \mu\left(\mathrm{u}_{1}+\mathrm{iu}_{2}\right) & =-\left(\phi_{, 1}+\mathrm{i} \phi_{, 2}\right)+\frac{2(\lambda+2 \mu)}{\lambda+\mu}\left(\mathrm{p}_{1}+\mathrm{i} \mathrm{p}_{2}\right) \\
& =-\left[\mathrm{G}(\mathrm{z})+\mathrm{z} \overline{\mathrm{G}^{\prime}(\mathrm{z})}+\overline{\mathrm{H}^{\prime}(\mathrm{z})}\right]+\frac{2(\lambda+2 \mu)}{\lambda+\mu} \mathrm{G}(\mathrm{z}) \\
& =\left[\frac{2(\lambda+2 \mu)}{\lambda+\mu}-1\right] \mathrm{G}(\mathrm{z})-\mathrm{z} \overline{\mathrm{G}^{\prime}(\mathrm{z})}-\overline{\mathrm{H}^{\prime}(\mathrm{z})} \\
& =\mathrm{k}_{0} \mathrm{G}(\mathrm{z})-\mathrm{z} \overline{\mathrm{G}^{\prime}(\mathrm{z})}-\overline{\mathrm{H}^{\prime}(\mathrm{z})}, \tag{28}
\end{align*}
$$

where

$$
\begin{equation*}
\mathrm{k}_{0}=\frac{\lambda+3 \mu}{\lambda+\mu}=3-4 \sigma . \tag{29}
\end{equation*}
$$

The quantity $\mathrm{k}_{0}$ is called the Muskhelishvile's constant.
The formulas given by (13), (14) and (28) are called Kolosov-Muskhelishvilli formulas. This result corresponds to the state of plane strain.
Remark :- In the generalized plane-stress problem, $\lambda$ must be replaced by $\bar{\lambda}=$ $\frac{2 \lambda \mu}{\lambda+2 \mu}$ and if $\mathrm{k}_{0}$ is the corresponding value of in (29). We find

$$
\mathrm{k}_{0}=\frac{\bar{\lambda}+3 \mu}{\bar{\lambda}+\mu}=\frac{5 \lambda+6 \mu}{3 \lambda+2 \mu}=\frac{3-\sigma}{1+\sigma} .
$$

We note that both k and $\mathrm{k}_{0}$ are greater than 1 .

## First and Second Boundary Values Problem of Plane Elasticity

We first consider the first boundary-value problem in which the stress component $\tau_{\alpha \beta}$ must be such that

$$
\begin{equation*}
\tau_{\alpha \beta} v_{\beta}=\mathrm{T}_{\alpha}(\mathrm{s}), \alpha, \beta=1,2, \tag{1}
\end{equation*}
$$

where the stress vector $T_{\alpha}(s)$ is specified on the boundary.
In terms of Airy stress function $\phi=\phi\left(\mathrm{x}_{1}, \mathrm{x}_{2}\right)$, the condition (1) is equivalent to

$$
\begin{equation*}
\phi_{, 1}(\mathrm{~s})=-\int \mathrm{T}_{2}(\mathrm{~s}) \mathrm{ds}, \tag{2}
\end{equation*}
$$

$$
\phi_{, 2}(\mathrm{~s})=\int \mathrm{T}_{1}(\mathrm{~s}) \mathrm{ds},
$$

on C. Now, we write

$$
\begin{equation*}
\phi_{, 1}+\mathrm{i} \phi_{, 2}=\mathrm{i} \int\left[\mathrm{~T}_{1}(\mathrm{~s})+\mathrm{iT}_{2}(\mathrm{~s})\right] \text { ds on } \mathrm{C} . \tag{3}
\end{equation*}
$$

we know that

$$
\begin{equation*}
\phi_{, 1}+\mathrm{i} \phi_{, 2}=\mathrm{G}(\mathrm{z})+\mathrm{z} \overline{\mathrm{G}^{\prime}(\mathrm{z})}+\overline{\mathrm{H}^{\prime}(\mathrm{z})}, \tag{4}
\end{equation*}
$$

where G and H are analytic functions. Thus, the boundary condition in terms of G and H is

$$
\begin{equation*}
\mathrm{G}(\mathrm{z})+\mathrm{z} \overline{\mathrm{G}^{\prime}(\mathrm{z})}+\overline{\mathrm{H}^{\prime}(\mathrm{z})}=\mathrm{i} \int\left[\mathrm{~T}_{1}(\mathrm{~s})+\mathrm{i}_{2}(\mathrm{~s})\right] \text { ds on } \mathrm{C} . \tag{5}
\end{equation*}
$$

The determination of the corresponding boundary conditions in the second boundary-value problem is as follows.

In this type of boundary value problem, boundary conditions are

$$
\begin{equation*}
\mathrm{u}_{\alpha}=\mathrm{g}_{\alpha}(\mathrm{s}), \text { on } \mathrm{C}, \tag{6}
\end{equation*}
$$

where the functions $\mathrm{g}_{\alpha}(\mathrm{s})$ are know functions. Equation (6) yields

$$
2 \mu\left(\mathrm{u}_{1}+\mathrm{i} \mathrm{u}_{2}\right)=2 \mu\left[\mathrm{~g}_{1}(\mathrm{~s})+\mathrm{ig}_{2}(\mathrm{~s})\right], \text { on } \mathrm{C},
$$

This implies
$\Rightarrow \quad \mathrm{k}_{0} \mathrm{G}(\mathrm{z})-\mathrm{z} \overline{\mathrm{G}^{\prime}(\mathrm{z})}-\overline{\mathrm{H}^{\prime}(\mathrm{z})}=2 \mu\left[\mathrm{~g}_{1}(\mathrm{~s})+\mathrm{i} \mathrm{g}_{2}(\mathrm{~s})\right]$, on C .

## Chapter-7 Torsion of Bars

### 7.1 TORSION OF A CIRCULAR SHAFT

Let us consider an elastic right circular beam of length $l$. We choose the $z$-axis along the axis of the beam so that its ends lie in the planes $\mathrm{z}=0$ and $\mathrm{z}=1$, respectively. The end $\mathrm{z}=0$ is fixed in the xy -plane and a couple of vector moment $\overrightarrow{\mathrm{M}}=\mathrm{Me} \hat{e}_{3}$ about the z -axis is applied at the end $\mathrm{z}=l$. The lateral surface of the circular beam is stress-free and body forces are neglected.


The problem is to compute the displacements, strains and stresses developed in the beam because of the twist (or torsion) it experiences due to the applied couple.

Because of the symmetry of the cross-section of the beam by planes normal to the z-axis, these sections will remain planes even after deformation.

That is, if ( $\mathbf{u}, \mathbf{v}, \mathbf{w}$ ) are the displacements, then

$$
\begin{equation*}
\mathrm{w}=0, \tag{1}
\end{equation*}
$$

However, these plane sections will get rotated about the z -axis through some angle $\theta$. The rotation $\theta$ will depend upon the distance of the cross-section from the fixed end $\mathrm{z}=0$ of the beam. For small deformations, we assume that

$$
\begin{equation*}
\theta=\alpha \mathbf{z} \tag{2}
\end{equation*}
$$

where $\alpha$ is a constant.

The constant $\alpha$ represents the twist per unit length ( for $\mathrm{z}=1$ ).
Now, we consider a cross-section of the right circular beam. Let $P(x, y)$ be a point on it before deformations and $P^{\prime}\left(x^{\prime}, y^{\prime}\right)$ be the same material point after deformation.

(Rotation of a section of the circular shaft)
Then

$$
\begin{align*}
u & =x^{\prime}-x \\
& =r[\cos (\theta+\beta)-r \cos \beta] \\
& =r[\cos \theta \cos \beta-\sin \theta \sin \beta]-r \cos \beta \\
& =r \cos \beta(\cos \theta-1)-y \sin \theta \\
& =x(\cos \theta-1)-y \sin \theta . \tag{3}
\end{align*}
$$

Since $\theta$ is small, so

$$
\cos \theta \approx 1, \sin \theta \approx \theta
$$

Therefore

$$
\begin{equation*}
\mathbf{u}=-\mathbf{y} \theta=-\alpha \mathbf{y z} . \tag{4}
\end{equation*}
$$

Similarly

$$
\begin{align*}
v & =y^{\prime}-y \\
& =r \sin (\beta+\theta)-r \sin \beta \\
& =x \sin \theta+\mathbf{y}(\cos \theta-\mathbf{1}) . \tag{5}
\end{align*}
$$

For small deformations, $\theta$ is small, so

$$
\begin{equation*}
\mathbf{v}=\mathbf{x} \theta=\alpha \mathbf{x} \mathbf{z} \tag{6}
\end{equation*}
$$

Thus, the displacement components at any point ( $x, y, z$ ) of the beam due to twisting are

$$
\begin{equation*}
\mathbf{u}=-\alpha \mathbf{y z}, \quad \mathbf{v}=\alpha \mathbf{x z}, \quad \mathbf{w}=\mathbf{0}, \tag{7}
\end{equation*}
$$

where $\alpha$ is the twist per unit length. Therefore, the displacement vector $\bar{u}$ is

$$
\begin{equation*}
\overline{\mathrm{u}}=-\alpha z(\mathrm{y} \hat{\mathrm{i}}-\mathrm{x} \hat{\mathrm{j}}) . \tag{8}
\end{equation*}
$$

Then

$$
\begin{equation*}
\overline{\mathrm{u}} \cdot \overline{\mathrm{r}}=-\alpha \mathrm{z}(y \hat{\mathrm{i}}-x \hat{\mathrm{j}}) \cdot(x \hat{\mathrm{i}}+y \hat{\mathrm{j}})=0, \tag{9}
\end{equation*}
$$

and

$$
\begin{equation*}
|\overline{\mathrm{u}}|=\alpha \mathrm{z} \sqrt{\mathrm{x}^{2}+\mathrm{y}^{2}}=\alpha \mathrm{zr} \tag{10}
\end{equation*}
$$

in polar coordinates ( $\mathbf{r}, \theta$ ).

That is, the displacement vector is in the tangential direction and is of magnitude $\alpha$ rz.

The corresponding strains are

$$
\begin{align*}
& \mathrm{e}_{\mathrm{xx}}=0, \mathrm{e}_{\mathrm{yy}}=0, \mathrm{e}_{\mathrm{zz}}=0, \\
& \mathrm{e}_{\mathrm{xy}}=0, \mathrm{e}_{\mathrm{yz}}=\frac{1}{2} \alpha \mathrm{x}, \mathrm{e}_{\mathrm{zx}}=-\frac{1}{2} \alpha \mathrm{y} . \tag{11}
\end{align*}
$$

The stress-strain relations

$$
\tau_{\mathrm{ij}}=\lambda \delta_{\mathrm{ij}} \mathrm{e}_{\mathrm{kk}}+2 \mu \mathrm{e}_{\mathrm{ij}},
$$

yield

$$
\begin{align*}
& \tau_{\mathrm{xx}}=\tau_{\mathrm{yy}}=\tau_{\mathrm{zz}}=\tau_{\mathrm{xy}}=0, \\
& \tau_{\mathrm{yz}}=\mu \alpha \mathrm{x}, \quad \tau_{\mathrm{xz}}=-\mu \alpha \mathrm{y} . \tag{12}
\end{align*}
$$

This system of stresses clearly satisfies the equations of equilibrium (for zero body)

$$
\tau_{\mathrm{ij}, \mathrm{j}}=0 .
$$

Also the Beltrami-Michell compatibility conditions for zero body force

$$
\begin{equation*}
\nabla^{2} \tau_{\mathrm{ij}}+\frac{1}{1+\sigma} \tau_{\mathrm{kk}, \mathrm{ij}}=0 \tag{13}
\end{equation*}
$$

are obviously satisfied.

Since the lateral surface is stress-free, so the boundary conditions to be satisfied on the lateral surface are

$$
\begin{equation*}
\tau_{\mathrm{ij}} v_{\mathrm{j}=0}, \quad \text { for } \mathrm{i}=1,2,3 . \tag{14}
\end{equation*}
$$

Since the normal is perpendicular to z -axis, so

$$
\begin{equation*}
v_{3}=0, \tag{15}
\end{equation*}
$$

on the lateral surface of the beam. The first two conditions are identically satisfied. The third condition becomes

$$
\begin{equation*}
\tau_{\mathrm{xz}} v_{\mathrm{x}}+\tau_{\mathrm{yz}} v_{\mathrm{y}}=0 \tag{16}
\end{equation*}
$$

Let ' $a$ ' be the radius of the cross-section. Let ( $x, y$ ) be a point on the boundary C of the cross-section. Then

$$
x^{2}+y^{2}=a^{2}
$$

and

$$
\begin{aligned}
& v_{x}=\cos (\hat{v}, x)=x / a \\
& v_{y}=\cos (\hat{v}, y)=y / a
\end{aligned}
$$

Therefore, on the boundary C,

$$
\begin{equation*}
\tau_{\mathrm{xz}} v_{\mathrm{x}}+\tau_{\mathrm{yz}} v_{\mathrm{y}}=-\mu \alpha \mathrm{y} .(\mathrm{x} / \mathrm{a})+\mu \alpha \mathrm{x} .(\mathrm{y} / \mathrm{a})=0 \tag{17}
\end{equation*}
$$

That is, the lateral surface of the circular beam is stress-free.
On the base $\mathrm{z}=l$, let

$$
\overline{\mathrm{F}}=\left(\mathrm{F}_{\mathrm{x}}, \mathrm{~F}_{\mathrm{y}}, \mathrm{~F}_{\mathrm{z}}\right),
$$

be the resultant force. Then

$$
\begin{align*}
\mathrm{F}_{\mathrm{x}} & =\iint_{\mathrm{R}} \tau_{\mathrm{zx}} d x d y \\
& =-\mu \alpha \iint_{\mathrm{R}} \mathrm{y} d x d y \\
& =\mathbf{0} \tag{18}
\end{align*}
$$

since the $y$ coordinate of the $\mathbf{C}$. G. is zero. Also,

$$
\begin{equation*}
\mathrm{F}_{\mathrm{y}}=\iint_{\mathrm{R}} \tau_{\mathrm{zy}} \mathrm{dx} d y=\alpha \mathrm{x} \iint_{\mathrm{R}} \mathrm{x} d x \mathrm{dy}=0 \tag{19}
\end{equation*}
$$

$$
\begin{equation*}
\mathrm{Fz}=\iint_{\mathrm{R}} \tau_{\mathrm{zz}} \mathrm{dx} \mathrm{dy}=0 \tag{20}
\end{equation*}
$$

because the $x$ coordinate of the $C$. $G$. is zero as the $C$. $G$. lies on the z-axis.
Thus, on the base $z=l$, the resultant force $\overline{\mathbf{F}}$ is zero.
Let $\bar{M}=\left(\mathrm{M}_{\mathrm{x}}, \mathrm{M}_{\mathrm{y}}, \mathrm{M}_{\mathrm{z}}\right)$ be the resultant couple on the base $\mathrm{z}=l$. Then

$$
\begin{equation*}
\bar{M}=\iint_{\mathrm{R}}\left(\mathrm{x} \hat{\mathrm{e}}_{1}+\mathrm{y} \hat{\mathrm{e}}_{2}+l \hat{\mathrm{e}}_{3}\right) \times\left(\tau_{\mathrm{xz}} \hat{\mathrm{e}}_{1}+\tau_{\mathrm{zy}} \hat{\mathrm{e}}_{2}+\tau_{\mathrm{zz}} \hat{\mathrm{e}}_{3}\right) \mathrm{dxdy} \tag{21}
\end{equation*}
$$

This gives

$$
\begin{aligned}
\mathrm{M}_{\mathrm{x}} & =\iint_{\mathrm{R}}\left(\mathrm{y} \tau_{\mathrm{zz}}-l \tau_{\mathrm{zy}}\right) \mathrm{dxdy}=\mu \alpha l \iint_{\mathrm{R}} \mathrm{yx} \mathrm{dxdy}=0 \\
\mathrm{M}_{\mathrm{y}} & =\iint_{\mathrm{R}}\left(l \tau_{\mathrm{xz}}-\mathrm{x} \tau_{\mathrm{zz}}\right) \mathrm{dxdy}=-\mu \alpha l \iint_{\mathrm{R}} \mathrm{y} \mathrm{dxdy}=0 \\
\mathrm{M}_{\mathrm{z}} & =\iiint_{\mathrm{R}}\left(\mathrm{x} \tau_{\mathrm{zy}}-\mathrm{y} \tau_{\mathrm{zx}}\right) \mathrm{dxdy} \\
& =\mu \mathrm{x} \iint_{\mathrm{R}}\left(\mathrm{x}^{2}+\mathrm{y}^{2}\right) \mathrm{dxdy}
\end{aligned}
$$

$=\mu \alpha x$ [moment of inertia of the cross-section $\mathrm{z}=l$ about the z -axis]

$$
\begin{equation*}
=\mu \alpha \frac{\pi \mathrm{a}^{4}}{2} \tag{24}
\end{equation*}
$$

As $\overline{\mathrm{M}}=\mathrm{M} \hat{\mathrm{e}}_{3}$, so we write

$$
\begin{equation*}
\mathrm{M}=\mathrm{M}_{\mathrm{z}}=\frac{\pi \mu \alpha \mathrm{a}^{4}}{2} \tag{25}
\end{equation*}
$$

where $a$ is the radius of the circular cross-section at $\mathrm{z}=\ell$.
This gives

$$
\begin{equation*}
\alpha=\frac{2 \mathrm{M}}{\mu \pi \mathrm{a}^{4}}, \tag{26}
\end{equation*}
$$

which determines the constant $\alpha$ in term of moment of the applied couple M , radius of the cross-section ' $a$ ' and rigidity $\mu$ of the medium of the beam.

The constant $\alpha$ is the twist per unit length. With $\alpha$ given by (26), the displacements, strains and stresses at any point of the beam due to applied twist became completely known by equations (7), (11), (12).

## Stress -Vector

The stress vector at any point $\mathrm{P}(\mathrm{x}, \mathrm{y})$ in any cross-section ( $\mathrm{z}=$ constant $)$ is given by

$$
\begin{equation*}
\stackrel{\mathrm{z}}{\mathrm{~T}}=\tau_{\mathrm{zx}} \hat{\mathrm{i}}+\tau_{\mathrm{zy}} \hat{\mathrm{j}}+\tau_{\mathrm{zz}} \hat{\mathrm{k}} . \tag{27}
\end{equation*}
$$

Using (12), we find

$$
\begin{align*}
\underset{\sim}{\mathrm{T}} & =\tau_{\mathrm{zx}} \hat{\mathrm{i}}+\tau_{\mathrm{zy}} \hat{\mathrm{j}} \\
& =\mu \alpha(-y \hat{\mathrm{i}}+x \hat{\mathrm{j}}), \tag{28}
\end{align*}
$$

which lies in the cross-section itself, i.e., the stress-vector is tangential.
Moreover, we note that the stress-vector $\underset{\sim}{\mathrm{T}}$ is perpendicular to the radius vector $r=x \hat{i}+y \hat{j}$ as

$$
\stackrel{\mathrm{z}}{\mathrm{~T}} \cdot \mathrm{r}=0 .
$$

The magnitude $\tau$ of the shier vector $\underset{\sim}{T}$ is given by

$$
\begin{equation*}
\tau=\sqrt{\tau_{\mathrm{zx}}^{2}+\tau_{\mathrm{zy}}^{2}}=\mu \alpha \sqrt{\mathrm{x}^{2}+\mathrm{y}^{2}}=\mu \alpha, \tag{29}
\end{equation*}
$$

which is maximum when $r=a$, and

$$
\begin{equation*}
\tau_{\max }=\mu \alpha \mathrm{a}=\frac{2 \mathrm{M}}{\pi \mathrm{a}^{3}} . \tag{30}
\end{equation*}
$$

Note:- The torsinal rigidity of the beam, denoted by $\mathrm{D}_{0}$, is defined by

$$
\begin{equation*}
\mathrm{D}_{0}=\frac{\mathrm{M}}{\alpha}=\mu \cdot \frac{\pi \mathrm{a}^{4}}{2} . \tag{31}
\end{equation*}
$$

The constant $D_{0}$ provides a measure of the rigidity of the beam. It depends on the modulus of rigidity $\mu$ and the shape of the cross-section of the beam only. The constant $D_{0}$ (which is equal to $\frac{M}{\alpha}$ ) represents the moment of the couple required to produce a unit angle of twist per unit of length.

It is also called the torsinal stiffness of the beam.

Example: Consider a circular shaft of length $l$, radius a, and shear modulus $\mu$, twisted by a couple $M$. Show that the greatest angle of twist $\theta$ and the maximum shear stress

$$
\tau=\sqrt{\tau_{\mathrm{zx}}^{2}+\tau_{\mathrm{zy}}^{2}}
$$

are given by

$$
\theta_{\max }=\frac{2 \mathrm{M} l}{\pi \cdot \mu \mathrm{a}^{4}}, \quad \quad \tau_{\max }=\frac{2 \mathrm{M}}{\pi \mathrm{a}^{3}}
$$

Solution: We know that

$$
\begin{equation*}
\theta=\alpha \mathrm{z}, \quad \mathrm{M}=\frac{\mu \alpha \pi \mathrm{a}^{4}}{2}, \quad \mathrm{~T}=\mu \alpha \mathrm{r} \tag{1}
\end{equation*}
$$

We find

$$
\begin{equation*}
\alpha=\frac{2 \mathrm{M}}{\mu \pi \mathrm{a}^{4}} \tag{2}
\end{equation*}
$$

Now

$$
\theta_{\max }=\alpha \cdot(\mathrm{z})_{\max }=\alpha l=\frac{2 \mathrm{M}}{\mu \pi \mathrm{a}^{4}},
$$

and

$$
\tau_{\max }=\mu \alpha(\mathrm{r})_{\max }=\mu \alpha \mathrm{a}=\frac{2 \mathrm{M}}{\pi \mathrm{a}^{3}} .
$$

### 7.2. TORSION OF BARS OF ARBITRARY CROSS-SECTION

Consider an elastic beam of length $l$ of uniform but arbitrary cross-section. The lateral surface is stress free and body forces are absent. Suppose that a couple of vector moment $\overrightarrow{\mathrm{M}}$ about the axis of the beam is applied at one end and the other end is fixed.

The problem is to compute the displacements, strains and stresses developed in the beam.

We choose z -axis along the axis of the beam. Let the end $\mathrm{z}=0$ be fixed and the end $\mathrm{z}=\ell$ be applied the vector couple of moment $\overrightarrow{\mathrm{M}}$.

(Torsion of a beam of arbitrary cross-section)

Let ( $u, v, w$ ) be the components of the displacement. In the case of a circular cross-section, a plane cross-section remains plane even after the deformation. However, when the cross-section is arbitrary, a plane cross-section will not remain plane after deformation, it gets warped (curved surface).

This phenomenon is known as warping.
We assume that each section is warped in the same way, i.e., warping is independent of $z$. We write

$$
\begin{equation*}
\mathbf{w}=\alpha \phi(\mathbf{x}, \mathbf{y}) \tag{1}
\end{equation*}
$$

which is the same for all sections. Here $\alpha$ is the angle of twist per unit length of the beam.

The function $\phi=\phi(x, y)$ is called the Saint-Venant's warping function or torsion function.

Due to twisting, a plane section get rotated about the axis of the beam. The angle $\theta$ of rotation depends upon the distance of the section from the fixed end. For small deformations, we assume that

$$
\begin{equation*}
\theta=\alpha z \tag{2}
\end{equation*}
$$

Let $\mathbf{P}(x, y)$ be any point on a cross section before deformations $\mathbf{P}^{\prime}\left(x^{\prime}, y^{\prime}\right)$ be the same material point after deformation.


Then

$$
\begin{align*}
u=x^{\prime}-x & =r \cos (\theta+\beta)-r \cos \beta \\
& =r \cos \theta \cos \beta-r \sin \theta \sin \beta-r \cos \beta \\
& =x(\cos \theta-1)-y \sin \theta \\
& =-y \theta . \tag{3}
\end{align*}
$$

Similarly

$$
\begin{equation*}
\mathrm{v}=\mathrm{x} \theta . \tag{4}
\end{equation*}
$$

Hence, for small deformations, the displacement components are given by the relations

$$
\begin{align*}
& \mathbf{u}=-\alpha \mathbf{y}, \\
& \mathbf{v}=\alpha \mathbf{x z}, \\
& \mathbf{w}=\alpha \phi(\mathbf{x}, \mathbf{y}) . \tag{5}
\end{align*}
$$

The strain components are

$$
\begin{align*}
& \mathrm{e}_{\mathrm{xx}}=\frac{\partial \mathrm{u}}{\partial \mathrm{x}}=0, \\
& \mathrm{e}_{\mathrm{yy}}=\frac{\partial \mathrm{v}}{\partial \mathrm{y}}=0, \\
& \mathrm{e}_{\mathrm{zz}}=\frac{\partial \mathrm{w}}{\partial \mathrm{z}}=0, \tag{6}
\end{align*}
$$

and

$$
\begin{align*}
& \quad \mathrm{e}_{\mathrm{xy}}=\frac{1}{2}\left(\frac{\partial \mathrm{u}}{\partial \mathrm{y}}+\frac{\partial \mathrm{v}}{\partial \mathrm{x}}\right)=\frac{1}{2} \leftarrow \alpha \mathrm{z}+\alpha \mathrm{z}=0, \\
& \mathrm{e}_{\mathrm{yz}}=\frac{1}{2} \alpha\left[\mathrm{x}+\frac{\partial \phi}{\partial \mathrm{y}}\right] \\
& \mathrm{e}_{\mathrm{zx}}=\frac{1}{2} \alpha\left[-\mathrm{y}+\frac{\partial \phi}{\partial \mathrm{x}}\right] . \tag{7}
\end{align*}
$$

The stresses are to be found by using the Hooke's law

$$
\begin{equation*}
\tau_{\mathrm{ij}}=\lambda \delta_{\mathrm{ij}} \mathrm{e}_{\mathrm{kk}}+2 \mu \mathrm{e}_{\mathrm{ij}} . \tag{8}
\end{equation*}
$$

We find

$$
\begin{align*}
& \tau_{\mathrm{xx}}=\tau_{\mathrm{yy}}=\tau_{\mathrm{zz}}=\tau_{\mathrm{xy}}=0,  \tag{9}\\
& \tau_{\mathrm{yz}}=\mu \alpha\left(\mathrm{x}+\frac{\partial \phi}{\partial \mathrm{y}}\right),  \tag{...}\\
& \tau_{\mathrm{zx}}=\mu \alpha\left(-\mathrm{y}+\frac{\partial \phi}{\partial \mathrm{x}}\right) . \tag{11}
\end{align*}
$$

These stresses must satisfy the following equilibrium equations for zero body force

$$
\begin{equation*}
\tau_{\mathrm{ij}, \mathrm{j}}=0, \tag{12}
\end{equation*}
$$

for $\mathrm{i}=1,2,3$ in R. First two equations are identically satisfied. The third equation gives

$$
\begin{array}{ll}
\tau_{\mathrm{zx}, \mathrm{x}}+\tau_{\mathrm{zy}, \mathrm{y}}=0 & \text { in } \mathrm{R} \\
\text { or } & \mu \alpha \frac{\partial}{\partial \mathrm{x}}\left[-\mathrm{y}+\frac{\partial \phi}{\partial \mathrm{x}}\right]+\mu \alpha \frac{\partial}{\partial \mathrm{y}}\left[\mathrm{x}+\frac{\partial \phi}{\partial \mathrm{y}}\right]=0 \\
\text { or } & \frac{\partial^{2} \phi}{\partial \mathrm{x}^{2}}+\frac{\partial^{2} \phi}{\partial \mathrm{y}^{2}}=0, \tag{13}
\end{array}
$$

## in the region $\mathbf{R}$ of the cross-section.

This shows that the torsion function $\phi$ is a harmonic function.
Let C be the boundary curve of the region R representing the cross-section of the beam. Since the lateral surface is stress free, so the boundary conditions to be satisfied are

$$
\begin{equation*}
\tau_{\mathrm{ij}} v_{\mathrm{j}}=0, \quad \text { on } \mathrm{C} \tag{14}
\end{equation*}
$$

for $\mathrm{i}=1,2,3$. The first two conditions are identically satisfied as $v_{3}=0$. The third condition is

$$
\tau_{\mathrm{zx}} v_{\mathrm{x}}+\tau_{\mathrm{zy}} v_{\mathrm{y}}=0 \quad \text { on } \mathrm{C}
$$

or

$$
\left(-\mathrm{y}+\frac{\partial \phi}{\partial \mathrm{x}}\right) v_{\mathrm{x}}+\left(\mathrm{x}+\frac{\partial \phi}{\partial \mathrm{y}}\right) v_{\mathrm{y}}=0 \quad \text { on } \mathrm{C}
$$

or

$$
\frac{\partial \phi}{\partial \mathrm{x}} v_{\mathrm{x}}+\frac{\partial \phi}{\partial \mathrm{y}} v_{\mathrm{y}}=\mathrm{y} v_{\mathrm{x}}-\mathrm{x} v_{\mathrm{y}} \quad \text { on } \mathrm{C}
$$

or

$$
\begin{equation*}
\frac{\mathrm{d} \phi}{\mathrm{dv}}=\mathrm{y} v_{\mathrm{x}}-\mathrm{x} v_{\mathrm{y}} \quad \text { on } \mathrm{C} . \tag{15}
\end{equation*}
$$

since

$$
\begin{align*}
\frac{\mathrm{d} \phi}{\mathrm{~d} v} & =\text { normal derivative of } \phi \\
& =\nabla \phi \cdot \hat{v} \\
& =\frac{\partial \phi}{\partial \mathrm{x}} \mathrm{v}_{\mathrm{x}}+\frac{\partial \phi}{\partial \mathrm{y}} v_{\mathrm{y}} . \tag{16}
\end{align*}
$$

Thus, the torsion function/ warping functions $\phi$ must be a solution of the two dimensional following Neumann boundary value problem.
$\nabla^{2} \phi=0 \quad, \quad$ in $\mathbf{R}$,
$\frac{\mathbf{d} \phi}{\mathbf{d v}}=$ normal derivations of $\phi=\left(\mathbf{y} \cdot \mathbf{v}_{\mathbf{x}}-\mathbf{x} \cdot \mathbf{v}_{\mathbf{y}}\right) \quad$ on $\mathbf{C}$.
So, we solve the torsion problem as a Neummann problem (17).
On the base $\mathrm{z}=l$, let $\overline{\mathrm{F}}$ be the resultant force. Then

$$
\begin{align*}
\mathrm{F}_{\mathrm{x}} & =\iint_{\mathrm{R}} \tau_{\mathrm{zx}} \mathrm{dx} d y=\mu \alpha \iint_{\mathrm{R}}\left(\frac{\partial \phi}{\partial \mathrm{x}}-\mathrm{y}\right) \mathrm{dx} d \mathrm{y} \\
& =\mu \alpha \iint_{\mathrm{R}} \int\left[\frac{\partial}{\partial \mathrm{x}}\left\{\mathrm{x}\left(\frac{\partial \phi}{\partial \mathrm{x}}-\mathrm{y}\right)\right\}+\frac{\partial}{\partial \mathrm{y}}\left\{\mathrm{x}\left(\frac{\partial \phi}{\partial \mathrm{y}}+\mathrm{x}\right)\right\}\right] \mathrm{dx} d \mathrm{y} \\
& =\mu \alpha \int_{\mathrm{C}}\left[-\mathrm{x}\left(\frac{\partial \phi}{\partial \mathrm{y}}+\mathrm{x}\right) \mathrm{dx}+\mathrm{x}\left(\frac{\partial \phi}{\partial \mathrm{x}}-\mathrm{y}\right) \mathrm{dy}\right], \tag{18}
\end{align*}
$$

using Green's theorem

$$
\begin{equation*}
\mathcal{L}^{P} \boldsymbol{P d x}+\mathrm{Qdy}=\iint_{\mathrm{R}}\left(\frac{\partial \mathrm{Q}}{\partial \mathrm{x}}-\frac{\partial \mathrm{P}}{\partial \mathrm{y}}\right) \mathrm{dx} \mathrm{dy}, \tag{19}
\end{equation*}
$$

which converts surface integral into a line integral.
In case of two-dimensional curve C , directions cosines of the tangent are $\left.<\frac{\mathrm{dx}}{\mathrm{ds}}, \frac{\mathrm{dy}}{\mathrm{ds}}\right\rangle$, and therefore, d. c.'s of the normal are $\left.<+\frac{\mathrm{dy}}{\mathrm{ds}},-\frac{\mathrm{dx}}{\mathrm{ds}}\right\rangle$, i. e.,

$$
\begin{equation*}
v_{\mathrm{x}}=\mathrm{dy} / \mathrm{ds}, \quad v_{\mathrm{y}}=-\mathrm{dx} / \mathrm{ds} . \tag{20}
\end{equation*}
$$

From equations (18) and (20), we write

$$
\begin{align*}
\mathrm{F}_{\mathrm{x}} & =\mu \alpha \int_{\mathrm{x}} \mathrm{x}\left[\left(-\mathrm{x}-\frac{\partial \phi}{\partial \mathrm{y}}\right) \frac{\mathrm{dx}}{\mathrm{ds}}+\left(\frac{\partial \phi}{\partial \mathrm{x}}-\mathrm{y}\right) \frac{\mathrm{dy}}{\mathrm{ds}}\right] \mathrm{ds} \\
& =\mu \alpha \int_{\mathrm{x}} \mathrm{x}\left[\frac{\partial \phi}{\partial \mathrm{x}} \cdot v_{\mathrm{x}}+\frac{\partial \phi}{\partial \mathrm{y}} v_{\mathrm{y}}+\boldsymbol{v} v_{\mathrm{y}}-\mathrm{y} v_{\mathrm{x}}\right] \mathrm{ds} \\
& =\mu \alpha \int_{\mathrm{x}}\left[\frac{\partial \phi}{\partial v}+v_{\mathrm{y}}-\mathrm{y} v_{\mathrm{x}}-\mathrm{ds}\right. \\
& =\mathbf{0} \tag{21}
\end{align*}
$$

since, on the boundary $C$, the integrand is identically zero.
Similarly

$$
\mathrm{F}_{\mathrm{y}}=\mathrm{F}_{\mathrm{z}}=0 .
$$

Thus, the resultant force $\overline{\mathrm{F}}$ on the cross-section R of the beam vanishes.
On the base $\mathrm{z}=l$, let $\overline{\mathrm{M}}$ be the resultant couple. Then

$$
\begin{align*}
\mathrm{M}_{\mathrm{x}} & =\iint_{\mathrm{R}} \mathrm{y} \tau_{\mathrm{zz}}-\mathrm{z} \tau_{\mathrm{zy}} \underline{-} \mathrm{d} \mathrm{x} \mathrm{dy} \\
& =-l \iint_{\mathrm{R}} \tau_{\mathrm{zy}} \mathrm{dxdy} ; \quad(\because \mathrm{z}=l) \\
& =-l \mathrm{~F}_{\mathrm{y}} \\
& =\mathbf{0} \tag{23}
\end{align*}
$$

Similarly

$$
\begin{equation*}
\mathrm{M}_{\mathrm{y}}=0 \tag{24}
\end{equation*}
$$

Now
$M_{z}=\iint_{R} \tau_{z y}-y \tau_{z x} \underline{\bar{d}} x d y=\mu \alpha \iint_{R}\left[x^{2}+y^{2}+x \frac{\partial \phi}{\partial y}-y \frac{\partial \phi}{\partial x}\right] d x d y$,
This gives

$$
\begin{equation*}
M=\mu \alpha \iint_{R}\left[x^{2}+y^{2}+x \frac{\partial \phi}{\partial y}+y \frac{\partial \phi}{\partial z}\right] d x d y \tag{25}
\end{equation*}
$$

as it is given that $M$ is the moments of the torsion couple about $z$-axis.

