INSTITUTE OF AERONAUTICAL ENGINEERING
(Autonomous)
Dundigal, Hyderabad -500 043

## ELECTRICAL AND ELECTRONICS ENGINEERING

COURSE LECTURE NOTES

| Course Name | COMPLEX ANALYSIS AND PROBABILITY DISTRIBUTIONS |
| :--- | :--- |
| Course Code | AHSB06 |
| Programme | B.Tech |
| Semester | IIVB |
| Course <br> Coordinator | Mr. CH. Soma Shekar, Assistant Professor |
| Course Faculty | Mr. CH. Soma Shekar, Assistant Professor |
| Lecture Numbers | $1-63$ |
| Topic Covered | All |

## COURSE OBJECTIVES:

The course should enable the students to:

| I | Understand the basic theory of complex functions to express the power series. |
| :---: | :--- |
| II | Evaluate the contour integration using Cauchy residue theorem. |
| III | Enrich the knowledge of probability on single random variables and probability <br> distributions. |

COURSE LEARNING OUTCOMES (CLOs):
Students, who complete the course, will have demonstrated the ability to do the following:

| AHSB06.01 | Recall continuity, differentiability, analyticity of a function using limits. |
| :---: | :--- |
| AHSB06.02 | Interpret the conditions for a complex variable to be analytic and/or entire <br> function. |
| AHSB06.03 | Interpret the concepts of Cauchy-Riemann relations and harmonic functions. |
| AHSB06.04 | Analyze the Bilinear transformation by cross ratio property. |
| AHSB06.05 | Identify the conditions of fixed and critical point of Bilinear Transformation. |
| AHSB06.06 | Demonstrate the area under a curve using the concepts of indefinite <br> integration. |


| AHSB06.07 | Interpret the concepts of the Cauchy's integral formula and the generalized <br> Cauchy's integral formula. |
| :---: | :--- |
| AHSB06.08 | Demonstrate complex functions as power series and radius of convergence of <br> power series. |
| AHSB06.09 | Interpret the concept of complex integration to the real-world problems of <br> flow with circulation around a cylinder. |
| AHSB06.10 | Asses the Taylor's and Laurent series expansion of complex functions. |
| AHSB06.11 | Interpret the concept of different types of singularities for analytic function. |
| AHSB06.12 | Identify the poles, residues and solve integrals using Cauchy's residue <br> theorem. |
| AHSB06.13 | Interpret the concept of Cauchy's residue theorem to the real-world problems <br> of Quantum Mechanical scattering and Quantum theory of atomic collisions. |
| AHSB06.14 | Demonstrate an understanding of the basic concepts of probability and <br> random variables. |
| AHSB06.15 | Classify the types of random variables and calculate mean, variance. <br> AHSB06.16Estimate moment about origin, central moments, moment generating <br> function of probability distribution. |
| AHSB06.17 | Recognize where the Binomial distribution could be appropriate model of the <br> distributions. |
| AHSB06.18 | Recognize where the Poisson distribution could be appropriate model of the <br> distributions. |
| AHSB06.19 | Recognize where the Binomial distribution and Poisson distribution could be <br> appropriate to find mean, variance of the distributions. |
| AHSB06.20 | Apply the inferential methods relating to the means of normal distributions. <br> AHSB06.21Interpret Binomial distribution to the phenomena of real-world problem like <br> sick versus healthy. |
| AHSB06.22 | Identify the mapping of Normal distribution in real-world problem to analyze <br> the stock market. |
| AHSB06.23 | Use Poisson distribution in real-world problem to predict soccer scores. <br> Possess the knowledge and skills for employability and to succeed in national <br> and international level competitive examinations. |
| AHSB06.24 |  |

## SYLLABUS

| Module-I | COMPLEX FUNCTIONS AND DIFFERENTIATION | Classes: 09 |
| :--- | :--- | :--- |
| Complex functions differentiation and integration: Complex functions and its representation <br> on argand plane, concepts of limit, continuity, differentiability, analyticity, Cauchy-Riemann <br> conditions and harmonic functions; Milne-Thomson method. Bilinear Transformation. |  |  |


| Module-III | COMPLEX INTEGRATION | Classes: 09 |
| :--- | :--- | :--- |

Line integral: Evaluation along a path and by indefinite integration; Cauchy's integral theorem; Cauchy's integral formula; Generalized integral formula; Power series expansions of complex functions and contour Integration: Radius of convergence.

## Module -IIII POWER SERIES EXPANSION OF COMPLEX FUNCTION

Classes: 09
Expansion in Taylor's series, Maclaurin's series and Laurent series. Singular point; Isolated singular point; Pole of order m; Essential singularity; Residue: Cauchy Residue Theorem. Evaluation of Residue by Laurent Series and Residue Theorem.
Evaluation of integrals of the type $\int_{0}^{2 \pi} f(\cos \theta, \sin \theta) d \theta$ and $\int_{-\infty}^{\infty} f(x) d x$

| Module -IV | SINGLE RANDOM VARIABLES | Classes: 09 |
| :--- | :--- | :---: |
| Random variables: Discrete and continuous, probability distributions, mass function-density <br> function of a probability distribution. Mathematical expectation. Moment about origin, central <br> moments, moment generating function of probability distribution. |  |  |
| Module -V | PROBABILITY DISTRIBUTIONS | Classes: 09 |
| Binomial, Poisson and normal distributions and their properties. |  |  |

## Text Books:

1. Kreyszig, "Advanced Engineering Mathematics", John Wiley \& Sons Publishers, $10^{\text {th }}$ Edition, 2010
2. B. S. Grewal, "Higher Engineering Mathematics", Khanna Publishers, $4{ }^{\text {rd }}$ Edition, 2015.

Reference Books:

1. T.K.V Iyengar, B.Krishna Gandhi, "Engineering Mathematics - III", S. Chand \& Co., $12{ }^{\text {th }}$ Edition, 2015.
2. T.K.V Iyengar, B.Krishna Gandhi, "Probability and Statistics", S. Chand \& Co., $7^{\text {th }}$ Edition, 2015.
3. Churchill, R.V. and Brown, J.W, "Complex Variables and Applications", Tata Mc Graw-Hill, $8^{\text {th }}$ Edition, 2012.

## Web References:

1. http://www.efunda.com/math/math_home/math.cfm
2. http://www.ocw.mit.edu/resourcs/\#Mathematics
3. http://www.sosmath.com
4. http://www.mathworld.wolfram.com

## E-Text Books:

1. http://www.keralatechnologicaluniversity.blogspot.in/2015/06/erwin-kreyszig-advanced-engineering- mathematics-ktu-ebook-download.html
2. http://www.faadooengineers.com/threads/13449-Engineering-Maths-II-eBooks

## MODULE-I

## COMPLEX ANALYSIS AND DIFFERENTIATION

## COMPLEX FUNCTIONS

## Complex number

For a complex number $z=x+i y$, the number $\operatorname{Re} z=x$ is called the real part of $z$ and the number $\operatorname{Im} \mathrm{z}=\mathrm{y}$ is said to be the its imaginary part. If $\mathrm{x}=0, \mathrm{z}$ is said to be a purely imaginary number.

Definition : Let $\mathrm{z}=\mathrm{x}+\mathrm{iy} \in \mathrm{C}$. The complex number $\mathrm{z}=\mathrm{x}-\mathrm{iy}$ is called the complex conjugate of $z$ and $|z|=\sqrt{x^{2}+y^{2}}$ is said to be the absolute value or the modulus of the complex number z.

## Functions of a Complex Variable :

Let D be a nonempty set in C . A single-valued complex function or, simply, a complex function $\mathrm{f}: \mathrm{D} \rightarrow \mathrm{C}$ is a map that assigns to each complex argument $\mathrm{z}=\mathrm{x}+$ iy in D a unique complex number $w=u+i v$. We write $w=f(z)$.

The set $D$ is called the domain of the function $f$ and the set $f(D)$ is the range or the image of f . So, a complex-valued function f of a complex variable z is a rule that assigns to each complex number $z$ in a set $D$ one and only one complex number $w$. We call w the image of $z$ under $f$.

If $z=x+i y \in D$, we shall write $f(z)=u(x, y)+i v(x, y)$ or $f(z)=u(z)+i v(z)$. The real functions $u$ and $v$ are called the real and, respectively, the imaginary part of the complex function f. Therefore, we can describe a complex function with the aid of two real functions depending on two real variables.

Example 1.The function $\mathrm{f}: \mathrm{C} \rightarrow \mathrm{C}$, defined by $\mathrm{f}(\mathrm{z})=\mathrm{z}^{3}$, can be written as $\mathrm{f}(\mathrm{z})=\mathrm{u}(\mathrm{x}, \mathrm{y})+\mathrm{iv}(\mathrm{x}$, $y$, with $u, v: R^{2} \rightarrow R$ given by $u(x, y)=x^{3}-3 x y^{2}, v(x, y)=3 x^{2} y-y^{3}$.

Example 2. For the function $f: C \rightarrow C$, defined by $f(z)=e^{z}$, we have $u(x, y)=e^{x} \cos y, v(x, y)$ $=e^{x} \sin y$, for any $(x, y) \in R^{2}$.

Limits of Functions: Let $\mathrm{D} \subseteq \mathrm{C}, \mathrm{a} \in \mathrm{D}^{\prime}$ and $\mathrm{f}: \mathrm{D} \rightarrow \mathrm{C}$. A number $\mathrm{l} \in \mathrm{C}$ is called a limit of the function $f$ at the point a if for any $V \in V(1)$, there exists $U \in V(a)$ such that, for any $z \in U \cap D \backslash$ $\{\mathrm{a}\}$, it follows that $\mathrm{f}(\mathrm{z}) \in \mathrm{V}$. We shall use the notation $\mathrm{l}=\lim _{z \rightarrow z_{0}} \mathrm{f}(\mathrm{z})$.

Remark : If a complex function $\mathrm{f}: \mathrm{D} \rightarrow \mathrm{C}$ possesses a limit 1 at a given point a, then this limit is unique.

Exercise 1: Prove that $\lim _{z \rightarrow z_{0}} \frac{z}{z}$ does not exist.
Solution : To prove that the above limit does not exist, we compute this limit as $z \rightarrow 0$ on the real and on the imaginary axis, respectively. In the first situation, i.e. for $z=x \in R$, the value of the limit is 1 . In the second situation,
i.e. for $\mathrm{z}=\mathrm{i} y$, with $\mathrm{y} \in \mathrm{R}$, the limit is -1 . Thus, the limit depends on the direction from which we approach 0 , which implies that the limit does not exist.