We write

$$
\begin{equation*}
\mathbf{M}=\alpha \mathbf{D} \tag{26}
\end{equation*}
$$

where M is applied moment, $\alpha$ is the twist per unit length and D is the torsinal rigidity given by

$$
\begin{equation*}
\mathrm{D}=\mu \iint_{\mathrm{R}}\left[\mathrm{x}^{2}+\mathrm{y}^{2}+\mathrm{x} \frac{\partial \phi}{\partial \mathrm{y}}-\mathrm{y} \frac{\partial \phi}{\partial \mathrm{x}}\right] \mathrm{dx} d \mathrm{y} \tag{27}
\end{equation*}
$$

which depends upon $\mu$ (i. e., the material of the beam) and the shape of the cross-section $R$ (i. e., the geometry of the beam).

For given $\mathrm{M}, \alpha$ can be determined from equation (26).
However, when $\alpha$ is given, then we can calculate the required moment M , from equation (27), to produce the twist $\alpha$ per unit length.

After finding $\phi$ (by solving the Neumann boundary value problem) the torsion function $\phi$ becomes know and consequently the torsinal rigidity D becomes known.

### 7.3. DIRICHLET BOUNDARY VALUE PROBLEM

Let $\psi(\mathrm{x}, \mathrm{y})$ be the conjugate harmonic function of the harmonic function $\phi(x, y)$. Then

$$
\begin{equation*}
\frac{\partial \phi}{\partial x}=\frac{\partial \psi}{\partial y}, \frac{\partial \phi}{\partial y}=-\frac{\partial \psi}{\partial x}, \tag{28}
\end{equation*}
$$

and the function $\phi+\mathrm{i} \psi$ is an analytic function of the complex variable $\mathrm{x}+\mathrm{iy}$.
Now

$$
\begin{aligned}
\frac{\mathrm{d} \phi}{\mathrm{~d} \nu} & =\frac{\partial \phi}{\partial \mathrm{x}} v_{\mathrm{x}}+\frac{\partial \phi}{\partial \mathrm{y}} v_{\mathrm{y}} \\
& =\frac{\partial \phi}{\partial \mathrm{x}} \frac{\mathrm{dy}}{\mathrm{ds}}-\frac{\partial \phi}{\partial \mathrm{y}} \frac{\mathrm{dx}}{\mathrm{ds}}
\end{aligned}
$$

$$
\begin{align*}
& =\frac{\partial \psi}{\partial \mathrm{y}} \frac{\mathrm{dy}}{\mathrm{ds}}+\frac{\partial \psi}{\partial \mathrm{x}} \frac{\mathrm{dx}}{\mathrm{ds}} \\
& =\frac{\mathrm{d} \psi}{\mathrm{ds}} \tag{29}
\end{align*}
$$

On the boundary C, the boundary condition (15) now becomes

$$
\frac{\mathrm{d} \psi}{\mathrm{ds}}=\mathrm{y} v_{\mathrm{x}}-\mathrm{x} v_{\mathrm{y}} \quad \text { on } \mathrm{C}
$$

or

$$
\frac{\mathrm{d} \psi}{\mathrm{ds}}=\mathrm{y} \frac{\mathrm{dy}}{\mathrm{ds}}+\mathrm{x} \frac{\mathrm{dx}}{\mathrm{ds}} \quad \text { on } \mathrm{C}
$$

or

$$
\begin{equation*}
\psi=\frac{1}{2} \mathbf{t}^{2}+y^{2}+\text { constt, } \quad \text { on } C . \tag{30}
\end{equation*}
$$

Thus, the determination of the function $\psi=\psi(\mathbf{x}, \mathbf{y})$ is the problem of solving the Dirichlet problem

$$
\begin{align*}
& \frac{\partial^{2} \psi}{\partial \mathbf{x}^{2}}+\frac{\partial^{2} \psi}{\partial \mathbf{y}^{2}}=0 \quad \text { in } \mathbf{R} \\
& \psi=\frac{1}{2} \mathbf{1}^{2}+\mathbf{y}^{2}+\text { constt. on } \mathbf{C} \tag{31}
\end{align*}
$$

whose solution is unique.

Once $\psi$ becomes known, the torsion function $\phi=\phi(\mathrm{x}, \mathrm{y})$ can be obtained from relations in (28).

We generally, take, constt. $=0$ in (31).
The Dirichlet problem of potential theory can be solved by standard techniques.

### 7.4. STRESS- FUNCTION

Stress-Function: The stress function, denoted by $\Psi$, is defined as

$$
\begin{equation*}
\Psi(\mathrm{x}, \mathrm{y})=\psi(\mathrm{x}, \mathrm{y})-\frac{1}{2} \mathrm{t}^{2}+\mathrm{y}^{2} \tag{32}
\end{equation*}
$$

where $\psi(x, y)$ is the solution of the Dirichlet problem (31).

We know that $\psi(\mathrm{x}, \mathrm{y})$ is also the conjugate function of the harmonic torsion function $\phi(\mathrm{x}, \mathrm{y})$.
The stress $\Psi(x, y)$ was introduced by Prandtl.
From equation (32), we find

$$
\frac{\partial \Psi}{\partial \mathrm{x}}=\frac{\partial \Psi}{\partial \mathrm{x}}-\mathrm{x}
$$

or

$$
\begin{equation*}
\frac{\partial \Psi}{\partial \mathrm{x}}=\frac{\partial \Psi}{\partial \mathrm{x}}+\mathrm{x} \tag{33}
\end{equation*}
$$

and

$$
\frac{\partial \Psi}{\partial y}=\frac{\partial \psi}{\partial y}-y
$$

or

$$
\begin{equation*}
\frac{\partial \psi}{\partial y}=\frac{\partial \Psi}{\partial y}+y . \tag{..}
\end{equation*}
$$

Further, we find

$$
\begin{align*}
\nabla^{2} \Psi & =\frac{\partial^{2} \Psi}{\partial \mathrm{x}^{2}}+\frac{\partial^{2} \Psi}{\partial \mathrm{y}^{2}} \\
& =\left(\frac{\partial^{2} \Psi}{\partial \mathrm{x}^{2}}-1\right)+\left(\frac{\partial^{2} \psi}{\partial \mathrm{y}^{2}}-1\right) \\
& =-\mathbf{2} \quad \text { in } \mathbf{R} . \tag{35}
\end{align*}
$$

From equations (31) and (32), we find

$$
\begin{equation*}
\Psi=\text { constt on } \mathrm{C} . \tag{..}
\end{equation*}
$$

The differential equation in (35) is called Poisson's equation.

Thus, the stress function $\Psi(\mathrm{x}, \mathrm{y})$ is a solution of the boundary value problem consists of equations (35) and (36).

The shear stresses $\tau_{\mathrm{zx}}$ and $\tau_{\mathrm{zy}}$, given in (11), can also be expressed in terms of stress function $\Psi$. We find

$$
\tau_{z x}=\mu \alpha\left(\frac{\partial \phi}{\partial x}-y\right)
$$

$$
\begin{align*}
& =\mu \alpha\left(\frac{\partial \Psi}{\partial y}-y\right) \\
& =\mu \alpha \frac{\partial \Psi}{\partial y}  \tag{37}\\
\tau_{z y} & =\mu \alpha\left(\frac{\partial \phi}{\partial y}+x\right) \\
& =\mu \alpha\left(-\frac{\partial \psi}{\partial x}+x\right) \\
& =-\mu \alpha \frac{\partial \Psi}{\partial x} . \tag{38}
\end{align*}
$$

The torsional rigidity D in term of stress function $\Psi(\mathrm{x}, \mathrm{y})$ can be found. We know that

$$
\begin{align*}
& D=\frac{1}{\alpha} \iint_{\mathrm{R}} \mathrm{x} \tau_{\mathrm{zy}}-\mathrm{y} \tau_{\mathrm{zx}} \underline{-\bar{d} x} \mathrm{dy} \\
& =\frac{1}{\alpha} \iint_{\mathrm{R}}\left[\mathrm{x}\left\{-\mu \alpha \frac{\partial \Psi}{\partial \mathrm{x}}\right\}-\mathrm{y}\left\{\mu \alpha \frac{\partial \Psi}{\partial \mathrm{y}}\right\}\right] \mathrm{dx} \mathrm{dy} \\
& =-\mu \iint_{\mathrm{R}}\left[\mathrm{x} \frac{\partial \Psi}{\partial \mathrm{x}}+\mathrm{y} \frac{\partial \Psi}{\partial \mathrm{y}}\right] \mathrm{dx} \mathrm{dy} \\
& =-\mu \iint_{\mathrm{R}}\left[\frac{\partial}{\partial \mathrm{x}} \times \Psi+\frac{\partial}{\partial \mathrm{y}} \mathbf{y} \Psi\right] \mathrm{dx} \mathrm{~d} y+2 \mu \iint_{\mathrm{R}} \Psi \mathrm{~d} x \mathrm{dy} \\
& =-\mu \int_{\mathrm{C}} \Psi \mathrm{x} v_{\mathrm{x}}+\mathrm{y} v_{\mathrm{y}}^{\underline{\mathrm{d}} \mathrm{-}}+2 \mu \int_{\mathrm{R}} \int_{\mathrm{C}} \Psi \mathrm{dx} \mathrm{dy}, \tag{39}
\end{align*}
$$

using the Green theorem for plane.
On the boundary, C ,

$$
\Psi=\text { constt }
$$

We choose

$$
\begin{equation*}
\Psi=0 \text { on } \mathbf{C} \tag{40}
\end{equation*}
$$

Then

$$
\mathrm{D}=2 \mu \int_{\mathrm{R}} \int_{\mathrm{D}} \Psi \mathrm{dxdy}
$$

which is the expression for torsinal rigidity $D$ in terms of stress function $\Psi$

### 7.5. LINES OF SHEARING SHEER

Consider the family of curves in the xy-plane given by

$$
\begin{equation*}
\Psi=\text { constt. } \tag{42}
\end{equation*}
$$

Then

$$
\frac{\partial \Psi}{\partial x}+\frac{\partial \Psi}{\partial y} \frac{d y}{d x}=0
$$

This gives

$$
\begin{align*}
& \frac{d y}{d x}=-\frac{\Psi x}{\Psi y} . \\
& =\frac{\mu \alpha \Psi x}{\mu \alpha \Psi y} \\
& =\frac{\tau_{z y}}{\tau_{z x}} . \tag{..}
\end{align*}
$$

This relation shows that at each of the curve, $\Psi=$ constt, the stress vector

$$
\begin{equation*}
\bar{\tau}=\tau_{\mathrm{zx}} \hat{\mathrm{i}}+\tau_{\mathrm{zy}} \hat{\mathrm{j}} \tag{44}
\end{equation*}
$$


is defined along the tangent to the curve $\Psi=$ constt. at that point.

The curves

$$
\Psi=\text { constt }
$$

are called lines of shearing stress.

### 7.6. SPECIAL CASES OF BEAMS : TORSION OF AN ELLIPTIC BEAM

Let the boundary C of the cross-section be

$$
\begin{equation*}
\frac{x^{2}}{\mathrm{a}^{2}}+\frac{\mathrm{y}^{2}}{\mathrm{~b}^{2}}=1 \tag{A1}
\end{equation*}
$$

We assume that the solution of the Dirichlet boundary values problem (31) is of the type

$$
\begin{equation*}
\psi(x, y)=c^{2}\left(x^{2}-y^{2}\right)+k^{2} \tag{A2}
\end{equation*}
$$

for constants $\mathrm{c}^{2}$ and $\mathrm{k}^{2}$.
It is obvious that

$$
\begin{equation*}
\nabla^{2} \psi=0 \tag{A3}
\end{equation*}
$$

At each point ( $\mathrm{x}, \mathrm{y}$ ) on the boundary C , we must have

$$
c^{2}\left(x^{2}-y^{2}\right)+k^{2}=\frac{1}{2}\left(x^{2}+y^{2}\right)
$$

or

$$
\begin{equation*}
x^{2}\left(\frac{1}{2}-c^{2}\right)+y^{2}\left(\frac{1}{2}+c^{2}\right)=k^{2} \tag{A4}
\end{equation*}
$$

which becomes the ellipse (A1) if

$$
c^{2}<\frac{1}{2}
$$

and

$$
\mathrm{a}=\frac{\mathrm{k}}{\sqrt{\frac{1}{2}-\mathrm{c}^{2}}}, \quad \mathrm{~b}=\frac{\mathrm{k}}{\sqrt{\frac{1}{2}+\mathrm{c}^{2}}},
$$

or

$$
\mathrm{a}^{2}=\frac{\mathrm{k}^{2}}{2-\mathrm{c}^{2}}, \quad \mathrm{~b}^{2}=\frac{\mathrm{k}^{2}}{\frac{1}{2}+\mathrm{c}^{2}}
$$

or

$$
\begin{equation*}
c^{2}=\frac{1}{2}\left(\frac{a^{2}-b^{2}}{a^{2}+b^{2}}\right), \quad k^{2}=\frac{a^{2} b^{2}}{a^{2}+b^{2}} . \tag{A5}
\end{equation*}
$$

Therefore, solution of the Dirichlet problem (31) for this particular type of elliptic beam is

$$
\begin{equation*}
\psi(x, y)=\frac{\mathrm{a}^{2}-\mathrm{b}^{2}}{2 \mathrm{a}^{2}+\mathrm{b}^{2}}-\left(\mathrm{x}^{2}-\mathrm{y}^{2}\right)+\frac{\mathrm{a}^{2} \mathrm{~b}^{2}}{\mathrm{a}^{2}+\mathrm{b}^{2}} . \tag{A6}
\end{equation*}
$$

The torsion function $\phi$ is given by the formula

$$
\begin{align*}
\phi(x, y) & =\int\left(\frac{\partial \phi}{\partial x} d x+\frac{\partial \phi}{\partial y} d y\right) \\
& =\int\left[\frac{\partial \psi}{\partial y} d x-\frac{\partial \psi}{\partial x} d y\right] \\
& =-\frac{a^{2}-b^{2}}{a^{2}+b^{2}} \int y d x+x d y \\
& =-\frac{a^{2}-b^{2}}{a^{2}+b^{2}} x y . \tag{A7}
\end{align*}
$$

The stress function $\Psi(\mathrm{x}, \mathrm{y})$ is given by the formula

$$
\begin{align*}
\Psi & =\psi(\mathrm{x}, \mathrm{y})-\frac{1}{2}\left(\mathrm{x}^{2}+\mathrm{y}^{2}\right) \\
& =\frac{-\mathrm{a}^{2} \mathrm{~b}^{2}}{\mathrm{a}^{2}+\mathrm{b}^{2}}\left(\frac{\mathrm{x}^{2}}{\mathrm{a}^{2}}+\frac{\mathrm{y}^{2}}{\mathrm{~b}^{2}}-1\right) \tag{A8}
\end{align*}
$$

using (A6).
To calculate the torsional rigidity $\mathrm{D}_{\mathrm{e}}$ and twist $\alpha$.
The non-zero shear stresses are

$$
\begin{align*}
\tau_{z x} & =\mu \alpha\left(\frac{\partial \psi}{\partial y}-y\right) \\
& =\mu \alpha\left[-\frac{a^{2}-b^{2}}{a^{2}+b^{2}} y-y\right] \\
& =\frac{-2 \mu \alpha a^{2} y}{a^{2}+b^{2}} \tag{A9}
\end{align*}
$$

and

$$
\begin{align*}
\tau_{z y} & =\mu \alpha\left(-\frac{\partial \psi}{\partial x}+x\right) \\
& =\mu \alpha\left[-\frac{a^{2}-b^{2}}{a^{2}+b^{2}} \mathrm{x}+\mathrm{x}\right] \\
& =\frac{2 \mu \alpha \mathrm{~b}^{2} \mathrm{x}}{\mathrm{a}^{2}+\mathrm{b}^{2}} \tag{A10}
\end{align*}
$$

Let M be the torsion moment of the couple about z -axis. Then

$$
\begin{aligned}
\mathrm{M} & =\iint_{\mathrm{R}} \mathrm{x} \tau_{\mathrm{zy}}-\mathrm{y} \tau_{\mathrm{zx}} \underset{-}{\bar{d} x d y} \\
& =\frac{2 \mu \alpha}{a^{2}+\mathrm{b}^{2}}\left[\mathrm{~b}^{2} \iint_{\mathrm{R}} \int^{2} \mathrm{x}^{2} \mathrm{dxdy}+\mathrm{a}^{2} \iint_{\mathrm{R}} \mathrm{y}^{2} \mathrm{dxdy}\right]
\end{aligned}
$$

This implies

$$
\begin{equation*}
\mathrm{M}=\frac{2 \mu \alpha}{\mathrm{a}^{2}+\mathrm{b}^{2}} \mathrm{~b}^{2} \mathrm{I}_{\mathrm{y}}+\mathrm{a}^{2} \mathrm{I}_{\mathrm{x}} \tag{A11}
\end{equation*}
$$

where $I_{x}$ and $I_{y}$ are the moments of inertia of the elliptic cross-section about $x-$ and $y$-axis, respectively.

We know that

$$
\begin{equation*}
\mathrm{I}_{\mathrm{x}}=\frac{\pi \mathrm{ab}^{3}}{4}, \quad \mathrm{I}_{\mathrm{y}}=\frac{\pi \mathrm{a}^{3} \mathrm{~b}}{4} \tag{A12}
\end{equation*}
$$

Putting these values in (A11), we find

$$
\begin{equation*}
\mathrm{M}=\frac{\pi \mu \mathrm{a}^{3} \mathrm{~b}^{3}}{\mathrm{a}^{2}+\mathrm{b}^{2}} \tag{A13}
\end{equation*}
$$

The torsional rigidity $D_{e}$ is given by the formula

$$
\begin{equation*}
\mathrm{M}=\mathrm{D}_{\mathrm{e}} \alpha . \tag{A14}
\end{equation*}
$$

Hence, we find

$$
\begin{align*}
& D_{e}=\frac{\pi \mu a^{3} b^{3}}{a^{2}+b^{2}}  \tag{A15}\\
& \alpha=\frac{M\left(a^{2}+b^{2}\right)}{\pi \mu a^{3} b^{3}} \tag{A16}
\end{align*}
$$

Equations (A9) , (A10) and (A16) imply

$$
\begin{align*}
& \tau_{z \mathrm{x}}=\frac{-2 \mathrm{My}}{\pi \mathrm{ab}^{3}}  \tag{A17}\\
& \tau_{\mathrm{zy}}=\frac{-2 \mathrm{Mx}}{\pi \mathrm{a}^{3} \mathrm{~b}} \tag{A18}
\end{align*}
$$

Lines of shearing stress: The line of shearing stress are given by the formula

$$
\begin{equation*}
\Psi=\text { constt. } \tag{A19}
\end{equation*}
$$

Equations (A8) and (A19) imply

$$
\frac{x^{2}}{a^{2}}+\frac{y^{2}}{b^{2}}=\text { constt }
$$

as the lines the shearing stress.
This is a family of concentric ellipses, similar to the given ellipse $\frac{x^{2}}{a^{2}}+\frac{y^{2}}{b^{2}}=1$.


Displacements. The displacements at any point of the given elliptic beam are now given by the formulas

$$
\begin{equation*}
u=-\alpha y z, \quad v=\alpha x z, \quad w=\alpha \phi=-\left(\frac{a^{2}-b^{2}}{a^{2}+b^{2}}\right) \alpha x y \tag{A21}
\end{equation*}
$$

with the twist per unit length, $\alpha$, as given by equation (A16).
Maximum shear stress: The shear stress $\tau$ is given by

$$
\begin{align*}
& \tau=\sqrt{\tau_{\mathrm{zx}}^{2}+\tau_{\mathrm{zy}}^{2}} . \\
& =\frac{2 \mu \alpha}{\mathrm{a}^{2}+\mathrm{b}^{2}} \sqrt{\mathrm{a}^{4} \mathrm{y}^{2}+\mathrm{b}^{4} \mathrm{x}^{2}} . \tag{A22}
\end{align*}
$$

We know that the maximum shear stress occurs on the boundary C

$$
\mathrm{C}: \quad \frac{\mathrm{x}^{2}}{\mathrm{a}^{2}}+\frac{\mathrm{y}^{2}}{\mathrm{~b}^{2}}=1 .
$$

So, on the boundary C , the shear stress becomes

$$
\begin{align*}
\tau & =\frac{2 \mu \alpha}{a^{2}+b^{2}} \sqrt{a^{4} b^{2}\left(1-\frac{x^{2}}{a^{2}}\right)+b^{4} x^{2}} \\
& =\frac{2 \mu \alpha a b}{a^{2}+b^{2}} \sqrt{a^{2}-x^{2}+b^{2} x^{2} / a^{2}} \\
& =\frac{2 \mu \alpha b}{a^{2}+b^{2}} \sqrt{a^{2}-x^{2}\left(1-b^{2} / a^{2}\right)} \\
(\tau) & =\frac{2 \mu \alpha a b}{a^{2}+b^{2}} \sqrt{a^{2}-e^{2} x^{2}}, \tag{A23}
\end{align*}
$$

where

$$
\begin{equation*}
\mathrm{e}^{2}=1-\mathrm{b}^{2} / \mathrm{a}^{2} \tag{A24}
\end{equation*}
$$

is the eccentricity of the ellipse.

We note that the maximum value of shearing stress (A23) occurs when $x=$ 0.

Thus,

$$
\begin{equation*}
\tau_{\max }=\frac{2 \mu \alpha \mathrm{a}^{2} \mathrm{~b}}{\mathrm{a}^{2}+\mathrm{b}^{2}}=\frac{2 \mathrm{M}}{\pi \mathrm{ab}^{2}} \quad \text { at the points }(0, \pm \mathrm{b}) \tag{A25}
\end{equation*}
$$

Thus, the maximum shear stress occurs at the extremities of the minor axis

of the ellipse and maximum shear stress is given by (A25).
Warping curves: The warping curves are given by the relation
w = constt.
or
xy = constt.,
using (A21).
These are rectangular hyperbolas. We know that

$$
\mathrm{w}=-\frac{\mathrm{a}^{2}-\mathrm{b}^{2}}{\mathrm{a}^{2}+\mathrm{b}^{2}} \alpha \mathrm{xy} .
$$

In the first quadrant, $x$ and $y$ are both positive so $\mathbf{w}<0$.

In the third quadrant $x<0, y<0$, so $w<0$.
Therefore, in I and III quadrants, the curves becomes concave. In figure, -ve sign represents concave curves.


In the II and IV quadrants, the curves become convex which are represented by + sign in the figure.

Remark 1. Let A be the area of the ellipse and I be the moment of inertia about z -axis. Then

$$
\mathrm{A}=\pi \mathrm{ab},
$$

$$
\begin{align*}
& \mathrm{I}=\mathrm{I}_{\mathrm{x}}+\mathrm{I}_{\mathrm{y}}=\frac{\pi \mathrm{ab}}{4}\left(\mathrm{a}^{2}+\mathrm{b}^{2}\right) \\
& \mathrm{D}_{\mathrm{e}}=\frac{\mu \mathrm{A}^{4}}{4 \pi^{2} \mathrm{I}} \tag{A27}
\end{align*}
$$

Remark 2. The results for a circular bar can be derived as a particular case of the above-a bar with elliptic cross section, on taking $b=a$. We find

$$
\begin{aligned}
& \phi(x, y)=0 \\
& w(x, y)=0 \\
& \Psi(x, y)=0
\end{aligned}
$$

### 7.7. TORSION OF BEAMS WITH TRIANGULAR CROSS-SECTION

Consider a cylinder of length $l$ whose cross section is a triangular prism. Let zaxis lie along the central line of the of the beam and one end of the beam lying in the plane $\mathrm{z}=0$ is fixed at the origin and the other end lie in the plane $\mathrm{z}=l$. A couple of moment M is applied at the centriod $(0,0, l)$ of end.

## We shall determine the resulting deformation.

Let

$$
\begin{equation*}
\phi+\mathrm{i} \psi=\mathrm{ic}(\mathrm{x}+\mathrm{iy})^{3}+\mathrm{ik}, \tag{B1}
\end{equation*}
$$

where c and k are constants. We find

$$
\begin{align*}
& \phi=c\left(-3 x^{2} y+y^{3}\right)  \tag{B2}\\
& \psi=c\left(x^{3}-3 x y^{2}\right)+k \tag{B3}
\end{align*}
$$

We shall be solving the Dirichlet problem in $\psi$ :
(i) $\nabla^{2} \psi=0 \quad$ in R ,
(ii) $\psi=\frac{1}{2}\left(\mathrm{x}^{2}+\mathrm{y}^{2}\right) \quad$ on C ,
where C is the boundary of a triangular cross section occupying the region R .
We note that $\psi$, given by (B3), is a harmonic function in the two-dimensional region R .
On the boundary C, we must have

$$
\begin{equation*}
c\left(x^{3}-3 x y^{2}\right)+k=\frac{1}{2}\left(x^{2}+y^{2}\right) \tag{B6}
\end{equation*}
$$

By altering the value of constant c and k , we obtain various cross-sections of the beam.

In particularly, if we set

$$
\begin{equation*}
\mathrm{C}=-\frac{1}{6 \mathrm{a}} \quad \text { and } \quad \mathrm{k}=\frac{2}{3} \mathrm{a}^{2} \tag{B7}
\end{equation*}
$$

in equation (B6), we have

$$
-\frac{1}{6 a}\left(x^{3}-3 x y^{2}\right)+\frac{2}{3} a^{2}=\frac{1}{2}\left(x^{2}+y^{2}\right)
$$

or

$$
x^{3}-3 x y^{2}+3 a x^{2}+3 a y^{2}-4 a^{3}=0
$$

or

$$
\begin{equation*}
(x-a)(x-y \sqrt{3}+2 a)(x+y \sqrt{3}+2 a)=0 . \tag{B8}
\end{equation*}
$$

This shows that the boundary C consists of an equilateral triangle ABC

formed by the st. lines

$$
\begin{align*}
& x=a \\
& x+\sqrt{3} y+2 a=0 \\
& x-\sqrt{3} y+2 a=0 \tag{B9}
\end{align*}
$$

The altitude of this cross-section is 3 a and each length of the triangle is $2 \sqrt{3}$ a. The centroid of this cross-section is origin , lying on the central line of the beam.

The unique solution of the Dirichlet problem $(\mathrm{B} 4,5)$ is

$$
\begin{equation*}
\psi=-\frac{1}{6 a}\left(x^{3}-3 \mathrm{xy}^{2}\right)+\frac{2}{3} \mathrm{a}^{2} . \tag{B10}
\end{equation*}
$$

The torsion function becomes

$$
\begin{equation*}
\phi=-\frac{1}{6 a}\left(y^{3}-3 x^{2} y\right) \tag{B11}
\end{equation*}
$$

The stress function $\Psi$ is given by

$$
\begin{align*}
\Psi & =\psi(x, y)-\frac{1}{2}\left(x^{2}+y^{2}\right) \\
& =-\frac{1}{6 a}\left[x^{3}-3 x y^{2}+3 x^{2} a+3 y^{2}-4 a^{3}\right] \tag{B12}
\end{align*}
$$

The non-zero shear stresses are

$$
\begin{align*}
\tau_{\mathrm{zx}} & =\mu \alpha\left(\frac{\partial \psi}{\partial \mathrm{y}}-\mathrm{y}\right) \\
& =\frac{\mu \alpha}{\mathrm{a}}(\mathrm{x}-\mathrm{a}) \mathrm{y}  \tag{B13}\\
\tau_{\mathrm{zy}} & =\mu \alpha\left(-\frac{\partial \psi}{\partial \mathrm{x}}+\mathrm{x}\right) \\
& =\frac{\mu \alpha}{2 \mathrm{a}}\left(\mathrm{x}^{2}+2 \mathrm{ax}-\mathrm{y}^{2}\right) . \tag{B14}
\end{align*}
$$

The displacement at any points of the triangular beam are given by

$$
\begin{align*}
& u=-\alpha y z \\
& v=\alpha x z \\
& w=\alpha \phi=\frac{\alpha}{6 a}\left(3 x^{2} y-y^{3}\right)
\end{align*}
$$

where $\alpha$ is the twist per unit length.
We know that the maximum value of the shearing sheers $\tau$ occurs at the boundary C .
On the boundary $\mathrm{x}=\mathrm{a}$, we find

$$
\begin{equation*}
\tau_{\mathrm{zx}}=0, \quad \tau=\tau_{\mathrm{zy}}=\frac{\mu \alpha}{2 \mathrm{a}}\left(3 \mathrm{a}^{2}-\mathrm{y}^{2}\right) \tag{B16}
\end{equation*}
$$

which is maximum where $y=0$. Thus

$$
\begin{equation*}
\tau_{\max }=\frac{3}{2} \mu \alpha \mathrm{a} \tag{B17}
\end{equation*}
$$

at the point $\mathrm{D}(\mathrm{a}, 0)$, which is the middle point of BC .

Also, $\tau=0$ at the points (corner pts B \& C) where $\mathrm{y}= \pm \sqrt{3} \mathrm{a}$.
At the point $\mathrm{A}(-2 \mathrm{a}, 0), \tau$ is zero. Thus, the stress $\tau$ is zero at the corner points $\mathrm{A}, \mathrm{B}, \mathrm{C}$. We note that at the centroid $O(0,0)$, the shear stress is also zero.

Similarly, we may check that the shearing stress $\tau$ is maximum at the middle points of the sides AC and AB , and the corresponding maximum shear stress is each equal to $\frac{3}{2} \mu \alpha \mathrm{a}$.

Torsional rigidity: We know that M is the moment of the applied couple along z-axis. Therefore

$$
\begin{align*}
M & =\iint_{R} t \tau_{z y}-y \tau_{z x} \frac{\bar{d} x d y}{} \\
& =\frac{\mu \alpha}{2 a} \iint_{R} t^{3}+2 a x^{2}-x y^{2}-2 y^{2} x+2 y^{2} a \bar{d} x d y \\
& =2 \cdot \frac{\mu \alpha}{2 a} \int_{x=-2 a}^{a} \int_{y=0}^{y=\frac{x+2 a}{\sqrt{3}}} x^{3}+2 a x^{2}-x y^{2}-2 y^{2} x+2 a y^{2}-\bar{d} y d x \tag{B18}
\end{align*}
$$



This gives

$$
\mathrm{M}=\frac{9 \sqrt{3}}{5} \mu \alpha \mathrm{a}^{4}
$$

Consequently, we obtain

$$
\begin{align*}
& \mathrm{D}=\frac{9 \sqrt{3}}{5} \mu \mathrm{a}^{4},  \tag{B20}\\
& \alpha=\frac{5 \mathrm{M}}{9 \sqrt{3} \mu \mathrm{a}^{4}}, \tag{B21}
\end{align*}
$$

with $\mathrm{a}=\frac{1}{3} \mathrm{rd}$ of the altitude of the equilateral triangle, each side being equal to $2 \sqrt{3} a$.
Equation (B21) determines the constant $\alpha$ when the moment M and the crosssection are known. Equation (B13), (B14), and (B21) yield

$$
\begin{align*}
& \tau_{\mathrm{zx}}=\frac{5 \mathrm{M}}{9 \sqrt{3 a^{5}}} y(x-a)  \tag{B22}\\
& \tau_{\mathrm{zy}}=\frac{5 \mathrm{M}}{18 \sqrt{3} \mathrm{a}^{5}}\left(\mathrm{x}^{2}+2 a \mathrm{ax}-\mathrm{y}^{2}\right)
\end{align*}
$$

Equations (B17) and (B21) imply

$$
\tau_{\max }=\frac{15 \mathrm{M}}{18 \sqrt{3} \mathrm{a}^{4}}
$$

at the point $\mathrm{D}(\mathrm{a}, 0)$.
Theorem. Show that the points at which the shearing stress is maximum lie on the boundary C of the cross-section of the beam.

Proof. To prove this theorem, we use the following theorem from analysis.
"Let a function $f(x, y)$ be continuous and has continuous partial derivatives of the first and second orders and not identically equal to a constant and satisfy the in equality $\nabla^{2} f \geq 0$ in the region $R$. Then $f(x, y)$ attains its maximum value on the boundary $C$ of the region $R$ ".

We knows that the shear stress $\tau$ is given by

$$
\begin{equation*}
\tau^{2}=\mu^{2} \alpha^{2}\left(\Psi_{\mathrm{x}}^{2}+\Psi_{\mathrm{y}}^{2}\right) \tag{1}
\end{equation*}
$$

where $\Psi$ is the stress function. Now

$$
\begin{align*}
\frac{\partial}{\partial \mathrm{x}} \tau^{2} & =\mu^{2} \alpha^{2} 2 \mid \Psi_{\mathrm{x}} \Psi_{\mathrm{xx}}+2 \Psi_{\mathrm{y}} \Psi_{\mathrm{yx}} \\
& =2 \mu^{2} \alpha^{2}\left[\Psi_{\mathrm{x}} \Psi_{\mathrm{xx}+} \Psi_{\mathrm{y}} \Psi_{\mathrm{yy}}\right] \tag{2}
\end{align*}
$$

and

$$
\frac{\partial^{2}}{\partial \mathrm{x}^{2}} \tau^{2}=2 \mu^{2} \alpha^{2} \Psi_{\mathrm{zz}}^{2}+\Psi_{\mathrm{x}} \Psi_{\mathrm{xxx}}+\Psi_{\mathrm{yx}}^{2}+\Psi_{\mathrm{y}} \Psi_{\mathrm{yxx}}
$$

Similarly

$$
\begin{equation*}
\frac{\partial^{2}}{\partial \mathrm{y}^{2}} \tau^{2}=2 \mu^{2} \alpha^{2} \Psi_{\mathrm{yy}}^{2}+\Psi_{\mathrm{y}} \Psi_{\mathrm{yyy}}+\Psi_{\mathrm{xy}}^{2}+\Psi_{\mathrm{x}} \Psi_{\mathrm{xyy}} \tag{4}
\end{equation*}
$$

Adding (3) and (4),

$$
\begin{align*}
\nabla^{2} \tau^{2} & =2 \mu^{2} \alpha^{2}\left[\Psi_{\mathrm{xx}}^{2}+\Psi_{\mathrm{yy}}^{2}+\Psi_{\mathrm{x}} \frac{\partial}{\partial \mathrm{x}} \Psi_{\mathrm{xx}}+\Psi_{\mathrm{yy}}+\Psi_{\mathrm{y}} \frac{\partial}{\partial \mathrm{y}} \Psi_{\mathrm{xx}}+\Psi_{\mathrm{yy}}+2 \Psi_{\mathrm{xy}}^{2}\right] \\
& =2 \mu^{2} \propto^{2} \Psi_{\mathrm{xx}}^{2}+\Psi_{\mathrm{yy}}^{2}+2 \Psi_{\mathrm{xy}}^{2}  \tag{5}\\
& \geq 0 \quad \text { in R. }
\end{align*}
$$

Therefore, by above result, $\tau^{2}$ (and hence stress $\tau$ ) attains the maximum value on the boundary C of the region R .
Question. Let $D_{0}$ be the torsional rigidity of a circular cylinder, $D_{e}$ that of an elliptic cylinder, and $D_{t}$ that of a beam whose cross-section is an equilateral triangle. Show that for cross-section of equal areas

$$
\mathrm{D}_{\mathrm{e}}=\mathrm{kD}_{0}, \quad \mathrm{D}_{\mathrm{t}}=\frac{2 \pi \sqrt{3}}{15} \mathrm{D}_{0}, \text { where } \mathrm{k}=\frac{2 \mathrm{ab}}{\mathrm{a}^{2}+\mathrm{b}^{2}} \leq 1
$$

and $\mathrm{a}, \mathrm{b}$ are the semi-axis of the elliptical section.
Solution. We know that for a circular cylinder of radius r ,

$$
\begin{equation*}
\mathrm{D}_{0}=\frac{\pi}{2} \mu \mathrm{r}^{4} \tag{1}
\end{equation*}
$$

We know that for an elliptic cylinder $\frac{x^{2}}{a^{2}}+\frac{y^{2}}{b^{2}}=1$,

$$
\begin{equation*}
\mathrm{D}_{\mathrm{e}}=\frac{\pi \mu \mathrm{a}^{3} \mathrm{~b}^{3}}{\mathrm{a}^{2}+\mathrm{b}^{2}} \tag{2}
\end{equation*}
$$

We know that for an equilateral triangle (with each side of length $2 \sqrt{3 x}$ )

$$
\begin{equation*}
D_{t}=\frac{9 \sqrt{3}}{5} \mu x^{4} \tag{3}
\end{equation*}
$$

Since the areas of all the cross-section are equal, so

$$
\begin{equation*}
\pi r^{2}=\pi \mathrm{ab}=3 \sqrt{3} \mathrm{x}^{2} . \tag{4}
\end{equation*}
$$



Now

$$
\frac{\mathrm{D}_{\mathrm{e}}}{\mathrm{D}_{0}}=\frac{\pi \mu \mathrm{a}^{3} \mathrm{~b}^{3}}{\mathrm{a}^{2}+\mathrm{b}^{2}} \times \frac{2}{\pi \mu \mathrm{r}^{2}}
$$

$$
\begin{align*}
& =\frac{2 a^{3} b^{3}}{r^{4}\left(a^{2}+b^{2}\right)} \\
& =\frac{2 a^{3} b^{3}}{a^{2} b^{2}\left(a^{2}+b^{2}\right)} \\
& =\frac{2 a b}{a^{2}+b^{2}} \\
& =k \tag{5}
\end{align*}
$$

Also

$$
\begin{align*}
\frac{D_{t}}{D_{0}} & =\frac{a \sqrt{3} \mu x^{4}}{5} \times \frac{2}{\pi \mu r^{4}} \\
& =\frac{18 \sqrt{3} x^{4}}{5 \pi r^{4}} \\
& =\frac{18 \sqrt{3} x^{4}}{5 \pi\left(27 \mathrm{x}^{4} / \pi\right)} \\
& =\frac{18 \sqrt{3} \pi}{5 \times 27}=\frac{2 \pi \sqrt{3}}{15} \tag{6}
\end{align*}
$$

Hence, the result.