## Differentiability of complex function :

Let $\mathrm{w}=\mathrm{f}(\mathrm{z})$ be a given function defined for all z in a neighbourhood of $\mathrm{z}_{0}$. If $\lim _{\Delta z \rightarrow 0} \frac{f\left(z_{0}+\Delta z\right)-f\left(z_{0}\right)}{\Delta z}$ exists, the function $\mathrm{f}(\mathrm{z})$ is said to be derivable at $\mathrm{z}_{0}$ and the limit is denoted by $f^{\prime}\left(z_{0}\right) \cdot f^{\prime}\left(z_{0}\right)$ if exists is called the derivative of $\mathrm{f}(\mathrm{z})$ at $\mathrm{z}_{0}$.

Exercise : $\mathrm{f}(\mathrm{z})=\left|z^{2}\right|$ is a function which is continuous at all z but not derivable at any $\mathrm{z} \neq 0$
Solution: Let $\mathrm{f}(\mathrm{z})=\left|z^{2}\right|=z \bar{z}$
Then $\mathrm{f}(\mathrm{z})=z_{0} \bar{z}_{0}$

We have to prove that $\underset{z \rightarrow z_{0}}{\operatorname{lt}} \mathrm{z}=\mathrm{z}_{0}$ and $\underset{z \rightarrow z_{0}}{\operatorname{lt}} \bar{z}=\bar{z}_{0} \quad$ Thus $\underset{z \rightarrow z_{0}}{l t} z \bar{z}=z_{0} \bar{z}_{0}$
$\therefore \operatorname{lt}_{z \rightarrow z_{0}} \mathrm{f}(\mathrm{z})=\mathrm{f}\left(\mathrm{z}_{0}\right)$
$\therefore$ The function is continuous at all z
$\therefore \mathrm{f}\left(\mathrm{z}_{0}+\Delta z\right)=\left(z_{0}+\Delta z\right)(\bar{z}+\Delta \bar{z})=z_{0} \bar{z}_{0}+z_{0} \Delta \bar{z}+\Delta z \bar{z}_{0}+\Delta z \Delta \bar{z}$
Now $\frac{f\left(z_{0}+\Delta z\right)-f\left(z_{0}\right)}{\Delta z}=\frac{z_{0} \Delta \bar{z}+\Delta z \bar{z}_{0}+\Delta z \Delta \bar{z}}{\Delta z}$
Consider the limit as $\Delta z \rightarrow 0$

Case 1: let $\Delta z \rightarrow 0$ along x-axis then $\Delta x=\Delta x, \Delta y=0 \Rightarrow \Delta z=\Delta x$
$\therefore \operatorname{lt}_{\Delta z \rightarrow 0} \frac{f\left(z_{0}+\Delta z\right)-f(z)}{\Delta z}=\operatorname{lt}_{\Delta x \rightarrow 0} \frac{z_{0} \Delta x+\Delta x \bar{z}_{0}+\Delta x \Delta x}{\Delta x}=z_{0}+\bar{z}_{0}$
Case 2: Let $\Delta z \rightarrow 0$ along $y$-axis then $\Delta x=0, \Delta y=\Delta y \Rightarrow \Delta z=i \Delta y$

$$
\operatorname{lt}_{\Delta z \rightarrow 0} \frac{f\left(z_{0}+\Delta z\right)-f(z)}{\Delta z}=\operatorname{lt}_{\Delta y \rightarrow 0} \frac{z_{0}(-i \Delta y)+i \Delta y \bar{z}_{0}+(i \Delta y)(-i \Delta y)}{i \Delta y}=-z_{0}+\bar{z}_{0} \quad \rightarrow(2
$$

Thus, from (1) and (2) for $f^{\prime}\left(z_{0}\right)$ to exists
i.e., $z_{0}=-z_{0} \Rightarrow 2 z_{0}=0 \Rightarrow z_{0} \neq 0$
$\therefore \mathrm{f}^{\prime}(\mathrm{z})$ does not exists though $\mathrm{f}(\mathrm{z})=\left|z^{2}\right|$ is continuous at all z .

## polar form of Cauchy-Riemann equation:

## Theorem:

If $f(z)=f\left(r e^{i \theta}\right)=u(r, \theta)+i v(r, \theta)$ and $\mathrm{f}(\mathrm{z})$ is derivable at $z_{0}=r_{0} e^{i \theta_{0}}$ then
$\frac{\partial u}{\partial r}=\frac{1}{r} \frac{\partial v}{\partial \theta}, \frac{\partial v}{\partial r}=-\frac{1}{r} \frac{\partial u}{\partial \theta}$
Proof: Let $z=r e^{i \theta}$ Then $f(z)=f\left(r e^{i \theta}\right)=u(r, \theta)+i v(r, \theta)$
Differentiating it with respect to r partially,
$\frac{\partial}{\partial r} f(z)=f^{\prime}(z) \frac{\partial z}{\partial r}=f^{\prime}(z) e^{i \theta}$
$\therefore \mathrm{f}^{\mathrm{\prime}}(\mathrm{z})=\frac{1}{e^{i \theta}} \frac{\partial f}{\partial r}=\frac{1}{e^{i \theta}}\left(u_{r}+i v_{r}\right)$

Similarly differentiating partially with respect to $\theta$

$$
\begin{align*}
& \frac{\partial f}{\partial \theta}=f^{\prime}(z) \frac{\partial z}{\partial \theta}=f^{\prime}(z) \cdot r i e^{i \theta} \\
& \therefore f^{\prime}(z)=\frac{1}{r i e^{i \theta}}\left(u_{\theta}+i v_{\theta}\right) \tag{2}
\end{align*}
$$

From (1) and (2) we have
$\frac{1}{e^{i \theta}}\left(u_{r}+i v_{r}\right)=\frac{1}{r i e^{i \theta}}\left(u_{\theta}+i v_{\theta}\right)$
$\therefore u_{r}+i v_{r}=\frac{1}{r} \frac{\partial v}{\partial \theta}-i \frac{1}{r} \frac{\partial u}{\partial \theta}$
Equating real and imaginary parts, we get

$$
\frac{\partial u}{\partial r}=\frac{1}{r} \frac{\partial v}{\partial \theta} \text { and } \frac{\partial v}{\partial r}=-\frac{1}{r} \frac{\partial u}{\partial \theta}
$$

## Analytic function:

A complex function is said to be analytic on a region $R$ if it is complex differentiable at every point in R. The terms holomorphic function, differentiable function, and complex differentiable function are sometimes used interchangeably with "analytic function". Many mathematicians prefer the term "holomorphic function" (or "holomorphic map") to "analytic function" .

If a complex function is analytic on a region $R$, it is infinitely differentiable in $R$.

## Singularities:

A complex function may fail to be analytic at one or more points through the presence of singularities, or along lines or line segments through the presence of branch cuts.

Eg. $\mathrm{f}(\mathrm{z})=\frac{1}{z}$ is analytic every where except at $\mathrm{z}=0$.

At $\mathrm{z}=0 \quad f^{\prime}(z)$ does not exist.
So $\mathrm{z}=0$ is an isolated singular point.

## Entire function:

A complex function that is analytic at all finite points of the complex plane is said to be entire. A single-valued function that is analytic in all but possibly a discrete subset of its domain, and at those singularities goes to infinity like a polynomial (i.e., these exceptional points must be poles and not essential singularities), is called a meromorphic function.

## Cauchy-Riemann equations:

The Cauchy-Riemann equations on a pair of real-valued functions of two real variables $u(x, y)$ and $v(x, y)$ are the two equations:

1. $\frac{\partial u}{\partial x}=\frac{\partial v}{\partial y}$
2. $\frac{\partial u}{\partial y}=-\frac{\partial v}{\partial x}$

Typically $u$ and $v$ are taken to be the real and imaginary parts respectively of a complex-valued function of a single complex variable $z=x+i y, f(x+\mathrm{i} y)=u(x, y)+\mathrm{i} v(x, y)$

## Relation with harmonic functions :

Analytic functions are intimately related to harmonic functions. We say that a real-valued function $\mathrm{h}(\mathrm{x}, \mathrm{y})$ on the plane is harmonic if it obeys Laplace's equation:

$$
\frac{\partial^{2} h}{\partial^{2} x}+\frac{\partial^{2} h}{\partial^{2} y}=0
$$

In fact, as we now show, the real and imaginary parts of an analytic function are harmonic. Let $f$ $=u+i v$ be analytic in some open set of the complex plane.

$$
\text { Then, } \begin{aligned}
\frac{\partial^{2} u}{\partial^{2} x}+\frac{\partial^{2} u}{\partial^{2} y} & =\frac{\partial}{\partial x} \frac{\partial u}{\partial x}+\frac{\partial}{\partial y} \frac{\partial u}{\partial y} \\
& =\frac{\partial}{\partial x} \frac{\partial u}{\partial y}-\frac{\partial}{\partial y} \frac{\partial u}{\partial x} \quad \text { (using Cauchy-Riemann) } \\
& =\frac{\partial^{2} u}{\partial x \partial y}-\frac{\partial^{2} u}{\partial y \partial x} \\
& =0
\end{aligned}
$$

A similar calculation shows that v is also har monic. This result is important in applications because it shows that one can obtain solutions of a second order partial differential equation by solving a system of first order partial differential equations. It is particularly important in this case because we will be able to obtain solutions of the Cauchy-Riemann equations without really solving these equations.

Given a harmonic function $u$ we say that another harmonic function $v$ is its harmonic conjugate if the complex-valued function $f=u+i v$ is analytic.

## Conjugate harmonic function:

If two harmonic functions $u$ and $v$ satisfy the Cauchy-Reimann equations in a domain $D$ and they are real and imaginary parts of an analytic function $f$ in $D$ then $v$ is said to be a conjugate
harmonic function of $u$ in D.If $f(z)=u+i v$ is an analytic function and if $u$ and $v$ satisfy Laplace's equation ,then $u$ and $v$ are called conjugate harmonic functions.

## Polar form of cauchys Riemann equations:

The Cauchy-Riemann equations can be written in other coordinate systems. For instance, it is not difficult to see that in the system of coordinates given by the polar representation $z=r e^{i \theta}$ these equations take the following form:

$$
\begin{aligned}
& \frac{\partial u}{\partial r}=\frac{1}{r} \frac{\partial v}{\partial \theta} \\
& \frac{\partial v}{\partial r}=-\frac{1}{r} \frac{\partial u}{\partial \theta}
\end{aligned}
$$

Problem: Show that the function $\mathrm{f}: \mathrm{C} \rightarrow \mathrm{C}$, defined by $\mathrm{f}(\mathrm{z})=\bar{z}$ does not satisfy the CauchyRiemann equations.

Solution: Indeed, since $u(x, y)=x, v(x, y)=-y$, it follows that $\partial u / \partial x=1$, while $\partial v / \partial y=-1$. So, this function, despite the fact that it is continuous everywhere on C , it is R differentiable on C , is nowhere C -derivable.