## Chapter-8 <br> Variational Methods

### 8.1. INTRODUCTION

We shall be using the minimum principles in deriving the equilibrium and compatibility equations of elasticity.

### 8.2. DEFLECTION OF AN ELASTIC STRING

Let a stretched string, with the end points fixed at $(0,0)$ and $(\ell, 0)$, be deflected by a distributed transverse load $f(\mathrm{x})$ per unit length of the string. We suppose that the transverse deflection $\mathrm{y}(\mathrm{x})$ is small and the change in the stretch force T produced by deflection is neglible.

These are the usual assumption used in deriving the equation for $\mathrm{y}(\mathrm{x})$ from considerations of static equilibrium.

We shall deduce this equation from the Principle of Minimum Potential Energy.

## We know that the potential energy V is defined by the formula

$$
\begin{equation*}
\mathrm{V}=\mathrm{U}-\int_{0}^{\ell} f(\mathrm{x}) \mathrm{ydx}, \tag{1}
\end{equation*}
$$

where the strain energy $U$ is equal to the product of the tensile force $T$ by the total stretch e of the string. That is

$$
\begin{equation*}
\mathrm{U}=\mathrm{Te} \mathrm{e}, \tag{2}
\end{equation*}
$$

where

$$
\begin{align*}
\mathrm{e} & =\int_{0}^{\ell} \mathbb{d} \mathrm{s}-\mathrm{dx}^{-} \\
& =\int_{0}^{\ell} \sqrt{1+\left(\mathrm{y}^{\prime}\right)^{2}}-1 \mathrm{dx} . \tag{3}
\end{align*}
$$

Since, we are dealing with the linear theory, so

$$
\begin{equation*}
\left(y^{\prime}\right)^{2} \ll 1, \tag{4}
\end{equation*}
$$

and equation (3) can be written as

$$
\begin{equation*}
\mathrm{e}=\frac{1}{2} \int_{0}^{e} \mathrm{~s}^{z} \mathrm{dx} \tag{5}
\end{equation*}
$$

From equations (1), (2) and (5), finally we write

$$
\begin{equation*}
\mathrm{V}=\int_{0}^{e}\left[\frac{1}{2} \mathrm{~T} \mathrm{y}_{-}^{\mathrm{z}}-\mathrm{f}(\mathrm{x}) \mathrm{y}\right] \mathrm{dx} . \tag{6}
\end{equation*}
$$

The appropriate Euler's equation of the functional (6) is (left as an exercise)

$$
\begin{equation*}
\mathrm{T} \frac{\mathrm{~d}^{2} \mathrm{y}}{\mathrm{dx}^{2}}+\mathrm{f}(\mathrm{x})=0 \tag{7}
\end{equation*}
$$

This equation is the familiar / well known equation for the transverse deflection of the string under the load $f(\mathbf{x})$.

### 8.3. DEFLECTION OF THE CENTRAL LINE OF A BEAM

Let the axis of a beam of constant cross-section coincides with the x -axis. Suppose that the beam is bent by a transverse load

$$
\begin{equation*}
\mathrm{p}=f(\mathrm{x}), \tag{1}
\end{equation*}
$$

estimated per unit length of the beam.
As per theory of deformation of beams, we suppose that the shearing stresses are negligible in comparison with the tensile stress

$$
\begin{equation*}
\tau_{\mathrm{xx}}=\frac{\mathrm{My}}{\mathrm{I}}, \tag{2}
\end{equation*}
$$

where $M$ is the magnitude of the moment about the $x$-axis and $I$ is the moment of inertia of the cross-section about $x$-axis.

The strain $\mathrm{e}_{\mathrm{xx}}$ is then given by

$$
\begin{align*}
\mathrm{e}_{\mathrm{xx}} & =\frac{\tau_{\mathrm{xx}}}{\mathrm{E}} \\
& =\frac{\mathrm{My}}{\mathrm{IE}}, \tag{4}
\end{align*}
$$

where $E$ is the Young's modulus.

## The strain-energy function $W$ is given by

$$
\begin{align*}
\mathrm{W} & =\frac{1}{2} \tau_{\mathrm{xx}} \mathrm{e}_{\mathrm{xx}} \\
& =\frac{\mathrm{M}^{2} \mathrm{y}^{2}}{2 \mathrm{EI}^{2}} \tag{5}
\end{align*}
$$

The strain energy per unit length of the beam is found by integrating W over the cross-section of the beam, and we get

$$
\begin{align*}
\int_{\mathrm{R}} \mathrm{Wd} d \sigma & =\frac{\mathrm{M}^{2}}{2 E I^{2}} \int_{\mathrm{R}} \mathrm{y}^{2} \mathrm{~d} \sigma \\
& =\frac{\mathrm{M}^{2}}{2 \mathrm{EI}} \tag{6}
\end{align*}
$$

## The well known Bernoulli- Euler law is

$$
\begin{equation*}
\mathrm{M}=-\mathrm{EI} \frac{\mathrm{~d}^{2} \mathrm{y}}{\mathrm{dx}^{2}} \tag{7}
\end{equation*}
$$

The total strain energy $U$ obtained by integrating the expression (6) over the length of the beam, and using (7), we find

$$
\begin{equation*}
\mathrm{U}=\frac{1}{2} \int_{0}^{\ell} \mathrm{EI}\left(\frac{\mathrm{~d}^{2} \mathrm{y}}{\mathrm{dx}^{2}}\right) \mathrm{dx} . \tag{8}
\end{equation*}
$$

We suppose that the ends of the beam are clamped, hinged, or free, so that the supporting forces do not work and contribute nothing to potential energy V .

If we neglect the weight of the beam, the only external load is $p=f(x)$, then the formula

$$
\begin{equation*}
\mathrm{V}=\mathrm{U}-\int_{0}^{\ell} \mathrm{f}(\mathrm{x}) \mathrm{ydx} \tag{9}
\end{equation*}
$$

for the potential energy gives

$$
\begin{equation*}
\mathrm{V}=\int_{0}^{\ell}\left[\frac{1}{2} \mathrm{EI}\left(\frac{\mathrm{~d}^{2} \mathrm{y}}{\mathrm{dx}^{2}}\right)-\mathrm{f}(\mathrm{x}) \mathrm{y}\right] \mathrm{dx} . \tag{10}
\end{equation*}
$$

The Euler's equation of the above function $V$ is (left as an exercise)

$$
\begin{equation*}
\frac{\mathrm{d}^{2}}{\mathrm{dx}^{2}}\left(\mathrm{EI} \frac{\mathrm{~d}^{2} \mathrm{y}}{\mathrm{dx}^{2}}\right)-f(\mathrm{x})=0 \tag{11}
\end{equation*}
$$

### 8.4. DEFLECTION OF AN ELASTIC MEMBRANE

Let the membrane, with fixed edges, occupy some region in the xy-plane. We suppose that the membrane is stretched so that the tension T is uniform and that T is so great that it is not appreciably changed when the membrane is deflected by a distributed normal load of intensity $f(\mathrm{x}, \mathrm{y})$.

## We first compute the strain energy $U$. The total stretch e of the surface

$$
\begin{equation*}
\mathbf{z}=\mathbf{u}(\mathbf{x}, \mathbf{y}), \tag{1}
\end{equation*}
$$

is

$$
\begin{align*}
\mathrm{e} & =\iint_{\mathrm{R}} \mathrm{~d} \sigma-\mathrm{dxdy} \\
& =\iint_{\mathrm{R}} \sqrt{\mathrm{u}_{\mathrm{x}}^{2}+\mathrm{u}_{\mathrm{y}}^{2}+1}-1-\overline{d x} d y \tag{2}
\end{align*}
$$

where

$$
\begin{equation*}
d \sigma=\sqrt{u_{x}^{2}+u_{y}^{2}+1} d x d y \tag{3}
\end{equation*}
$$

is the element of area of the membrane in the deformed state.
As usual, it is assumed that the displacement $u$ and its first derivatives are small. Then, we can write (2) as

$$
\begin{equation*}
\mathrm{e}=\frac{1}{2} \iint_{\mathrm{R}} \int_{\mathrm{x}}^{2}+\mathrm{u}_{\mathrm{y}}^{2} \mathrm{dx} d \mathrm{y} \tag{4}
\end{equation*}
$$

Hence, the strain energy $U$ is given by

$$
\begin{align*}
\mathbf{U} & =\mathbf{T} \mathbf{e} \\
& =\frac{\mathrm{T}}{2} \iint_{\mathrm{R}} \mathrm{u}_{\mathrm{x}}^{2}+\mathrm{u}_{\mathrm{y}}^{2} \mathrm{dx} \mathrm{dy} . \tag{5}
\end{align*}
$$

We know that the potential energy is given by the formula

$$
\begin{equation*}
\mathrm{V}=\mathrm{U}-\iint_{\mathrm{R}} \int_{\mathrm{f}} \mathrm{f}(\mathrm{x}, \mathrm{y}) \mathrm{udx} \mathrm{dy} . \tag{6}
\end{equation*}
$$

Equations (5) and (6) give the potential energy as

$$
\begin{equation*}
V=\iint_{R}\left\{\frac{T}{2} \mathbf{t}_{x}^{2}+u_{y}^{2}-f(x, y)\right\} u d x d y . \tag{7}
\end{equation*}
$$

The equilibrium state is characterized by the condition

$$
\begin{equation*}
\delta \mathrm{V}=0 . \tag{8}
\end{equation*}
$$

This gives (left as an exercise)

$$
\begin{equation*}
\mathrm{T} \nabla^{2} \mathrm{u}+f(\mathrm{x}, \mathrm{y})=0 . \tag{9}
\end{equation*}
$$

### 8.5. TORSION OF CYLINDERS

We consider the Saint -Venant torsion problem for a cylinder of arbitrary cross-section.

We shall use the Principle of Minimum Complementary Energy to deduce the appropriate Compatibility equation.

## We know that the displacement components in the cross-section

are

$$
\begin{equation*}
u(x, y, z)=-\alpha z y, \quad v(x, y, z)=\alpha z x, \tag{1}
\end{equation*}
$$

where $\alpha$ is the twist per unit length of the cylinder.
We assume with Saint-Venant principle that nonvanishing stresses are $\tau_{\text {zx }}$ and $\tau_{\mathrm{zy}}$.

## The formula for complementary energy is

$$
\begin{equation*}
\mathrm{V}^{*}=\mathrm{U}-\int_{\Sigma^{\mathrm{u}}} \mathrm{~T}_{\mathrm{i}} \mathrm{u}_{\mathrm{i}} \mathrm{~d} \sigma, \tag{2}
\end{equation*}
$$

where the surface integral is evaluated over the ends of the cylinder, and strain energy U is given by

$$
\begin{equation*}
\mathrm{U}=\int_{\tau} \mathrm{W} \mathrm{~d} \tau, \tag{3}
\end{equation*}
$$

where

$$
\begin{align*}
\mathrm{W} & =\frac{1}{2} \tau_{\mathrm{ij}} \mathrm{e}_{\mathrm{ij}} \\
& =\left(\tau_{\mathrm{zx}} \mathrm{e}_{\mathrm{zx}}+\tau_{\mathrm{zy}} \mathrm{e}_{\mathrm{zy}}\right) \tag{4}
\end{align*}
$$

## From shear strain relations, we have

$$
\begin{align*}
& \tau_{\mathrm{zx}}=2 \mu \mathrm{e}_{\mathrm{zx}}, \\
& \tau_{\mathrm{zy}}=2 \mu \mathrm{e}_{\mathrm{zy}} . \tag{5}
\end{align*}
$$

## From equations (4) and (5), we find

$$
\begin{equation*}
\mathrm{W}=\frac{1}{2 \mu}\left(\tau_{z x}^{2}+\tau_{z y}^{2}\right), \tag{6}
\end{equation*}
$$

and hence

$$
\begin{equation*}
\mathrm{U}=\frac{1}{2 \mu} \int_{\tau} \tau_{\mathrm{zx}}^{2}+\tau_{\mathrm{zy}}^{2}{ }_{-}^{-} \mathrm{d} \tau \tag{7}
\end{equation*}
$$

Now, we shall compute the surface integral in (2).
i). For the end $z=0$, we have

$$
u+v=0
$$

so

$$
\begin{equation*}
\int_{\mathrm{R}} \mathrm{~T}_{\mathrm{i}} \cdot \mathrm{u}_{\mathrm{i}} \mathrm{~d} \sigma^{\prime}=0 . \tag{8}
\end{equation*}
$$

ii). On the end $\mathrm{z}=\ell$, we have

$$
\begin{equation*}
\int_{\mathrm{R}} \mathrm{~T}_{\mathrm{i}} \cdot \mathrm{u}_{\mathrm{i}} \mathrm{~d} \sigma=\alpha \iint_{\mathrm{R}} \int \ell \mathrm{y} \tau_{\mathrm{zx}}+\ell \mathrm{x} \tau_{\mathrm{zy}} \overline{\mathrm{dx}} \mathrm{dy} \tag{9}
\end{equation*}
$$

Thus, we find

$$
\begin{equation*}
\mathrm{V}^{*}=\frac{\ell}{2 \mu} \iint_{\mathrm{R}} \boldsymbol{\tau}_{\mathrm{zx}}^{2}+\tau_{\mathrm{zy}}^{2} \overline{\mathrm{~d}} \mathrm{dx} \mathrm{dy}-\alpha \ell \iint_{\mathrm{R}} \mathrm{k} \tau_{\mathrm{zy}}-\mathrm{y} \tau_{\mathrm{zx}} \overline{-} \mathrm{dx} \mathrm{dy} . \tag{10}
\end{equation*}
$$

In this case, the admissible stresses satisfy the equilibrium equation (left as an exercise)

$$
\begin{equation*}
\frac{\partial \tau_{\mathrm{zx}}}{\partial \mathrm{x}}+\frac{\partial \tau_{\mathrm{zy}}}{\partial \mathrm{y}}=0, \quad \text { in } \mathrm{R}, \tag{11}
\end{equation*}
$$

and the boundary condition

$$
\begin{equation*}
\tau_{\mathrm{zx}} \cos (\mathrm{x}, v)+\tau_{\mathrm{zy}} \cos (\mathrm{y}, v)=0, \quad \text { on } \mathrm{C} . \tag{12}
\end{equation*}
$$

Equilibrium equation (11) will clearly be satisfied identically if we introduce the stress function.

$$
\Psi=\Psi(\mathrm{x}, \mathrm{y})
$$

such that

$$
\begin{align*}
& \tau_{\mathrm{zx}}=\mu \alpha \frac{\partial \Psi}{\partial y} \\
& \tau_{\mathrm{zy}}=-\mu \alpha \frac{\partial \Psi}{\partial \mathrm{x}} . \tag{13}
\end{align*}
$$

The boundary condition (12) then requires that

$$
\mu \alpha\left(\frac{\partial \Psi}{\partial y} \frac{d y}{d s}+\frac{\partial \Psi}{\partial x} \frac{d x}{d s}\right)=0
$$

or

$$
\frac{\mathrm{d} \Psi}{\mathrm{ds}}=0
$$

or

$$
\begin{equation*}
\Psi=\text { constant }, \quad \text { on } \mathbf{C} \tag{14}
\end{equation*}
$$

On substituting expressions for stresses from (13) into equation (10), we get

$$
\begin{equation*}
\mathrm{V}^{*}=\left(\frac{\mu \alpha^{2} \ell}{2}\right) \iint_{\mathrm{R}} \int_{\mathrm{x}} \Psi_{-}^{-2}+\Psi \Psi_{\mathrm{y}}^{\overline{2}}+2 \boldsymbol{x _ { \mathrm { x } }}+\mathrm{y} \Psi_{\mathrm{y}}^{\overline{-}} \mathrm{dx} d y \tag{15}
\end{equation*}
$$

The corresponding Euler equation (exercise) is

$$
\begin{equation*}
\nabla^{2} \Psi=-2, \quad \text { in } \mathrm{R}, \tag{16}
\end{equation*}
$$

which is precisely the equation for the Prandtl stress function.
Remark: The formula (15) for $\mathrm{V}^{*}$ can be written in a simple form which we shall find useful in subsequent considerations.

We note that

$$
\begin{equation*}
\iint_{\mathrm{R}} \int_{\mathrm{x}}+\Psi_{\mathrm{y}} \Psi_{\mathrm{y}}^{-\mathrm{dx}} \mathrm{dy}=\iint_{\mathrm{R}}\left[\frac{\partial}{\partial \mathrm{x}} \Psi \pm \frac{\partial}{\partial \mathrm{y}} \mathrm{y} \Psi\right] \mathrm{dxdy}-2 \iint_{\mathrm{R}} \Psi \mathrm{dxdy} \tag{17}
\end{equation*}
$$

But

$$
\begin{equation*}
\int_{\mathrm{R}} \int_{[ }\left[\frac{\partial}{\partial \mathrm{x}} \Psi \pm \frac{\partial}{\partial \mathrm{y}} \Psi \Psi_{-} \mathrm{dx} \mathrm{dy}=\int_{\mathrm{C}} \Psi \mathrm{x} \cos (\mathrm{x}, v)+\mathrm{y} \cos (\mathrm{y}, v) \underline{\mathrm{d}} \mathrm{~s}\right. \tag{18}
\end{equation*}
$$

## So that we obtain

$\mathrm{V}^{*}=\left(\frac{\mu \alpha^{2} \ell}{2}\right) \iint \Sigma \Psi_{-}^{-2}-\mathrm{u} \Psi_{-}^{-}{ }_{-}^{\mathrm{dxdy}}+2 \int_{\mathrm{C}} \Psi \mathrm{Y} \cos (\mathrm{x}, v)+\mathrm{y} \cos (\mathrm{y}, v) \underline{\overline{\mathrm{d}} \mathrm{s}, .}$
where

$$
\begin{equation*}
(\nabla \Psi)^{2}=\left(\Psi_{\mathrm{x}}\right)^{2}+\left(\Psi_{\mathrm{y}}\right)^{2} \tag{20}
\end{equation*}
$$

If the region $R$ is simply connected, we can take

$$
\begin{equation*}
\Psi=0, \quad \text { on } C \tag{21}
\end{equation*}
$$

and for the determination of $\Psi$, we have the functional

$$
\begin{equation*}
\mathrm{V}^{*}=\left(\frac{\mu \alpha^{2} \ell}{2}\right) \iint_{\mathrm{R}} \nabla \Psi^{-2}-4 \Psi_{-}^{-} \mathrm{dx} \mathrm{dy} \tag{22}
\end{equation*}
$$

This functional (22) is to be minimized on the set of all functions of class $C^{2}$ vanishing on the boundary $C$ of the simply connected region $R$.

### 8.6. VARIATIONAL PROBLEM RELATED TO THE BIHARMONIC EQUATION

Consider the variational problem

$$
\begin{equation*}
I(u)=\iint_{R}\left[\nabla^{2} u^{2}-2 f u\right] d x d y=\min , \tag{1}
\end{equation*}
$$

where the admissible functions $\mathrm{u}(\mathrm{x}, \mathrm{y})$ satisfy on the boundary C of the region R the conditions

$$
\begin{align*}
\mathbf{u} & =\phi(\mathbf{s}),  \tag{2}\\
\frac{\partial u}{\partial \eta} & =\mathrm{h}(\mathrm{~s}), \tag{3}
\end{align*}
$$

Suppose that the set $\{\mathrm{u}(\mathrm{x}, \mathrm{y})\}$ of all admissible functions includes the minimizing function $\mathrm{u}(\mathrm{x}, \mathrm{y})$. We represent an arbitrary function $\overline{\mathrm{u}}(\mathrm{x}, \mathrm{y})$ of this set in the form

$$
\begin{equation*}
\overline{\mathrm{u}}(\mathrm{x}, \mathrm{y})=\mathrm{u}(\mathrm{x}, \mathrm{y})+\in \mathrm{\eta}(\mathrm{x}, \mathrm{y}) \tag{4}
\end{equation*}
$$

where $\in$ is a small real parameter.

## We know that the necessary condition that u minimize the

 integral (1) is$$
\begin{equation*}
\delta \mathrm{I}=\left[\frac{\mathrm{d}}{\mathrm{~d} \in} \mathrm{I}+\in \eta\right]_{\in=0}=0 \tag{5}
\end{equation*}
$$

Using (1), we write

$$
\begin{equation*}
\mathrm{I}(\mathrm{u}+\in \eta)=\iint_{\mathrm{R}}\left[\nabla^{2} u+\in \eta_{-}^{-2}-2 \mathrm{f}(\mathrm{u}+\in \eta)\right] \mathrm{dx} \mathrm{dy} \tag{6}
\end{equation*}
$$

Therefore equations (5) and (6) give

$$
\begin{equation*}
\iint_{\mathrm{R}} \int^{2} u \nabla^{2} \eta-f \eta d x d y=0 . \tag{7}
\end{equation*}
$$

We note that

$$
\begin{aligned}
& \nabla^{2} u \nabla^{2} \eta,=\left(\nabla^{2} u\right)\left(\frac{\partial^{2} \eta}{\partial x^{2}}+\frac{\partial^{2} \eta}{\partial y^{2}}\right) \\
& =\left[\frac{\partial}{\partial x}\left(\nabla^{2} u \frac{\partial \eta}{\partial x}\right)+\frac{\partial}{\partial y}\left(\nabla^{2} u \frac{\partial \eta}{\partial y}\right)\right]-\left[\frac{\partial}{\partial x}\left(\eta \frac{\partial}{\partial x} \nabla^{2} u\right)+\frac{\partial}{\partial y}\left(\eta \frac{\partial}{\partial y} \nabla^{2} u\right)\right]
\end{aligned}
$$

## Applying Gauss Diwergence \& Stoke's theorems we get

$$
\begin{align*}
& \int_{\mathrm{R}} \int_{[ }\left[\frac{\partial}{\partial \mathrm{x}}\left(\nabla^{2} \mathbf{u} \frac{\partial \eta}{\partial \mathrm{x}}\right)+\frac{\partial}{\partial y}\left(\nabla^{2} \mathbf{u} \frac{\partial \eta}{\partial y}\right)\right] \mathrm{dx} \mathrm{dy}=\int_{\mathrm{C}} \nabla^{2} \mathbf{u} \frac{\partial \eta}{\partial v} \mathrm{ds},  \tag{9}\\
& \int_{\mathrm{R}}\left[\int_{\partial} \frac{\partial}{\partial \mathrm{x}}\left(\eta \frac{\partial}{\partial \mathrm{x}} \nabla^{2} \mathbf{u}\right)+\frac{\partial}{\partial y}\left(\eta \frac{\partial}{\partial y} \nabla^{2} \mathbf{u}\right)\right] \mathrm{dx} \mathrm{dy}=\int_{C} \eta \frac{\partial}{\partial v} \nabla^{2} \mathbf{u} \overline{\mathrm{ds}} . \tag{10}
\end{align*}
$$

For equations (7) to (10), we obtain
$\delta I=\iint_{R} \nabla^{2} \nabla^{2} u-f-\bar{\eta} d x d y+\int_{C} \nabla^{2} u \frac{\partial \eta}{\partial v} d s-\int_{C} \eta \frac{\partial}{\partial v} \boldsymbol{\nabla}^{2} u \bar{d} s=0$
From equations (2) to (4), we have

$$
\begin{equation*}
\eta=0 \text { and } \frac{\partial \eta}{\partial v}=0 \quad \text { on } C . \tag{12}
\end{equation*}
$$

## Using result (12), equation (11) gives

$$
\begin{equation*}
\nabla^{4} u=f(x, y) \quad \text { in the region } R \text {. } \tag{13}
\end{equation*}
$$

This is the same differential equation which arises in the study of the transverse deflection of thin elastic plates.

We have assumed in the foregoing that the admissible functions in the set $\{u(x, y)\}$ satisfy the boundary conditions (2) and (3).

If we consider a larger set $S$ of all functions $u$ belonging to class $C^{4}$, then (11) must hold for every $u$ in this set. But the set $S$ includes functions that satisfy the boundary conditions (2) and (3), and thus we must have,

$$
\begin{equation*}
\iint_{\mathrm{R}} \boldsymbol{\nabla}^{4} \mathrm{u}-\mathrm{f}-\eta \mathrm{dx} \mathrm{dy}=0 \tag{14}
\end{equation*}
$$

Since $\eta$ is arbitrary, it follows that the minimizing function $u(x, y)$ again satisfy (13) and we conclude from (11) that

$$
\begin{equation*}
\int_{\mathrm{C}} \nabla^{2} \mathbf{u} \frac{\partial \eta}{\partial v} \mathrm{ds}-\int_{\mathrm{C}} \eta \frac{\partial}{\partial v} \nabla^{2} \mathbf{u}-\overline{d s}=0 \tag{15}
\end{equation*}
$$

for every $\eta$ of class $C^{4}$.
Now if we consider first all $\eta$ such that

$$
\begin{equation*}
\eta=0 \quad \text { on } C \text { and } \frac{\partial \eta}{\partial v} \neq 0 \quad \text { on } C, \tag{16}
\end{equation*}
$$

it follows from (15) that

$$
\begin{equation*}
\nabla^{2} \mathrm{u}=0, \quad \text { on } \mathrm{C} . \tag{17}
\end{equation*}
$$

On the other hand, if we consider only those $\eta$ such that

$$
\begin{equation*}
\eta \neq 0, \quad \frac{\partial \eta}{\partial v}=0, \quad \text { on } C \tag{18}
\end{equation*}
$$

## We get the condition

$$
\frac{\partial}{\partial v}\left(\nabla^{2} u\right)=0, \quad \text { on } C .
$$

Hence if the functional in (1) is minimized on the set $S$ of all $u$ of class $C^{4}$, the minimizing function will be found among those functions of $S$ which satisfy the conditions (17) and (19) on the boundary of the region.

Remark:- For this variational problem (1), we shall obtain the same differential equation (13) when the minimizing function $u=u(x, y)$, instead of the system given in (2) and (3), satisfies the following boundary condition.

$$
\begin{aligned}
\nabla^{2} \mathrm{u} & =0, \\
\frac{\partial}{\partial v} \nabla^{2} \mathrm{u} & =0, \quad \text { on } \mathrm{C} .
\end{aligned}
$$

### 8.7. RITZ METHOD :- ONE DIMENSIONAL CASE

## Consider the variational problem

$$
\begin{equation*}
\mathrm{I}(\mathrm{y})=\int_{\mathrm{x}_{0}}^{\mathrm{x}_{1}} \mathrm{~F} \mathbf{\wedge}, \mathrm{y}, \mathrm{y}^{\prime} \underline{-\mathrm{d}} \mathrm{x}, \tag{1}
\end{equation*}
$$

in which all admissible function $\mathrm{y}=\mathrm{y}(\mathrm{x})$ are such that

$$
\begin{equation*}
\mathrm{y}\left(\mathrm{x}_{0}\right)=\mathrm{y}_{1}, \quad \mathrm{y}\left(\mathrm{x}_{1}\right)=\mathrm{y}_{1} . \tag{2}
\end{equation*}
$$

We know that such a function $y$ is a solution of the Euler's equation

$$
\begin{equation*}
\mathrm{F}_{\mathrm{y}}-\frac{\mathrm{d}}{\mathrm{dx}} \mathrm{~F}_{\mathrm{y}}^{\prime}=0 . \tag{3}
\end{equation*}
$$

A direct method to obtain the desired function was proposed by W. Ritz in 1911.

In this method, we construct a sequence of functions which converge to desired solution of the Euler's equation (3).

Outlines of the Ritz Method:
Let $\mathrm{y}=\mathrm{y}^{*}(\mathrm{x})$ be the exact solution of the given variational problem. Let $\mathrm{I}\left(\mathrm{y}^{*}\right)=$ m be the minimum value of the functional in (1).

In this method, one tries to find a sequence $\left\{\bar{y}_{\mathrm{n}}(\mathrm{x})\right\}$ of admissible functions such that

$$
\begin{equation*}
\lim _{n \rightarrow \infty} I\left(\bar{y}_{n}(x)\right)=m, \tag{4}
\end{equation*}
$$

so that

$$
\begin{equation*}
\lim _{n \rightarrow \infty} \bar{y}_{n}(x)=y^{*}(x), \tag{5}
\end{equation*}
$$

## is the required function.

Ritz proposed to construct a function $\overline{\mathrm{y}}(\mathrm{x})$ by choosing a family of functions

$$
\begin{equation*}
y(x)=\phi\left(x, a_{1}, a_{2} \ldots \ldots \ldots \ldots, a_{k}\right) \tag{6}
\end{equation*}
$$

depending on $k$ real parameters $a_{i}$, where $\phi$ is such that for all values of the $a_{i}$, the end conditions given in (2) are satisfied.

Then, on putting the value of $y$ from (6) in (1), one obtains

$$
\mathrm{I}\left(\mathrm{a}_{1}, \mathrm{a}_{2}, a_{3}, \ldots \ldots \ldots, \mathrm{a}_{\mathrm{k}}\right) .
$$

This functional can be minimized by determining the values of the parameters $\mathrm{a}_{\mathrm{i}}$ from the following system of equations:

$$
\begin{equation*}
\frac{\partial \mathrm{I}\left(\mathrm{a}_{1}, \mathrm{a}_{2}, \ldots \ldots \ldots . \mathrm{a}_{\mathrm{k}}\right)}{\partial \mathrm{a}_{\mathrm{r}}}=0 \quad \text { for } \mathrm{r}=1,2, \ldots \ldots \ldots \mathrm{k} . \tag{7}
\end{equation*}
$$

Let this system has solution $\overline{\mathrm{a}}_{1}, \overline{\mathrm{a}}_{2}, \ldots . . . . ., \overline{\mathrm{a}}_{\mathrm{k}}$. Then, the minimizing function, say $\bar{y}(x)$, is

$$
\begin{equation*}
\overline{\mathrm{y}}(\mathrm{x})=\phi\left(\mathrm{x}, \overline{\mathrm{a}}_{1}, \overline{\mathrm{a}}_{2}, \ldots . . . . . . . ., \overline{\mathrm{a}}_{\mathrm{k}}\right) . \tag{8}
\end{equation*}
$$

It is expected that $\bar{y}(x)$ will be a fair approximation to the minimizing function $y^{*}(x)$ when the number of parameters in (3) is made sufficiently large. Now, we construct a sequence $\left\{\overline{\mathrm{y}}_{\mathrm{n}}(\mathrm{x})\right\}$ of functions $\overline{\mathrm{y}}_{\mathrm{n}}(\mathrm{x})$ such that

$$
\begin{equation*}
\lim _{n \rightarrow \infty} I\left(\bar{y}_{n}\right)=m . \tag{9}
\end{equation*}
$$

Consider a sequence of families of functions of the type (6), namely,

$$
\left.\begin{array}{l}
y_{1}(x)=\phi_{1}\left(x, a_{1}\right) \\
y_{2}(x)=\phi_{2}\left(x, a_{1}, a_{2}\right) \\
y_{3}(x)=\phi_{3}\left(x, a_{1}, a_{2}, a_{3}\right)  \tag{10}\\
------------- \\
y_{n}(x)=\phi_{n}\left(x, a_{1}, a_{2}, \ldots \ldots . ., a_{n}\right)
\end{array}\right\},
$$

in which the family $y_{k}(x)=\phi_{k}\left(x, a_{1}, \ldots \ldots \ldots, a_{k}\right)$ includes in it all functions in the families with subscript less than k .

The parameters $a_{i}$ in each function $y_{k}$ can be determined so as to minimize the integral $I\left(y_{k}\right)$. We denote the values of the parameters thus obtained by $\bar{a}_{\mathrm{i}}$, so that the minimizing functions are

$$
\begin{equation*}
\overline{\mathrm{y}}_{\mathrm{n}}(\mathrm{x})=\phi_{\mathrm{n}}\left(\mathrm{x}, \overline{\mathrm{a}}_{1}, \overline{\mathrm{a}}_{2}, \ldots \ldots \ldots ., \overline{\mathrm{a}}_{\mathrm{n}}\right), \tag{11}
\end{equation*}
$$

for $n=1,2, \ldots$. Since each family $y_{k}(x)$ includes the families $y_{k-1}(x)$ for special values of parameters $a_{i}$, the successive minima $I\left(\bar{y}_{k}\right)$ are nonincreasing, therefore,
$\mathrm{I}\left(\overline{\mathrm{y}}_{1}\right) \geq \mathrm{I}\left(\overline{\mathrm{y}}_{2}\right) \geq \mathrm{I}\left(\overline{\mathrm{y}}_{3}\right) \geq \ldots \ldots \ldots \geq \mathrm{I} \overline{\mathrm{h}}_{\mathrm{n}-1} \geq \mathrm{I}\left(\overline{\mathrm{y}}_{\mathrm{n}}\right) \geq \ldots \ldots \ldots$.
Since the sequence $\left\{\mathrm{I}\left(\overline{\mathrm{y}}_{\mathrm{n}}\right)\right\}$ of real numbers is bounded below by m and is non-increasing ( m being the exact minima), therefore, it is a convergent sequence.