Problem: Show that the function $f(z)=e^{z}$ satisfies the Cauchy-Riemann equations.

## Solution:

$$
\text { since } e^{z}=e^{x}(\text { cos } y+i \sin y)
$$

Indeed it follows that

$$
u(x, y)=e^{x} \cos y, \quad v(x, y)=e^{x} \sin y
$$

and $\quad \frac{\partial u}{\partial x}=\mathrm{e}^{\mathrm{x}} \cos \mathrm{y}=\frac{\partial v}{\partial y} ; \quad \frac{\partial u}{\partial y}=\mathrm{e}^{\mathrm{x}} \sin \mathrm{y}=-\frac{\partial v}{\partial x} ;$
Moreover, $\mathrm{e}^{\mathrm{z}}$ is complex derivable and it follows immediately that its complex derivative is $\mathrm{e}^{\mathrm{z}}$.

## Holomorphic functions:

Holomorphic functions are complex functions, defined on an open subset of the complex plane, that are differentiable. In the context of complex analysis, the derivative of $f$ at $z_{0}$ is defined to be $f^{\prime}\left(z_{0}\right)=\lim _{z \rightarrow z_{0}} \frac{f(z)-f\left(z_{0}\right)}{z-z_{0}}, z \in c$.

## Construction of analytic function whose real or imaginary part is known:

Suppose $f(z)=u+i v$ is an analytic function, whose real part $u$ is known.We can find $v$, the imaginary part and also the function $f(z)$.

Problem: Showthat $\left(\frac{\partial^{2}}{\partial x^{2}}+\frac{\partial^{2}}{\partial y^{2}}\right) \log \left|f^{\prime}(z)\right|=0$ where $\mathrm{f}(\mathrm{z})$ is an analytic function.
Solution: Taking $x=\frac{z+\bar{z}}{2}, y=\frac{z-\bar{z}}{2}=\frac{-i}{2}(z-\bar{z})$
We have $\frac{\partial}{\partial z}=\frac{\partial}{\partial x}\left(\frac{\partial x}{\partial z}\right)+\frac{\partial}{\partial y}\left(\frac{\partial y}{\partial z}\right)=\frac{1}{2}\left(\frac{\partial}{\partial x}-i \frac{\partial}{\partial y}\right)$
And $\frac{\partial}{\partial \bar{z}}=\frac{1}{2}\left(\frac{\partial}{\partial x}+i \frac{\partial}{\partial y}\right)$
$\therefore \frac{\partial^{2}}{\partial z \partial \bar{z}}=\frac{1}{2}\left(\frac{\partial}{\partial x}-i \frac{\partial}{\partial y}\right) \cdot \frac{1}{2}\left(\frac{\partial}{\partial x}+i \frac{\partial}{\partial y}\right)=\frac{1}{4}\left(\frac{\partial^{2}}{\partial x^{2}}+\frac{\partial^{2}}{\partial y^{2}}\right)$
Hence $\left(\frac{\partial^{2}}{\partial x^{2}}+\frac{\partial^{2}}{\partial y^{2}}\right)\left(\log \left|f^{\prime}(z)\right|=4 \frac{\partial^{2}}{\partial z \partial \bar{z}}\left(\frac{1}{2} \log \left|f^{\prime}(z)\right|^{2}\right)\right.$

$$
\begin{aligned}
& =2 \frac{\partial^{2}}{\partial z \partial \bar{z}}\left[\left(\log f^{\prime}(z) f^{\prime}(\bar{z})\right)\right] \quad\left(\because|z|^{2}=z \bar{z}\right) \\
& =2 \frac{\partial^{2}}{\partial z \partial \bar{z}}\left[\left(\log f^{\prime}(z)+f^{\prime}(\bar{z})\right)\right] \\
& =2\left[\frac{\partial}{\partial \bar{z}} \frac{f^{\prime \prime}(z)}{f^{\prime}(z)}+\frac{\partial}{\partial z} \frac{f^{\prime \prime}(z)}{f^{\prime}\left(z_{-}\right)}\right] \\
& =2(0+0)=0
\end{aligned}
$$

Since $f(z)$ is analytic, $f(z)$ is analytic, $f^{\prime}(\bar{z})$ is also analytic and $\frac{\partial f^{\prime}(z)}{\partial \bar{z}}=0, \frac{\partial f^{\prime}(\bar{z})}{\partial z}=0$
Problem: Show that $\mathrm{f}(\mathrm{z})=\left\{\begin{array}{ll}\frac{x y(x+i y)}{x^{2}+y^{4}} & , z \neq 0 \\ 0 & , z=0\end{array} \quad\right.$ is not analytic at $\mathrm{z}=0$ although C-R equations satisified at origin.

Solution: $\quad \frac{f(z)-f(0)}{z-0}=\frac{f(z)-0}{z}=\frac{f(z)}{z}$

$$
=\frac{x y^{2}(x+i y)}{\left(x^{2}+y^{4}\right) \cdot z}=\frac{x y^{2}(z)}{\left(x^{2}+y^{4}\right) \cdot z}=\frac{x y^{2}}{\left(x^{2}+y^{4}\right)}
$$

Clearly $\underset{y \rightarrow 0}{\lim } 0 \frac{x y^{2}}{\left(x^{2}+y^{4}\right)}=\underset{y \rightarrow 0}{\lim } 0 \frac{x y^{2}}{\left(x^{2}+y^{4}\right)}=0$
Along path $\mathrm{y}=\mathrm{mx}$

$$
z \rightarrow 0 \frac{\lim _{\rightarrow}}{} \frac{f(z)-f(0)}{z-0}=x \xrightarrow{\lim } 0 \frac{x\left(m^{2} \cdot x^{2}\right)}{x^{2}+m^{4} \cdot x^{4}}=x \xrightarrow{\lim } 0 \frac{m^{2} \cdot x^{2}}{1+m^{4} \cdot x^{2}}=0
$$

Along path $x=m y^{2}$
$z \rightarrow 0 \frac{f(z)-f(0)}{z-0}=y \xrightarrow{\lim } 0 \frac{y^{2}\left(m \cdot y^{2}\right)}{y^{4}+m^{2} \cdot y^{4}}=y \xrightarrow{\lim } 0 \frac{m}{1+m^{2}} \neq 0$
Limit value depends on $m$ i.e on the path of approach and its different for the different paths Followed and therefore limit does not exists.

Hence $f(z)$ is not differentiable at $z=0$. Thus $f(z)$ is not analytic at $z=0$
To prove that C-R conditions are satisified at origin
Let $f(z)=u+i v=\frac{x y^{2}(x+i y)}{\left(x^{2}+y^{4}\right)}$
Then $\mathrm{u}(\mathrm{x}, \mathrm{y})=\frac{x^{2} y^{2}}{\left(x^{2}+y^{4}\right)}$ and $\mathrm{v}(\mathrm{x}, \mathrm{y})=\frac{x y^{3}}{\left(x^{2}+y^{4}\right)}$ for $\mathrm{z} \neq 0$
Also $\mathrm{u}(0,0)=0$ and $\mathrm{v}(0,0)=0 \quad[\because \mathrm{f}(\mathrm{z})=0$ at $\mathrm{z}=0]$
Now $\frac{\partial u}{\partial x}=x \rightarrow \lim _{\rightarrow} 0 \frac{u(x, 0)-u(0,0)}{x}=x \rightarrow 0 \frac{0}{x}=0$

$$
\begin{aligned}
& \frac{\partial u}{\partial y}=y \xrightarrow{\lim _{n} 0 \frac{u(0, y)-u(0,0)}{y}=x \rightarrow 0 \frac{0}{y}=0} \\
& \frac{\partial v}{\partial x}=x \rightarrow 0 \frac{\lim _{\rightarrow}}{} 0 \frac{v(x, 0)-v(0,0)}{x}=x \rightarrow 0 \frac{0}{x}=0
\end{aligned}
$$

$$
\frac{\partial v}{\partial y}=y \xrightarrow{\lim } 0 \frac{v(0, y)-v(0,0)}{y}=x \rightarrow 0 \frac{0}{y}=0
$$

Thus C-R equations are satisified are satisified at the origin
Hence $f(z)$ is not analytic at $z=0$ even C-R equations are satisified at origin.

## Milne Thomson method:

Problem : Find the regular function whose imaginary part is $\log \left(x^{2}+y^{2}\right)+x-2 y$.
Solution: Given $v=\log \left(x^{2}+y^{2}\right)+x-2 y$

$$
\begin{gathered}
\therefore \frac{\partial v}{\partial x}=\frac{2 x}{x^{2}+y^{2}}+1 \quad-\rightarrow(1) \quad \text { and } \frac{\partial v}{\partial y}=\frac{2 y}{x^{2}+y^{2}}-2 \rightarrow(2) \\
\mathrm{f}^{\prime}(\mathrm{z})=\frac{\partial u}{\partial x}+i \frac{\partial v}{\partial x}=\frac{\partial v}{\partial y}+i \frac{\partial v}{\partial x} \quad \text { (Using C-R equation) } \\
\quad=\frac{2 y}{x^{2}+y^{2}}-2+\left(\frac{2 x}{x^{2}+y^{2}}+1\right) \quad \text { (using (1),(2)) }
\end{gathered}
$$

By Milne Thomson method, $\mathrm{f}^{\prime}(\mathrm{z})$ is expressed in terms of z by replacing x z and y by 0 .
Hence $f^{\prime}(z)=-2+i\left(\frac{2 z}{z^{2}}+1\right)=-2+i\left(\frac{2}{z}+1\right)$
On integrating, $f(z)=\int\left[-2+i\left(\frac{2}{z}+1\right)\right] d z+c$

$$
=-2 z+i(2 \log z+z)+c=2 i \log z-(2-i) z+c .
$$

Problem: Show that the function $u=4 x y-3 x+2$ is harmonic .construct the corresponding analytic function $f(z)=u+i v$ in terms of $z$.