In order to ensure the convergence of this sequence to $\mathrm{I}\left(\mathrm{y}^{*}\right)$, one must impose some conditions on the choice of functions $\phi_{i}$ in (10). We take the set of functions in (10) to be relatively complete.

Then for each $\in \rightarrow 0$, there exists (by definition of relatively complete) in the family (10), a function

$$
y_{n}^{*}(x)=y_{n}\left(x, a_{1}^{*}, a^{*}{ }_{2}, \ldots \ldots . . a_{n}^{*}\right)
$$

such that

$$
\begin{equation*}
\left|\mathrm{y}_{\mathrm{n}}^{*}-\mathrm{y}^{*}\right|<\in \tag{13}
\end{equation*}
$$

and

$$
\begin{equation*}
\left|y_{n}^{* \prime}-y^{* \prime}\right|<\in, \tag{14}
\end{equation*}
$$

for all $\mathrm{x} \in\left(\mathrm{x}_{0}, \mathrm{x}_{1}\right)$. But, it is known that $\mathrm{F}\left(\mathrm{x}, \mathrm{y}, \mathrm{y}^{\prime}\right)$ is a continuous function of its arguments, therefore,

$$
\begin{equation*}
\left|\mathrm{F}\left(\mathrm{x}, \mathrm{y}^{*}{ }_{\mathrm{n}}, \mathrm{y}_{\mathrm{n}}^{* \prime}\right)-\mathrm{f}\left(\mathrm{x}, \mathrm{y}^{*}, \mathrm{y}^{*}{ }^{*}\right)\right|<\in, \tag{15}
\end{equation*}
$$

for all x in ( $\mathrm{x}_{0}, \mathrm{x}_{1}$ ). Consequently

$$
\mathrm{I}\left(\mathrm{y}_{\mathrm{n}}^{*}\right)-\mathrm{I}\left(\mathrm{y}^{*}\right)=\int_{\mathrm{x}_{0}}^{\mathrm{x}_{1}} \mathrm{~F}\left(\mathrm{x}, \mathrm{y}_{\mathrm{n}}^{*}, \mathrm{y}_{\mathrm{n}}^{* \prime}\right)-\mathrm{F}\left(\mathrm{x}, \mathrm{y}^{*}, \mathrm{y}^{* '}\right) \mathrm{dx}
$$

$$
\begin{aligned}
& =\int_{x_{0}}^{x_{1}}\left|\mathrm{~F}\left(\mathrm{x}, \mathrm{y}_{\mathrm{n}}^{*}, \mathrm{y}_{\mathrm{n}}^{*}\right)-\mathrm{F}\left(\mathrm{x}, \mathrm{y}^{*}, \mathrm{y}^{* \prime}\right)\right| \mathrm{dx} \\
& <\in^{\prime}, \text { say. }
\end{aligned}
$$

## This gives

$$
\begin{equation*}
\mathrm{I}\left(\mathrm{y}_{\mathrm{n}}^{*}\right)<\mathrm{I}\left(\mathrm{y}^{*}\right)+\epsilon^{\prime} . \tag{1}
\end{equation*}
$$

As $y^{*}{ }_{n}$ is a function of the set (10) and $\mathrm{I}\left(\overline{\mathrm{y}}_{\mathrm{n}}\right)$ is a minimum of $\mathrm{I}(\mathrm{y})$ on the family $y_{n}$, therefore,

$$
\begin{equation*}
I\left(y^{*}{ }_{n}\right) \geq I\left(\bar{y}_{n}\right) . \tag{17}
\end{equation*}
$$

As $y^{*}$ is the exact solution of the problem and $\bar{y}_{\mathrm{n}}$ is an approximation of the same, therefore,

$$
\begin{equation*}
\mathrm{I}\left(\mathrm{y}^{*}\right) \leq\left(\overline{\mathrm{y}}_{\mathrm{n}}\right) . \tag{18}
\end{equation*}
$$

Combining the in equalities (16) to (18), we find

$$
\begin{equation*}
\mathrm{I}\left(\mathrm{y}^{*}\right) \leq \mathrm{I}\left(\overline{\mathrm{y}}_{\mathrm{n}}\right) \leq \mathrm{I}\left(\overline{\mathrm{y}}_{\mathrm{n}}\right)<\mathrm{I}\left(\mathrm{y}^{*}\right)+\in, \tag{19}
\end{equation*}
$$

but $\epsilon^{\prime}$ can be made as small as we wish, therefore, we get

$$
\begin{equation*}
\operatorname{Lt}_{n \rightarrow \infty} I\left(\bar{y}_{n}\right)=\underset{n \rightarrow \infty}{\operatorname{Lt}} I\left(y_{n}^{*}\right)=I\left(y^{*}\right) . \tag{20}
\end{equation*}
$$

This completes the proof.

Definition. Let $\mathrm{y}(\mathrm{x})$ be an admissible function satisfying the end conditions

$$
\mathrm{y}\left(\mathrm{x}_{0}\right)=\mathrm{y}_{0}, \mathrm{y}\left(\mathrm{x}_{1}\right)=\mathrm{y}_{1} .
$$

If , for each $\in>0$, their exists in the family (6) a function

$$
\mathrm{y}_{\mathrm{n}}^{*}(\mathrm{x})=\mathrm{y}_{\mathrm{n}}\left(\mathrm{x}, \mathrm{a}^{*}{ }_{1}, \mathrm{a}^{*}{ }_{2}, \mathrm{a}^{*}{ }_{3}, \ldots, \ldots, \mathrm{a}_{\mathrm{n}}^{*}\right)
$$

such that

$$
\left|y^{*}{ }_{n}-y^{*}\right|<\epsilon \quad \text { and } \quad\left|y_{n}^{*}-y^{*}\right|<\epsilon
$$

for all x in $\left(\mathrm{x}_{0}, \mathrm{x}_{1}\right)$, then the family of functions in (6) is said to be relatively complete.

Remark:- Among useful, relatively complete sets of functions in the interval $(0, \ell)$ are
(i) trigonometrically polynomials: $\sum_{\mathrm{k}=1}^{\mathrm{n}} \mathrm{a}_{\mathrm{k}} \sin \left(\frac{\mathrm{k} \pi \mathrm{x}}{\ell}\right)$,
(ii) algebraic polynomials : $\sum_{\mathrm{k}=1}^{\mathrm{n}} \mathrm{a}_{\mathrm{k}} \mathrm{x}^{\mathrm{k}}-\mathrm{x}$.

Question:- Show that the system of equations

$$
\frac{\partial}{\partial a_{j}} I\left(y_{n}\right)=0(\text { for } j=1,2,3, \ldots \ldots, n)
$$

for the coefficients in the approximate solution

$$
y_{n}(x)=\sum_{k=1}^{n} a_{k} \phi_{k}(x)
$$

of the variational problem

$$
\begin{aligned}
& I(y)=\int\left[p y^{\prime 2}+q y^{2}+2 f y\right] d x=\min \\
& y(0)=y(\ell)=0
\end{aligned}
$$

by the Ritz method is

$$
\iint^{b} y_{n}^{\prime} \phi_{j}^{\prime}+q y_{n} \phi_{j}+f \phi_{j} d x=0, j=1,2, \ldots \ldots \ldots, n
$$

Solution:- We have

$$
\begin{equation*}
\mathrm{I}\left(\mathrm{y}_{\mathrm{n}}\right)=\int\left[\mathrm{p} \mathrm{y}_{\mathrm{n}}{ }^{\prime{ }^{2}}+\mathrm{q} \mathrm{y}_{\mathrm{n}}^{2}+2 \mathrm{f} \mathrm{y}_{\mathrm{n}}\right] \mathrm{dx} . \tag{1}
\end{equation*}
$$

Therefore, we find

$$
\begin{align*}
\frac{\partial}{\partial \mathrm{a}_{\mathrm{j}}} \mathrm{I}\left(\mathrm{y}_{\mathrm{n}}\right) & =\int_{0}^{\ell}\left[2 \mathrm{p} \mathrm{y}_{\mathrm{n}} \prime^{\prime} \cdot \frac{\partial}{\partial \mathrm{a}_{\mathrm{j}}} \cdot\left(\mathrm{y}_{\mathrm{n}}^{\prime}\right)+2 \mathrm{q} \mathrm{y}_{\mathrm{n}} \frac{\partial}{\partial \mathrm{a}_{\mathrm{j}}} \cdot \mathrm{y}_{\mathrm{n}}+2 \mathrm{f} \frac{\partial}{\partial \mathrm{a}_{\mathrm{j}}} \mathrm{y}_{\mathrm{n}}\right] \mathrm{dx} \\
& =2 \int_{0}^{\ell} \mathrm{p} \mathrm{y}_{\mathrm{n}} \prime^{\prime} \phi_{\mathrm{j}}^{\prime}+\mathrm{q} \mathrm{y}_{\mathrm{n}} \phi_{\mathrm{j}}+\mathrm{f} \phi_{\mathrm{j}} d x . \tag{2}
\end{align*}
$$

## Hence, the system of equations

$$
\begin{equation*}
\frac{\partial}{\partial \mathrm{a}_{\mathrm{j}}} \mathrm{I}\left(\mathrm{y}_{\mathrm{n}}\right)=0 \tag{3}
\end{equation*}
$$

becomes, using (2),

$$
\begin{equation*}
\int_{0}^{\ell} \mathrm{d} \mathrm{y}_{\mathrm{n}}{ }^{\prime} \phi_{\mathrm{j}}^{\prime}+\mathrm{q} \mathrm{y}_{\mathrm{n}} \phi_{\mathrm{j}}+\mathrm{f} \phi_{\mathrm{j}} \mathrm{dx}=0 \tag{4}
\end{equation*}
$$

for $\mathrm{j}=1,2, \ldots \ldots, n$.
This completes the solution.

### 8.8. RITZ METHOD :- TWO-DIMENSIONAL CASE

## Consider the functional in the form

$$
\begin{equation*}
I(u)=\iint_{R} F\left(x, y, u, u_{x}, u_{y}\right) d x d y \tag{1}
\end{equation*}
$$

We suppose that the admissible functions in the variational problem,

$$
\begin{equation*}
\mathrm{I}(\mathrm{u})=\text { minimum }, \tag{2}
\end{equation*}
$$

satisfy the condition

$$
u=\varphi(s)
$$

## on the boundary $C$ of the region $R$.

Let $u^{*}(x, y)$ be an exact solution of the variational problem (obtained by solving the corresponding Euler's equation) and let

$$
\begin{equation*}
\mathrm{I}\left(\mathrm{u}^{*}\right)=\mathrm{m} \tag{3}
\end{equation*}
$$

be the minimum value of the functional (1).
We now introduce a sequence $\left\{u_{n}(x, y)\right\}$ of families of admissible functions

$$
\begin{equation*}
\mathrm{u}_{\mathrm{n}}(\mathrm{x}, \mathrm{y})=\phi_{\mathrm{n}}\left(\mathrm{x}, \mathrm{y}, \mathrm{a}_{1}, \mathrm{a}_{2}, \ldots \ldots ., \mathrm{a}_{\mathrm{n}}\right), \tag{4}
\end{equation*}
$$

with parameters $a_{i}$, and suppose that each family $u_{i}(x, y)$ includes in it families with subscripts less than $i$.

We further assume that the set (4) is relatively complete.
Then, for each $\in>0$, there exists a function

$$
\mathrm{u}_{\mathrm{n}}^{*}(\mathrm{x}, \mathrm{y})=\phi_{\mathrm{n}}\left(\mathrm{x}, \mathrm{y}, \mathrm{a}^{*}{ }_{1}, \mathrm{a}_{2}^{*}, \ldots \ldots \ldots, \mathrm{a}_{\mathrm{n}}^{*}\right)
$$

belonging to the set (4) such that
$\left|\mathrm{u}^{*}{ }_{\mathrm{n}}-\mathrm{u}^{*}\right|<\epsilon, \quad\left|\frac{\partial \mathrm{u}_{\mathrm{n}}^{*}}{\partial \mathrm{x}}-\frac{\partial \mathrm{u}^{*}}{\partial \mathrm{x}}\right|<\epsilon,\left|\frac{\partial \mathrm{u}^{*}{ }_{\mathrm{n}}}{\partial \mathrm{y}}-\frac{\partial \mathrm{u}^{*}}{\partial \mathrm{y}}\right|<\epsilon$,
for all ( $x, y$ ) $\in R$. With the help of (4), we form $I\left(u_{n}\right)$ and determine the parameters $a_{i}$ so that $I\left(u_{n}\right)$ is a minimum.

Let $\bar{a}_{i}$ be the values of the $a_{i}$ obtained by solving the system of equations (known as Ritz's equations)

$$
\begin{equation*}
\frac{\partial \mathrm{I}\left(\mathrm{u}_{\mathrm{n}}\right)}{\partial \mathrm{a}_{\mathrm{j}}}=0, \tag{6}
\end{equation*}
$$

for $\mathrm{j}=1,2, \ldots \ldots, \mathrm{n}$. We write

$$
\begin{equation*}
\overline{\mathrm{u}}_{\mathrm{n}}(\mathrm{x}, \mathrm{y})=\varphi_{\mathrm{n}}\left(\mathrm{x}, \mathrm{y}, \overline{\mathrm{a}}_{1}, \overline{\mathrm{a}}_{2}, \ldots \ldots \ldots . . \ldots . . \overline{\mathrm{a}}_{\mathrm{n}}\right) . \tag{7}
\end{equation*}
$$

The sequence $\left\{I\left(\bar{u}_{n}\right)\right\}$ of real numbers then converges to $I\left(u^{*}\right)=m$, where $u^{*}$ $(\mathrm{x}, \mathrm{y})$ is the function that minimizes (1).

The remaining proof is similar to the proof for one dimensional case.
Illustration. Find an approximate solution to the problem of extremising the functional

$$
\begin{equation*}
I(z)=\iint_{D} z_{x}^{2}+z_{y}^{2}-2 z d x d y \tag{1}
\end{equation*}
$$

where the region R is a sequence, $-\mathrm{a} \leq \mathrm{x} \leq \mathrm{a},-\mathrm{a} \leq \mathrm{y} \leq \mathrm{a}$ and $\mathrm{z}=0$ on the boundary of the sequence $D$.

Solution: We shall seek an approximate solution in the form

$$
\begin{equation*}
\mathrm{z}_{1}=\mathrm{z}=\alpha_{1}\left(\mathrm{x}^{2}-\mathrm{a}^{2}\right)\left(\mathrm{y}^{2}-\mathrm{a}^{2}\right), \tag{2}
\end{equation*}
$$

in which $\alpha_{1}$ is a constant to be determined.
It is clear that this function $\mathrm{z}_{1}$ satisfies the boundary condition. Putting the value of $z$ from (2) into (1), we find

$$
\begin{aligned}
\mathrm{I}\left(\mathrm{z}_{1}\right) & \left.=\iint_{D} 4 \alpha_{1}^{2} \mathrm{x}^{2} y^{2}-\mathrm{a}^{2}-4 \alpha_{1}^{2} \mathrm{y}^{2} \mathfrak{t}^{2}-\mathrm{a}^{2}-2 \alpha_{1} \mathrm{t}^{2}-\mathrm{a}^{2} y^{2}-\mathrm{a}^{2}\right] d x d y \\
& =4 \alpha_{1}^{2}
\end{aligned}
$$

$$
\int_{-a}^{a} x^{2} d x \int_{-a}^{a} y^{2}-a^{2}{ }_{-}^{2} d y+4 \alpha_{1}^{2} \int_{-a}^{a} y^{2} d y \cdot \int_{-a}^{a} t^{2}-a^{2}{ }_{-}^{2} d x-2 \alpha_{1} \int_{-a}^{a} t^{2}-a^{2}-x_{-a}^{a} \int_{-a} y^{2}-a^{2} \underline{d} y
$$

$$
\begin{align*}
& =4 \alpha_{1}^{2}\left[\frac{x^{3}}{3}\right]_{-a}^{a}\left[\frac{y^{5}}{5}-\frac{2 a^{3} y^{3}}{3}+a^{4} y\right]_{-a}^{a} \\
& +4 \alpha_{1}^{2}\left[\frac{y^{3}}{3}\right]_{-a}^{a}\left[\frac{x^{5}}{5}-2 a^{2} \frac{x^{3}}{3}+a^{4} x\right]_{-a}^{a} \\
& -2 \alpha_{1}\left[\frac{x^{3}}{3}-a^{2} x\right]_{-a}^{a}\left[\frac{y^{3}}{3}-a^{2} y\right]_{-a}^{a} \\
& =8 \alpha_{1}^{2} \frac{1}{3} a^{3}+a^{3}\left(-\frac{a^{5}}{5}-\frac{2 a^{5}}{3}+a^{5}\right)(2)-2 \alpha_{1}\left(\frac{a^{3}}{3}-a^{3}\right) \cdot 4 \\
& =\frac{32}{3} \alpha_{1}^{2} a^{3}\left(\frac{8 a^{5}}{15}\right)-8 \alpha_{1}^{2}\left(\frac{2 a^{3}}{3}\right)^{2}=\frac{32 \times 8}{3 \times 15} \alpha_{1}^{2} a^{8}-\frac{32}{9} \alpha_{1} a^{6} . \tag{3}
\end{align*}
$$

For an extremum value, we have

$$
\begin{equation*}
\frac{\mathrm{dI} \varkappa_{1}^{-}}{\mathrm{d} \alpha_{1}}=0 \tag{4}
\end{equation*}
$$

This gives

$$
\begin{equation*}
\alpha_{1}=\frac{5}{16 \mathrm{a}^{2}} \tag{5}
\end{equation*}
$$

Thus, an approximate solution is

$$
\begin{equation*}
\mathrm{z}=\frac{5}{16 \mathrm{a}^{2}}\left(\mathrm{x}^{2}-\mathrm{a}^{2}\right)\left(\mathrm{y}^{2}-\mathrm{a}^{2}\right) \tag{6}
\end{equation*}
$$

Question. Show that the system of equations $\frac{\partial}{\partial a_{j}} I\left(u_{n}\right)=0,(j=1,2, \ldots \ldots, n)$ for the determination of coefficients in the minimizing function

$$
\mathrm{u}_{\mathrm{n}}=\sum_{\mathrm{i}=1}^{\mathrm{n}} \mathrm{a}_{\mathrm{i}} \phi_{\mathrm{i}}(\mathrm{x}, \mathrm{y})
$$

for the problem

$$
I(u)=\iint_{\mathrm{R}} \mathrm{u}_{\mathrm{x}}^{2}+\mathrm{u}_{\mathrm{y}}^{2}+2 \mathrm{fu} \quad \mathrm{dx} d \mathrm{y}=\mathrm{min}, \quad \mathrm{u}=0 \quad \text { on } \mathrm{C}
$$

is

$$
\int_{\mathrm{R}}\left[\int_{\mathrm{D}} \frac{\partial \mathrm{u}_{\mathrm{n}}}{\partial \mathrm{x}} \frac{\partial \phi_{\mathrm{j}}}{\partial \mathrm{x}}+\frac{\partial \mathrm{u}_{\mathrm{n}}}{\partial \mathrm{y}} \frac{\partial \phi_{\mathrm{i}}}{\partial \mathrm{y}}+\mathrm{f} \phi_{\mathrm{j}}\right] \mathrm{dx} d y=0, \quad \mathrm{j}=1,2,3, \ldots \ldots, \mathrm{n} .
$$

Solution:- The minimizing function is

$$
\begin{equation*}
\mathrm{u}_{\mathrm{n}}=\sum \mathrm{a}_{\mathrm{i}} \phi_{\mathrm{i}} \tag{1}
\end{equation*}
$$

We form

$$
\begin{equation*}
I\left(u_{n}\right)=\iint_{R} u_{n, x}^{2}+u_{n, y}^{2}+2 f u_{n} d x d y \tag{2}
\end{equation*}
$$

Therefore, we find

$$
\begin{gather*}
\frac{\partial}{\partial \mathrm{a}_{\mathrm{j}}} \mathrm{I}\left(\mathrm{u}_{\mathrm{n}}\right) \\
\iint_{\mathrm{R}}\left[2 \mathrm{u}_{\mathrm{n}, \mathrm{x}} \frac{\partial}{\partial \mathrm{a}_{\mathrm{j}}} \mathbf{u}_{\mathrm{n}}, \mathrm{x}-2 \mathrm{u}_{\mathrm{x}, \mathrm{y}} \frac{\partial}{\partial \mathrm{a}_{\mathrm{j}}} \mathbf{a}_{\mathrm{n}}, \mathrm{y}-2 \mathrm{f} \frac{\partial}{\partial \mathrm{a}_{\mathrm{j}}} \mathrm{u}_{\mathrm{n}}\right] \mathrm{dx} \text { dy } \tag{3}
\end{gather*}
$$

We have

$$
\begin{equation*}
\mathrm{u}_{\mathrm{n}, \mathrm{x}}=\sum_{\mathrm{i}=1}^{\mathrm{n}} \mathrm{a}_{\mathrm{i}} \frac{\partial \phi_{\mathrm{i}}}{\partial \mathrm{x}}, \tag{4}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{\partial}{\partial \mathrm{a}_{\mathrm{j}}}\left(\mathrm{u}_{\mathrm{n}, \mathrm{x}}\right)=\frac{\partial \phi_{\mathrm{j}}}{\partial \mathrm{x}} \tag{5}
\end{equation*}
$$

Similarly,

$$
\begin{equation*}
\frac{\partial}{\partial \mathrm{a}_{\mathrm{j}}}\left(\mathrm{u}_{\mathrm{n}}, \mathrm{y}\right)=\frac{\partial \phi_{\mathrm{j}}}{\partial \mathrm{y}}, \tag{6}
\end{equation*}
$$

$$
\begin{equation*}
\frac{\partial}{\partial \mathrm{a}_{\mathrm{j}}} \mathrm{u}_{\mathrm{n}}=\phi_{\mathrm{j}} . \tag{7}
\end{equation*}
$$

Using (5) to (7) in relation (3), we write

$$
\begin{equation*}
\left.\frac{\partial}{\partial \mathrm{a}_{\mathrm{j}}} \mathrm{I}\left(\mathrm{u}_{\mathrm{n}}\right)=2 \int_{\mathrm{R}} \int_{\mathrm{R}}^{[ } \frac{\partial \mathrm{u}_{\mathrm{n}}}{\partial \mathrm{x}} \frac{\partial \phi_{\mathrm{j}}}{\partial \mathrm{x}}+\frac{\partial \mathrm{u}_{\mathrm{n}}}{\partial \mathrm{y}} \frac{\partial \phi_{\mathrm{j}}}{\partial \mathrm{y}}+\mathrm{f} \phi_{\mathrm{j}}\right] \mathrm{dx} d \mathrm{dy} \tag{8}
\end{equation*}
$$

Therefore, the system of equations

$$
\begin{equation*}
\frac{\partial}{\partial \mathrm{a}_{\mathrm{j}}} \mathrm{I}\left(\mathrm{u}_{\mathrm{m}}\right)=0 \tag{9}
\end{equation*}
$$

becomes

$$
\begin{equation*}
\iint_{\mathrm{R}}\left[\frac{\partial \mathrm{u}_{\mathrm{n}}}{\partial \mathrm{x}} \frac{\partial \phi_{\mathrm{j}}}{\partial \mathrm{x}}+\frac{\partial \mathrm{u}_{\mathrm{n}}}{\partial \mathrm{y}} \frac{\partial \phi_{\mathrm{j}}}{\partial \mathrm{y}}+\mathrm{f} \phi_{\mathrm{j}}\right] \mathrm{dx} \mathrm{dy}=0 \tag{10}
\end{equation*}
$$

for $\mathrm{j}=1,2, \ldots ., \mathrm{n}$.
Question:- Show that the system of equations $\frac{\partial \mathrm{I}\left(\mathrm{u}_{\mathrm{n}}\right)}{\partial \mathrm{a}_{\mathrm{j}}}=0,(\mathrm{j}=1,2, \ldots \ldots, n)$, for determining the coefficients in the approximate solution

$$
\mathrm{u}_{\mathrm{x}}=\sum_{\mathrm{i}=1}^{\mathrm{n}} \mathrm{a}_{\mathrm{i}} \phi_{\mathrm{i}}
$$

for the problem

$$
\begin{aligned}
& I(u)=\iint_{R}\left[V^{2} u^{2}-f u\right] d x d y=\min , \\
& u=0, \quad \frac{\partial u}{\partial \eta}=0 \quad \text { on } C,(C \text { being the boundary of } R),
\end{aligned}
$$

is

$$
\iint_{\mathrm{R}} \sum_{\mathrm{i}=1}^{\eta} \mathrm{a}_{\mathrm{i}} \nabla^{2} \phi_{\mathrm{i}} \nabla^{2} \phi_{\mathrm{j}}-\mathrm{f} \phi_{\mathrm{j}} . \mathrm{dx} d y=0
$$

for $\mathrm{j}=1,2, \ldots ., \mathrm{n}$.
Solution:- We form I ( $u_{n}$ ). We obtain

$$
\begin{equation*}
I\left(u_{n}\right)=\iint_{R} \int\left[\nabla^{2} u_{n}^{2}-2 f u_{n}\right] d x d y . \tag{1}
\end{equation*}
$$

We find

$$
\begin{align*}
\frac{\partial}{\partial \mathrm{a}_{\mathrm{j}}} \mathrm{I}\left(\mathrm{u}_{\mathrm{n}}\right) & \left.=2 \iint_{\mathrm{R}} \int_{\mathrm{R}}^{\infty} \nabla^{2} \mathrm{u}_{\mathrm{n}} \cdot \frac{\partial}{\partial \mathrm{a}_{\mathrm{J}}} \nabla^{2} \mathrm{u}_{\mathrm{n}}-\mathrm{f} \frac{\partial}{\partial \mathrm{a}_{\mathrm{J}}} \mathrm{u}_{\mathrm{n}}\right] \mathrm{dx} d y \\
& =2 \iint_{\mathrm{R}} \nabla^{2} \mathrm{u}_{\mathrm{n}} \cdot \nabla^{2} \phi_{\mathrm{j}}-\mathrm{f} \phi_{\mathrm{j}} . \mathrm{dx} d y \\
& =2 \iint_{\mathrm{R}} \sum_{\mathrm{i}=1}^{\eta} \nabla^{2} \phi_{\mathrm{i}} \nabla^{2} \phi_{\mathrm{j}}-\mathrm{f} \phi_{\mathrm{j}} \cdot \mathrm{dx} \text { dy. } \tag{2}
\end{align*}
$$

Therefore, the system of equations

$$
\begin{equation*}
\frac{\partial}{\partial \mathrm{a}_{\mathrm{j}}} \mathrm{I}\left(\mathrm{u}_{\mathrm{n}}\right)=0 \tag{3}
\end{equation*}
$$

becomes

$$
\begin{equation*}
\iint_{\mathrm{R}} \int_{\mathrm{i}=1}^{\eta} \nabla^{2} \phi_{\mathrm{i}} \nabla^{2} \phi_{\mathrm{j}}-\mathrm{f} \phi_{\mathrm{j}} \cdot \mathrm{dx} \mathrm{dy}=0 \tag{4}
\end{equation*}
$$

for $\mathrm{j}=1,2, \ldots \ldots \ldots, n$. This completes the solution.

### 8.9. GALERKIN METHOD

In 1915, Galerkin proposed a method of finding an approximate solution of the boundary value problems in mathematical physics. This method shall have wider scope than the method of Ritz.

## Method : Let it be required to solve a linear differential equation

$$
\begin{equation*}
\mathbf{L}(\mathbf{u})=\mathbf{0} \quad \text { in } \mathbf{R}, \tag{1}
\end{equation*}
$$

subject to some homogeneous boundary conditions, L being a linear differential operator.

It is assumed, for simplicity that the domain R is two-dimensional.
We seek an approximate solution of the problem of the type

$$
\begin{equation*}
\mathrm{u}_{\mathrm{n}}(\mathrm{x}, \mathrm{y})=\sum_{\mathrm{i}=1}^{\mathrm{n}} \mathrm{a}_{\mathrm{i}} \phi_{\mathrm{i}}(\mathrm{x}, \mathrm{y}) \tag{2}
\end{equation*}
$$

where the $\phi_{\mathrm{i}}$ are suitable coordinate functions and $\mathrm{a}_{\mathrm{i}}$ are constant.

We suppose that the functions $\phi_{i}$ satisfy the same boundary conditions as the exact solution $u(x, y)$. We further suppose that the set $\left\{\phi_{i}\right)$ is complete in the sense that every piecewise continues function $f(\mathrm{x}, \mathrm{y})$, say, can be approximated in $R$ by the sum $\sum_{i=1}^{n} c_{i} \phi_{i}$ in such a way that

$$
\begin{equation*}
\delta_{N}=\iint_{R}\left(f-\sum_{i=1}^{N} c_{i} \phi_{i}\right) d x d y \tag{3}
\end{equation*}
$$

can be made as small as we wish.

Ordinarily, $\mathrm{u}_{\mathrm{n}}$ given in (2) will not satisfy (1). Let

$$
\begin{equation*}
\mathrm{L}\left(\mathrm{u}_{\mathrm{n}}\right)=\epsilon_{\mathrm{n}}(\mathrm{x}, \mathrm{y}), \quad \text { where } \epsilon_{\mathrm{x}}(\mathrm{x}, \mathrm{y}) \neq 0, \quad \text { in } \mathrm{R} . \tag{4}
\end{equation*}
$$

If maximum of $\epsilon_{\mathrm{n}}(\mathrm{x}, \mathrm{y})$ is small, we can consider $\mathrm{u}_{\mathrm{n}}(\mathrm{x}, \mathrm{y})$ given is (2) as a satisfactory approximation to the exact solution $\mathrm{u}(\mathrm{x}, \mathrm{y})$.

Thus, to get a good approximation, we have to choose the constants $\mathrm{a}_{\mathrm{i}}$ so as to minimize the error function $\in_{\mathrm{n}}(\mathrm{x}, \mathrm{y})$.

A reasonable minimization technique is suggested by the following:
If one represents $u(x, y)$ by the serious $u(x, y)=\sum_{i=1}^{\infty} a_{i} \phi_{i}$, with suitable properties and consider $\quad u_{n}=\sum_{i=1}^{\eta} c_{i} \phi_{i}$ as the nth partial sum, then, the orthogality condition,

$$
\begin{equation*}
\iint_{\mathrm{R}} \int_{\mathrm{S}} \mathrm{~L}\left(\mathrm{u}_{\mathrm{n}}\right) \phi_{\mathrm{i}}(\mathrm{x}, \mathrm{y}) \mathrm{dx} \mathrm{dy}=0, \tag{5}
\end{equation*}
$$

as $n \rightarrow \infty$ is equivalent to the statement

$$
\begin{equation*}
\mathbf{L}(\mathbf{u})=\mathbf{0}, \tag{6}
\end{equation*}
$$

by virtue of (1).
This led Galerkin to impose on the error function $\in_{\mathrm{n}}$ a set of orthogality conditions (now called Galerkin conditions)

$$
\begin{equation*}
\iint_{\mathrm{R}} \mathrm{~S}^{\mathrm{L}}\left(\mathrm{u}_{\mathrm{n}}\right) \phi_{\mathrm{i}}(\mathrm{x}, \mathrm{y}) \mathrm{dx} \mathrm{dy}=0 \tag{7}
\end{equation*}
$$

for $\mathrm{i}=1,2, \ldots \ldots \ldots, \mathrm{n}$. This yields the set of equations

$$
\begin{equation*}
\iint_{R} \mathrm{~L}\left(\sum_{\mathrm{j}=1}^{\eta} \mathrm{a}_{\mathrm{j}} \phi_{\mathrm{j}}\right) \phi_{\mathrm{i}} \mathrm{dxdy}=0, \tag{8}
\end{equation*}
$$

for $i=1,2, \ldots \ldots \ldots, n$.
This set of $n$ equations determine the constants $a_{i}$ in the approximate solution (2).

Remark 1. When the differential equation and the boundary conditions are self-adjoint and the corresponding functional $\mathrm{I}(\mathrm{u})$ in the problem

$$
\begin{equation*}
\mathrm{I}(\mathrm{u})=\min , \tag{9}
\end{equation*}
$$

is positive definite, then the system of Galerkin equation in (8) is equivalent to the Ritz system

$$
\begin{equation*}
\frac{\partial}{\partial \mathrm{a}_{\mathrm{j}}} \mathrm{I}\left(\mathrm{u}_{\mathrm{n}}\right)=0 \tag{10}
\end{equation*}
$$

Remark 2. It is important to the note that in Galerkin's formulation, there is no reference to any connection of equation (1) with a variational problem. Indeeds, the Galerkin method can be applied to a wider class of problems phrased in terms of integrals and other types of functional equations.

Question. Solve the variational problem

$$
\begin{array}{r}
\int_{0}^{1}\left[y^{\prime 2}-y^{2}-2 x y\right] d x=\min \\
y(0)=y(1)=0 \tag{1}
\end{array}
$$

## by the Galerkin method.

Solution:- Here

$$
\mathrm{F}=\mathrm{y}^{1^{2}}-\mathrm{y}^{2}-2 \mathrm{xy},
$$

and Euler's equation is

$$
\begin{gather*}
y^{\prime \prime}+y-x=0, \quad \text { in } 0<x<1 \\
L[y]=\frac{d^{2} y}{d x^{2}}+y-x=0 \tag{2}
\end{gather*}
$$

We consider an approximate solution of the problem of the form

$$
\begin{equation*}
y_{n}=(1-x)\left[a_{1} x+a_{2} x^{2}+\ldots \ldots \ldots .+a_{n} x^{n}\right], \tag{3}
\end{equation*}
$$

which satisfy the boundary conditions.

The first approximation is

$$
\begin{align*}
y_{1} & =a_{1} x(1-x) \\
& =a_{1}\left(x-x^{2}\right) . \tag{4}
\end{align*}
$$

Here,

$$
\begin{equation*}
\phi_{1}=x-x^{2} . \tag{5}
\end{equation*}
$$

We find

$$
\begin{equation*}
\mathrm{L}\left[\mathrm{y}_{1}\right]=\mathrm{a}_{1}(-2)+\mathrm{a}_{1}\left(\mathrm{x}-\mathrm{x}^{2}\right)+(-\mathrm{x})=\mathrm{a}_{1}\left(\mathrm{x}-\mathrm{x}^{2}-2\right)-\mathrm{x}, \tag{6}
\end{equation*}
$$

and the coefficient $\mathrm{a}_{1}$ is determined from the Galerkin's equation

$$
\begin{equation*}
\int_{0}^{1} \mathrm{~L}\left(\mathrm{u}_{1}\right) \phi_{1} \mathrm{dx}=0 \tag{7}
\end{equation*}
$$

This yields

$$
\begin{aligned}
& \int_{0}^{1} a_{1} t-x^{2}-2-x \cdot\left(x-x^{2}\right) d x=0 \\
& \int_{0}^{1} a_{1}-x^{2}-2-x^{2}-t^{2}-x^{3}-d x=0 \\
& \int_{0}^{1} a_{1} 4 x+3 x^{2}-2 x^{3}+x^{4}-t^{2}-x^{3} d x=0 .
\end{aligned}
$$

This gives

$$
\begin{equation*}
a_{1}=\frac{-5}{18} . \tag{8}
\end{equation*}
$$

Thus, an approximate solution of the given variational problem, using Galerkin method, is

$$
y \equiv y_{1}(x)=\frac{-5}{18}\left(x-x^{2}\right)
$$

### 8.10 APPLICATION OF GALERKIN METHOD TO THE PROBLEM OF TORSION OF BEAMS

Consider a cylindrical bar subjected to no body forces and free external forces on its lateral surface. One end of the bar is fixed in the plane $z=0$ while the other end is in the plane $\mathrm{z}=\ell$ (say). The bar is twisted by a couple of magnitude M whose moment is directed along the axis of the bar (i.e., z -axis). Prandtl introduced a function $\Psi(x, y)$, known as Prandtl stress function, such that

$$
\begin{equation*}
\tau_{\mathrm{zx}}=\mu \alpha \frac{\partial \Psi}{\partial \mathrm{y}}, \quad \tau_{\mathrm{zy}}=-\mu \alpha \frac{\partial \Psi}{\partial \mathrm{x}} \tag{1}
\end{equation*}
$$

and stress function $\Psi$ is determined from the system

$$
\begin{align*}
\nabla^{2} \Psi & =-2 & & \text { in } \mathrm{R}  \tag{2}\\
\Psi & =\mathbf{0} & & \text { on } \mathbf{C} \tag{3}
\end{align*}
$$

where R is the region of the cross-section of the bar and C its boundary. Let $R$ be the rectangle $|x| \leq A,|y| \leq B$.


Now we have to solve the system consisting of equations (2) and (3) by using the Galerkin method. We write (2) as

$$
\begin{equation*}
L(\Psi)=0, \tag{4}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathrm{L} \equiv \nabla^{2}+2 \tag{5}
\end{equation*}
$$

## We take an approximate solution in the form

$\Psi_{n}(x, y)=\left(x^{2}-A^{2}\right)\left(y^{2}-B^{2}\right)\left(a_{1}+a_{2} x^{2}+a_{3} y^{2}+\ldots \ldots . . a_{n} x^{2 k} y^{2 k}\right)$.
This approximate solution satisfies the boundary conditions in (3). Here $a_{1}$, $\mathrm{a}_{2}, \ldots \ldots, \mathrm{a}_{\mathrm{n}}$ are constants to be determined by using Galerkin method.

## The first approximation is

$$
\begin{equation*}
\Psi_{1}=\mathrm{a}_{1} \phi_{1}=\mathrm{a}_{1}\left(\mathrm{x}^{2}-\mathrm{A}^{2}\right)\left(\mathrm{y}^{2}-\mathrm{B}^{2}\right), \tag{7}
\end{equation*}
$$

with

$$
\begin{equation*}
\phi_{1}=\left(\mathrm{x}^{2}-\mathrm{A}^{2}\right)\left(\mathrm{y}^{2}-\mathrm{B}^{2}\right) . \tag{8}
\end{equation*}
$$

The coefficient $a_{1}$ is determined from the Galerkin equation

$$
\begin{equation*}
\int_{-\mathrm{B}}^{\mathrm{B}} \int_{-\mathrm{A}}^{\mathrm{A}} \mathrm{~L} \Psi_{1} \bar{\Phi}_{1} \mathrm{dxdy}=0 \tag{9}
\end{equation*}
$$

## This implies

$$
\int_{-\mathrm{B}}^{\mathrm{B}} \int_{-\mathrm{A}}^{\mathrm{A}} \nabla^{2} \Psi_{1}+2 \phi_{-1} \mathrm{dxdy}=0
$$

$$
\int_{-B}^{B} \int_{-A}^{A} 2 a_{1}\left(y^{2}-B^{2}\right)+2 a_{1}\left(x^{2}-A^{2}\right)+2 \overline{-}^{2}-A^{2} \underline{s}^{2}-B^{2} \underline{d} x d y=0
$$

$$
2 a_{1} \int_{-A}^{A} t^{2}-A^{2} \underline{d} x \int_{-B}^{B} y^{2}-B^{2}{ }_{-}^{2} d y+2 a_{1} \int_{-A}^{A} t^{2}-A^{2} \underline{d} x \int_{-B}^{B} y^{2}-B^{2}{ }_{-}^{2} d y
$$

$$
+2 \int_{-A}^{A} x^{2}-A^{2} \underline{d} x \int_{-B}^{B} y^{2}-B^{2}{ }_{-}^{2} d y=0
$$

## Integration yields

Hence

$$
\begin{equation*}
\mathrm{a}_{1}=\frac{5}{4}\left(\frac{1}{\mathrm{~A}^{2}+\mathrm{B}^{2}}\right) . \tag{..}
\end{equation*}
$$

Therefore, an approximate solution, by Galerkin method, to the given boundary value problem is

$$
\begin{equation*}
\Psi_{1}=\frac{5}{4}\left(\frac{1}{\mathrm{~A}^{2}+\mathrm{B}^{2}}\right)\left(\mathrm{x}^{2}-\mathrm{A}^{2}\right)\left(\mathrm{y}^{2}-\mathrm{B}^{2}\right) \tag{11}
\end{equation*}
$$

Note:- Approximate values of the torsinal rigidity D and maximum shear stress $\tau_{\text {max }}$ can also be computed with the help of (7). We recall that

$$
\begin{equation*}
D=2 \mu \iint_{R} \Psi \mathrm{dx} \mathrm{dy} \tag{12}
\end{equation*}
$$

and the maximum shear stress $\tau_{\max }$ occurs at the mid points of the longer sides of the bar.

If $B>A$, then

$$
\begin{align*}
\tau_{\max } & =\tau_{\mathrm{yz}}^{\substack{\mathrm{x}=\mathrm{A} \\
\mathrm{y}=0}} \mathrm{~m}_{\substack{ \\
\mathrm{y}=0}} . \\
& =-\mu \alpha\left[\frac{\partial \Psi}{\partial \mathrm{x}=\mathrm{A}}\right]_{\substack{ \\
\mathrm{y}=0}} . \tag{13}
\end{align*}
$$

Inserting the values of $\Psi=\Psi_{1}$ given in (7) into relations (12) and (13), we find that

$$
\begin{align*}
& \mathrm{D}_{1}=\frac{5}{18} \mu \mathrm{a}^{2} \mathrm{~b}\left[\frac{\mathrm{~b} / \mathrm{a}^{2}}{1+\mathrm{b} / \mathrm{a}^{2}}\right]  \tag{14}\\
& \tau_{\max }=\frac{5}{4} \mu \alpha \mathrm{a}\left[\frac{\mathrm{~b} / \mathrm{a}^{2}}{1+\mathrm{b} / \mathrm{a}^{2}}\right] \tag{15}
\end{align*}
$$

with $\mathrm{a}=2 \mathrm{~A}, \mathrm{~b}=2 \mathrm{~B}$.

### 8.11. METHOD OF KANTOROVICH

In 1932, Kantorvich proposed a generalization of the Ritz method. In the present method, the integration of partial differential equation (Euler's equation) reduces to the integration of a system of ordinary differential equations.

## In the application of the Ritz method to the problem

$$
\begin{equation*}
I(u)=\iint_{R} F\left(x, y, u, u_{x}, u_{y}\right) d x d y=\min \tag{1}
\end{equation*}
$$

we consider approximate solutions in the form

$$
\begin{equation*}
\mathrm{u}_{\mathrm{n}}=\sum_{\mathrm{k}=1}^{\eta} \mathrm{a}_{\mathrm{k}} \phi_{\mathrm{k}}(\mathrm{x}, \mathrm{y}), \tag{2}
\end{equation*}
$$

where the functions $\varphi_{k}(x, y)$ satisfy the same boundary conditions as those imposed on the exact solution $u(x, y)$ and $a_{k}$ are constants. We then determined the coefficients $a_{k}$ so as to minimize $I\left(u_{n}\right)$.

In the method of Kantorvich, the $\mathrm{a}_{\mathrm{k}}$ in (2) are no longer constants but are unknown functions of $x$ such that the product

$$
a_{k}(x) \phi_{k}(x, y)
$$

satisfies the same boundary conditions as $u(x, y)$.
This led to minimize

$$
\begin{equation*}
\mathrm{I}\left(\mathrm{u}_{\mathrm{n}}\right)=\mathrm{I}\left(\sum_{\mathrm{k}=1}^{\eta} \mathrm{a}_{\mathrm{k}} \backslash \bar{\Phi}_{\mathrm{k}}(\mathrm{x}, \mathrm{y})\right) . \tag{3}
\end{equation*}
$$

Since the functions $\phi_{\mathrm{k}}(\mathrm{x}, \mathrm{y})$ are known functions, we can perform integration w. r. t. y in (1) and then obtain a functional of the type

$$
\begin{equation*}
I\left(u_{n}\right)=\int_{x_{0}}^{x_{1}} f a_{k}(x), a_{k}^{\prime}(x), x_{-}^{-} d x . \tag{4}
\end{equation*}
$$

Kantorvich proposed to determine the function $\mathrm{a}_{\mathrm{k}}(\mathrm{x})$ so that they minimize the functional (4). It is clear that $\mathrm{a}_{\mathrm{k}}(\mathrm{x})$ can be determined by solving the second order ordinary differential equations (Euler's equations)

$$
\begin{equation*}
\mathrm{f}_{\mathrm{a}_{\mathrm{k}}}-\frac{\mathrm{d}}{\mathrm{dx}} \mathrm{f}_{\left(\mathrm{a}_{\mathrm{k})}\right.}=0, \quad \text { for } \mathrm{k}=1,2, \ldots \ldots, \mathrm{n} . \tag{5}
\end{equation*}
$$

Once $a_{k}(x)$ are determined from (5), an approximate solution is known.

### 8.12. APPLICATION OF KANTORVICH METHOD TO THE TORSINAL PROBLEM

## The torsion boundary value problem is

$$
\left.\begin{array}{cc}
\nabla^{2} \Psi=-2 & \text { in } \mathrm{R} \\
\Psi=0 & \text { on } \mathrm{C} \tag{1}
\end{array}\right\}
$$

$C$ being the boundary of $R$, where $R$ is the rectangle $|x| \leq A,|y| \leq B$.
The variational problem associated with this boundary value problem is (a well known result)

$$
\begin{equation*}
I(\Psi)=\iint_{R} \Psi_{\mathrm{x}}^{2}+\Psi_{\mathrm{y}}^{2}-4 \Psi \mathrm{dx} d y=\min \tag{2}
\end{equation*}
$$

An approximate to its solution is

$$
\begin{equation*}
\Psi_{1}=\mathrm{a}_{1}(\mathrm{x})\left(\mathrm{y}^{2}-\mathrm{B}^{2}\right), \tag{3}
\end{equation*}
$$

with

$$
\begin{equation*}
\phi_{1}(x, y)=y^{2}-B^{2} . \tag{4}
\end{equation*}
$$

Here $\phi_{1}(x, y)$ vanishes on the part $y= \pm B$ of the boundary C. In order that

$$
\begin{equation*}
a_{1}(x)\left(y^{2}-B^{2}\right)=0 \quad \text { on } C, \tag{5}
\end{equation*}
$$

we shall determine $a_{1}(x)$ such that

$$
\begin{equation*}
\mathrm{a}_{1}(\mathrm{~A})=\mathrm{a}_{1}(-\mathrm{A})=0 . \tag{6}
\end{equation*}
$$

Inserting the value of $\Psi$ from (3) in (2), we get

$$
\left.I\left(a_{1}\right)=\iint a_{1}^{\prime}(x)^{2} y^{2}-B_{-}^{2}+4 y^{2} a_{1} x^{-2}-4 a_{1} y^{2}-B^{2}\right] d x d y
$$

$$
=\int_{-A}^{A}\left[a_{1}^{\prime}(x)^{2} \cdot \int_{-B}^{B} y^{4}+B^{4}-2 B^{2} y^{2}-d y+4 a_{1} x^{-2}\right.
$$

$$
\left.\times \int_{-B}^{B} y^{2} d y-4 a_{1} \cdot-\int_{-B}^{B} y^{2}-B^{2}-d y\right] d x
$$

$$
\begin{align*}
& =\int_{-A}^{A}\left[a_{1}{ }^{\prime}{ }^{-2}\left[\frac{y^{5}}{5}+B^{4} y-\frac{2 B^{2} y^{3}}{3}\right]_{-B}^{B}\right. \\
& \left.+4 a_{1}<-\left\{\frac{y^{3}}{3}\right\}_{-B}^{B}-4 a_{1}<\left\{\frac{y^{3}}{3}-B^{2} y\right\}_{-B}^{B}\right] d x \\
& =\int_{-A}^{A}\left[2 a_{1}^{\prime-2}\left(\frac{B^{5}}{5}+\frac{B^{5}}{1}-2 \frac{B^{5}}{3}\right)+4 a_{1}^{2}\left(2 \frac{B^{3}}{3}\right)-8 a_{1}\left(\frac{B^{3}}{3}-B^{3}\right)\right] d x \\
& =\left(\frac{2 B^{3}}{3}\right) \int_{-A}^{A}\left[\frac{8}{5} B^{2} \mathbf{a}_{1}{ }^{\prime 2}+4 a_{1}^{2}+8 a_{1}\right] d x . \tag{7}
\end{align*}
$$

Here

$$
\begin{align*}
& \mathrm{f}\left(\mathrm{a}_{1}, \mathrm{a}_{1}^{\prime}, \mathrm{x}\right)=\frac{8}{5} \mathrm{~B}^{2} a_{1^{\prime}}^{\prime-2}+4 \mathrm{a}_{1}^{2}+8 \mathrm{a}_{1} .  \tag{8}\\
& \mathrm{f}_{\mathrm{a}_{1}}=8 \mathrm{a}_{1}+8 \\
& \mathrm{f}_{\mathrm{a}_{1}}=\frac{16}{5} \mathrm{~B}^{2} \mathrm{a}_{1}^{\prime} . \tag{9}
\end{align*}
$$

## Euler's equation becomes

$$
\begin{align*}
& \left(8 a_{1}+8\right)-\frac{d}{d x}\left(\frac{16}{5} B^{2} a_{1}^{\prime}\right)=0 \\
& a_{1}{ }^{\prime \prime}(x) \frac{5}{2 B^{2}} a_{1}(x)-\frac{5}{2 B^{2}}=0 . \tag{10}
\end{align*}
$$

Its solution is

$$
\begin{equation*}
\mathrm{a}_{1}(\mathrm{x})=\mathrm{C}_{1} \cosh .\left(\frac{\mathrm{kx}}{\mathrm{~B}}\right)+\mathrm{C}_{2} \sinh \left(\frac{\mathrm{kx}}{\mathrm{~B}}\right)-1, \tag{11}
\end{equation*}
$$

with

$$
\begin{equation*}
\mathrm{k}=\sqrt{\frac{5}{2}} \tag{12}
\end{equation*}
$$

The function $\mathrm{a}_{1}(\mathrm{x})$ must be an even function, therefore $\mathrm{C}_{2}=0$.
Putting $\mathrm{a}_{1}(\mathrm{~A})=0$, we get,

5

$$
\begin{equation*}
\mathrm{C}_{1}=\frac{1}{\cosh (\mathrm{kx} / \mathrm{B})} \tag{13}
\end{equation*}
$$

Thus, we find

$$
\begin{equation*}
\Psi_{1}=\left(\mathrm{y}^{2}-\mathrm{B}^{2}\right)\left[\frac{\cosh (\mathrm{kx} / \mathrm{B})}{\cosh (\mathrm{kA} / \mathrm{B}}-1\right] \tag{14}
\end{equation*}
$$

as an approximate solution obtained by Kantorvich method.

## Chapter-9

## Waves in Elastic Solids

### 9.1. WAVES IN AN ISOTROPIC ELASTIC SOLID

## In the absence of body force, equations of motion are

$$
\begin{equation*}
\tau_{\mathrm{i}, \mathrm{j}, \mathrm{j}}=\rho \ddot{u}_{i} \tag{1}
\end{equation*}
$$

for $\mathrm{i}, \mathrm{j}=1,2,3$. Here, dot signifies the differentiation with respect to time t and $\rho$ is the density of the solid. $\tau_{\mathrm{ij}}$ is the stress tensor, $\mathrm{u}_{\mathrm{i}}$ is the displacement vector.

The left of (1) is a force in the $\mathrm{x}_{\mathrm{i}}$-direction due to stresses and the right of (1) is the inertia term - mass $\times$ acceleration.

We know that the generalized Hooke's law for an isotropic homogeneous elastic medium is

$$
\begin{equation*}
\tau_{\mathrm{ij}}=\lambda \delta_{\mathrm{ij}} \mathrm{u}_{\mathrm{k}, \mathrm{k}}+\mu\left(\mathrm{u}_{\mathrm{i}, \mathrm{j}}+\mathrm{u}_{\mathrm{j}, \mathrm{i}}\right) \tag{2}
\end{equation*}
$$

Here, homogeneity implies that $\rho, \lambda$ and $\mu$ are constants throughout the medium, $\lambda$ and $\mu$ being Lame' constants.

We put

$$
\begin{equation*}
\theta=\mathrm{u}_{\mathrm{k}, \mathrm{k}}=\operatorname{div} \overline{\mathrm{u}}, \tag{3}
\end{equation*}
$$

$$
\begin{equation*}
\bar{\Omega}=\operatorname{curl} \overline{\mathrm{u}} \tag{4}
\end{equation*}
$$

## $\theta$ is the cubical dilatation and $\bar{\Omega}$ is the rotation vector.

Putting (2) into (1) and using (3), we obtain the following Navier's equation of motion (exercise)

$$
\begin{equation*}
(\lambda+\mu) \operatorname{grad} \theta+\mu \nabla^{2} \overline{\mathrm{u}}=\rho \ddot{\overline{\mathrm{u}}} \tag{5}
\end{equation*}
$$

in which $\nabla^{2}$ is the Laplacian.
Taking divergence of both sides of (5) and using (3), we obtain

$$
(\lambda+\mu) \operatorname{div}(\operatorname{grad} \theta)+\mu \nabla^{2}(\theta)=\rho \ddot{\theta}
$$

or

$$
(\lambda+\mu) \nabla^{2} \theta+\mu \nabla^{2} \theta=\rho \ddot{\theta}
$$

or

$$
\begin{equation*}
\nabla^{2} \theta=\frac{1}{\alpha^{2}} \frac{\partial^{2} \theta}{\partial t^{2}} \tag{6}
\end{equation*}
$$

where

$$
\begin{equation*}
\alpha=\sqrt{\frac{\lambda+2 \mu}{\rho}}=\sqrt{\frac{k+\frac{4}{3} \mu}{\rho}} \tag{7}
\end{equation*}
$$

It shows that the changes in the cubical dilation $\theta$ propagates through the elastic isotropic solid with speed $\alpha$.

Here, k is the modulus of compressibility.
Taking curl of (5) both sides, we write

$$
\begin{equation*}
(\lambda+\mu) \operatorname{curl}(\operatorname{grad} \theta)+\mu \nabla^{2}(\operatorname{curl} \bar{u})=\rho \frac{\partial^{2}}{\partial t^{2}}(\operatorname{curl} \bar{u}) \tag{8}
\end{equation*}
$$

From, vector calculus, we have the identity

$$
\begin{equation*}
\text { curl } \operatorname{grad} \phi \equiv 0 \tag{9}
\end{equation*}
$$

Using (4) and (9) , equation (8) gives

$$
\begin{equation*}
\nabla^{2} \bar{\Omega}=\frac{1}{\beta^{2}} \frac{\partial^{2} \bar{\Omega}}{\partial t^{2}}, \tag{10}
\end{equation*}
$$

where

$$
\begin{equation*}
\beta=\sqrt{\frac{\mu}{\rho}} . \tag{11}
\end{equation*}
$$

Equation (10) shows that changes in rotation $\bar{\Omega}$ propagates with speed $\beta$.
Remark 1: The speed of both the waves depend upon the elastic parameters $\lambda$, $\mu$ and the density $\rho$ of the medium. Since $\lambda>0 \& \mu>0$, it can be seen that

$$
\alpha>\beta
$$

That is, the dilatational waves propagate faster than the shear waves.
Therefore, dilatational waves arrive first while rotational waves arrive after that on a seismogram. For this reason, dilatational waves are also called primary waves and rotational waves are called secondary waves.

Remark 2: For a Poisson's solid $(\lambda=\mu)$, we have

$$
\alpha=\sqrt{3} \beta .
$$

For most solids , particularly rocks in Earth, there is a small difference between $\lambda \& \mu$. So, we may take $\lambda=\mu$ and solid is then called a Poisson's solid.

Remark 3: In seismology, the dilatational waves are called P-waves and rotational waves are denoted by $\mathbf{S}$ - waves.

Remark 4: If $\mu=0$, then $\beta=0$. That is, there is no $S$ - wave in a medium with zero rigidity.

That is, in liquids, $\mathbf{S}$ - waves cann't exist. However, P - waves exists in a liquid medium.

Remark 5: Since

$$
\operatorname{div} \bar{\Omega}=\operatorname{div}(\operatorname{curl} \bar{u}) \equiv 0
$$

it follows that a rotational wave is free of expansion/compression of volume.
For this reason , the rotational wave $\bar{\Omega}$ is also called equivoluminal/dilatationless.

Remark 6: The dilatational wave $(\theta \neq 0)$ causes a change in volume of the material elements in the body. Rotational wave (when $\bar{\Omega} \neq 0$ ) produces a change in shape of the material element without changes in the volume of material elements.

## Rotational waves are also referred as shear waves or a wave of distortion.

Remark 7: For a typical metal like copper, the speeds of primary and secondary and waves are estimated as
$\alpha=4.36 \mathrm{~km} / \mathrm{sec}$ and $\beta=2.13 \mathrm{~km} / \mathrm{sec}$,
respectively.
Remark 8: At points far away from initially disturbed region, the waves are plane waves.

This suits seismology because the recording station of a disturbance during an earthquake is placed at a great distance in comparison to the dimensions of initial source.

Remark 9: We note that equations (6) and (10) both are forms of wave equations. Equation (6) shows that a disturbance $\theta$, called dilatation wave/compressive wave propagates through the elastic medium with velocity $\alpha$. Similarly, equation (10) shows that a disturbance $\bar{\Omega}$, called a rotational wave, propagates through the elastic medium with velocity $\beta$.

Thus, we conclude that any disturbance in an infinite homogeneous isotropic elastic medium can be propagated in the form of two types of these waves.

The speed $\alpha$ depends upon rigidity $\mu$ and modulus of compressibility k. On the other hand , $\beta$ depends upon rigidity $\mu$ only.

## Helmholtz's Theorem (P and S wave of Seismology)

Any vector point function $\overline{\mathrm{F}}$ which is finite, uniform and continuous and which vanishes at infinity, may be expressed as the sum of a gradient of a scalar function $\phi$ and curl of a zero divergence vector $\psi$.

## The function $\phi$ is called the scalar potential of $\overline{\mathbf{F}}$ and $\bar{\psi}$ is called the vector potential of $\overline{\mathbf{F}}$.

The equation of motion for an elastic isotropic solid with density $\rho$ (for zero body force) is

$$
\begin{equation*}
(\lambda+2 \mu) \nabla(\nabla \cdot \overline{\mathbf{u}})-\mu \nabla \times \nabla \times \overline{\mathrm{u}}=\rho \frac{\partial^{2} \bar{u}}{\partial t^{2}} . \tag{1}
\end{equation*}
$$

We write the decomposition for the displacement vector $u$ in the form

$$
\begin{equation*}
\mathbf{u}=\nabla \phi+\nabla \times \bar{\psi}, \tag{2}
\end{equation*}
$$

with

$$
\begin{equation*}
\operatorname{div} \bar{\psi}=0 \tag{3}
\end{equation*}
$$

The scalar point function $\phi$ and vector point function $\bar{\psi}$ are called Lame' potentials.

We know the following vector identities.

$$
\begin{align*}
& \nabla \times \nabla \phi \equiv \overline{0}  \tag{4}\\
& \nabla \cdot(\nabla \times \bar{\psi}) \equiv 0 .  \tag{5}\\
& \begin{aligned}
\nabla \cdot \overline{\mathbf{u}} & =\nabla \cdot(\nabla \phi+\nabla \times \bar{\psi}) \\
& =\nabla \cdot \nabla \phi \\
& =\nabla^{2} \phi
\end{aligned}
\end{align*}
$$

Now
and

$$
\begin{align*}
\nabla \times \overline{\mathrm{u}} & =\nabla \times(\nabla \phi+\nabla \times \bar{\psi}) \\
& =\nabla \times \nabla \times \bar{\psi} \\
& =\operatorname{grad}(\operatorname{div} \bar{\psi})-\nabla^{2} \bar{\psi} \\
& =-\nabla^{2} \bar{\psi} . \tag{7}
\end{align*}
$$

Using (6) and (7) in (1), we write
$(\lambda+2 \mu) \nabla\left\{\nabla^{2} \phi\right\}+\mu \nabla \times\left\{\nabla^{2} \bar{\psi}\right\}=\rho \frac{\partial^{2}}{\partial t^{2}}(\nabla \phi)+\rho \frac{\partial^{2}}{\partial t^{2}}(\nabla \times \bar{\psi})$
or
$(\lambda+2 \mu) \nabla\left[\nabla^{2} \phi-\left(\frac{\rho}{\lambda+2 \mu}\right) \frac{\partial^{2} \phi}{\partial t^{2}}\right]+\mu \nabla \times\left\{\left[\nabla^{2} \bar{\psi}-\left(\frac{\rho}{\mu}\right) \frac{\partial^{2} \bar{\psi}}{\partial t^{2}}\right]=\overline{0}\right.$,
which is satisfied if we take $\phi$ and $\Psi$ to be solutions of wave equations

$$
\begin{align*}
& \nabla^{2} \phi=\frac{1}{\alpha^{2}} \frac{\partial^{2} \phi}{\partial t^{2}}  \tag{9}\\
& \nabla^{2}-\bar{\psi}=\frac{1}{\beta^{2}} \frac{\partial^{2} \bar{\psi}}{\partial t^{2}} \tag{10}
\end{align*}
$$

where

$$
\begin{align*}
& \alpha=\sqrt{\frac{\lambda+2 \mu}{\rho}}  \tag{11}\\
& \beta=\sqrt{\frac{\mu}{\rho}} \tag{12}
\end{align*}
$$

Wave equation (9) is a Scalar wave equation and equation (10) is a vector wave equation.

Waves represented by (9) are P-waves while waves represented by (10) are Swaves.
$\phi$ and $\bar{\psi}$ are now called the scalar and vector potentials associated with Pand S-waves, respectively.

Note : We can write the displacement vector $\overline{\mathrm{u}}$ as

$$
\begin{equation*}
\overline{\mathrm{u}}=\overline{\mathrm{u}}_{\mathrm{P}}+\overline{\mathrm{u}}_{\mathrm{S}}, \tag{1}
\end{equation*}
$$

where

$$
\begin{align*}
& \overline{\mathrm{u}}_{\mathrm{P}}=\nabla \phi,  \tag{2}\\
& \overline{\mathrm{u}}_{\mathrm{S}}=\operatorname{curl} \quad \bar{\psi} . \tag{3}
\end{align*}
$$

Hence $\overline{\mathrm{u}}_{\mathrm{P}}$ is the displacement due to P-wave alone and $\overline{\mathrm{u}}_{\mathrm{S}}$ is the displacement due to S -wave alone.

## One - Dimensional Waves

(a) P-waves : Consider the solution where $\bar{\psi}=0$. Then the one-dimensional P -wave satisfies the equation

$$
\begin{equation*}
\frac{\partial^{2} \phi}{\partial x^{2}}=\frac{1}{\alpha^{2}} \frac{\partial^{2} \phi}{\partial t^{2}} \quad, \quad \alpha=\sqrt{\frac{\lambda+2 \mu}{\rho}}, \tag{1}
\end{equation*}
$$

whose solution is

$$
\begin{equation*}
\phi=\phi(x \pm \alpha \mathrm{t}) \tag{2}
\end{equation*}
$$

representing a wave in x -direction. The corresponding displacement vector is

$$
\overline{\mathrm{u}}_{\mathrm{P}}=\nabla \phi
$$

$$
\begin{equation*}
=\left(\frac{\partial \phi}{\partial x}\right) \hat{e}_{1} . \tag{3}
\end{equation*}
$$

That is, the displacement vector $\bar{u}_{P}$ is in the direction of propagation of P wave. Therefore, P -waves are longitudinal waves.

Since

$$
\begin{align*}
\operatorname{curl} \overline{\mathrm{u}}_{\mathrm{P}} & =\operatorname{curl}(\operatorname{grad} \phi) \\
& =0 \tag{4}
\end{align*}
$$

So , P-waves are irrotational/rotationless. P-waves do not cause any rotation of the material particles of the medium.

Since

$$
\begin{align*}
\operatorname{div} \overline{\mathrm{u}}_{\mathrm{P}} & =\operatorname{div}(\operatorname{grad} \phi) \\
& =\nabla^{2} \phi \\
& \neq 0 \tag{5}
\end{align*}
$$

therefore, P-waves are dilatational/compressional.
In this case, wavefronts are planes,

$$
\mathrm{x}=\mathrm{constant} .
$$

The particle motion, for P -waves, is perpendicular to the wave fronts(figure).

## If we look for a free oscillation of angular velocity $\omega$, then we take

$$
\begin{equation*}
\phi(x, t)=e^{i \omega t} \phi(x), \tag{6}
\end{equation*}
$$

where

Here,

$$
\begin{equation*}
\frac{\mathrm{d}^{2} \phi}{\mathrm{dx}^{2}}+\mathrm{h}^{2} \phi=0 \tag{7}
\end{equation*}
$$

$$
\begin{equation*}
\mathrm{h}=\frac{\omega}{\alpha} \tag{8}
\end{equation*}
$$

is the wave number.
(b) S-waves: The one-dimensional wave equation for $S$-wave is

$$
\begin{equation*}
\frac{\partial^{2} \bar{\psi}}{\partial \mathrm{x}^{2}}=\frac{1}{\beta^{2}} \frac{\partial^{2} \bar{\psi}}{\partial \mathrm{t}^{2}} \tag{1}
\end{equation*}
$$

where

$$
\begin{equation*}
\beta=\sqrt{\frac{\mu}{\rho}} \tag{2}
\end{equation*}
$$

The corresponding displacement vector is

$$
\begin{equation*}
\overline{\mathrm{u}}_{\mathrm{S}}=\operatorname{curl} \bar{\psi} . \tag{3}
\end{equation*}
$$

Solution of (1) is of the type

$$
\begin{equation*}
\bar{\psi}(\mathrm{x}, \mathrm{t})=\psi_{1}(\mathrm{x}-\beta \mathrm{t}) \hat{e}_{1}+\psi_{2}(\mathrm{x}-\beta \mathrm{t}) \hat{e}_{2}+\psi_{3}(\mathrm{x}-\beta \mathrm{t}) \hat{e}_{3} \tag{4}
\end{equation*}
$$

From (3) and (4), we find

$$
\begin{equation*}
\overline{\mathrm{u}}_{\mathrm{S}}=-\frac{\partial \psi_{3}}{\partial \mathrm{x}} \hat{e}_{2}+\frac{\partial \psi_{2}}{\partial x} \hat{e}_{3} . \tag{5}
\end{equation*}
$$

This shows that the displacement vector $\overline{\mathrm{u}}_{\mathrm{S}}$ lies in a plane parallel to yz-plane, which is perpendicular to the x -direction, representing the direction of S -wave propagation.

The wavefronts are planes having x -axis their normals.
Thus, particle motion due to a $S$-wave is parallel to the wavefronts.
So, S-wave are transverse waves (figure).


Since

$$
\begin{align*}
\operatorname{div} \overline{\mathrm{u}}_{\mathrm{S}} & =\operatorname{div}\left(\operatorname{curl} \bar{\psi}^{\prime}\right) \\
& =0 \quad, \tag{6}
\end{align*}
$$

so, $S$-waves are dilationless/equivoluminal.

## Since

$$
\begin{align*}
\operatorname{curl} \overline{\mathrm{u}}_{\mathrm{s}} & =\operatorname{curl} \operatorname{curl} \bar{\psi}_{\mathrm{s}} \\
& =-\nabla^{2} \bar{\psi} \\
& \neq 0 \quad, \tag{7}
\end{align*}
$$

so, S -waves are not irrotational, they are rotational.
As strains are not zero, S-waves are shear/distortional waves.
Each component $\psi_{\mathrm{i}}$ satisfy the scalar wave equation for S -waves.

### 9.2. SV- AND SH-WAVES

We shall now consider the study of waves as related to Earth, for example , waves generated by earthquakes. We consider the surface of the earth (taken as plane , approximately) as horizontal. Let z-axis be taken vertically downwards and xy-plane as horizontal.

To determine $x$ - and $y$-directions in the horizontal plane, we proceed as below.

Let $\overline{\mathrm{p}}$ be the direction of propagation of a S-wave. Let the plane made by the vector $\bar{p}$ and the $z$-axis be the xz-plane. Then $x$-axis lies in the horizontal plane $(\mathrm{z}=0)$ bounding the earth and the propagation vector $\overline{\mathrm{p}}$ lies in the vertical xzplane (figure).

(A S-wave)
We choose y -axis as the direction perpendicular to the xz-plane so that x -, $y$ - and $z$-axis form a right handed system.

The displacement vector $\overline{\mathrm{u}}_{\mathrm{S}}$ corresponding to a S-wave propagating in the $\overline{\mathrm{p}}$ direction, is perpendicular to $\overline{\mathrm{p}}$-direction.

We resolve $\overline{\mathrm{u}}_{\mathrm{s}}$ into two components - the first component $\overline{\mathrm{u}}_{\mathrm{SV}}$ lying in the vertical xz-plane and the second component $\bar{u}_{S H}$ parallel to $y$-axis(i.e. $\perp$ to xzplane).

A S-wave representing the motion corresponding to the first component of the displacement vector is known as a SV-wave.

For a SV-wave, the particle motion is perpendicular to the direction $\overline{\mathrm{p}}$ of propagation of wave and lies in the vertical xz-plane which is normal to the horizontal bounding surface. Let $\mathrm{u}_{\mathrm{Sv}}$ denote the corresponding displacement vector. Then, we write

$$
\overline{\mathrm{u}}_{\mathrm{SV}}=(\mathrm{u}, 0, \mathrm{w}),
$$

contains both horizontal component $u$ in the $x$-direction and vertical component w in the z -direction. A $S V$-wave is a vertically polarized shear wave. SV stands for vertical shear.

A S-wave representing the second displacement component, parallel to the y -axis, lying in the horizontal plane, is known as SH-wave. Let $\bar{u}_{\text {SH }}$ denote the corresponding displacement. Then

$$
\bar{u}_{S H}=(0, \mathrm{v}, 0),
$$

is parallel to $y$-axis. The motion due to a SH -wave is perpendicular to
p-direction (i.e., in a transverse direction) and along a horizontal direction. A SH-wave is a horizontally polarized shear wave. SH stands for horizontal shear.

When a P-wave propagates in the $\overline{\mathrm{p}}$-direction, then the corresponding displacement, denoted by $\bar{u}_{P}$, is given by

$$
\bar{u}_{P}=(u, 0, w) .
$$

The displacement $\bar{u}_{P}$ contains both horizontal component in the $x$-direction and vertical component in the $z$-direction(figure)

(A P-wave)

### 9.3. WAVE PROPAGATION IN TWO DIMENSIONS

We assume that waves are propagating in planes parallel to the vertical xzplane containing the wave propagation vector

$$
\hat{p}=l \hat{e}_{1}+\mathrm{n} \hat{e}_{2}
$$

with

$$
1^{2}+n^{2}=1
$$

Then the wave motion will be the same in all planes parallel to xz-plane (figure) and independent of y so that $\frac{\partial}{\partial y} \equiv 0$.