Solution: Given $u=4 x y-3 x+2 \rightarrow(1)$
Differentiating (1) partially w.r.t .x, $\frac{\partial u}{\partial x}=4 y-3$

Again differentiating $\frac{\partial^{2} u}{\partial x^{2}}=0$
Again differentiating (1) partially w.r.t .y, $\frac{\partial u}{\partial x}=4 x$
Again differentiating $\frac{\partial^{2} u}{\partial y^{2}}=0$
$\therefore \frac{\partial^{2} u}{\partial x^{2}}+\frac{\partial^{2} u}{\partial y^{2}}=0$
Hence $u$ is Harmonic.
Now $f^{\prime}(z)=\frac{\partial u}{\partial x}+\mathrm{i} \frac{\partial v}{\partial y}=\frac{\partial u}{\partial x}-i \frac{\partial u}{\partial y} \Rightarrow f^{\prime}(z)=4 y-3-i .4 x$
Using Milne Thomson method
$f^{\prime}(z)=-3-i 4 z($ putting $\mathrm{x}=\mathrm{z}$ and $\mathrm{y}=0)$
Integrating, $f(z)=-3 z-i 2 z^{2}+c$

Problem : Find the imaginary part of an analytic function whose real part is $e^{x}(x \cos y-y \sin y)$

Solution: Let $f(z)=u+i v$ where $u=e^{x}(x \cos y-y \sin y)$

$$
\begin{aligned}
f^{\prime}(z)= & \frac{\partial u}{\partial x}+\mathrm{i} \frac{\partial v}{\partial x}=\frac{\partial u}{\partial x}-\mathrm{i} \frac{\partial u}{\partial y} \quad \text { (using C-R equ) } \\
& =\left[e^{x}(x \cos y-y \sin y)+e^{x} \cos y\right]-i\left[e^{x}(-x \sin y-\sin y-y \cos y)\right]
\end{aligned}
$$

By milne's method $f^{\prime}(z)=\left(z e^{z}+e^{z}\right)-i(0)=z e^{z}+e^{z}$
Integrating, We get

$$
\mathrm{f}(\mathrm{z})=\int\left(z e^{z}+e^{z}\right) d z+c=(z-1) e^{z}+e^{z}+c=z e^{z}+c
$$

i.e., $u+i v=(x+i y) e^{x+i y}+c$

$$
\begin{aligned}
& =(x+i y) e^{x} \cdot e^{i y}+c \\
& =e^{x}(x+i y)(\cos y+i \sin y)+c \\
& =e^{x}(x \cos y+i x \sin y+i y \cos y-y \sin y)+c \\
= & e^{x}[(x \cos y-y \sin y)+i(x \sin y+y \cos y)]+c
\end{aligned}
$$

## Bilinear Transformation-Mobius Transformations:

Another important class of elementary mappings was studied by August Ferdinand Möbius (1790-1868). These mappings are conveniently expressed as the quotient of two linear expressions and are commonly known as linear fractional or bilinear transformations. They arise naturally in mapping problems involving the function $\arctan (\mathrm{z})$. In this section, we show how they are used to map a disk one-to-one and onto a half-plane. An important property is that these transformations are conformal in the entire complex plane except at one point. (see Section 10.1)

Let $a, b, c$, and $d$ denote four complex constants with the restriction that $a d \neq b c$. Then the function

$$
\begin{equation*}
w=S(z)=\frac{a z+b}{c z+d} \tag{10-13}
\end{equation*}
$$

is called a bilinear transformation, a Möbius transformation, or a linear fractional transformation.
If the expression for $S(z)$ in Equation (10-13) is multiplied through by the quantity $c z+d$ , then the resulting expression has the bilinear form $c w z-a z+d w-b=0$.
We collect terms involving $z$ and write $z(c w-a)=-d w+b$. Then, for values of $w \neq \frac{a}{c}$ the inverse transformation is given by

$$
\begin{equation*}
z=S^{-1}(w)=\frac{-d w+b}{c w-a} . \tag{10-14}
\end{equation*}
$$

We can extend $S(z)$ and $S^{-1}(w)$ to mappings in the extended complex plane. The value $S(\infty)$ should be chosen to equal the limit of $S(z)$ as $z \rightarrow \infty$. Therefore we define

$$
S(\infty)=\lim _{z \rightarrow \infty} S(z)=\lim _{z \rightarrow \infty} \frac{a+\frac{b}{z}}{c+\frac{d}{z}}=\frac{a}{c}
$$

and the inverse is $s^{-1}\left(\frac{a}{c}\right)=\infty$. Similarly, the value $s^{-1}(\infty)$ is obtained by

$$
S^{-1}(\infty)=\lim _{m \rightarrow \infty} S^{-1}(w)=\lim _{w \rightarrow \infty} \frac{-d+\frac{b}{w}}{c-\frac{2}{w}}=\frac{-d}{c},
$$

and the inverse is $s\left(\frac{-d}{c}\right)=\infty$. With these extensions we conclude that the transformation $w=S(z)$ is a one-to-one mapping of the extended complex z-plane onto the extended complex w-plane.

We now show that a bilinear transformation carries the class of circles and lines onto itself. If $S(z)$ is an arbitrary bilinear transformation given by Equation (10-13) and $c=0$ , then $S(z)$ reduces to a linear transformation, which carries lines onto lines and circles onto circles. If $c \neq 0$, then we can write $S(z)$ in the form

$$
\begin{align*}
S(z) & =\frac{a z+b}{c z+d} \\
& =\frac{c(a z+b)}{c(c z+d)} \\
& =\frac{a c z+b c}{c(c z+d)} \\
& =\frac{a c z+a d-a d+b c}{c(c z+d)} \\
& =\frac{a(c z+d)-a d+b c}{c(c z+d)} \\
& =\frac{a}{c}+\frac{b c-a d}{c} \frac{1}{c z+d} \tag{10-15}
\end{align*}
$$

The condition $a d \neq b c$ precludes the possibility that $S(z)$ reduces to $a$ constant. Equation (10-15) indicates that $S(z)$ can be considered as a composition of functions.

It is a linear mapping $\xi=c z+d$, followed by the reciprocal transformation

$$
z=\frac{1}{\xi}, \text { followed }
$$ by $w=\frac{a}{c}+\frac{b c-a d}{c} z$ . In Section 2.1 we showed that each function in this composition maps the class of circles and lines onto itself; it follows that the bilinear transformation $S(z)$ has this property. A half-plane can be considered to be a family of parallel lines and a disk as a family of circles. Therefore we conclude that a bilinear transformation maps the class of half-planes and disks onto itself. Example 10.3 illustrates this idea.

The general formula for a bilinear transformation (Equation (10-13)) appears to involve four independent coefficients: $a, b, c$, and $d$. But as $S(z)$ is not identically constant, either $a \neq 0$ or $c \neq 0$, we can express the transformation with three unknown coefficients and write either

$$
S(z)=\frac{z+\frac{b}{2}}{\frac{c}{2} z+\frac{d}{2}} \quad \text { or } \quad S(z)=\frac{\frac{2}{c} z+\frac{b}{c}}{z+\frac{d}{c}}
$$

respectively. Doing so permits us to determine a unique a bilinear transformation if three distinct image values $S\left(z_{1}\right)=w_{1}, \quad S\left(z_{z}\right)=w_{2}$, and $S\left(z_{3}\right)=w_{3}$ are specified. To determine such a mapping, we can conveniently use an implicit formula involving z and w .

Theorem 10.3 (The Implicit Formula). There exists a unique bilinear transformation that maps three distinct points $z_{1}, z_{2}$, and $z_{3}$ onto three distinct points $w_{1}, w_{2}$, and $w_{3}$, respectively. An implicit formula for the mapping is given by the equation

$$
\begin{equation*}
\frac{\left(z-z_{1}\right)\left(z_{2}-z_{3}\right)}{\left(z-z_{3}\right)\left(z_{2}-z_{1}\right)}=\frac{\left(w-w_{1}\right)\left(w_{2}-w_{3}\right)}{\left(w-w_{3}\right)\left(w_{2}-w_{1}\right)} . \tag{10-18}
\end{equation*}
$$

Example 1. Construct the bilinear transformation $w=S(z)$ that maps the points $z_{1}=-\dot{1}, z_{i}=1, z_{3}=\dot{I}$ onto the points $w_{1}=-1, w_{2}=0, w_{3}=1$, respectively.



Solution. We use the implicit formula, Equation (10-18), and write

$$
\begin{aligned}
& \frac{(z-(-\dot{I}))(1-\dot{I})}{(z-\dot{I})(1-(-\dot{I}))}=\frac{(w-(-1))(0-1)}{(w-1)(0-(-1))} \\
& \frac{(z+\dot{i})(1-\dot{I})}{(z-\dot{I})(1+\dot{I})}=\frac{(w+1)(0-1)}{(w-1)(0+1)} \\
& \frac{(z+\dot{i})(1-\dot{I})}{(z-\dot{I})(1+\dot{I})}=\frac{w+1}{-w+1} .
\end{aligned}
$$

Expanding this equation, collecting terms involving w and zw on the left and then simplify.

```
(z-\dot{I})(l+il)(w+l)=(z+\dot{I})(l-\dot{I})(-W+l)
```



```
=
    (-1+ i!) zw+(-1- i})w+(1-\dot{I}) z+(l+\dot{I}
zw+ i! zw+w- i}w+z+\dot{I}z+l-\dot{I
=
```



```
2zw+2w= - 2íz z+2i
    zw + W = - İ z + ì 
W(l+z)= íi(l-z)
```

Therefore the desired bilinear transformation is

$$
\mathrm{w}=S(z)=\frac{\dot{1}(1-z)}{1+z} .
$$

Example 2. Find the bilinear transformation $w=S(z)$ that maps the points $z_{1}=-2, \quad z_{2}=-1-i i_{,}, z_{3}=0$ onto the points $w_{1}=-1, \quad w_{2}=0, w_{3}=1$, respectively.



Solution. Again, we use the implicit formula, Equation (10-18), and write

$$
\begin{aligned}
& \frac{(z-(-2))((-1-\dot{I})-0)}{(z-0)((-1-\dot{I})-(-2))}=\frac{(w-(-1))(0-1)}{(w-1)(0-(-1))} \\
& \frac{(z+2)(-1-\dot{\text { i }})}{(z)(-1-\dot{I}+2)}=\frac{(w+1)(-1)}{(w-1)(1)} \\
& \frac{z+2}{z} \frac{-1-\dot{I}}{1-\dot{I}}=\frac{1+w}{1-w}
\end{aligned}
$$

Using the fact that $\frac{-1-\dot{\underline{i}}}{1-\dot{I}}=\frac{1}{\dot{I}}$, we rewrite this equation as

$$
\frac{z+2}{\dot{\text { i } z}}=\frac{1+w}{1-w} .
$$

We now expand the equation and obtain

$$
\begin{aligned}
& (z+2)(1-w)=\dot{\text { in }} z(1+w) \\
& z+2-z w-2 w=\dot{1} z+\dot{\text { in }} \mathbf{z} w \\
& z-\dot{1} z+2=z w+\dot{1} z w+2 w \\
& (1-\dot{\text { II }}) z+2=w(z+\dot{I} z+2) \\
& (1-\dot{I}) z+2=w((1+\dot{I}) z+2)
\end{aligned}
$$

which can be solved for $w$ in terms of z , giving the desired solution

$$
W=S(z)=\frac{(l-\dot{I}) z+2}{(1+\dot{\text { I}}) z+2}
$$

Corollary (The Implicit Formula with a point at Infinity). In equation (10-18) the point at infinity can be introduced as one of the prescribed points in either the z plane or the w plane.