Under this assumption, the Navier equation of motion for isotropic elastic materials

$$
\begin{equation*}
(\lambda+\mu) \operatorname{grad} \operatorname{div} \overline{\mathrm{u}}+\mu \nabla^{2} \overline{\mathrm{u}}=\rho \ddot{\overrightarrow{\mathrm{u}}} \tag{1}
\end{equation*}
$$

gives

$$
\begin{gather*}
(\lambda+\mu) \frac{\partial}{\partial x}\left(\frac{\partial u}{\partial x}+\frac{\partial w}{\partial z}\right)+\mu \nabla^{2} \mathrm{u}=\rho \ddot{u}  \tag{2}\\
(\lambda+\mu) \frac{\partial}{\partial z}\left(\frac{\partial u}{\partial z}+\frac{\partial w}{\partial x}\right)+\mu \nabla^{2} \mathrm{w}=\rho \ddot{w}  \tag{3}\\
\mu \nabla^{2} \mathrm{v}=\rho \ddot{\mathrm{v}} \tag{4}
\end{gather*}
$$

where the displacement vector $\bar{u}=(u, v, w)$ is independent of $y$,

$$
\begin{equation*}
\text { cubical dilatation } \equiv \theta=\frac{\partial u}{\partial x}+\frac{\partial w}{\partial z} \tag{5}
\end{equation*}
$$

and Laplacian is given by

$$
\begin{equation*}
\nabla^{2}=\frac{\partial^{2}}{\partial x^{2}}+\frac{\partial^{2}}{\partial z^{2}} . \tag{6}
\end{equation*}
$$

Let $\mathrm{z}=0$ be the boundary surface. The components of the stress acting on the surface $\mathrm{z}=0$ are given by

$$
\begin{align*}
& \tau_{\mathrm{zx}}=2 \mu \mathrm{e}_{\mathrm{zx}}=\mu\left(\frac{\partial u}{\partial z}+\frac{\partial w}{\partial x}\right)  \tag{7}\\
& \tau_{\mathrm{zz}}=\lambda \theta+2 \mu \mathrm{e}_{\mathrm{zz}}=(\lambda+2 \mu) \frac{\partial w}{\partial z}+\lambda \frac{\partial u}{\partial x}  \tag{8}\\
& \tau_{\mathrm{zy}}=2 \mu \mathrm{e}_{\mathrm{zy}}=\mu \frac{\partial v}{\partial z} \tag{9}
\end{align*}
$$

From above equations, we conclude that the general two - dimensional problem of propagation of plane elastic waves, parallel to the xz-plane, splits into two independent problems stated below -

Problem I - Consisting of equations (4) and (9).
Problem II - Consisting of equations (2), (3) (7) and (8).
Problem I (SH-problem) : In this problem the displacement components are

$$
\begin{equation*}
u=w \equiv 0, \quad, \mathrm{v}=\mathrm{v}(\mathrm{x}, \mathrm{z}, \mathrm{t}) \tag{10}
\end{equation*}
$$

and

$$
\begin{equation*}
\tau_{\mathrm{zx}}=\tau_{\mathrm{zz}}=0 \quad, \tau_{\mathrm{zy}}=\tau_{\mathrm{zy}}(\mathrm{x}, \mathrm{z}, \mathrm{t})=\mu \frac{\partial \mathrm{v}}{\partial \mathrm{z}} \tag{11}
\end{equation*}
$$

In this problem, the displacement component $\mathrm{v}(\mathrm{x}, \mathrm{z}, \mathrm{t})$ satisfies the scalar wave equation

$$
\begin{equation*}
\nabla^{2} v=\frac{1}{\beta^{2}} \frac{\partial^{2} v}{\partial t^{2}} \tag{12}
\end{equation*}
$$

This differential equation is independent of modulus of compressibility k. So the motion due to such waves is equivoluminal. Since the displacement component v is horizontal and is perpendicular to the direction $\hat{p}$ of propagation ( $\hat{p}$ lies in xz-plane), the waves represented by v are horizontally polarized waves or SH -waves.

Problem II (P-SV problem) : In this problem

$$
\begin{align*}
& \mathrm{v} \equiv 0, \mathrm{u}=\mathrm{u}(\mathrm{x}, \mathrm{z}, \mathrm{t}), \mathrm{w}=\mathrm{w}(\mathrm{x}, \mathrm{z}, \mathrm{t})  \tag{13}\\
& \tau_{\mathrm{zy}} \equiv 0, \tau_{\mathrm{zx}}=\mu\left(\frac{\partial u}{\partial z}+\frac{\partial w}{\partial x}\right), \tau_{\mathrm{zz}}=(\lambda+2 \mu) \frac{\partial w}{\partial z}+\lambda \frac{\partial u}{\partial x} \\
& \tau_{\mathrm{xx}}=\lambda\left(\frac{\partial u}{\partial x}+\frac{\partial w}{\partial z}\right)+2 \mu \frac{\partial u}{\partial x} \tag{14}
\end{align*}
$$

In this problem, the in-plane displacement components $u$ and $w$ satisfying the two simultaneous partial differential equations (2) and (3), and the in-plane stress components $\tau_{\mathrm{xx}}, \tau_{\mathrm{zx}}$ and $\tau_{\mathrm{zz}}$ also contain these two displacement components. The displacement component v plays no role in the solution of this problem and hence taken as identically zero.

The displacement components $u$ and $w$ can be expressed in terms of two scalar potentials

$$
\phi=\phi(\mathrm{x}, \mathrm{z}, \mathrm{t}) \text { and } \psi=\psi(\mathrm{x}, \mathrm{z}, \mathrm{t})
$$

through the relations

$$
\begin{align*}
& u=\frac{\partial \phi}{\partial x}+\frac{\partial \psi}{\partial z}  \tag{15}\\
& w=\frac{\partial \phi}{\partial z}-\frac{\partial \psi}{\partial x} \tag{16}
\end{align*}
$$

with help of Helmholtz's theorem (on taking $\bar{\psi}=-\psi \hat{e}_{2}$ ).
Using (15) and (16) in the equations of motion (2) and (3), we obtain

$$
\begin{align*}
& \frac{\partial}{\partial x}\left[\alpha^{2} \nabla^{2} \phi-\frac{\partial^{2} \phi}{\partial t^{2}}\right]+\frac{\partial}{\partial z}\left[\beta^{2} \nabla^{2} \psi-\frac{\partial^{2} \psi}{\partial t^{2}}\right]=0  \tag{17}\\
& \frac{\partial}{\partial z}\left[\alpha^{2} \nabla^{2} \phi-\frac{\partial^{2} \phi}{\partial t^{2}}\right]-\frac{\partial}{\partial x}\left[\beta^{2} \nabla^{2} \psi-\frac{\partial^{2} \psi}{\partial t^{2}}\right]=0 \tag{18}
\end{align*}
$$

These equations are identically satisfied when the scalar potentials $\phi$ and $\psi$ are solutions of following scalar wave equations

$$
\begin{align*}
& \nabla^{2} \phi=\frac{1}{\alpha^{2}} \frac{\partial^{2} \phi}{\partial t^{2}}  \tag{19}\\
& \nabla^{2} \psi=\frac{1}{\beta^{2}} \frac{\partial^{2} \psi}{\partial t^{2}} \tag{20}
\end{align*}
$$

where

$$
\begin{equation*}
\alpha=\sqrt{\frac{(\lambda+2 \mu)}{\rho}} \quad, \quad \beta=\sqrt{\frac{\mu}{\rho}} . \tag{21}
\end{equation*}
$$

The plane wave solutions of (19) represents P-waves and those of (20) represent S-waves.

The potentials $\phi$ and $\psi$ are called displacement potentials.
In the displacements (15) and (16), the contribution from $\phi$ is due to P -waves and that from $\psi$ is due S waves.

P-SV wave is a combination of P-wave and the SV-wave. The displacement vector ( $u, 0$, w) lies in a vertical plane and $S$-waves represented by the potential $\psi$ are also propagated in a vertical plane, so these $S$-waves are vertically polarized shear waves or simple SV-waves.

The stress $\tau_{\mathrm{zx}}$ and $\tau_{\mathrm{zz}}$ in terms of potentials $\phi$ and $\psi$ are

$$
\begin{align*}
\tau_{\mathrm{zx}} & =\mu\left(\frac{\partial u}{\partial z}+\frac{\partial w}{\partial x}\right) \\
& =\mu\left(2 \frac{\partial^{2} \phi}{\partial x \partial z}+\frac{\partial^{2} \psi}{\partial z^{2}}-\frac{\partial^{2} \psi}{\partial x^{2}}\right),  \tag{22}\\
\tau_{\mathrm{zz}} & =\lambda \operatorname{div} \overline{\mathrm{u}}+2 \mu \cdot \frac{\partial w}{\partial z} \\
& =\lambda \nabla^{2} \phi+2 \mu\left(\frac{\partial^{2} \phi}{\partial z^{2}}-\frac{\partial^{2} \psi}{\partial x \partial z}\right) \\
& =\left(\lambda \nabla^{2}+2 \mu \frac{\partial^{2}}{\partial z^{2}}\right) \phi-2 \mu \frac{\partial^{2} \psi}{\partial x \partial z} . \tag{23}
\end{align*}
$$

Since,

$$
\begin{equation*}
\frac{\lambda}{\mu}=\frac{\lambda+2 \mu-2 \mu}{\mu}=\frac{\lambda+2 \mu}{\mu}-2=\frac{\alpha^{2}}{\beta^{2}}-2 \tag{24}
\end{equation*}
$$

therefore, we write

$$
\begin{equation*}
\tau_{z z}=\mu\left[\left(\frac{\alpha^{2}}{\beta^{2}}-2\right) \nabla^{2} \phi+2 \frac{\partial^{2} \phi}{\partial z^{2}}-2 \frac{\partial^{2} \psi}{\partial x \partial z}\right] . \tag{25}
\end{equation*}
$$

Note 1: In the SH-problem, other stresses are

$$
\begin{equation*}
\tau_{\mathrm{xx}}=\tau_{\mathrm{yy}}=0, \quad \tau_{\mathrm{xy}}=\mu \frac{\partial \mathrm{v}}{\partial \mathrm{x}} \tag{26}
\end{equation*}
$$

Note 2: In the P-SV problem, other stresses are (in terms of $\phi$ an $\psi$ )

$$
\begin{align*}
\tau_{\mathrm{xx}} & =\lambda \nabla^{2} \phi+2 \mu \mathrm{e}_{\mathrm{xx}} \\
& =\lambda \nabla^{2} \phi+2 \mu\left(\frac{\partial^{2} \phi}{\partial \mathrm{x}^{2}}+\frac{\partial^{2} \psi}{\partial \mathrm{x} \partial \mathrm{z}}\right),  \tag{27}\\
\tau_{\mathrm{yy}} & =\lambda \nabla^{2} \phi+2 \mu \mathrm{e}_{\mathrm{yy}} \\
& =\lambda \nabla^{2} \phi  \tag{28}\\
\tau_{\mathrm{xy}} & =2 \mu \mathrm{e}_{\mathrm{xy}}=0 \tag{29}
\end{align*}
$$

Note 3: In the case of two - dimensional wave propagation, the SH motion is decoupled from the P-SV motion. The displacement vector due to P-SV type motion is

$$
\begin{aligned}
\overline{\mathrm{u}} & =\nabla \phi+\nabla \times\left(\psi \hat{e}_{2}\right) \\
& =(\mathrm{u}, \mathrm{o}, \mathrm{w}) .
\end{aligned}
$$

### 9.4. PLANE WAVES

A geometric surface of all points in space over which the phase of a wave is constant is called a wavefronts.

Wavefronts can have many shapes. For example, wavefronts can be planes or spheres or cylinders.

A line normal to the wave fronts, indicating the direction of motion of wave, is called a ray.

If the waves are propagated in a single direction, the waves are called plane waves, and the wavefronts for plane waves are parallel planes with normal along the direction of propagation of the wave.

Thus, a plane wave is a solution of the wave equation in which the disturbance/displacement varies only in the direction of wave propagation and is constant in all the directions orthogonal to propagation direction.

The rays for plane waves are parallel straight lines.

For spherical waves, the disturbance is propagated out in all directions from a point source of waves. The wavefronts are concentric spheres and the rays are radial lines leaving the point source in all directions.

Since seismic energy is usually radiated from localized sources , seismic wavefronts are always curved to some extent. However, at sufficiently large distances from the source, the wavefronts become flat enough that a plane wave approximation become locally valid.

Illustration (1) : For one -dimensional wave equation

$$
\begin{equation*}
\frac{\partial^{2} \phi}{\partial x^{2}}=\frac{1}{c^{2}} \frac{\partial^{2} \phi}{\partial t^{2}}, \tag{1}
\end{equation*}
$$

a progressive wave travelling with speed c in the positive x -direction is represented by

$$
\begin{equation*}
\phi(\mathrm{x}, \mathrm{t})=\mathrm{A} \mathrm{e}^{\mathrm{ik}(\mathrm{x}-\mathrm{ct})}=\mathrm{A} \mathrm{e}^{\mathrm{i}(\mathrm{kx}-\mathrm{wt})} \tag{2}
\end{equation*}
$$

where

$$
\begin{aligned}
& \mathrm{k}=\text { wave number }, \\
& \omega=\mathrm{c} \mathrm{k}=\text { angular frequency }, \\
& \lambda=\text { wavelength }=2 \pi / \mathrm{k} \\
& \mathrm{~A}=\text { amplitude of the wave. }
\end{aligned}
$$

Let

$$
\begin{equation*}
\overline{\mathrm{x}}=\mathrm{x} \hat{e}_{1} \quad, \overline{\mathrm{k}}=\mathrm{k} \hat{e}_{1} . \tag{3}
\end{equation*}
$$

Then $\overline{\mathrm{k}}$ is called the propagation vector and (2) can be written in the form

$$
\begin{equation*}
\phi(\mathrm{x}, \mathrm{t})=\mathrm{A} \mathrm{e}^{\mathrm{i}(\overline{\mathrm{k}} \cdot \overline{\mathrm{x}}-\omega \mathrm{t})} . \tag{4}
\end{equation*}
$$

In this type of waves, wavefronts are planes

$$
\begin{equation*}
\mathrm{x}=\mathrm{constant}, \tag{5}
\end{equation*}
$$

which are parallel to yz - plane.
Illustration 2: A two - dimensional wave equation with speed c is

$$
\begin{equation*}
\frac{\partial^{2} \phi}{\partial x^{2}}+\frac{\partial^{2} \phi}{\partial y^{2}}=\frac{1}{c^{2}} \frac{\partial^{2} \phi}{\partial t^{2}} \tag{1}
\end{equation*}
$$

If we take

$$
\begin{align*}
& \mathrm{u}=l \mathrm{x}+\mathrm{my}-\mathrm{ct},  \tag{2}\\
& \mathrm{v}=l \mathrm{x}+\mathrm{my}+\mathrm{ct} \tag{3}
\end{align*}
$$

where $l, \mathrm{~m}$ are constants, then equation (1) is reduced to (exercise)

$$
\begin{equation*}
\frac{\partial^{2} \phi}{\partial \mathrm{u} \partial \mathrm{v}}=0 \tag{4}
\end{equation*}
$$

provided

$$
\begin{equation*}
l^{2}+\mathrm{m}^{2}=1 \tag{5}
\end{equation*}
$$

## Integrating (4), we get

$$
\begin{align*}
\phi & =\mathrm{f}(\mathrm{u})+\mathrm{g}(\mathrm{v}) \\
& =\mathrm{f}(l \mathrm{x}+\mathrm{my}-\mathrm{ct})+\mathrm{g}(l \mathrm{x}+\mathrm{my}+\mathrm{ct}) \tag{6}
\end{align*}
$$

where $f$ and $g$ are arbitrary functions.
Let

$$
\begin{equation*}
\hat{v}=l \hat{e}_{1}+\mathrm{m} \hat{e}_{2} \tag{7}
\end{equation*}
$$

Then $\hat{v}$ is a unit vector perpendicular to the system of straight lines

$$
\begin{equation*}
l \mathrm{x}+\mathrm{my}=\mathrm{constant} \tag{8}
\end{equation*}
$$

## in two-dimensional xy-plane.



At any instant, say $t=t_{0}$, the disturbance $\phi$ is constant for all points ( $x, y$ ) lying on the line (8). Therefore, $\phi$ represents a plane propagating with speed c in the direction $\hat{v}$. The wavefronts are straight lines given by equation (8).

Illustration 3: Three-dimensional wave equation is

$$
\begin{equation*}
\frac{\partial^{2} \phi}{\partial x^{2}}+\frac{\partial^{2} \phi}{\partial y^{2}}+\frac{\partial^{2} \phi}{\partial z^{2}}=\frac{1}{c^{2}} \frac{\partial^{2} \phi}{\partial t^{2}}, \tag{1}
\end{equation*}
$$

which propagates with speed c .
As discussed earlier, solution of (1) is

$$
\begin{equation*}
\phi(\mathrm{x}, \mathrm{y}, \mathrm{z}, \mathrm{t})=\mathrm{f}(l \mathrm{x}+\mathrm{my}+\mathrm{nz}-\mathrm{ct})+\mathrm{g}(l \mathrm{x}+\mathrm{my}+\mathrm{nz}+\mathrm{ct}) \tag{2}
\end{equation*}
$$

where the constants $l, \mathrm{~m}, \mathrm{n}$ satisfy the relation

$$
\begin{equation*}
l^{2}+\mathrm{m}^{2}+\mathrm{n}^{2}=1 . \tag{3}
\end{equation*}
$$

Let

$$
\begin{equation*}
\hat{v}=l \hat{e}_{1}+\mathrm{m} \hat{e}_{2}+\mathrm{n} \hat{e}_{3} . \tag{4}
\end{equation*}
$$

Then $\hat{v}$ is a unit vector which is normal to the system of parallel planes

$$
\begin{equation*}
l \mathrm{x}+\mathrm{my}+\mathrm{nz}=\text { constant } . \tag{5}
\end{equation*}
$$

The wavefronts are parallel planes given by equation (5), which travel with speed c in the direction $\hat{v}$. So, the wave (2) is a three-dimensional progressive plane wave.

Question : Determine the wavelength and velocity of a system of plane waves given by

$$
\phi=\mathrm{a} \sin (\mathrm{~A} x+\mathrm{By}+\mathrm{Cz}-\mathrm{D} \mathrm{t}),
$$

where $\mathrm{a}, \mathrm{A}, \mathrm{B}, \mathrm{C}, \mathrm{D}$ are constants.
Solution : Let $\langle l, \mathrm{~m}, \mathrm{n}\rangle$ be the direction cosines of the direction of wave propagation with speed c . Then

$$
\begin{equation*}
l^{2}+\mathrm{m}^{2}+\mathrm{n}^{2}=1 \tag{1}
\end{equation*}
$$

## Equation of a wave front is

$$
\begin{equation*}
l \mathrm{x}+\mathrm{my}+\mathrm{nz}=\mathrm{p} \tag{2}
\end{equation*}
$$

where $p$ is the length of perpendicular from the origin to the wave front.


$$
\langle l, \mathrm{~m}, \mathrm{n}\rangle
$$

## The plane wave $\phi$, therefore, must be of the type

$$
\begin{equation*}
\phi=\mathrm{a} \sin \mathrm{k}(l \mathrm{x}+\mathrm{my}+\mathrm{n} \mathrm{z}-\mathrm{ct}), \tag{3}
\end{equation*}
$$

where

$$
\begin{aligned}
& \mathrm{k}=\text { wave number and }, \\
& \mathrm{c}=\text { wave velocity } .
\end{aligned}
$$

Comparing (3) with the given form

$$
\begin{equation*}
\phi=\mathrm{a} \sin (\mathrm{Ax}+\mathrm{By}+\mathrm{Cz}-\mathrm{D} \mathrm{t}), \tag{4}
\end{equation*}
$$

we find
$l=\frac{A}{\sqrt{A^{2}+B^{2}+C^{2}}}, \quad \mathrm{~m}=\frac{B}{\sqrt{A^{2}+B^{2}+C^{2}}}, \quad \mathrm{n}=\frac{C}{\sqrt{A^{2}+B^{2}+C^{2}}}$,
and

$$
\begin{align*}
& \mathrm{k}=\sqrt{A^{2}+B^{2}+C^{2}},  \tag{6}\\
& \mathrm{c}=\frac{D}{\sqrt{A^{2}+B^{2}+C^{2}}} . \tag{7}
\end{align*}
$$

Therefore ,

$$
\begin{equation*}
\text { wave length }=\frac{2 \pi}{k}=\frac{2 \pi}{\sqrt{A^{2}+B^{2}+C^{2}}}, \tag{8}
\end{equation*}
$$

and

$$
\begin{equation*}
\text { wave velocity }=\frac{D}{\sqrt{A^{2}+B^{2}+C^{2}}} \tag{9}
\end{equation*}
$$

## Chapter-10

## Surface Waves

### 10.1. INTRODUCTION

In an elastic body, it is possible to have another type of waves (other than body waves) which are propagated over the surface and penetrate only a little into the interior of the body.

Such waves are similar to waves produced on a smooth surface of water when a stone is thrown into it.

These type of waves are called SURFACE WAVES.
Surface waves are "tied" to the surface and diminish exponentially as they get farther from the surface.

The criterion for surface waves is that the amplitude of the displacement in the medium dies exponentially with the increasing distance from the surface.

In seismology, the interfaces are, in the ideal case, horizontal and so the plane of incidence is vertical. Activity of surface waves is restricted to the neighbourhood of the interface(s) or surface of the medium.

Under certain conditions, such waves can propagated independently along the surface/interface. For surface waves, the disturbance is confined to a depth equal to a few wavelengths.

Let us take xz - plane as the plane of incidence with z - axis vertically downwards. Let $\mathrm{z}=0$ be the surface of a semi-infinite elastic medium (Figure)


For a surface wave, its amplitude will be a function of $z$ (rather than an exponential function) which tends to zero as $\mathrm{z} \rightarrow \infty$. For such surface - waves , the motion will be two - dimensional, parallel to $x z$ - plane, so that $\frac{\partial}{\partial y}=0$.

The existence of surface waves raises the question of whether they might (under certain conditions) be able to travel freely along the plane (horizontal) as a guided wave.

### 10.2. RAYLEIGH WAVES

Rayleigh (1885) discussed the existence of a simplest surface wave propagating on the free - surface of a homogeneous isotropic elastic half space.

Let the half - space occupies the region $\mathrm{z} \geq 0$ with z - axis taken as vertically downwards. Let $\rho$ be the density and $\lambda, \mu$ be the Lame' elastic moduli and $\mathrm{z}=$ 0 be the stress - free boundary of the half - space (figure).


Suppose that a train of plane waves is propagating in the media in the positive x - direction such that
(i) the plane of incidence is the vertical plane ( $x z$ - plane) so that the motion, i.e., disturbance is independent of y and hence $\frac{\partial}{\partial y} \equiv 0$.
(ii) the amplitude of the surface wave decreases exponentially as we move away from the surface in the z - direction.

This problem is a plane strain problem and the displacement vector $\bar{u}$ is of the type

$$
\begin{equation*}
\overline{\mathrm{u}}=(\mathrm{u}, 0, \mathrm{w}), \quad \mathrm{u}=\mathrm{u}(\mathrm{x}, \mathrm{z}, \mathrm{t}), \quad \mathrm{w}=\mathrm{w}(\mathrm{x}, \mathrm{z}, \mathrm{t}) . \tag{1}
\end{equation*}
$$

The displacement components $u$ and $w$ are given in term of potential $\phi$ and $\psi$ through the relations

$$
\begin{align*}
& \mathrm{u}=\frac{\partial \phi}{\partial x}+\frac{\partial \psi}{\partial z}  \tag{2}\\
& \mathrm{w}=\frac{\partial \phi}{\partial z}-\frac{\partial \psi}{\partial x} \tag{3}
\end{align*}
$$

where potentials $\phi$ and $\psi$ satisfy the scalar wave equations

$$
\begin{align*}
& \nabla^{2} \phi=\frac{\partial^{2} \phi}{\partial x^{2}}+\frac{\partial^{2} \phi}{\partial z^{2}}=\frac{1}{\alpha^{2}} \frac{\partial^{2} \phi}{\partial t^{2}}  \tag{4}\\
& \nabla^{2} \psi=\frac{\partial^{2} \psi}{\partial x^{2}}+\frac{\partial^{2} \psi}{\partial z^{2}}=\frac{1}{\beta^{2}} \frac{\partial^{2} \psi}{\partial t^{2}} \tag{5}
\end{align*}
$$

and wave velocities $\alpha$ and $\beta$ are given by

$$
\begin{align*}
& \alpha=\sqrt{(\lambda+2 \mu) / \rho}  \tag{6}\\
& \beta=\sqrt{\mu / \rho} \tag{7}
\end{align*}
$$

This wave motion is of a P-SV type wave travelling along the stress - free surface of the half space and such a motion takes place in the xz - plane.

Now, we seek solutions of wave equations (4) and (5) of the form

$$
\begin{equation*}
\phi(\mathrm{x}, \mathrm{z}, \mathrm{t})=\mathrm{f}(\mathrm{z}) \mathrm{e}^{\mathrm{ik}(\mathrm{ct}-\mathrm{x})} \tag{8}
\end{equation*}
$$

$$
\begin{equation*}
\psi(\mathrm{x}, \mathrm{z}, \mathrm{t})=\mathrm{g}(\mathrm{z}) \mathrm{e}^{\mathrm{ik}(\mathrm{ct}-\mathrm{x})} \tag{9}
\end{equation*}
$$

where k is the wave number, c is the speed of surface wave travelling in the + ve x -direction $\quad$ and $\omega=\mathrm{ck}$ is the angular frequency.

The substitution of (8) and (9) into wave equations (4) and (5) leads to two ordinary differential equations (exercise)

$$
\begin{equation*}
\frac{d^{2} f}{d z^{2}}-\mathrm{k}^{2} \mathrm{a}^{2} \mathrm{f}=0 \tag{10}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{d^{2} g}{d z^{2}}-\mathrm{k}^{2} \mathrm{~b}^{2} \mathrm{~g}=0 \tag{11}
\end{equation*}
$$

with
and

$$
\begin{align*}
& \mathrm{a}=\sqrt{1-c^{2} / \alpha^{2}}>0,  \tag{12}\\
& \mathrm{~b}=\sqrt{1-\mathrm{c}^{2} / \beta^{2}}>0 . \tag{13}
\end{align*}
$$

From equations (8) to (11); we find (exercise)

$$
\begin{align*}
& \phi(x, z, t)=\left(A e^{-a k z}+A_{1} e^{a k z}\right) e^{i k(c t-x)}  \tag{14}\\
& \psi(x, z, t)=\left(B e^{-b k z}+B_{1} e^{b k z}\right) e^{i k(c t-x)} \tag{15}
\end{align*}
$$

where $\mathrm{A}, \mathrm{A}_{1} \mathrm{~B}$ and $\mathrm{B}_{1}$ are constants.
Since the disturbance due to surface waves must die rapidly as $\mathrm{z} \rightarrow \infty$, we must have

$$
\begin{equation*}
\mathrm{A}_{1}=\mathrm{B}_{1}=0 \tag{16}
\end{equation*}
$$

Thus, for Rayleigh waves, we get

$$
\begin{equation*}
\phi(\mathrm{x}, \mathrm{z}, \mathrm{t})=\mathrm{Ae}^{-\mathrm{akz}} \mathrm{e}^{\mathrm{i}(\omega \mathrm{t}-\mathrm{kx})} \tag{17}
\end{equation*}
$$

and

$$
\begin{equation*}
\psi(\mathrm{x}, \mathrm{z}, \mathrm{t})=\mathrm{B}^{-\mathrm{bkz}} \mathrm{e}^{\mathrm{i}(\omega \mathrm{t}-\mathrm{kx})} \tag{18}
\end{equation*}
$$

with

$$
\begin{equation*}
c<\beta<\alpha \tag{19}
\end{equation*}
$$

The equation (19) gives the condition for the existence of surface waves on the surface of a semi - infinite isotropic elastic media with velocity c in the positive x - direction.

To find frequency equation for the velocity c and ratio of A and B ; we use the stress - free boundary conditions. This conditions are

$$
\begin{equation*}
\tau_{\mathrm{zx}}=\tau_{\mathrm{zy}}=\tau_{\mathrm{zz}}=0 \quad \text { at } \mathrm{z}=0 \tag{20}
\end{equation*}
$$

From relation in (2) and (3), we find

$$
\begin{align*}
\operatorname{div} \overline{\mathrm{u}} & =\frac{\partial u}{\partial x}+\frac{\partial w}{\partial z} \\
& =\frac{\partial^{2} \phi}{\partial x^{2}}+\frac{\partial^{2} \phi}{\partial z^{2}} \\
& =\nabla^{2} \phi \tag{21}
\end{align*}
$$

The stresses in terms of potentials $\phi$ and $\psi$ are given by

$$
\begin{gather*}
\tau_{\mathrm{zx}}=\mu\left(\frac{\partial u}{\partial z}+\frac{\partial w}{\partial x}\right) \\
=\mu\left[2 \frac{\partial^{2} \phi}{\partial \mathrm{x} \partial \mathrm{z}}+\frac{\partial \psi}{\partial \mathrm{z}^{2}}-\frac{\partial^{2} \psi}{\partial \mathrm{x}^{2}}\right]  \tag{22}\\
\tau_{\mathrm{zz}}=\lambda \operatorname{div} \overline{\mathrm{u}}+\mu \frac{\partial w}{\partial z} \\
=\mu\left[\left(\frac{\alpha^{2}}{\beta^{2}}-2\right) \frac{\partial^{2} \phi}{\partial x^{2}}++\frac{\alpha^{2}}{\beta^{2}} \frac{\partial^{2} \phi}{\partial z^{2}}-2 \frac{\partial^{2} \psi}{\partial x \partial z}\right],  \tag{23}\\
\tau_{\mathrm{zy}}=\mu\left(\frac{\partial \mathrm{v}}{\partial \mathrm{z}}+\frac{\partial \mathrm{w}}{\partial \mathrm{y}}\right) \\
=0 \tag{24}
\end{gather*}
$$

The boundary condition $\tau_{\mathrm{zx}}=0$ at $\mathrm{z}=0$ in (20), gives (exercise)
or

$$
2 A(-i k)(-a k)-B(-i k)^{2}+B(-b k)^{2}=0
$$

$$
2 i a k^{2} A+B\left(1+b^{2}\right) k^{2}=0
$$

or

$$
2 i a A=-B\left(1+b^{2}\right)
$$

or

$$
\begin{equation*}
\frac{A}{B}=\frac{i\left(1+b^{2}\right)}{2 a} \tag{25}
\end{equation*}
$$

The boundary condition $\tau_{z \mathrm{z}}=0$ at $\mathrm{z}=0$, gives (exercise)
or

$$
\left(\frac{\alpha^{2}}{\beta^{2}}-2\right) \cdot \mathrm{A} \cdot\left(-\mathrm{k}^{2}\right)+\frac{\alpha^{2}}{\beta^{2}} \cdot \mathrm{~A} \cdot\left(\alpha^{2} \mathrm{k}^{2}\right)-2 \mathrm{~B}(-\mathrm{ik})(-\mathrm{bk})=0
$$

$$
\mathrm{A}\left[-\frac{\alpha^{2}}{\beta^{2}}+2+\frac{\alpha^{2} a^{2}}{\beta^{2}}\right] k^{2}-2 \text { i b B k }{ }^{2}=0
$$

or

$$
\mathrm{A}\left(2-\frac{\alpha^{2}}{\beta^{2}}+\frac{\alpha^{2} a^{2}}{\beta^{2}}\right)=2 \mathrm{ibB}
$$

$$
\begin{equation*}
\frac{A}{B}=\left(\frac{2 i b}{2-\frac{\alpha^{2}}{\beta^{2}}+\frac{\alpha^{2} a^{2}}{\beta^{2}}}\right) \tag{26}
\end{equation*}
$$

Eliminating $A / B$ from equations (25) and (26) and substituting the value of a and $b$ from equation (12) and (13), we obtain
or

$$
\begin{aligned}
& \left(1+b^{2}\right)\left(2-\frac{\alpha^{2}}{\beta^{2}}+\frac{\alpha^{2} a^{2}}{\beta^{2}}\right)=4 \mathrm{ab} \\
& \left(2-\frac{c^{2}}{\beta^{2}}\right)\left[2-\frac{\alpha^{2}}{\beta^{2}}+\frac{\alpha^{2}}{\beta^{2}}\left(1-\frac{\mathrm{c}^{2}}{\alpha^{2}}\right)\right] \\
& =4 \sqrt{1-\frac{c^{2}}{\alpha^{2}}} \sqrt{1-\frac{c^{2}}{\beta^{2}}}
\end{aligned}
$$

or

$$
\begin{align*}
& \left(2-\frac{c^{2}}{\beta^{2}}\right)\left[2-\frac{\alpha^{2}}{\beta^{2}}+\frac{\alpha^{2}}{\beta^{2}}-\frac{c^{2}}{\beta^{2}}\right]=4 \sqrt{1-\frac{c^{2}}{\alpha^{2}}} \sqrt{1-\frac{c^{2}}{\beta^{2}}} \\
& \left(2-\frac{c^{2}}{\beta^{2}}\right)^{2}-4\left(1-\frac{c^{2}}{\alpha^{2}}\right)^{\frac{1}{2}}\left(1-\frac{c^{2}}{\beta^{2}}\right)^{\frac{1}{2}}=0 \tag{27}
\end{align*}
$$

This equation contains only one unknown c.
This equation determines the speed $\mathbf{c}$ for Rayleigh surface waves in an uniform half - space.

Equation (27) is called the RAYLEIGH EQUATION for Rayleigh waves. It is also called the Rayleigh frequency equation. It is also called the Rayleigh wave - velocity equation.

This equation is the period equation for Rayleigh waves.
In order to prove that Rayleigh waves really exists , we must show that frequency equation (27) has at least one real root for c. To show this, we proceed as follows :

From equation (27), we write
or

$$
\begin{aligned}
& \left(2-\frac{c^{2}}{\beta^{2}}\right)^{4}=16\left(1-\frac{c^{2}}{\alpha^{2}}\right)\left(1-\frac{c^{2}}{\beta^{2}}\right) \\
& 16-32 \frac{c^{2}}{\beta^{2}}+24 \frac{c^{4}}{\beta^{4}}-8 \frac{c^{6}}{\beta^{6}}+\frac{c^{8}}{\beta^{8}} \\
& =16\left[1-\frac{c^{2}}{\alpha^{2}}-\frac{c^{2}}{\beta^{2}}+\frac{c^{4}}{\alpha^{2} \beta^{2}}\right] \\
& =16\left[1-\frac{c^{2}}{\beta^{2}} \frac{\beta^{2}}{\alpha^{2}}-\frac{c^{2}}{\beta^{2}}+\frac{c^{4}}{\beta^{4}} \cdot \frac{\beta^{2}}{\alpha^{2}}\right]
\end{aligned}
$$

or
$-32+24 \frac{c^{2}}{\beta^{2}}-8 \frac{c^{4}}{\beta^{4}}+\frac{c^{6}}{\beta^{6}}+16 \frac{\beta^{2}}{\alpha^{2}}+16-16 \frac{c^{2}}{\beta^{2}} \cdot \frac{\beta^{2}}{\alpha^{2}}=0$.
Put

$$
\begin{equation*}
\mathrm{s}=\mathrm{c}^{2} / \beta^{2} \tag{29}
\end{equation*}
$$

then equation (28) gives

$$
\begin{equation*}
f(s)=s^{3}-8 s^{2}+\left(24-16 \frac{\beta^{2}}{\alpha^{2}}\right) s-16\left(1-\frac{\beta^{2}}{\alpha^{2}}\right)=0 . \tag{30}
\end{equation*}
$$

It is polynomial equation of degree $\mathbf{3}$ in $\mathbf{s}$ with real coefficients.
Hence the frequency equation (30) has at least one real root since complex roots occur in conjugate pairs.