## Proof.

Case 1. If $z_{3}=\infty$, then we can write $\frac{\left(z_{\varepsilon}-z_{3}\right)}{\left(z-z_{3}\right)}=\frac{\left(z_{2}-\infty\right)}{(z-\infty)}=1$ and substitute this expression into Equation (10-18) to obtain $\frac{\left(z-z_{1}\right)\left(z_{2}-\infty\right)}{(z-\infty)\left(z_{2}-z_{1}\right)}=\frac{\left(w-w_{1}\right)\left(w_{2}-w_{3}\right)}{\left(w-w_{3}\right)\left(w_{2}-w_{1}\right)}$ which can be rewritten as $\frac{\left(z-z_{1}\right)\left(z_{\varepsilon}-\infty\right)}{\left(z_{\varepsilon}-z_{1}\right)(z-\infty)}=\frac{\left(w-w_{1}\right)\left(w_{\varepsilon}-w_{3}\right)}{\left(w-w_{3}\right)\left(w_{2}-w_{1}\right)}$ and simplifies to obtain

$$
\frac{z-z_{1}}{z_{2}-z_{1}}=\frac{\left(w-w_{1}\right)\left(w_{2}-w_{3}\right)}{\left(w-w_{3}\right)\left(w_{2}-w_{1}\right)} .
$$

Case 2. If $w_{3}=\infty$, then we can write $\frac{\left(w_{2}-w_{3}\right)}{\left(w^{2}-w_{3}\right)}=\frac{\left(w_{2}-\infty\right)}{\left(w_{2}-\infty\right)}=1$ and substitute this expression into Equation (10-18) to obtain $\frac{\left(z-z_{1}\right)\left(z_{2}-z_{3}\right)}{\left(z-z_{3}\right)\left(z_{2}-z_{1}\right)}=\frac{\left(w-w_{1}\right)\left(w_{2}-\infty\right)}{(w-\infty)\left(w_{2}-w_{1}\right)}$ which can be rewritten as $\frac{\left(z-z_{1}\right)\left(z_{2}-z_{3}\right)}{\left(z-z_{3}\right)\left(z_{2}-z_{1}\right)}=\frac{\left(w-w_{1}\right)\left(w_{2}-\infty\right)}{\left(w_{2}-w_{1}\right)(w-\infty)}$ and simplifies to obtain $\frac{\left(z-z_{1}\right)\left(z_{2}-z_{3}\right)}{\left(z-z_{3}\right)\left(z_{2}-z_{1}\right)}=\frac{w-w_{1}}{w_{2}-w_{1}}$.

Example 1: Find the bilinear transformation $\mathrm{w}=\mathrm{S}(\mathrm{z})$ that maps the points

$$
z_{1}=0, z_{i}=\dot{1} \text {, and } z_{3}=-\dot{1} \text { onto } w_{1}=-1, w_{2}=1 \text {, and } w_{3}=0 \text {, respectively. }
$$

Solution. Method I. Use the implicit formula $\frac{\left(z-z_{1}\right)\left(z_{2}-z_{3}\right)}{\left(z-z_{3}\right)\left(z_{2}-z_{1}\right)}=\frac{\left(w-w_{1}\right)\left(w_{2}-w_{3}\right)}{\left(w-w_{3}\right)\left(w_{2}-w_{1}\right)}$.
Substitute the values given above and get

$$
\frac{(z-0)(\dot{I}+\dot{\text { II }})}{(z+\dot{\text { II}})(\dot{I}-0)}=\frac{(w+1)(1-0)}{(w-0)(1+1)},
$$

then simplify and get

$$
\frac{2 z}{z+\dot{I}}=\frac{W+1}{2 w} .
$$

Solving for w we obtain

$$
\begin{aligned}
& (2 z)(2 w)=(z+\dot{I})(w+1) \\
& 4 W z=W(z+\dot{I})+z+\dot{I} \\
& W(4 z-(z+\dot{\text { II}}))=z+\dot{\text { II }} \\
& W(3 z-\dot{I})=z+\dot{I} \\
& W=\frac{z+\dot{\text { II }}}{3 z-\dot{\text { II }}}
\end{aligned}
$$

Therefore,
$W=S(z)=\frac{z+\dot{\text { I }}}{3 z-\dot{\text { II }}}$
Solution. Method II. The general form of a bilinear transformation is

$$
w=S(z)=\frac{a z+b}{c z+d}, \quad \text { and it is not the case that both } a=0 \text { and } c=0
$$

So the desired formula must have one of the following two forms:
either

Let us assume that the first form

$$
S(z)=\frac{z+b}{c z+d} \quad \text { is the one that works out. }
$$

Then we can set up three equations to solve $\frac{z_{k}+b}{c z_{k}+d}=w_{k}$ for $k=1,2,3$ :

$$
\frac{0+\mathrm{b}}{\mathrm{c} * 0+\mathrm{d}}=-1, \quad \frac{\dot{\operatorname{li}}+\mathrm{b}}{\mathrm{c} * \dot{\operatorname{l}}+\mathrm{d}}=1, \quad \frac{-\dot{\operatorname{li}}+\mathrm{b}}{\mathrm{c} *(-\dot{\mathrm{I}})+\mathrm{d}}=0
$$

then simplify these equations get

$$
\mathrm{b}=-\mathrm{d}, \quad \dot{\text { I }}+\mathrm{b}=\dot{\text { II }} \mathrm{c}+\mathrm{d}, \quad-\dot{\text { I }}+\mathrm{b}=0 .
$$

The last equation is easy to solve and we get $b=\dot{I}$ and then the first equation
yields $\quad d=-b=-\dot{\text { II }}$.
Use these values to rewrite the second equation as $\dot{\mathbf{I}}+\dot{I}=\dot{I} \mathbf{c}-\dot{I}$ and then obtain $\mathrm{c}=3$.

Substituting these into $S(z)=\frac{z+b}{c z+d}$ produces the desired result:

$$
w=S(z)=\frac{z+\dot{I}}{3 z-\dot{I}} .
$$

Example 2: Find the bilinear transformation $w=S(z)$ that maps the points

$$
z_{1}=-\dot{1}, z_{2}=0 \text {, and } z_{3}=\dot{I} \quad \text { onto } w_{1}=-1, w_{2}=\dot{1} \text {, and } w_{3}=1 \text {, respectively. }
$$

Solution. Method I. Use the implicit formula $\frac{\left(z-z_{1}\right)\left(z_{2}-z_{3}\right)}{\left(z-z_{3}\right)\left(z_{2}-z_{1}\right)}=\frac{\left(w-w_{1}\right)\left(w_{2}-w_{3}\right)}{\left(w-w_{3}\right)\left(w_{2}-w_{1}\right)}$.
Substitute the values given above and get

$$
\begin{aligned}
& \frac{(z-(-\dot{I}))(0-\dot{I})}{(z-\dot{I})(0-(-\dot{I}))}=\frac{(w-(-1))(\dot{I}-1)}{(w-1)(\dot{I}-(-1))}, \\
& \frac{(z+\dot{I})(0-\dot{\text { I }})}{(z-\dot{I})(0+\dot{I})}=\frac{(w+1)(\dot{I}-1)}{(w-1)(\dot{I}+1)},
\end{aligned}
$$

then simplify and get

$$
-\frac{z+\dot{I}}{z-\dot{I}}=\frac{\dot{1} w+\dot{I}}{w-1} .
$$

Solving for w we obtain

$$
\begin{aligned}
& -(z+\dot{I})(w-1)=(z-\dot{I})(\dot{I} w+\dot{I}) \\
& -w z-\dot{I} w+z+\dot{I}=w+\dot{I} w z+\dot{I} z+1 \\
& -w z-\dot{I} w z-w-\dot{I} w=-z+\dot{I} z+1-\dot{I} \\
& w(-z-\dot{I} z-1-\dot{I})=-z+\dot{I} z+1-\dot{I}
\end{aligned}
$$

$$
\begin{aligned}
& W=\frac{-z+\dot{\text { II }} z+1-\dot{\text { in }}}{-z-\dot{\text { II }} z-1-\dot{\text { in }}} \\
& W=\frac{z(-1+\dot{\text { II }})+1-\dot{\text { II }}}{z(-1-\dot{\text { II }})-1-\dot{\text { II }}} \\
& W=\frac{(-1+\dot{\text { II }})(z-1)}{(-1-\dot{I})(z+1)} \\
& W=\left(\frac{-1+\dot{\text { II }}}{-1-\dot{\text { II }}}\right) \frac{(z-1)}{(z+1)} \\
& W=(-\dot{I}) \frac{(z-1)}{(z+1)} \\
& W=\frac{-\dot{\text { II }}(z-1)}{z+1}
\end{aligned}
$$

Therefore,

$$
w=S(z)=\frac{-\dot{\text { I }} z+\dot{\text { İ }}}{z+1} .
$$

Solution. Method II. The general form of a bilinear transformation is

$$
w=S(z)=\frac{a z+b}{c z+d}, \quad \text { and it is not the case that both } a=0 \text { and } c=0
$$

So the desired formula must have one of the following two forms:
either

$$
S(z)=\frac{1 z+\frac{b}{2}}{\frac{c}{2} z+\frac{d}{2}}=\frac{z+b_{1}}{c_{1} z+d_{1}} \quad \text { or } \quad S(z)=\frac{\frac{2}{c} z+\frac{b}{c}}{\frac{c}{c} z+\frac{d}{c}}=\frac{a_{\varepsilon} z+b_{\varepsilon}}{z+d_{\varepsilon}}
$$

Let us assume that the first form $S(z)=\frac{z+b}{c z+d}$ is the one that works out.
Then we can set up three equations to solve $\frac{z_{k}+b}{c z_{k}+d}=w_{k}$ for $k=1,2,3$ :

$$
\frac{-\dot{I}+b}{c *(-\dot{I})+d}=-1, \quad \frac{0+b}{c * 0+d}=\dot{\text { in }}, \quad \frac{\dot{I}+b}{c * \dot{I}+d}=1
$$

then simplify these equations get the system of equations

```
b - \dot{I}\textrm{C}}+\textrm{d}=\dot{\mathrm{ I}
b - \dot{I}d=0
b - \dot{I}c-d=-\dot{I}
```

Add row 1 to row 3 and get

$$
\begin{aligned}
\mathrm{b}-\dot{\text { in }}+\mathrm{d} & =\dot{\mathbf{I}} \\
\mathrm{b}-\dot{\text { I }} \mathrm{d} & =0 \\
2 \mathrm{~b}-\dot{\text { in } 2 \mathrm{c}} & =0
\end{aligned}
$$

Divide row 2 by 1 and subtract it from row 1 to get

| $d$ | $=\dot{\text { i }}$ |
| ---: | :--- |
| $b-$ ì $d$ | $=0$ |
| $b-\dot{I} c$ | $=0$ |

Use $d=\dot{I}$ to rewrite the second equation as $b-\dot{I} * \dot{I}=0$ and then obtain $b=-1$.
Use $b=-1$ to rewrite the third equation as $-1-\dot{\text { in }} \mathbf{c}=0$ and then obtain $\mathrm{c}=$ in .

Substituting these into $S(z)=\frac{z+b}{c z+d}$ produces the desired result:

$$
w=S(z)=\frac{z-1}{\dot{\underline{I}} z+\dot{\text { I }}}=\frac{-\dot{\operatorname{I}} z+\dot{\mathrm{I}}}{z+1} .
$$

Example 3: Find the bilinear transformation $w=S(z)$ that maps the points $z_{1}=0, z_{2}=1$, and $z_{3}=2$ onto $w_{1}=0, w_{2}=1$, and $w_{3}=\infty$, respectively.
Solution. Method I. Use the implicit formula $\frac{\left(z-z_{1}\right)\left(z_{2}-z_{3}\right)}{\left(z-z_{3}\right)\left(z_{2}-z_{1}\right)}=\frac{\left(w-w_{1}\right)}{\left(w_{2}-w_{1}\right)}$.
Substitute the values given above and get $\frac{(z-0)(1-2)}{(z-2)(1-0)}=\frac{(w-0)}{(1-0)}$,
then simplify and obtain $\frac{z}{2-z}=\frac{w}{1}$.
Therefore, ${ }^{w}=S(z)=\frac{z}{2-z}$.
Solution. Method II. The general form of a bilinear transformation is

$$
w=S(z)=\frac{a z+b}{c z+d}, \quad \text { and it is not the case that both } a=0 \text { and } c=0 .
$$

So the desired formula must have one of the following two forms:
either

$$
S(z)=\frac{l z+\frac{b}{2}}{\frac{c}{2} z+\frac{d}{2}}=\frac{z+b_{1}}{c_{1} z+d_{1}} \quad \text { or } \quad S(z)=\frac{\frac{2}{c} z+\frac{b}{c}}{\frac{c}{c} z+\frac{d}{c}}=\frac{a_{\varepsilon} z+b_{2}}{z+d_{\varepsilon}} .
$$

Let us assume that the first form $s(z)=\frac{z+b}{c z+d}$ is the one that works out.
Then we can set up three equations to solve $\frac{z_{k}+b}{c z_{k}+d}=w_{k}$ for $k=1,2,3$ :

$$
\frac{0+b}{c * 0+d}=0, \frac{l+b}{c * l+d}=1, \frac{2+b}{c * 2+d}=\infty
$$

In the third equation we will take reciprocals and write it as $\frac{c * 2+d}{2+b}=0$, then we have

$$
\frac{0+b}{c * 0+d}=0, \frac{l+b}{c * l+d}=1, \frac{c * 2+d}{2+b}=0
$$

then simplify these equations get

$$
\mathrm{b}=0,1+\mathrm{b}=\mathrm{c}+\mathrm{d}, 2 \mathrm{c}+\mathrm{d}=0
$$

Use $b=0$ to rewrite the second equation as $l=c+d$ then solve the system of two equations

$$
\begin{aligned}
c+d & =1 \\
2 c+d & =0
\end{aligned}
$$

Subtracting the first equation from the second equation and get $\quad c=-1$.
Use $c=-1$ in the first equation and get $d=2$.

Substituting these into

$$
s(z)=\frac{z+b}{c z+d} \text { produces the desired result: }
$$

$$
w=s(z)=\frac{z+0}{-z+2}=\frac{z}{2-z} .
$$

Example 4: Find the bilinear transformation $w=S(z)$ that maps the points

$$
z_{1}=1, z_{2}=\dot{1} \text {, and } z_{3}=-1 \text { onto } w_{1}=0, w_{2}=1 \text {, and } w_{3}=\infty \text {, respectively. }
$$

Solution. Method I. Use the implicit formula $\frac{\left(z-z_{1}\right)\left(z_{\varepsilon}-z_{3}\right)}{\left(z-z_{3}\right)\left(z_{2}-z_{1}\right)}=\frac{\left(w-w_{1}\right)}{\left(w_{2}-w_{1}\right)}$.

Substitute the values given above and get

$$
\begin{aligned}
& \frac{(z-1)(\dot{I}-(-1))}{(z-(-1))(\dot{I}-1)}=\frac{(w-0)}{(1-0)} \\
& \frac{(z-1)}{(z+1)} \frac{(\dot{I}+1)}{(\dot{I}-1)}=\frac{W}{1} \\
& \frac{(z-1)}{(z+1)}(-i)=w \\
& \frac{(-\dot{1})(z-1)}{(z+1)}=W \\
& \frac{\dot{I}(1-z)}{1+z}=W \\
& \text { Therefore, } w=S(z)=\frac{\dot{\text { i }}-\dot{\text { I }} z}{1+z} \text {. }
\end{aligned}
$$

Solution. Method II. The general form of a bilinear transformation is

$$
w=S(z)=\frac{a z+b}{c z+d}, \quad \text { and it is not the case that both } a=0 \text { and } c=0 .
$$

So the desired formula must have one of the following two forms:
either

$$
S(z)=\frac{1 z+\frac{b}{2}}{\frac{c}{2} z+\frac{d}{2}}=\frac{z+b_{1}}{c_{1} z+d_{1}} \quad \text { or } \quad S(z)=\frac{\frac{2}{c} z+\frac{b}{c}}{\frac{c}{c} z+\frac{d}{c}}=\frac{a_{\varepsilon} z+b_{2}}{z+d_{\varepsilon}} .
$$

Let us assume that the first form $S(z)=\frac{z+b}{c z+d}$ is the one that works out.
Then we can set up three equations to solve $\frac{z_{k}+b}{c z_{k}+d}=w_{k}$ for $k=1,2,3$ :

$$
\frac{1+b}{c * l+d}=0, \frac{\dot{I}+b}{c * \dot{I}+d}=1, \frac{-1+b}{c *(-1)+d}=\infty
$$

In the third equation we will take reciprocals and write it as $\frac{-c+d}{-1+b}=0$, then we have

$$
\frac{l+b}{c * l+d}=0, \frac{\dot{I}+b}{c * \dot{I}+d}=1, \frac{-c+d}{-1+b}=0,
$$

then simplify these equations get

$$
\mathrm{l}+\mathrm{b}=0, \quad \dot{\mathrm{I}}+\mathrm{b}=\mathrm{c} * \dot{\mathrm{I}}+\mathrm{d}, \quad-\mathrm{c}+\mathrm{d}=0 .
$$

The first equation is easy to solve and we get $b=-1$.
Use $b=-1$ to rewrite the second equation as $\dot{1}-1=\dot{I} c+d$ then solve the system of two equations

$$
\begin{aligned}
& \text { ï } c+d=-1+\text { 프 } \\
& -c+d=0
\end{aligned}
$$

Subtract the second equation from the first equation obtain ( $1+\dot{\mathrm{i}}) \mathrm{c}=-1+\dot{\mathrm{i}}$ and get $\mathrm{c}=\mathrm{i}$. Use $c=\dot{1}$ in the second equation and get $d=$ it .

Substituting these into

$$
s(z)=\frac{z+b}{c z+d} \text { produces the desired result: }
$$

$$
w=S(z)=\frac{z-1}{\dot{i} z+\dot{\text { in }}}=\frac{\dot{\text { in }}-\dot{\text { in }} z}{1+z} .
$$

## Fixed Point:

A fixed point of a mapping $w=f(z)$ is a point $z_{0}$ such that $f\left(z_{0}\right)=z_{0}$.
Example 1: Show that a bilinear transformation, $\quad w=f(z)=\frac{a z+b}{c z+d}$, can have at most two fixed points.

## Solution .

The equation $z=\frac{a z+b}{c z+d}$ can be written as

$$
c z^{2}+(d-a) z-b=0
$$

and this quadratic equation has, at most, two distinct solutions:

$$
z=\frac{a-d+\sqrt{(d-a)^{2}+4 b c}}{2 c}
$$

and

$$
z=\frac{a-d-\sqrt{(d-a)^{2}+4 b c}}{2 c}
$$

Example 2: Find the fixed points of (a). $\quad w=S(z)=\frac{z-1}{z+1} .(b) . \quad w=S(z)=\frac{4 z+3}{2 z-1}$.

Solution. (a). Solve the equation $z=\frac{z-1}{z+1}$ for $z$ and get

$$
\begin{aligned}
& z=\frac{z-1}{z+1} \\
& z(z+1)=z-1 \\
& z^{2}+z=z-1 \\
& z^{2}=-1 \\
& z= \pm \dot{1}
\end{aligned}
$$

Therefore, the fixed points of $S(z)=\frac{z-1}{z+1}$ are $z=+\dot{i}$ and $z=-\dot{\text { in }}$.
Just for fun, we can substitute $z= \pm \dot{\text { in }}$ into the formula $S(z)=\frac{z-1}{z+1}$.

$$
\begin{aligned}
& \frac{-\dot{\text { I }}-1}{-\dot{I}+1}=\frac{-1-\dot{I}}{1-\dot{I}}=\frac{(-1-\dot{\text { I }})(1+\dot{\text { I }})}{(1-\dot{\text { I }})(1+\dot{\text { I }})}=\frac{-2 \dot{\text { I }}}{2}=-\dot{\text { II }}, \quad \text { and } \\
& \frac{\dot{\text { i }}-1}{\dot{I}+1}=\frac{-1+\dot{I}}{1+\dot{I}}=\frac{(-1+\dot{I})(1-\dot{I})}{(1+\dot{I})(1-\dot{I})}=\frac{2 \dot{\text { i }}}{2}=\dot{I} .
\end{aligned}
$$

(b). Solution. Solve the equation $z=\frac{4 z+3}{2 z-1}$ for $z$ and get

$$
\begin{aligned}
& z=\frac{4 z+3}{2 z-1} \\
& z(2 z-1)=4 z+3 \\
& 2 z^{x}-z=4 z+3 \\
& 2 z^{2}-5 z-3=0 \\
& (z z+1 y\{z-3)=0 \\
& 2\left(z+\frac{1}{2}\right)(z-3)=0
\end{aligned}
$$