Moreover

$$
\begin{equation*}
\mathrm{f}(0)=-16\left(1-\frac{\beta^{2}}{\alpha^{2}}\right)<0 \tag{31}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathrm{f}(1)=1-8+24-\frac{16 \beta^{2}}{\alpha^{2}}+\frac{16 \beta^{2}}{\alpha^{2}}=1>0 \tag{32}
\end{equation*}
$$

since $\beta<\alpha$.
Hence, by the intermediate theorem of calculus, the frequency equation (30) has at least one root in the range $(0,1)$.

That is, there exists $s \in(0,1)$ which is a root of the Rayleigh frequency equation.

Further, $0<\mathrm{s}<1$ implies

$$
0<\frac{c^{2}}{\beta^{2}}<1
$$

or

$$
\begin{equation*}
0<c<\beta . \tag{33}
\end{equation*}
$$

This establishes that Rayleigh surface waves propagating with speed $c=c_{R}$, where

$$
0<c_{R}<\beta<\alpha,
$$

always exists.
This proves the existence of Rayleigh surface waves.
Special Case (Poisson's solid) : When the semi-infinite elastic medium is a Poisson's solid,
then $\lambda=\mu$ and $\sigma=\frac{1}{4}$ and

$$
\begin{equation*}
\frac{\alpha^{2}}{\beta^{2}}=\frac{\lambda+2 \mu}{\mu}=3 \tag{34}
\end{equation*}
$$

For this type of elastic medium , frequency equation (30) becomes
or

$$
\begin{aligned}
& 3 s^{3}-24 s^{2}+56 s-32=0 \\
& (s-4)\left(3 s^{2}-12 s+8\right)=0
\end{aligned}
$$

giving

$$
\mathrm{s}=4,2+\frac{2}{\sqrt{3}}, 2-\frac{2}{\sqrt{3}}
$$

or

$$
\begin{equation*}
\frac{c^{2}}{\beta^{2}}=4,2+\frac{2}{\sqrt{3}}, 2-\frac{2}{\sqrt{3}} . \tag{35}
\end{equation*}
$$

But c must satisfy the inequality $0<\mathrm{c}<\beta$, therefore, the only possible value for $c=c_{R}$ is

$$
\frac{c_{R}^{2}}{\beta^{2}}=2-\frac{2}{\sqrt{3}}
$$

or

$$
\begin{equation*}
c_{\mathrm{R}}=\left[\sqrt{2-\frac{2}{\sqrt{3}}}\right] \cdot \beta \cong 0.9195 \beta \tag{36}
\end{equation*}
$$

giving the speed of propagation of Rayleigh surface waves along the stress free boundary of the Poissonian half-space in the x-direction.

Thus, the order of arrival of P, SV-and Rayleigh waves is


From equations (12) and (13) ; we write

$$
\begin{aligned}
\mathrm{a}^{2} & =1-\frac{c_{R}^{2}}{\alpha^{2}} \\
& =1-\frac{\mathrm{c}_{\mathrm{R}}{ }^{2}}{\beta^{2}} \cdot \frac{\beta^{2}}{\alpha^{2}} \\
& =\frac{1}{9}(3+2 \sqrt{3}),
\end{aligned}
$$

or

$$
\begin{equation*}
\mathrm{a}=\frac{1}{3} \sqrt{3+2 \sqrt{3}} \cong 0.8475 \approx 0.848 \tag{37}
\end{equation*}
$$

and

$$
\begin{aligned}
\mathrm{b}^{2} & =1-\frac{c_{R}^{2}}{\beta^{2}} \\
& =\frac{1}{3}[2 \sqrt{3}-3],
\end{aligned}
$$

or

$$
\begin{equation*}
\mathrm{b} \cong 0.3933 \cong 0.393 . \tag{38}
\end{equation*}
$$

Further , from equations (25), (37) and (38), we find

$$
\begin{equation*}
\frac{A}{B}=1.468 \mathrm{i} \tag{39}
\end{equation*}
$$

## Remark 1: Displacements due to Rayleigh waves (Particle motion)

## From equations (2) , (3) , (17) , (18) and (26) ; we find

$$
\begin{align*}
u(x, z, t) & =\left[A(-i k) e^{-a k z}+B(-b k) e^{-b k z}\right] \mathrm{e}^{\mathrm{i}(\omega t-k x)} \\
& =\mathrm{k}\left[-\mathrm{i} A \mathrm{e}^{-\mathrm{akz}}-\mathrm{A}\left(\frac{2-c^{2} / \beta^{2}}{2 i}\right) \mathrm{e}^{-\mathrm{bkz}}\right] \mathrm{e}^{\mathrm{i}(\omega \mathrm{t}-\mathrm{kx})} \\
& =-\mathrm{ikA}\left[\mathrm{e}^{-\mathrm{akz}}-\left(1-\frac{c^{2}}{2 \beta^{2}}\right) \mathrm{e}^{-\mathrm{bkz}}\right] \mathrm{e}^{\mathrm{i}(\omega \mathrm{t}-\mathrm{kx})} . \tag{40}
\end{align*}
$$

Similarly, we shall find

$$
\begin{align*}
& \mathrm{w}(\mathrm{x}, \mathrm{z}, \mathrm{t})=\left[(-\mathrm{ak}) \mathrm{Ae}^{-\mathrm{akz}}-(-\mathrm{ik}) \mathrm{B} \mathrm{e}^{-\mathrm{bkz}}\right] \mathrm{e}^{\mathrm{i}(\omega \mathrm{t}-\mathrm{kx})} \\
= & \mathrm{k}\left[-\mathrm{a} \mathrm{~A}^{-\mathrm{akz}}+\frac{i e^{b \mathrm{kz}}(-2 i A a)}{\left(2-c^{2} / \beta^{2}\right)}\right] \mathrm{e}^{\mathrm{i}(\omega \mathrm{t}-\mathrm{kx})} \\
= & \mathrm{kA}\left[-\mathrm{a}^{-\mathrm{akz}}+\left(\frac{a}{1-c^{2} / 2 \beta^{2}}\right) \mathrm{e}^{-\mathrm{bkz}}\right] \mathrm{e}^{\mathrm{i}(\omega \mathrm{t}-\mathrm{kx})} . \tag{41}
\end{align*}
$$

Then , taking the real parts of equations (40) and (41) and using equation (42) to (44) ; we find

$$
\begin{align*}
& \mathrm{u}(\mathrm{x}, \mathrm{z}, \mathrm{t})=\operatorname{AkU}(\mathrm{z}) \sin \theta,  \tag{45}\\
& \mathrm{w}(\mathrm{x}, \mathrm{z}, \mathrm{t})=\operatorname{AkW}(\mathrm{z}) \cos \theta, \tag{46}
\end{align*}
$$

we remember that A is the potential amplitude. Eliminating $\theta$ from equations (45) and (46) ,

> we obtain

$$
\begin{equation*}
\frac{u^{2}}{[A k U(z)]^{2}}+\frac{w^{2}}{[A k W(z)]^{2}}=1 \tag{47}
\end{equation*}
$$

which is an equation of an ellipse in the vertical xz-plane.
Equation (47) shows that particles, during the propagation of Rayleigh Surface Waves, describe ellipses.

Remark 2: Particle Motion at the surface ( $\mathrm{z}=0$ )

On the surface $z=0$,we find, at $z=0$,

$$
\mathrm{U}(0)=\frac{c^{2}}{2 \beta^{2}} \quad, \quad \mathrm{~W}(0)=\mathrm{a}\left(\frac{\frac{c^{2}}{2 \beta^{2}}}{1-c^{2} / 2 \beta^{2}}\right)
$$

Since $0<c<\beta$, so

$$
\begin{equation*}
\mathrm{U}(0)>0 \quad, \quad \mathrm{~W}(0)>0 \tag{49}
\end{equation*}
$$

Let

$$
\begin{equation*}
\mathrm{a}_{1}=\mathrm{AkW}(0), \quad \mathrm{b}_{1}=\operatorname{AkU}(0) . \tag{50}
\end{equation*}
$$

Then

$$
\begin{equation*}
a_{1}>0, b_{1}>0, \tag{51}
\end{equation*}
$$

and equation (47) reduces to

$$
\begin{equation*}
\frac{u^{2}}{b_{1}^{2}}+\frac{w^{2}}{a_{1}^{2}}=1 \tag{52}
\end{equation*}
$$

## Remark 3: Particular case (Poisson's Solid) :

In this case, we find, at $\mathrm{z}=0$,

$$
\begin{equation*}
\frac{a_{1}}{b_{1}}=\frac{W(0)}{U(0)}=\sqrt{3} \mathrm{a} \cong 1.5 \tag{53}
\end{equation*}
$$

using (37). So

$$
\begin{equation*}
a_{1}>b_{1} \tag{54}
\end{equation*}
$$

Thus, the surface particle motion ( $\mathrm{on} \mathrm{z}=0$ ) is an ellipse with a vertical major axis.

The horizontal and vertical displacement components are out of phase by $\frac{\pi}{2}$.
The resulting surface particle motion is Retrograde (opposite to that of wave propagation).

Remark 4: The dependencies of the displacement components $u$ and $w$ on depth(z) are given by equation (43) to (46) .

There is a value of z for which $\mathrm{u}=0$ (For Poisson's solid, at $\mathrm{z}=+0.19 \lambda, \lambda=$ $\left.\frac{2 \pi}{k}, \mathrm{u}=0\right)$, whereas w is never zero.

At the depth, where $u$ is zero, its amplitude changes sign.
For greater depths, the particle motion is Prograde.
With increasing depth , the amplitudes of $u$ and $w$ decrease exponentially , with $w$ always larger than $u$.

Thus, the elliptic motion changes from retrograde at the surface to prograde at depth, passing through a node at which there is no horizontal motion.

So , for the propagation of Rayleigh surface waves , a surface particle describes an ellipse, about its mean position, in the retrograde sense.


Particle motion for the fundamental Rayleigh mode in a uniform half-space, propagating from left to right. One horizontal wavelength ( $\wedge$ ) is shown; the dots are plotted at a fixed time point. Motion is counter clockwise (retrograde) at the surface, changing to purely vertical motion at a depth of about $N 5$, and becoming clockwise (prograde) at greater depths. Note that the time behavior at a fixed distance is given by looking from right to left in this plot.

Remark 5: We see that frequency equation (27) for Rayleigh surface waves is independent of $\omega$. Therefore, the velocity $c_{R}$ of Rayleigh surface waves is constant and fixed.

This phenomenon is called nondispessive.
That is , Rayleigh waves are undispresed.
Remark 6: Maximum displacement parallel to the direction of Rayleigh waves

$$
\begin{aligned}
& =(\mathrm{u})_{\max .} \\
& =\mathrm{b}_{1} \\
& =\frac{2}{3} \mathrm{a}_{1} \quad, \quad \text { for a Poisson solid. } \\
& =\text { two-third of the maximum displacement in the } \\
\text { vertical direction } & \\
& \text { for a Poisson solid. }
\end{aligned}
$$

Note (1) : Rayleigh waves are important because the largest disturbances caused by an earthquake recorded on a distant seismogram are usually those of Rayleigh waves.

GROUND ROLL is the term commonly used for Rayleigh waves.
Note (2) : Although a "a free surface" means contact with a vacuum , the elastic constants and density of air are so lows in comparison with values for rocks that the surface of the earth is approximately a free surface.

Note (3): The boundary conditions $\tau_{z x}=\tau_{z z}=0$ at $z=0$ require that these two conditions must be satisfied, and so we require two parameters than can be adjusted. Therefore, we assume that both P-and SV-components exist and adjust their amplitude to satisfy the boundary conditions.

Exercise : Show that the displacement components at the surface of an elastic Poisson solid due to Rayleigh waves are

$$
\begin{aligned}
& u(x, t)=-0.423 k A \sin k\left(x-c_{R} t\right) \\
& w(x, t)=0.620 k A \cos k\left(x-c_{R} t\right), \quad v \equiv 0,
\end{aligned}
$$

with usual notation.

### 10.3. SURFACE WAVES OF SH-TYPE IN A HALF-SPACE

We consider first the possibility of the propagation of SH type surface waves (called Love waves) in a homogeneous semi-infinite isotropic elastic medium occupying the half-space $\mathrm{z} \geq 0$. The horizontal boundary $\mathrm{z}=0$ of the medium is assumed to be stress free. Let $\rho$ be the density of the medium and $\lambda, \mu$ Lame' constants(figure) .

(Elastic isotropic half-space)
Let the two - dimensional SH-wave motion takes place in the xz-plane. The basic equations for SH - wave motion are

$$
\begin{align*}
& u=w=0, \quad v=v(x, z, t)  \tag{1}\\
& \frac{\partial^{2} v}{\partial x^{2}}+\frac{\partial^{2} v}{\partial z^{2}}=\frac{1}{\beta^{2}} \frac{\partial^{2} v}{\partial t^{2}}  \tag{2}\\
& \beta^{2}=\frac{\mu}{\rho} \tag{3}
\end{align*}
$$

We try a plane wave solution of wave equation (2) of the form

$$
\begin{equation*}
v(x, z, t)=B \cdot e^{-b k z} \cdot e^{i(\omega t-k x)} \tag{4}
\end{equation*}
$$

where $\omega$ is the angular frequency of wave, $\mathrm{k}=$ wave number and $\mathrm{c}=\omega / \mathrm{k}$ is the speed with which surface waves are travelling in the x -direction on the surface $\mathrm{z}=0 ; \mathrm{b}>0$ and B is an arbitrary constant.

The amplitude of surface wave is $B e^{-b k z}$ which die exponentially as $z$ increases.

Substituting the value of $v(x, z, t)$ from equation (4) into equation (2), we find

$$
\begin{equation*}
\mathrm{b}^{2}=1-\frac{c^{2}}{\beta^{2}} . \tag{5}
\end{equation*}
$$

Since $\mathrm{b}>0$,so

$$
\begin{equation*}
c<\beta \tag{6}
\end{equation*}
$$

Using the stress-displacement relations, we find

$$
\begin{align*}
& \tau_{31}=\tau_{33}=0, \\
& \tau_{32}=\mu \frac{\partial \mathrm{v}}{\partial \mathrm{z}}=-\mu \mathrm{bk} \cdot \mathrm{~B} \cdot \mathrm{e}^{\mathrm{i}(\omega \mathrm{t}-\mathrm{kx})} \cdot \mathrm{e}^{-\mathrm{bkz}} \tag{7}
\end{align*}
$$

Hence, the stress-free condition of the boundary $\mathrm{z}=0$ implies that, using (7),

$$
\mu \mathrm{bkB} \mathrm{e}{ }^{\mathrm{i}(\omega t-\mathrm{kx})}=0
$$

or

$$
\begin{equation*}
\mathrm{B}=0, \tag{8}
\end{equation*}
$$

as $\mu \neq 0, b \neq 0, k \neq 0$.
This implies that in the case of a homogeneous isotropic elastic half space, Love waves do not exist at all.

### 10.4. PROPAGATION OF LOVE WAVES

Surface waves of the SH-type are observed to occur on the earth's surface. Love (1911) showed that if the earth is modelled as an isotropic elastic layer of finite thickness lying over a homogeneous elastic isotropic halfspace (rather than considering earth as a purely uniform half-space) then SH type waves occur in the stress-free surface of a layered half-space.

Now, we consider the possibility of propagation of surface waves of SH-type (Love waves) in a semi-infinite elastic isotropic medium consisting of a horizontal elastic layer of uniform thickness H lying over a half-space.

It is assumed that two elastic isotropic media are welded together and the horizontal boundary $\mathrm{z}=0$ of the semi-infinite medium is stress - free (see , figure).


$$
\beta_{1}<c_{L}<\beta_{2}
$$

Z
Let the layer and the half-space have different densities $\rho_{1}, \rho_{2}$ and different shear moduli $\mu_{1}, \mu_{2}$ respectively. Let two-dimensional SH-motion takes place parallel to xz-plane. The basic equation for SH -wave motion are

$$
\begin{align*}
& \mathrm{u}_{1}=\mathrm{w}_{1} \equiv 0, \mathrm{v}_{1}=\mathrm{v}_{1}(\mathrm{x}, \mathrm{z}, \mathrm{t})  \tag{1}\\
& \frac{\partial^{2} \mathrm{v}_{1}}{\partial \mathrm{x}^{2}}+\frac{\partial^{2} \mathrm{v}_{1}}{\partial \mathrm{z}^{2}}=\frac{1}{\beta_{1}^{2}} \frac{\partial^{2} \mathrm{v}_{1}}{\partial \mathrm{t}^{2}}  \tag{2}\\
& \beta_{1}^{2}=\frac{\mu_{1}}{\rho_{1}} \tag{3}
\end{align*}
$$

in the layer $0 \leq \mathrm{z} \leq \mathrm{H}$, and

$$
\begin{align*}
& u_{2}=w_{2}=0, \quad v_{2}=v_{2}(x, z, t)  \tag{4}\\
& \frac{\partial^{2} v_{2}}{\partial x^{2}}+\frac{\partial^{2} v_{2}}{\partial \mathrm{z}^{2}}=\frac{1}{\beta_{2}{ }^{2}} \frac{\partial^{2} \mathrm{v}_{2}}{\partial \mathrm{t}^{2}}  \tag{5}\\
& \beta_{2}^{2}=\frac{\mu_{2}}{\rho_{2}} \tag{6}
\end{align*}
$$

in the half-space $(\mathrm{z} \geq \mathrm{H})$.
Suitable plane wave solutions of wave equations (2) and (5) are (exercise), as discussed in detail already ,

$$
\begin{equation*}
\mathrm{v}_{1}(\mathrm{x}, \mathrm{z}, \mathrm{t})=\left(\mathrm{A}_{1} e^{-b_{1} k z}+\mathrm{B}_{1} e^{b_{1} k z}\right) \mathrm{e}^{\mathrm{i}(\omega \mathrm{t}-\mathrm{kx})} \tag{7}
\end{equation*}
$$

in the layer $0 \leq \mathrm{z} \leq \mathrm{H}$, and

$$
\begin{equation*}
\mathrm{v}_{2}(\mathrm{x}, \mathrm{z}, \mathrm{t})=\mathrm{A}_{2} e^{-b_{2} k z} \mathrm{e}^{\mathrm{i}(\omega \mathrm{t}-\mathrm{kx})} \tag{8}
\end{equation*}
$$

in the half-space $(\mathrm{z} \geq \mathrm{H}) . \omega$ is the angular wave frequency and k is the wave number, and $c=\omega / k$ is the speed of propagation of surface wave (if it exists) in the positive $x$-direction. $A_{1}, B_{1}$ are constants, $\mathrm{b}_{1}$ and $\mathrm{b}_{2}$ are real numbers with $\mathrm{b}_{2}>0$. However, $\mathrm{b}_{1}$ is unrestricted because z is finite in the layer.

Substituting for $\mathrm{v}_{2}$ from (8) into (5) yields the relation

$$
\begin{equation*}
\mathrm{b}_{2}^{2}=\left(1-\frac{c^{2}}{\beta_{2}^{2}}\right) \tag{9}
\end{equation*}
$$

and, therefore ,

$$
\begin{equation*}
c<\beta_{2}, \tag{10}
\end{equation*}
$$

otherwise

$$
\mathrm{v}_{2} \rightarrow \infty \text { as } \mathrm{z} \rightarrow \infty .
$$

From equations (7) and (2), we find

$$
\begin{equation*}
\mathrm{b}_{1}^{2}=\left(1-\frac{c^{2}}{\beta_{1}^{2}}\right) \tag{11}
\end{equation*}
$$

The stress-displacement relations imply

$$
\begin{equation*}
\tau_{31}=\tau_{33} \equiv 0 \tag{12}
\end{equation*}
$$

in the layer as well as in the half space. Also

$$
\begin{equation*}
\tau_{32}=\mu_{1} \mathrm{k}\left(-\mathrm{b}_{1} \mathrm{~A}_{1} e^{-b_{1} k z}+\mathrm{b}_{1} \mathrm{~B}_{1} e^{b_{1} k z}\right) \mathrm{e}^{\mathrm{i}(\omega \mathrm{t}-\mathrm{kx})} \tag{13}
\end{equation*}
$$

in the layer $0 \leq \mathrm{z} \leq \mathrm{H}$, and

$$
\begin{equation*}
\tau_{32}=\mathrm{k} \mu_{2} \mathrm{~A}_{2}\left(-\mathrm{b}_{2}\right) e^{-b_{2} k z} \cdot \mathrm{e}^{\mathrm{i}(\omega t-\mathrm{kx})} \tag{14}
\end{equation*}
$$

in the half space $\mathrm{z} \geq \mathrm{H}$.
The stress-free boundary $\mathrm{z}=0$ implies that

$$
\begin{equation*}
\tau_{32}=0, \tag{15}
\end{equation*}
$$

at $\mathrm{z}=0$. This gives

$$
\begin{equation*}
\mathrm{B}_{1}=\mathrm{A}_{1} . \tag{16}
\end{equation*}
$$

Since, there is a welded contact between the layer and the half-space at the interface $\mathrm{z}=\mathrm{H}$, so the displacement and the tractions must be continuous across the interface $\mathrm{z}=\mathrm{H}$.

Thus, the boundary conditions at $\mathrm{z}=\mathrm{H}$ are

$$
\begin{align*}
& \mathrm{v}_{1}=\mathrm{v}_{2}  \tag{17}\\
& \left.\tau_{32}\right|_{\text {layer }}=\left.\tau_{32}\right|_{\text {half-space }} . \tag{18}
\end{align*}
$$

From equations (13), (14) , (17) and (18), we find

$$
\begin{equation*}
\mathrm{A}_{1} e^{-k b_{1} H}+\mathrm{B}_{1} e^{k k_{1} H}=\mathrm{A}_{2} e^{-k b_{2} H}, \tag{19}
\end{equation*}
$$

and

$$
\begin{equation*}
\mu_{1} \mathrm{k}_{1}\left[-\mathrm{A}_{1} e^{-b_{1} k H}+\mathrm{B}_{1} e^{b_{1} k H}\right]=-\mu_{2} \mathrm{~b}_{2} \mathrm{k} e^{-b_{2} k H} . \tag{20}
\end{equation*}
$$

Equation (16), (19) and (20) are three homogeneous in $\mathrm{A}_{2}, \mathrm{~A}_{1}, \mathrm{~B}_{1}$ and $\mathrm{A}_{2}$.
We shall now eliminate them. From equations (19) \& (20), we write (after putting $\mathrm{B}_{1}=\mathrm{A}_{1}$ )

$$
\frac{A_{1}\left[e^{-k b_{1} H}+e^{k b_{1} H}\right]}{A_{1} \cdot \mu_{1} b_{1}\left(e^{-b_{1} k H}-e^{b_{1} k H}\right) k}=\frac{1}{\mu_{2} b_{2} k}
$$

or

$$
\frac{e^{-b_{1} k H}-e^{b_{1} k H}}{e^{-k b_{1} H}+e^{k b_{1} H}}=\frac{\mu_{2} b_{2}}{\mu_{1} b_{1}}
$$

or

$$
\begin{equation*}
\tan \mathrm{h}\left(\mathrm{~b}_{1} \mathrm{kH}\right)+\frac{\mu_{2} b_{2}}{\mu_{1} b_{1}}=0 \tag{21}
\end{equation*}
$$

Equation (21) is known as period equation/frequency equation/Dispersion equation for surface Love waves.

Equation (21) can also be written as

$$
\begin{align*}
& \tan h\left[k h . \sqrt{1-\frac{c^{2}}{\beta_{1}{ }^{2}}}\right]=-\frac{\mu_{2}}{\mu_{1}} \cdot \frac{\sqrt{1-\frac{c^{2}}{\beta_{2}{ }^{2}}}}{\sqrt{1-\frac{c^{2}}{\beta_{1}{ }^{2}}}} \\
& \tanh \left[\frac{\omega \mathrm{H}}{\mathrm{c}} \sqrt{1-\frac{\mathrm{c}^{2}}{\beta_{1}{ }^{2}}}\right]=-\frac{\mu_{2}}{\mu_{1}} \cdot \frac{\sqrt{1-\frac{c^{2}}{\beta_{2}{ }^{2}}}}{\sqrt{1-\frac{c^{2}}{\beta_{1}{ }^{2}}}} \tag{22}
\end{align*}
$$

Equation (22) is a transcendental equation.

For given $\omega$, we can find the speed c for surface Love waves. We note the value of $\mathbf{c}$ depends upon $\omega$. This means that waves of different frequencies will, in general , propagate with different phase velocity.

## This phenomenon is known as dispersion.

It is caused by the inhomogeneity of the medium (layered medium) due to some abrupt discontinuities within the medium (or due to continuous change of the elastic parameters which is not the present case).

Thus, Love waves are dispressed.
We consider now the following two possibilities between c and $\beta_{1}$.
(i) Either
$\mathrm{c} \leq \beta_{1}$,
(ii) or $\mathrm{c}>\beta_{1}$.

When $\mathrm{c} \leq \beta_{1}$ : In this case $\mathrm{b}_{1}$ is real (see, equation (11)) and left side of (22) becomes real and positive and right side of (22) is real and negative. Therefore , equation (22) can not possess any real solution for c .

Therefore, in this case, Love waves do not exist.
So , for the existence of surface Love waves, we must have

$$
\begin{equation*}
\mathbf{c}>\beta_{1} . \tag{24}
\end{equation*}
$$

In this case (24), $b_{1}$ is purely imaginary and we may write

$$
\begin{equation*}
\mathrm{b}_{1}==\sqrt{1-\frac{c^{2}}{\beta_{1}^{2}}}=\mathrm{i}\left(\sqrt{\frac{c^{2}}{\beta_{1}}-1}\right) . \tag{25}
\end{equation*}
$$

## Then equation (22) becomes

$$
\begin{equation*}
\tan \left(k H \sqrt{\frac{c^{2}}{\beta_{1}{ }^{2}}-1}\right)=\frac{\mu_{2}}{\mu_{1}}\left(\frac{\sqrt{1-\frac{c^{2}}{\beta_{2}{ }^{2}}}}{\sqrt{\frac{c^{2}}{\beta_{1}{ }^{2}}-1}}\right) . \tag{26}
\end{equation*}
$$

From equations (10) and (24) ; we find

$$
\begin{equation*}
\beta_{1}<c<\beta_{2} . \tag{27}
\end{equation*}
$$

Equation (26) is a transcendental equation that yields infinitely many roots for c.

Thus, the possible speeds of the Love waves are precisely the roots of equations (26) that lie in the interval ( $\beta_{1}, \beta_{2}$ ).

This indicates that the shear velocity in the layer must be less than the shear velocity in the half-space for the possible existence of Love waves.

This gives the upper and lower bounds for the speed of Love waves.
Remark 1: If the layer and the half - space are such that $\beta_{1} \leq \beta$, then existence of Love waves are not possible

Remark 2: In the limiting case when the layer is absent, we have

$$
\mu=\mu_{1} \text { and } \rho=\rho_{1}
$$

and therefore

$$
\beta=\beta_{1}
$$

Equation (22) leads to the impossible condition

$$
0=-1 .
$$

Hence, in this case, the wave considered can not exist.
Remark 3: When k or $\omega \rightarrow 0$, we get $\mathrm{c} \rightarrow \beta_{1}$.
The dispersion curve is given in the following figure.


## Here, if we assume

$$
\mu_{1} / \mu=1.8, \beta=3.6 \mathrm{~km} / \mathrm{sec}, \beta_{1}=4.6 \mathrm{~km} / \mathrm{sec}
$$

then

$$
c_{L}=\text { speed of Love waves }=4.0 \mathrm{~km} / \mathrm{sec} .
$$

## SURFACE WAVES

### 10.1. INTRODUCTION

In an elastic body, it is possible to have another type of waves (other than body waves) which are propagated over the surface and penetrate only a little into the interior of the body.

Such waves are similar to waves produced on a smooth surface of water when a stone is thrown into it.

These type of waves are called SURFACE WAVES.

## Surface waves are "tied" to the surface and diminish exponentially as they get farther from the surface.

The criterion for surface waves is that the amplitude of the displacement in the medium dies exponentially with the increasing distance from the surface.

In seismology, the interfaces are, in the ideal case, horizontal and so the plane of incidence is vertical. Activity of surface waves is restricted to the neighbourhood of the interface(s) or surface of the medium.

Under certain conditions, such waves can propagated independently along the surface/interface. For surface waves, the disturbance is confined to a depth equal to a few wavelengths.

Let us take xz - plane as the plane of incidence with z - axis vertically downwards. Let $\mathrm{z}=0$ be the surface of a semi-infinite elastic medium (Figure)


For a surface wave, its amplitude will be a function of z (rather than an exponential function) which tends to zero as $\mathrm{z} \rightarrow \infty$. For such surface - waves , the motion will be two - dimensional, parallel to $\mathrm{xz}-$ plane, so that $\frac{\partial}{\partial y}=0$.

The existence of surface waves raises the question of whether they might (under certain conditions) be able to travel freely along the plane (horizontal) as a guided wave.

### 10.2. RAYLEIGH WAVES

## Rayleigh (1885) discussed the existence of a simplest surface wave propagating on the free - surface of a homogeneous isotropic elastic half space.

Let the half - space occupies the region $\mathrm{z} \geq 0$ with z - axis taken as vertically downwards. Let $\rho$ be the density and $\lambda, \mu$ be the Lame' elastic moduli and $\mathrm{z}=$ 0 be the stress - free boundary of the half - space (figure).


Suppose that a train of plane waves is propagating in the media in the positive x - direction such that
(i) the plane of incidence is the vertical plane ( xz - plane) so that the motion, i.e., disturbance is independent of y and hence $\frac{\partial}{\partial y} \equiv 0$.
(ii) the amplitude of the surface wave decreases exponentially as we move away from the surface in the z - direction.

This problem is a plane strain problem and the displacement vector $\bar{u}$ is of the type

$$
\begin{equation*}
\overline{\mathrm{u}}=(\mathrm{u}, 0, \mathrm{w}), \quad \mathrm{u}=\mathrm{u}(\mathrm{x}, \mathrm{z}, \mathrm{t}), \quad \mathrm{w}=\mathrm{w}(\mathrm{x}, \mathrm{z}, \mathrm{t}) . \tag{1}
\end{equation*}
$$

The displacement components $u$ and $w$ are given in term of potential $\phi$ and $\psi$ through the relations

$$
\begin{align*}
& \mathrm{u}=\frac{\partial \phi}{\partial x}+\frac{\partial \psi}{\partial z}  \tag{2}\\
& \mathrm{w}=\frac{\partial \phi}{\partial z}-\frac{\partial \psi}{\partial x} \tag{3}
\end{align*}
$$

where potentials $\phi$ and $\psi$ satisfy the scalar wave equations

$$
\begin{align*}
& \nabla^{2} \phi=\frac{\partial^{2} \phi}{\partial x^{2}}+\frac{\partial^{2} \phi}{\partial z^{2}}=\frac{1}{\alpha^{2}} \frac{\partial^{2} \phi}{\partial t^{2}}  \tag{4}\\
& \nabla^{2} \psi=\frac{\partial^{2} \psi}{\partial x^{2}}+\frac{\partial^{2} \psi}{\partial z^{2}}=\frac{1}{\beta^{2}} \frac{\partial^{2} \psi}{\partial t^{2}} \tag{5}
\end{align*}
$$

and wave velocities $\alpha$ and $\beta$ are given by

$$
\begin{align*}
& \alpha=\sqrt{(\lambda+2 \mu) / \rho}  \tag{6}\\
& \beta=\sqrt{\mu / \rho} \tag{7}
\end{align*}
$$

This wave motion is of a P-SV type wave travelling along the stress - free surface of the half space and such a motion takes place in the xz - plane.

Now, we seek solutions of wave equations (4) and (5) of the form

$$
\begin{align*}
& \phi(x, z, t)=f(z) e^{i k(c t-x)}  \tag{8}\\
& \psi(x, z, t)=g(z) e^{i k(c t-x)} \tag{9}
\end{align*}
$$

where k is the wave number, c is the speed of surface wave travelling in the +ve x -direction $\quad$ and $\omega=\mathrm{ck}$ is the angular frequency.

The substitution of (8) and (9) into wave equations (4) and (5) leads to two ordinary differential equations (exercise)

$$
\begin{equation*}
\frac{d^{2} f}{d z^{2}}-\mathrm{k}^{2} \mathrm{a}^{2} \mathrm{f}=0 \tag{10}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{d^{2} g}{d z^{2}}-\mathrm{k}^{2} \mathrm{~b}^{2} \mathrm{~g}=0 \tag{11}
\end{equation*}
$$

with

$$
\begin{equation*}
\mathrm{a}=\sqrt{1-c^{2} / \alpha^{2}}>0 \tag{12}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathrm{b}=\sqrt{1-\mathrm{c}^{2} / \beta^{2}}>0 \tag{13}
\end{equation*}
$$

From equations (8) to (11) ; we find (exercise)

$$
\begin{align*}
& \phi(x, z, t)=\left(A e^{-a k z}+A_{1} e^{a k z}\right) e^{i k(c t-x)}  \tag{14}\\
& \psi(x, z, t)=\left(B e^{-b k z}+B_{1} e^{b k z}\right) e^{i k(c t-x)} \tag{15}
\end{align*}
$$

where $\mathrm{A}, \mathrm{A}_{1} \mathrm{~B}$ and $\mathrm{B}_{1}$ are constants.
Since the disturbance due to surface waves must die rapidly as $\mathrm{z} \rightarrow \infty$, we must have

$$
\begin{equation*}
\mathrm{A}_{1}=\mathrm{B}_{1}=0 . \tag{16}
\end{equation*}
$$

Thus, for Rayleigh waves, we get

$$
\begin{equation*}
\phi(\mathrm{x}, \mathrm{z}, \mathrm{t})=\mathrm{Ae}^{-\mathrm{akz}} \mathrm{e}^{\mathrm{i}(\omega \mathrm{t}-\mathrm{kx})} \tag{17}
\end{equation*}
$$

and

$$
\begin{equation*}
\psi(\mathrm{x}, \mathrm{z}, \mathrm{t})=\mathrm{B}^{-\mathrm{bkz}} \mathrm{e}^{\mathrm{i}(\omega \mathrm{t}-\mathrm{kx})}, \tag{18}
\end{equation*}
$$

with

$$
\begin{equation*}
c<\beta<\alpha . \tag{19}
\end{equation*}
$$

The equation (19) gives the condition for the existence of surface waves on the surface of a semi - infinite isotropic elastic media with velocity $c$ in the positive x - direction.