Therefore, the fixed points of $S(z)=\frac{4 z+3}{2 z-1}$ are $\quad z=-\frac{1}{2}$ and $z=3$.
Just for fun, we can substitute $\quad z=-\frac{1}{2}$ and $z=3$ into the formula $S(z)=\frac{4 z+3}{2 z-1}$.

$$
\begin{aligned}
& \frac{4\left(-\frac{1}{2}\right)+3}{2\left(-\frac{1}{2}\right)-1}=\frac{-2+3}{-1-1}=-\frac{1}{2}, \quad \text { and } \\
& \frac{4(3)+3}{2(3)-1}=\frac{12+3}{6-1}=\frac{15}{5}=3 .
\end{aligned}
$$

## EXCERCISE PROBLEMS:

1)Show that the real part of an analytic function $\mathrm{f}(\mathrm{z})$ where $\mathrm{u}=e^{-2 x y} \sin \left(x^{2}-y^{2}\right)$ is a harmonic function. Hence find its harmonic conjugate.
2) Prove that the real part of analytic function $\mathrm{f}(\mathrm{z})$ where $\mathrm{u}=\log |z|^{2}$ is harmonic function. If so find the analytic function by Milne Thompson method.
3)Obtain the regular function $\mathrm{f}(\mathrm{z})$ whose imaginary part of an analytic function is $\frac{x-y}{x^{2}+y^{2}}$
4) Find an analytic function $\mathrm{f}(\mathrm{z})$ whose real part of an analytic function is $\mathrm{u}=\frac{\sin 2 x}{\cosh 2 y-\cos 2 x}$ by Milne-Thompson method.
5) Find an analytic function $f(z)=u+i v$ if the real part of an analic function is $u=a$ $(1+\cos \theta)$ using Cauchy-Riemann equations in polar form.
6) Prove that if $u=x^{2}-y^{2}, v=-\frac{y}{x^{2}+y^{2}}$ both $u$ and $v$ satisfy Laplace's equation, but $u+$ iv is not a regular (analytic) function of $z$.
7)Show that the function $\mathrm{f}(\mathrm{z})=\sqrt{|x y|}$ is not analytic at the origin although Cauchy Riemann equations are satisfied at origin.
8) If $w=\emptyset+i \varphi$ represents the complex potential for an electric field where $\varphi=x^{2}-y^{2}+\frac{x}{x^{2}+y^{2}}$ then determine the function $\varphi$.
9)State and Prove the necessary condition for $\mathrm{f}(\mathrm{z})$ to be an analytic function in Cartesian form.
10)If $u$ and $v$ are conjugate harmonic functions then show that $u v$ is also a harmonic function.
11)Find the orthogonal trajectories of the family of curves $r^{2} \cos 2 \theta=\mathrm{c}$
12)Find an analytic function whose real part is $\mathrm{u}=\frac{\sin 2 x}{\cosh 2 y-\cos 2 x}$
13)Find an analytic function whose imaginary part is $\mathrm{v}=e^{x}(x \sin y+y \cos y)$
14)Find an analytic function whose real part is (i) $\mathrm{u}=\frac{x}{x^{2}+y^{2}}$ (ii) $\mathrm{u}=\frac{y}{x^{2}+y^{2}}$
15)Find an analytic function whose imaginary part is $\mathrm{v}=\frac{2 \sin x \sin y}{\cosh 2 x+\cosh 2 y}$
16) Find an analytic function $f(z)=u+i v$ if $u=a(1+\cos \theta)$
17)Find the conjugate harmonic of $u=e^{x^{2}-y^{2}} \cos 2 x y$ and find $f(z)$ in terms of $z$.
18)If $f(z)$ is an analytic function of $z$ and if $u-v=e^{x}(\cos y-\sin y)$ find $f(z)$ in terms of $z$.
19)If $f(z)$ is an analytic function of $z$ and if $\left.u-v=(x-y)\left(x^{2}+4 x y+y^{2}\right)\right)$ find $f(z)$ in terms of $z$.
20) Find the orthogonal trajectories of the family of curves $x^{3} y-x y^{3}=\mathrm{C}=$ constant

## COMPLEX INTEGRATION

## LINE INTEGRAL:

Defination: In mathematics, a line integral is an integral where the function to be integrated is evaluated along a curve. The terms path integral, curve integral, and curvilinear integral are also used; contour integral as well, although that is typically reserved for line integrals in the complex plane.

The function to be integrated may be a scalar field or a vector field. The value of the line integral is the sum of values of the field at all points on the curve, weighted by some scalar function on the curve (commonly arc length or, for a vector field, the scalar product of the vector field with a differential vector in the curve). This weighting distinguishes the line integral from simpler integrals defined on intervals. Many simple formulae in physics (for example, $W=\mathbf{F} \cdot \mathbf{s}$ ) have natural continuous analogs in terms of line integrals ( $W=\int_{C} \mathbf{F} \cdot \mathrm{~d} \mathbf{s}$ ). The line integral finds the work done on an object moving through an atomic or gravitational field.

In complex analysis, the line integral is defined in terms of multiplication and addition of complex numbers.

Let us consider $\mathrm{F}(\mathrm{t})=\mathrm{u}(\mathrm{t})+\mathrm{i} \mathrm{v}(\mathrm{t}), a \leq t \leq b$. Where u and v are real valued continuous functions of $t$ in $[a, b]$.

$$
\text { we define } \int_{a}^{b} F(t) d t=\int_{a}^{b} u(t) d t+i \int_{a}^{b} v(t) d t
$$

Thus, $\int_{a}^{b} F(t) d t$ is a complex number such that real part of $\int_{a}^{b} F(t) d t$ is $\int_{a}^{b} u(t) d t$ and imaginary part of $\int_{a}^{b} F(t) d t$ is $\int_{a}^{b} v(t) d t$.

Problem: Evaluate $\int_{0}^{1+i}\left(x^{2}-i y\right) d z \quad$ along the paths 1$\left.) \mathrm{y}=\mathrm{x} \quad 2\right) \mathrm{y}=\mathrm{x}^{2}$
Solution: 1)along the line $y=x, d y=d x$ so that $d z=d x+i d x=(1+i) d x$

$$
\therefore \int_{0}^{1+i}\left(x^{2}-i y\right) d z=\int_{0}^{1}\left(x^{2}-i x\right)(1+i) d x, \quad \text { since } \mathrm{y}=\mathrm{x}
$$

$$
\begin{aligned}
& =(1+\mathrm{i})\left[\frac{x^{3}}{3}-i \frac{x^{2}}{2}\right]_{0}^{1} \\
& =(1+\mathrm{i})\left[\frac{1}{3}-\frac{1}{2} i\right] \\
& =\frac{5}{6}-\frac{1}{6} i
\end{aligned}
$$

2) Along the parabola $y=x^{2}, d y=2 x d x$ so that $d z=d x+2 i x d x$

$$
\mathrm{dz}=(1+2 \mathrm{ix}) \mathrm{dx} \text { and } \mathrm{xvaries} \text { from } 0 \text { to } 1
$$

$$
\begin{aligned}
\therefore \int_{0}^{1+i}\left(x^{2}-i y\right) d z & =\int_{0}^{1}\left(x^{2}-i x^{2}\right)(1+2 i x) d x \\
& =(1-i) \int_{0}^{1} x^{2}(1+2 i x) d x \\
& =(1-i)\left(\frac{1}{3}+\frac{1}{2} i\right) \\
& =\frac{(1-i)(2+3 i)}{6} \\
& =\frac{5}{6}+\frac{1}{6} i
\end{aligned}
$$

Problem: Evaluate $\int_{z=0}^{z=1+i}\left(x^{2}+2 x y+i\left(y^{2}-x\right)\right) d z$ along $\mathrm{y}=\mathrm{x}^{2}$
Solution: Given $\left.\mathrm{f}(\mathrm{z})=\mathrm{x}^{2}+2 x y+i\left(y^{2}-x\right)\right) d z$

$$
\mathrm{Z}=\mathrm{x}+\mathrm{iy}, \mathrm{dz}=\mathrm{dx}+\mathrm{idy}
$$

$\therefore$ thecurvey $=x^{2}, d y=2 x d x$

$$
\therefore d z=d x+2 x i d x=(1+2 i x) d x
$$

$$
f(z)=x^{2}+2 x\left(x^{2}\right)+i\left(x^{4}-x\right)
$$

$$
=x^{2}+2 x^{3}+i\left(x^{4}-x\right)
$$

$\left.f(z) d z=\left(x^{2}+2 x^{3}\right)+i\left(x^{4}-x\right)(1+2 i x)\right) d x$

$$
\begin{gathered}
=\mathrm{x}^{2}+2 \mathrm{x}^{3}+\mathrm{i}\left(\mathrm{x}^{4}-\mathrm{x}\right)+2 \mathrm{ix}^{3}+4 \mathrm{ix}^{4}-2 \mathrm{x}^{5}+2 \mathrm{x}^{2} \\
\therefore \int_{c} f(z) d z=\int_{z=0}^{1+i} x^{2}+2 x y+i\left(y^{2}-x\right) d z \\
\left.\quad=\int_{0}^{1}\left(-2 x^{5}+3 x^{2}+2 x^{3}+i\left(5 x^{4}-x+2 x^{3}\right)\right) d x\right)
\end{gathered}
$$

$$
=\left[-\frac{x^{6}}{3}+x^{3}+\frac{x^{4}}{2}+i\left(\frac{5 x^{5}}{5}-\frac{x^{2}}{2}+\frac{x^{4}}{2}\right]_{0}^{1}\right.
$$

$$
=\left[\left(\frac{-1}{3}+1+\frac{1}{2}\right)+\left(\frac{5}{5}-\frac{1}{2}+\frac{1}{2}\right)\right]-0
$$

$$
=\frac{7}{6}+\frac{5}{5} i=\frac{7}{6}+i
$$

$$
\int_{c} f(z) d z=\frac{7}{6}+i
$$

Cauchy-Goursat Theorem: Let $f(z)$ be analytic in a simply connected domain D. If C is a simple closed contour that lies in $D$, then

$$
\int_{\mathrm{C}} \mathrm{f}(\mathrm{z}) \mathrm{dlz}=0
$$

Let us recall that $\mathbb{E}^{x}, \cos (z)$, and $z^{n}$ (where n is a positive integer) are all entire functions and have continuous derivatives. The Cauchy-Goursat theorem implies that, for any simple closed contour,
(a)
(b)

$$
\int_{C} \mathbb{E}^{\mathrm{E}} \mathrm{dlz}=0
$$

(b)

$$
\int_{c} \cos (z) d z=0 \text {, and }
$$

(c)

$$
\int_{\mathrm{c}} z^{\mathrm{n}} d \mathrm{z}=0 .
$$

## Cauchy integral formula

STATEMENT : let $\mathrm{F}(\mathrm{z})=\mathrm{u}(\mathrm{x}, \mathrm{y})+\mathrm{iv}(\mathrm{x}, \mathrm{y})$ be analytic on and within a simple closed contour (or curve ) ' c ' and let f ' $(\mathrm{z})$ be continuous there, then $\int f(z) d z=0$