To find frequency equation for the velocity c and ratio of A and B ; we use the stress - free boundary conditions. This conditions are

$$
\begin{equation*}
\tau_{\mathrm{zx}}=\tau_{\mathrm{zy}}=\tau_{\mathrm{zz}}=0 \quad \text { at } \mathrm{z}=0 \tag{20}
\end{equation*}
$$

From relation in (2) and (3), we find

$$
\begin{align*}
\operatorname{div} \overline{\mathrm{u}} & =\frac{\partial u}{\partial x}+\frac{\partial w}{\partial z} \\
& =\frac{\partial^{2} \phi}{\partial x^{2}}+\frac{\partial^{2} \phi}{\partial z^{2}} \\
& =\nabla^{2} \phi \tag{21}
\end{align*}
$$

The stresses in terms of potentials $\phi$ and $\psi$ are given by

$$
\begin{align*}
& \tau_{\mathrm{zx}}=\mu\left(\frac{\partial u}{\partial z}+\frac{\partial w}{\partial x}\right) \\
&=\mu\left[2 \frac{\partial^{2} \phi}{\partial \mathrm{x} \partial \mathrm{z}}+\frac{\partial \psi}{\partial \mathrm{z}^{2}}-\frac{\partial^{2} \psi}{\partial \mathrm{x}^{2}}\right],  \tag{22}\\
& \tau_{\mathrm{zz}}=\lambda \operatorname{div} \overline{\mathrm{u}}+\mu \frac{\partial w}{\partial z} \\
&=\mu\left[\left(\frac{\alpha^{2}}{\beta^{2}}-2\right) \frac{\partial^{2} \phi}{\partial x^{2}}++\frac{\alpha^{2}}{\beta^{2}} \frac{\partial^{2} \phi}{\partial z^{2}}-2 \frac{\partial^{2} \psi}{\partial x \partial z}\right],  \tag{23}\\
& \tau_{\mathrm{zy}}=\mu\left(\frac{\partial \mathrm{v}}{\partial \mathrm{z}}+\frac{\partial \mathrm{w}}{\partial \mathrm{y}}\right) \\
&=0 . \tag{24}
\end{align*}
$$

The boundary condition $\tau_{\mathrm{zx}}=0$ at $\mathrm{z}=0$ in (20), gives (exercise)

$$
2 A(-i k)(-a k)-B(-i k)^{2}+B(-b k)^{2}=0
$$

or

$$
2 i a k^{2} A+B\left(1+b^{2}\right) k^{2}=0
$$

or

$$
2 i a A=-B\left(1+b^{2}\right)
$$

or

$$
\begin{equation*}
\frac{A}{B}=\frac{i\left(1+b^{2}\right)}{2 a} \tag{25}
\end{equation*}
$$

The boundary condition $\tau_{z z}=0$ at $\mathrm{z}=0$, gives (exercise)

$$
\left(\frac{\alpha^{2}}{\beta^{2}}-2\right) \cdot \mathrm{A} \cdot\left(-\mathrm{k}^{2}\right)+\frac{\alpha^{2}}{\beta^{2}} \cdot \mathrm{~A} \cdot\left(\alpha^{2} \mathrm{k}^{2}\right)-2 \mathrm{~B}(-\mathrm{ik})(-\mathrm{bk})=
$$

0
or

$$
\mathrm{A}\left[-\frac{\alpha^{2}}{\beta^{2}}+2+\frac{\alpha^{2} a^{2}}{\beta^{2}}\right] k^{2}-2 \text { i b B k }{ }^{2}=0
$$

or

$$
\mathrm{A}\left(2-\frac{\alpha^{2}}{\beta^{2}}+\frac{\alpha^{2} a^{2}}{\beta^{2}}\right)=2 \mathrm{ibB}
$$

or

$$
\begin{equation*}
\frac{A}{B}=\left(\frac{2 i b}{2-\frac{\alpha^{2}}{\beta^{2}}+\frac{\alpha^{2} a^{2}}{\beta^{2}}}\right) \tag{26}
\end{equation*}
$$

## Eliminating A/B from equations (25) and (26) and substituting the value of $a$ and $b$ from equation (12) and (13), we obtain

$$
\begin{aligned}
& \left(1+b^{2}\right)\left(2-\frac{\alpha^{2}}{\beta^{2}}+\frac{\alpha^{2} a^{2}}{\beta^{2}}\right)=4 \mathrm{ab} \\
& \left(2-\frac{c^{2}}{\beta^{2}}\right)\left[2-\frac{\alpha^{2}}{\beta^{2}}+\frac{\alpha^{2}}{\beta^{2}}\left(1-\frac{\mathrm{c}^{2}}{\alpha^{2}}\right)\right] \\
& =4 \sqrt{1-\frac{c^{2}}{\alpha^{2}}} \sqrt{1-\frac{c^{2}}{\beta^{2}}}
\end{aligned}
$$

or

$$
\begin{align*}
& \left(2-\frac{c^{2}}{\beta^{2}}\right)\left[2-\frac{\alpha^{2}}{\beta^{2}}+\frac{\alpha^{2}}{\beta^{2}}-\frac{c^{2}}{\beta^{2}}\right]=4 \sqrt{1-\frac{c^{2}}{\alpha^{2}}} \sqrt{1-\frac{c^{2}}{\beta^{2}}} \\
& \left(2-\frac{c^{2}}{\beta^{2}}\right)^{2}-4\left(1-\frac{c^{2}}{\alpha^{2}}\right)^{\frac{1}{2}}\left(1-\frac{c^{2}}{\beta^{2}}\right)^{\frac{1}{2}}=0 \tag{27}
\end{align*}
$$

This equation contains only one unknown c.
This equation determines the speed $\mathbf{c}$ for Rayleigh surface waves in an uniform half - space.

Equation (27) is called the RAYLEIGH EQUATION for Rayleigh waves. It is also called the Rayleigh frequency equation. It is also called the Rayleigh wave - velocity equation.

This equation is the period equation for Rayleigh waves.
In order to prove that Rayleigh waves really exists, we must show that frequency equation (27) has at least one real root for c . To show this, we proceed as follows :

From equation (27), we write

$$
\begin{aligned}
& \left(2-\frac{c^{2}}{\beta^{2}}\right)^{4}=16\left(1-\frac{c^{2}}{\alpha^{2}}\right)\left(1-\frac{c^{2}}{\beta^{2}}\right) \\
& 16-32 \frac{c^{2}}{\beta^{2}}+24 \frac{c^{4}}{\beta^{4}}-8 \frac{c^{6}}{\beta^{6}}+\frac{c^{8}}{\beta^{8}} \\
& =16\left[1-\frac{c^{2}}{\alpha^{2}}-\frac{c^{2}}{\beta^{2}}+\frac{c^{4}}{\alpha^{2} \beta^{2}}\right] \\
& =16\left[1-\frac{c^{2}}{\beta^{2}} \frac{\beta^{2}}{\alpha^{2}}-\frac{c^{2}}{\beta^{2}}+\frac{c^{4}}{\beta^{4}} \cdot \frac{\beta^{2}}{\alpha^{2}}\right]
\end{aligned}
$$

or
$-32+24 \frac{c^{2}}{\beta^{2}}-8 \frac{c^{4}}{\beta^{4}}+\frac{c^{6}}{\beta^{6}}+16 \frac{\beta^{2}}{\alpha^{2}}+16-16 \frac{c^{2}}{\beta^{2}} \cdot \frac{\beta^{2}}{\alpha^{2}}=0$.
Put

$$
\begin{equation*}
s=c^{2} / \beta^{2} \tag{29}
\end{equation*}
$$

then equation (28) gives

$$
\begin{equation*}
f(s)=s^{3}-8 s^{2}+\left(24-16 \frac{\beta^{2}}{\alpha^{2}}\right) s-16\left(1-\frac{\beta^{2}}{\alpha^{2}}\right)=0 . \tag{30}
\end{equation*}
$$

It is polynomial equation of degree $\mathbf{3}$ in s with real coefficients.
Hence the frequency equation (30) has at least one real root since complex roots occur in conjugate pairs.

Moreover

$$
\begin{equation*}
\mathrm{f}(0)=-16\left(1-\frac{\beta^{2}}{\alpha^{2}}\right)<0, \tag{31}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathrm{f}(1)=1-8+24-\frac{16 \beta^{2}}{\alpha^{2}}+\frac{16 \beta^{2}}{\alpha^{2}}=1>0, \tag{32}
\end{equation*}
$$

since $\beta<\alpha$.

Hence, by the intermediate theorem of calculus, the frequency equation (30) has at least one root in the range $(0,1)$.

That is , there exists $s \in(0,1)$ which is a root of the Rayleigh frequency equation.

Further, $0<\mathrm{s}<1$ implies

$$
\begin{align*}
& 0<\frac{c^{2}}{\beta^{2}}<1 \\
& \mathbf{0}<\mathbf{c}<\boldsymbol{\beta} . \tag{33}
\end{align*}
$$

or
This establishes that Rayleigh surface waves propagating with speed $c=c_{R}$, where

$$
0<c_{R}<\beta<\alpha,
$$

always exists.
This proves the existence of Rayleigh surface waves.
Special Case (Poisson's solid) : When the semi-infinite elastic medium is a Poisson's solid ,
then $\lambda=\mu$ and $\sigma=\frac{1}{4}$ and

$$
\begin{equation*}
\frac{\alpha^{2}}{\beta^{2}}=\frac{\lambda+2 \mu}{\mu}=3 . \tag{34}
\end{equation*}
$$

For this type of elastic medium , frequency equation (30) becomes
or

$$
\begin{aligned}
& 3 s^{3}-24 s^{2}+56 s-32=0 \\
& (s-4)\left(3 s^{2}-12 s+8\right)=0
\end{aligned}
$$

giving
or

$$
\begin{align*}
& \mathrm{s}=4,2+\frac{2}{\sqrt{3}}, 2-\frac{2}{\sqrt{3}} \\
& \frac{c^{2}}{\beta^{2}}=4,2+\frac{2}{\sqrt{3}}, 2-\frac{2}{\sqrt{3}} . \tag{35}
\end{align*}
$$

But c must satisfy the inequality $0<\mathrm{c}<\beta$, therefore, the only possible value for $\mathrm{c}=\mathrm{c}_{\mathrm{R}}$ is

$$
\frac{c_{R}{ }^{2}}{\beta^{2}}=2-\frac{2}{\sqrt{3}}
$$

or

$$
\begin{equation*}
c_{R}=\left[\sqrt{2-\frac{2}{\sqrt{3}}}\right] \cdot \beta \cong 0.9195 \beta \tag{36}
\end{equation*}
$$

giving the speed of propagation of Rayleigh surface waves along the stress free boundary of the Poissonian half-space in the x-direction.

Thus, the order of arrival of P, SV-and Rayleigh waves is


From equations (12) and (13) ; we write

$$
\begin{aligned}
\mathrm{a}^{2} & =1-\frac{c_{R}^{2}}{\alpha^{2}} \\
& =1-\frac{\mathrm{c}_{\mathrm{R}}^{2}}{\beta^{2}} \cdot \frac{\beta^{2}}{\alpha^{2}} \\
& =\frac{1}{9}(3+2 \sqrt{3}),
\end{aligned}
$$

or

$$
\begin{equation*}
\mathrm{a}=\frac{1}{3} \sqrt{3+2 \sqrt{3}} \cong 0.8475 \approx 0.848 \tag{37}
\end{equation*}
$$

and

$$
\begin{aligned}
\mathrm{b}^{2} & =1-\frac{c_{R}^{2}}{\beta^{2}} \\
& =\frac{1}{3}[2 \sqrt{3}-3],
\end{aligned}
$$

or

$$
\begin{equation*}
\mathrm{b} \cong 0.3933 \cong 0.393 . \tag{38}
\end{equation*}
$$

Further , from equations (25) , (37) and (38), we find

$$
\begin{equation*}
\frac{A}{B}=1.468 \mathrm{i} \tag{39}
\end{equation*}
$$

## Remark 1: Displacements due to Rayleigh waves (Particle motion)

## From equations (2) , (3) , (17) , (18) and (26) ; we find

$$
\begin{align*}
\mathrm{u}(\mathrm{x}, \mathrm{z}, \mathrm{t}) & =\left[\mathrm{A}(-\mathrm{ik}) \mathrm{e}^{-\mathrm{akz}}+\mathrm{B}(-\mathrm{bk}) \mathrm{e}^{-\mathrm{bkz}}\right] \mathrm{e}^{\mathrm{i}(\omega t-\mathrm{kx})} \\
& =\mathrm{k}\left[-\mathrm{iA} \mathrm{e}^{-\mathrm{akz}}-\mathrm{A}\left(\frac{2-c^{2} / \beta^{2}}{2 i}\right) \mathrm{e}^{-\mathrm{bkz}}\right] \mathrm{e}^{\mathrm{i}(\omega \mathrm{t}-\mathrm{kx})} \\
& =-\mathrm{ikA}\left[\mathrm{e}^{-\mathrm{akz}}-\left(1-\frac{c^{2}}{2 \beta^{2}}\right) \mathrm{e}^{-\mathrm{bkz}}\right] \mathrm{e}^{\mathrm{i}(\omega \mathrm{t}-\mathrm{kx})} . \tag{40}
\end{align*}
$$

Similarly, we shall find

$$
\begin{align*}
& \mathrm{w}(\mathrm{x}, \mathrm{z}, \mathrm{t})=\left[(-\mathrm{ak}) \mathrm{Ae}^{-\mathrm{akz}}-(-\mathrm{ik}) \mathrm{B} \mathrm{e}^{-\mathrm{bkz}}\right] \mathrm{e}^{\mathrm{i}(\omega \mathrm{t}-\mathrm{kx})} \\
&=\mathrm{k}[-\mathrm{a} \mathrm{~A} \mathrm{e} \\
&-\mathrm{akz}\left.\frac{i e^{b k z}(-2 i A a)}{\left(2-c^{2} / \beta^{2}\right)}\right] \mathrm{e}^{\mathrm{i}(\omega \mathrm{t}-\mathrm{kx})}  \tag{41}\\
&= \mathrm{kA}\left[-\mathrm{a} \mathrm{e}^{-\mathrm{akz}}+\left(\frac{a}{1-c^{2} / 2 \beta^{2}}\right) \mathrm{e}^{-\mathrm{bkz}}\right] \mathrm{e}^{\mathrm{i}(\omega \mathrm{t}-\mathrm{kx})} .
\end{align*}
$$

Let

$$
\begin{align*}
& \theta=\omega \mathrm{t}-\mathrm{kx}  \tag{42}\\
& \mathrm{U}(\mathrm{z})=\mathrm{e}^{-\mathrm{akz}}-\left(1-\frac{c^{2}}{2 \beta^{2}}\right) \mathrm{e}^{-\mathrm{bzk}},  \tag{43}\\
& \mathrm{~W}(\mathrm{z})=-\mathrm{a} \mathrm{e}^{-\mathrm{akz}}+\left(\frac{a}{1-c^{2} / 2 \beta^{2}}\right) \mathrm{e}^{-\mathrm{bkz}} . \tag{44}
\end{align*}
$$

Then , taking the real parts of equations (40) and (41) and using equation (42) to (44) ; we find

$$
\begin{align*}
& \mathrm{u}(\mathrm{x}, \mathrm{z}, \mathrm{t})=\mathrm{AkU}(\mathrm{z}) \sin \theta  \tag{45}\\
& \mathrm{w}(\mathrm{x}, \mathrm{z}, \mathrm{t})=\mathrm{AkW}(\mathrm{z}) \cos \theta \tag{46}
\end{align*}
$$

we remember that A is the potential amplitude. Eliminating $\theta$ from equations (45) and (46) ,

> we obtain

$$
\begin{equation*}
\frac{u^{2}}{[A k U(z)]^{2}}+\frac{w^{2}}{[A k W(z)]^{2}}=1 \tag{47}
\end{equation*}
$$

which is an equation of an ellipse in the vertical xz-plane.
Equation (47) shows that particles, during the propagation of Rayleigh Surface Waves, describe ellipses.

Remark 2: Particle Motion at the surface ( $\mathrm{z}=0$ )

On the surface $z=0$,we find, at $z=0$,

$$
\mathrm{U}(0)=\frac{c^{2}}{2 \beta^{2}} \quad, \quad \mathrm{~W}(0)=\mathrm{a}\left(\frac{\frac{c^{2}}{2 \beta^{2}}}{1-c^{2} / 2 \beta^{2}}\right)
$$

Since $0<c<\beta$, so

$$
\begin{equation*}
\mathrm{U}(0)>0 \quad, \quad \mathrm{~W}(0)>0 \tag{49}
\end{equation*}
$$

Let

$$
\begin{equation*}
\mathrm{a}_{1}=\mathrm{AkW}(0), \quad \mathrm{b}_{1}=\operatorname{AkU}(0) . \tag{50}
\end{equation*}
$$

Then

$$
\begin{equation*}
a_{1}>0, b_{1}>0, \tag{51}
\end{equation*}
$$

and equation (47) reduces to

$$
\begin{equation*}
\frac{u^{2}}{b_{1}^{2}}+\frac{w^{2}}{a_{1}^{2}}=1 \tag{52}
\end{equation*}
$$

## Remark 3: Particular case (Poisson's Solid) :

In this case, we find, at $\mathrm{z}=0$,

$$
\begin{equation*}
\frac{a_{1}}{b_{1}}=\frac{W(0)}{U(0)}=\sqrt{3} \mathrm{a} \cong 1.5 \tag{53}
\end{equation*}
$$

using (37). So

$$
\begin{equation*}
a_{1}>b_{1} \tag{54}
\end{equation*}
$$

Thus, the surface particle motion ( $\mathrm{on} \mathrm{z}=0$ ) is an ellipse with a vertical major axis.

The horizontal and vertical displacement components are out of phase by $\frac{\pi}{2}$.
The resulting surface particle motion is Retrograde (opposite to that of wave propagation).

Remark 4: The dependencies of the displacement components $u$ and $w$ on depth(z) are given by equation (43) to (46) .

There is a value of z for which $\mathrm{u}=0$ (For Poisson's solid, at $\mathrm{z}=+0.19 \lambda, \lambda=$ $\left.\frac{2 \pi}{k}, \mathrm{u}=0\right)$, whereas w is never zero.

At the depth, where $u$ is zero, its amplitude changes sign.
For greater depths, the particle motion is Prograde.
With increasing depth , the amplitudes of $u$ and $w$ decrease exponentially , with $w$ always larger than $u$.

Thus, the elliptic motion changes from retrograde at the surface to prograde at depth , passing through a node at which there is no horizontal motion.

So , for the propagation of Rayleigh surface waves, a surface particle describes an ellipse, about its mean position, in the retrograde sense.



Particle motion for the fundamental Rayleigh mode in a uniform half-space, propagating from left to right. One horizontal wavelength $(\Lambda)$ is shown; the dots are plotted at a fixed time point. Motion is counter clockwise (retrograde) at the surface, changing to purely vertical motion at a depth of about $\mathcal{N} 5$, and becoming clockwise (prograde) at greater depths. Note that the time behavior at a fixed distance is given by looking from right to left in this plot.

Remark 5: We see that frequency equation (27) for Rayleigh surface waves is independent of $\omega$. Therefore, the velocity $c_{R}$ of Rayleigh surface waves is constant and fixed.

This phenomenon is called nondispessive.
That is , Rayleigh waves are undispresed.
Remark 6: Maximum displacement parallel to the direction of Rayleigh waves

$$
\begin{aligned}
& =(u)_{\max } \\
& =\mathrm{b}_{1} \\
& =\frac{2}{3} \mathrm{a}_{1} \quad, \quad \text { for a Poisson solid. }
\end{aligned}
$$

$=$ two-third of the maximum displacement in the vertical direction

## for a Poisson solid.

Note (1) : Rayleigh waves are important because the largest disturbances caused by an earthquake recorded on a distant seismogram are usually those of Rayleigh waves.

GROUND ROLL is the term commonly used for Rayleigh waves.
Note (2) : Although a "a free surface" means contact with a vacuum , the elastic constants and density of air are so lows in comparison with values for rocks that the surface of the earth is approximately a free surface.

Note (3) : The boundary conditions $\tau_{z x}=\tau_{z z}=0$ at $\mathrm{z}=0$ require that these two conditions must be satisfied, and so we require two parameters than can be adjusted. Therefore, we assume that
both P-and SV-components exist and adjust their amplitude to satisfy the boundary conditions.

Exercise : Show that the displacement components at the surface of an elastic Poisson solid due to Rayleigh waves are

$$
\begin{aligned}
& u(x, t)=-0.423 k A \sin k\left(x-c_{R} t\right) \\
& w(x, t)=0.620 k A \cos k\left(x-c_{R} t\right), \quad v \equiv 0,
\end{aligned}
$$

with usual notation.

### 10.3. SURFACE WAVES OF SH-TYPE IN A HALF-SPACE

We consider first the possibility of the propagation of SH type surface waves (called Love waves) in a homogeneous semi-infinite isotropic elastic medium occupying the half-space $z \geq 0$. The horizontal boundary $z=0$ of the medium is assumed to be stress free. Let $\rho$ be the density of the medium and $\lambda, \mu$ Lame' constants(figure) .

(Elastic isotropic half-space)
Let the two - dimensional SH-wave motion takes place in the xz-plane. The basic equations for SH - wave motion are

$$
\begin{align*}
& u=w=0, \quad v=v(x, z, t)  \tag{1}\\
& \frac{\partial^{2} v}{\partial x^{2}}+\frac{\partial^{2} v}{\partial z^{2}}=\frac{1}{\beta^{2}} \frac{\partial^{2} v}{\partial t^{2}}  \tag{2}\\
& \beta^{2}=\frac{\mu}{\rho} \tag{3}
\end{align*}
$$

We try a plane wave solution of wave equation (2) of the form

$$
\begin{equation*}
\mathrm{v}(\mathrm{x}, \mathrm{z}, \mathrm{t})=\mathrm{B} \cdot \mathrm{e}^{-\mathrm{bkz}} \cdot \mathrm{e}^{\mathrm{i}(\omega \mathrm{t}-\mathrm{kx})} \tag{4}
\end{equation*}
$$

where $\omega$ is the angular frequency of wave, $k=$ wave number and $c=\omega / k$ is the speed with which surface waves are travelling in the x -direction on the surface $\mathrm{z}=0 ; \mathrm{b}>0$ and B is an arbitrary constant.

The amplitude of surface wave is $\mathrm{B} \mathrm{e}^{-\mathrm{bkz}}$ which die exponentially as z increases.

Substituting the value of $v(x, z, t)$ from equation (4) into equation (2), we find

$$
\begin{equation*}
\mathrm{b}^{2}=1-\frac{c^{2}}{\beta^{2}} \tag{5}
\end{equation*}
$$

Since $\mathrm{b}>0$,so

$$
\begin{equation*}
c<\beta . \tag{6}
\end{equation*}
$$

Using the stress-displacement relations, we find

$$
\begin{align*}
& \tau_{31}=\tau_{33}=0, \\
& \tau_{32}=\mu \frac{\partial \mathrm{v}}{\partial \mathrm{z}}=-\mu \mathrm{bk} \cdot \mathrm{~B} \cdot \mathrm{e}^{\mathrm{i}(\omega \mathrm{t}-\mathrm{kx})} \cdot \mathrm{e}^{-\mathrm{bkz}} \tag{7}
\end{align*}
$$

Hence, the stress-free condition of the boundary $\mathrm{z}=0$ implies that, using (7),

$$
\mu \mathrm{bkB} \mathrm{e} \quad{ }^{\mathrm{i}(\omega t-\mathrm{kx})}=0
$$

or

$$
\begin{equation*}
\mathrm{B}=0 \tag{8}
\end{equation*}
$$

as $\mu \neq 0, b \neq 0, k \neq 0$.
This implies that in the case of a homogeneous isotropic elastic half space, Love waves do not exist at all.

### 10.4. PROPAGATION OF LOVE WAVES

Surface waves of the SH-type are observed to occur on the earth's surface. Love (1911) showed that if the earth is modelled as an isotropic elastic layer of finite thickness lying over a homogeneous elastic isotropic halfspace (rather than considering earth as a purely uniform half-space) then SH type waves occur in the stress-free surface of a layered half-space.

Now, we consider the possibility of propagation of surface waves of SH-type (Love waves) in a semi-infinite elastic isotropic medium consisting of a horizontal elastic layer of uniform thickness H lying over a half-space.

It is assumed that two elastic isotropic media are welded together and the horizontal boundary $\mathrm{z}=0$ of the semi-infinite medium is stress - free (see , figure).


Let the layer and the half-space have different densities $\rho_{1}, \rho_{2}$ and different shear moduli $\mu_{1}, \mu_{2}$ respectively. Let two-dimensional SH-motion takes place parallel to xz-plane. The basic equation for SH-wave motion are

$$
\begin{align*}
& \mathrm{u}_{1}=\mathrm{w}_{1} \equiv 0, \mathrm{v}_{1}=\mathrm{v}_{1}(\mathrm{x}, \mathrm{z}, \mathrm{t})  \tag{1}\\
& \frac{\partial^{2} \mathrm{v}_{1}}{\partial \mathrm{x}^{2}}+\frac{\partial^{2} \mathrm{v}_{1}}{\partial \mathrm{z}^{2}}=\frac{1}{\beta_{1}^{2}} \frac{\partial^{2} \mathrm{v}_{1}}{\partial \mathrm{t}^{2}}  \tag{2}\\
& \beta_{1}^{2}=\frac{\mu_{1}}{\rho_{1}} \tag{3}
\end{align*}
$$

in the layer $0 \leq \mathrm{z} \leq \mathrm{H}$, and

$$
\begin{align*}
& u_{2}=w_{2}=0, \quad v_{2}=v_{2}(x, z, t)  \tag{4}\\
& \frac{\partial^{2} v_{2}}{\partial x^{2}}+\frac{\partial^{2} v_{2}}{\partial z^{2}}=\frac{1}{\beta_{2}{ }^{2}} \frac{\partial^{2} v_{2}}{\partial t^{2}}  \tag{5}\\
& \beta_{2}{ }^{2}=\frac{\mu_{2}}{\rho_{2}} \tag{6}
\end{align*}
$$

in the half-space $(\mathrm{z} \geq \mathrm{H})$.
Suitable plane wave solutions of wave equations (2) and (5) are (exercise), as discussed in detail already,

$$
\begin{equation*}
\mathrm{v}_{1}(\mathrm{x}, \mathrm{z}, \mathrm{t})=\left(\mathrm{A}_{1} e^{-b_{1} k z}+\mathrm{B}_{1} e^{b_{1} k z}\right) \mathrm{e}^{\mathrm{i}(\omega \mathrm{t}-\mathrm{kx})} \tag{7}
\end{equation*}
$$

in the layer $0 \leq \mathrm{z} \leq \mathrm{H}$, and

$$
\begin{equation*}
\mathrm{v}_{2}(\mathrm{x}, \mathrm{z}, \mathrm{t})=\mathrm{A}_{2} e^{-b_{2} k z} \mathrm{e}^{\mathrm{i}(\omega \mathrm{t}-\mathrm{kx})} \tag{8}
\end{equation*}
$$

in the half-space $(\mathrm{z} \geq \mathrm{H}) . \omega$ is the angular wave frequency and k is the wave number, and $c=\omega / k$ is the speed of propagation of surface wave (if it exists) in the positive x-direction. $A_{1}, B_{1}$ are constants, $b_{1}$ and $b_{2}$ are real numbers with $\mathrm{b}_{2}>0$. However, $\mathrm{b}_{1}$ is unrestricted because z is finite in the layer.

Substituting for $\mathrm{v}_{2}$ from (8) into (5) yields the relation

$$
\begin{equation*}
\mathrm{b}_{2}^{2}=\left(1-\frac{c^{2}}{\beta_{2}^{2}}\right) \tag{9}
\end{equation*}
$$

and, therefore ,

$$
\begin{equation*}
\mathrm{c}<\beta_{2}, \tag{10}
\end{equation*}
$$

otherwise

$$
\mathrm{v}_{2} \rightarrow \infty \text { as } \mathrm{z} \rightarrow \infty .
$$

From equations (7) and (2), we find

$$
\begin{equation*}
\mathrm{b}_{1}^{2}=\left(1-\frac{c^{2}}{\beta_{1}^{2}}\right) \tag{11}
\end{equation*}
$$

The stress-displacement relations imply

$$
\begin{equation*}
\tau_{31}=\tau_{33} \equiv 0 \tag{12}
\end{equation*}
$$

in the layer as well as in the half space. Also

$$
\begin{equation*}
\tau_{32}=\mu_{1} \mathrm{k}\left(-\mathrm{b}_{1} \mathrm{~A}_{1} e^{-b_{1} k z}+\mathrm{b}_{1} \mathrm{~B}_{1} e^{b_{1} k z}\right) \mathrm{e}^{\mathrm{i}(\omega \mathrm{t}-\mathrm{kx})} \tag{13}
\end{equation*}
$$

in the layer $0 \leq \mathrm{z} \leq \mathrm{H}$, and

$$
\begin{equation*}
\tau_{32}=\mathrm{k} \mu_{2} \mathrm{~A}_{2}\left(-\mathrm{b}_{2}\right) e^{-b_{2} k z} \cdot \mathrm{e}^{\mathrm{i}(\omega t-\mathrm{kx})} \tag{14}
\end{equation*}
$$

in the half space $\mathrm{z} \geq \mathrm{H}$.
The stress-free boundary $\mathrm{z}=0$ implies that

$$
\begin{equation*}
\tau_{32}=0 \tag{15}
\end{equation*}
$$

at $\mathrm{z}=0$. This gives

$$
\begin{equation*}
\mathrm{B}_{1}=\mathrm{A}_{1} . \tag{16}
\end{equation*}
$$

Since, there is a welded contact between the layer and the half-space at the interface $\mathrm{z}=\mathrm{H}$, so the displacement and the tractions must be continuous across the interface $\mathrm{z}=\mathrm{H}$.

Thus, the boundary conditions at $\mathrm{z}=\mathrm{H}$ are

$$
\begin{align*}
& \mathrm{v}_{1}=\mathrm{v}_{2}  \tag{17}\\
& \left.\tau_{32}\right|_{\text {layer }}=\left.\tau_{32}\right|_{\text {half-space }} . \tag{18}
\end{align*}
$$

From equations (13), (14), (17) and (18), we find

$$
\begin{equation*}
\mathrm{A}_{1} e^{-k b_{1} H}+\mathrm{B}_{1} e^{k k_{1} H}=\mathrm{A}_{2} e^{-k b_{2} H}, \tag{19}
\end{equation*}
$$

and

$$
\begin{equation*}
\mu_{1} \mathrm{k} \mathrm{~b}_{1}\left[-\mathrm{A}_{1} e^{-b_{1} k H}+\mathrm{B}_{1} e^{b_{1} k H}\right]=-\mu_{2} \mathrm{~b}_{2} \mathrm{k} e^{-b_{2} k H} \tag{20}
\end{equation*}
$$

Equation (16), (19) and (20) are three homogeneous in $\mathrm{A}_{2}, \mathrm{~A}_{1}, \mathrm{~B}_{1}$ and $\mathrm{A}_{2}$. We shall now eliminate them. From equations (19) \& (20) , we write (after putting $\mathrm{B}_{1}=\mathrm{A}_{1}$ )
or

$$
\frac{A_{1}\left[e^{-k b_{1} H}+e^{k b_{1} H}\right]}{A_{1} \cdot \mu_{1} b_{1}\left(e^{-b_{1} k H}-e^{b_{1} k H}\right) k}=\frac{1}{\mu_{2} b_{2} k}
$$

$$
\frac{e^{-b_{1} k H}-e^{b_{1} k H}}{e^{-k b_{1} H}+e^{k b_{1} H}}=\frac{\mu_{2} b_{2}}{\mu_{1} b_{1}}
$$

or

$$
\begin{equation*}
\tan \mathrm{h}\left(\mathrm{~b}_{1} \mathrm{kH}\right)+\frac{\mu_{2} b_{2}}{\mu_{1} b_{1}}=0 \tag{21}
\end{equation*}
$$

Equation (21) is known as period equation/frequency equation/Dispersion
equation for surface Love waves.
Equation (21) can also be written as

$$
\tan h\left[k h \cdot \sqrt{1-\frac{c^{2}}{\beta_{1}^{2}}}\right]=-\frac{\mu_{2}}{\mu_{1}} \cdot \frac{\sqrt{1-\frac{c^{2}}{\beta_{2}^{2}}}}{\sqrt{1-\frac{c^{2}}{\beta_{1}^{2}}}}
$$

or

$$
\begin{equation*}
\tanh \left[\frac{\omega \mathrm{H}}{\mathrm{c}} \sqrt{1-\frac{\mathrm{c}^{2}}{\beta_{1}^{2}}}\right]=-\frac{\mu_{2}}{\mu_{1}} \cdot \frac{\sqrt{1-\frac{c^{2}}{\beta_{2}^{2}}}}{\sqrt{1-\frac{c^{2}}{\beta_{1}^{2}}}} . \tag{22}
\end{equation*}
$$

Equation (22) is a transcendental equation.
For given $\omega$, we can find the speed c for surface Love waves. We note the value of $\mathbf{c}$ depends upon $\omega$. This means that waves of different frequencies will, in general , propagate with different phase velocity.

## This phenomenon is known as dispersion.

It is caused by the inhomogeneity of the medium (layered medium) due to some abrupt discontinuities within the medium (or due to continuous change of the elastic parameters which is not the present case).

Thus, Love waves are dispressed.
We consider now the following two possibilities between c and $\beta_{1}$.
(i) Either
$c \leq \beta_{1}$,
(ii) or $c>\beta_{1}$.

When $\mathrm{c} \leq \beta_{1}$ : In this case $\mathrm{b}_{1}$ is real (see, equation (11)) and left side of (22) becomes real and positive and right side of (22) is real and negative. Therefore , equation (22) can not possess any real solution for c .

Therefore, in this case, Love waves do not exist.
So , for the existence of surface Love waves, we must have

$$
\begin{equation*}
\mathbf{c}>\beta_{1} . \tag{24}
\end{equation*}
$$

In this case (24), $b_{1}$ is purely imaginary and we may write

$$
\begin{equation*}
\mathrm{b}_{1}==\sqrt{1-\frac{c^{2}}{\beta_{1}^{2}}}=\mathrm{i}\left(\sqrt{\frac{c^{2}}{\beta_{1}}-1}\right) . \tag{25}
\end{equation*}
$$

## Then equation (22) becomes

$$
\begin{equation*}
\tan \left(k H \sqrt{\frac{c^{2}}{\beta_{1}^{2}}-1}\right)=\frac{\mu_{2}}{\mu_{1}}\left(\frac{\sqrt{1-\frac{c^{2}}{\beta_{2}^{2}}}}{\sqrt{\frac{c^{2}}{\beta_{1}^{2}}-1}}\right) \tag{26}
\end{equation*}
$$

From equations (10) and (24); we find

$$
\begin{equation*}
\beta_{1}<c<\beta_{2} . \tag{27}
\end{equation*}
$$

Equation (26) is a transcendental equation that yields infinitely many roots for c.

Thus, the possible speeds of the Love waves are precisely the roots of equations (26) that lie in the interval $\left(\beta_{1}, \beta_{2}\right)$.

This indicates that the shear velocity in the layer must be less than the shear velocity in the half-space for the possible existence of Love waves.

This gives the upper and lower bounds for the speed of Love waves.
Remark 1: If the layer and the half - space are such that $\beta_{1} \leq \beta$, then existence of Love waves are not possible

Remark 2: In the limiting case when the layer is absent, we have

$$
\mu=\mu_{1} \text { and } \rho=\rho_{1}
$$

and therefore

$$
\beta=\beta_{1}
$$

Equation (22) leads to the impossible condition

$$
0=-1 .
$$

Hence, in this case, the wave considered can not exist.
Remark 3: When k or $\omega \rightarrow 0$, we get $\mathrm{c} \rightarrow \beta_{1}$.
The dispersion curve is given in the following figure.

$\omega$
Here, if we assume

$$
\mu_{1} / \mu=1.8, \beta=3.6 \mathrm{~km} / \mathrm{sec}, \beta_{1}=4.6 \mathrm{~km} / \mathrm{sec}
$$

then

$$
c_{L}=\text { speed of Love waves }=4.0 \mathrm{~km} / \mathrm{sec}
$$