Proof: $\mathrm{f}(\mathrm{z})=\mathrm{u}(\mathrm{x}, \mathrm{y})+\mathrm{iv}(\mathrm{x}, \mathrm{y})$
And dz=dx+idy
$\Rightarrow \mathrm{f}(\mathrm{z}) \cdot \mathrm{dz}=(\mathrm{u}(\mathrm{x}, \mathrm{y})+\mathrm{iv}(\mathrm{x}, \mathrm{y})) \mathrm{dx}+\mathrm{idy}$
$f(z) \cdot d z=u(x, y) d x+i u(x, y) d y+i v(x, y) d x+i^{2} v(x, y) d y$
$f(z) \cdot d z=u(x, y) d x-v(x, y) d y+i(u(x, y) d y+v(x, y) d x$
Integrate both sides, we get

$$
\int f(z) d z=\int(\mathrm{udx}-\mathrm{vdy})+i(\mathrm{udy}+\mathrm{vdx})
$$

By greens theorem, we have

$$
\int M d x+N d y=\iint \frac{\partial \mathrm{N}}{\partial \mathrm{x}}-\frac{\partial \mathrm{M}}{\partial \mathrm{Y}} \mathrm{dxdy}
$$

Now $\int f(z) d z=\iint\left(-\frac{\partial \mathrm{v}}{\partial \mathrm{x}}-\frac{\partial \mathrm{u}}{\partial \mathrm{Y}}\right) \mathrm{dxdy}+i\left(\frac{\partial \mathrm{u}}{\partial \mathrm{x}}-\frac{\partial \mathrm{v}}{\partial \mathrm{Y}}\right) \mathrm{dxdy}$
Since $f^{\prime}(z)$ is continuous \&four partial derivatives i.e $\frac{\partial u}{\partial x}, \frac{\partial u}{\partial Y}, \frac{\partial v}{\partial x}, \frac{\partial v}{\partial Y}$ are also continuous in the region R enclosed by C, Hence we can apply Green's Theorem.

Using Green's Theorem in plane , assuming that R is the region bounded by C .
It is given that $\mathrm{f}(\mathrm{z})=\mathrm{u}(\mathrm{x}, \mathrm{y})+\mathrm{iv}(\mathrm{x}, \mathrm{y})$ is analytic on and within c .
Hence $\frac{\partial u}{\partial x}=\frac{\partial v}{\partial y}, \frac{\partial u}{\partial y}=-\frac{\partial v}{\partial x}$
Using this we have

$$
\int_{c} f(z) d z=\iint_{R} 0 \quad d x d y+i \iint_{R} 0 \quad d x d y=0
$$

Hence the theorem.

## Cauchy's integral formula:

Cauchy's integral formula states that

$$
\begin{equation*}
f\left(z_{0}\right)=\frac{1}{2 \pi i} \oint_{\gamma} \frac{f(z) d z}{z-z_{0}}, \tag{1}
\end{equation*}
$$


It can be derived by considering the contour integral

$$
\begin{equation*}
\oint_{y} \frac{f(z) d z}{z-z_{0}}, \tag{2}
\end{equation*}
$$

defining a path $\gamma_{\text {r }}$ as an infinitesimal counterclockwise circle around the point $z_{0}$, and defining the path $\gamma_{0}$ as an arbitrary loop with a cut line (on which the forward and reverse contributions cancel each other out) so as to go around $z_{0}$. The total path is then

$$
\begin{equation*}
\gamma=\gamma_{0}+\gamma_{r}, \tag{3}
\end{equation*}
$$

so

$$
\begin{equation*}
\oint_{\gamma} \frac{f(z) d z}{z-z_{0}}=\oint_{\gamma_{0}} \frac{f(z) d z}{z-z_{0}}+\oint_{\gamma_{r}} \frac{f(z) d z}{z-z_{0}} . \tag{4}
\end{equation*}
$$

From the Cauchy integral theorem, the contour integral along any path not enclosing a pole is 0 . Therefore, the first term in the above equation is 0 since $\%$ does not enclose the pole, and we are left with

$$
\begin{equation*}
\oint_{\gamma} \frac{f(z) d z}{z-z_{0}}=\oint_{\gamma_{r}} \frac{f(z) d z}{z-z_{0}} . \tag{5}
\end{equation*}
$$

Now, let $z \equiv z_{0}+r e^{i \theta}$, so $d z=i r e^{i \theta} d \theta$. Then

$$
\begin{equation*}
\oint_{\gamma} \frac{f(z) d z}{z-z_{0}} \quad=\quad \oint_{\gamma_{r}} \frac{f\left(z_{0}+r e^{i \theta}\right)}{r e^{i \theta}} i r e^{i \theta} d \theta \tag{6}
\end{equation*}
$$

$$
\begin{equation*}
=\quad \oint_{r=1} f\left(z_{0}+r e^{i \theta}\right) i d \theta . \tag{7}
\end{equation*}
$$

But we are free to allow the radius $r$ to shrink to 0 , so

$$
\begin{align*}
\oint_{\gamma} \frac{f(z) d z}{z-z_{0}} & =\quad \lim _{r \rightarrow 0} \oint_{\gamma_{r}} f\left(z_{0}+r e^{i \theta}\right) i d \theta  \tag{8}\\
& =\quad \oint_{\gamma_{r}} f\left(z_{0}\right) i d \theta  \tag{9}\\
& =\quad i f\left(z_{0}\right) \oint_{y_{r}} d \theta  \tag{10}\\
& =\quad 2 \pi i f\left(z_{0}\right), \tag{11}
\end{align*}
$$

giving (1).
If multiple loops are made around the point $z_{0}$, then equation (11) becomes

$$
\begin{equation*}
n\left(\gamma, z_{0}\right) f\left(z_{0}\right)=\frac{1}{2 \pi i} \oint_{\gamma} \frac{f(z) d z}{z-z_{0}}, \tag{12}
\end{equation*}
$$

where $n\left(\gamma, z_{0}\right)$ is the contour winding number.
A similar formula holds for the derivatives of $f(z)$,

$$
\begin{align*}
f^{\prime}\left(z_{0}\right) \quad & =\lim _{h \rightarrow 0} \frac{f\left(z_{0}+h\right)-f\left(z_{0}\right)}{h}  \tag{13}\\
& =\lim _{h \rightarrow 0} \frac{1}{2 \pi i h}\left[\oint_{\gamma} \frac{f(z) d z}{z-z_{0}-h}-\oint_{\gamma} \frac{f(z) d z}{z-z_{0}}\right]  \tag{14}\\
& =\quad \lim _{h \rightarrow 0} \frac{1}{2 \pi i h} \oint_{\gamma} \frac{f(z)\left[\left(z-z_{0}\right)-\left(z-z_{0}-h\right)\right] d z}{\left(z-z_{0}-h\right)\left(z-z_{0}\right)}  \tag{15}\\
& =\quad \lim _{h \rightarrow 0} \frac{1}{2 \pi i h} \oint_{\gamma} \frac{h f(z) d z}{\left(z-z_{0}-h\right)\left(z-z_{0}\right)}  \tag{16}\\
& =\frac{1}{2 \pi i} \oint_{\gamma} \frac{f(z) d z}{\left(z-z_{0}\right)^{2}} . \tag{17}
\end{align*}
$$

Iterating again,

$$
\begin{equation*}
f^{\prime \prime}\left(z_{0}\right)=\frac{2}{2 \pi i} \oint_{y} \frac{f(z) d z}{\left(z-z_{0}\right)^{3}} . \tag{18}
\end{equation*}
$$

Continuing the process and adding the contour winding number $n$,

$$
n\left(\gamma, z_{0}\right) f^{(r)}\left(z_{0}\right)=\frac{r!}{2 \pi i} \oint_{\gamma} \frac{f(z) d z}{\left(z-z_{0}\right)^{\gamma+1}} .
$$

Problem: Evaluate using cauchy's integral formula $\int_{c} \frac{e^{2 z}}{(z-1)(z-2)} d z$ where c is the circle $|z|=3$
Solution: Given $\int_{c} \frac{e^{2 z}}{(z-1)(z-2)} d z$
Both the points $\mathrm{z}=1, \mathrm{z}=2$ line inside $|z|=3$
Resolving into partial fractions
$\frac{1}{(z-1)(z-2)}=\frac{A}{(z-1)}+\frac{B}{(z-2)}$
$\mathrm{A}=-1, \mathrm{~B}=1$
From(1)

$$
\begin{aligned}
& \int_{c} \frac{e^{2 z}}{(z-1)(z-2)} d z \\
&=\int_{c} \frac{-e^{2 z}}{(z-1)} d z+\int_{c} \frac{e^{2 z}}{(z-2)} d z \quad \text { (by cauchy's integral formula) } \\
&=-2 \pi \mathrm{if}(1)+2 \pi \mathrm{if}(2) \\
&=-2 \pi \mathrm{ie}^{2.1}+2 \pi \mathrm{ie}^{2.2} \\
&=-2 \pi \mathrm{ie}^{2}+2 \pi i \mathrm{e}^{4}=2 \pi i\left(e^{4}-\mathrm{e}^{2}\right)
\end{aligned}
$$

Problem: Using cauchy'sintegralformula to evaluate $\int_{c} \frac{\sin \pi z^{2}+\cos \pi z^{2}}{(z-1) z-2)} d z$, where c is the circle $|z|=3$

Solution: $\int_{c} \frac{f(z)}{(z-1) z-2)} d z=\left(\int_{c} \frac{1}{(z-2)} d z+\int_{c} \frac{1}{(z-1)} d z\right) \mathrm{f}(\mathrm{z}) \mathrm{dz}$

$$
\begin{aligned}
= & \int_{c} \frac{f(z)}{(z-2)} d z+\int_{c} \frac{f(z)}{(z-1)} d z \\
& =2 \Pi i f(2)-2 \Pi \mathrm{if}(1) \\
= & 2 \Pi \mathrm{i}(\sin 4 \Pi+\cos 4 \Pi)-(\sin \Pi+\cos \Pi)) \\
= & 2 \Pi \mathrm{i}(1-(-1))=4 \Pi \mathrm{i}
\end{aligned}
$$

$$
\int_{c} \frac{\sin \pi z^{2}+\cos \pi z^{2}}{(z-1) z-2)} d z=4 \Pi i
$$

Problem: Evaluate $\int_{c} \frac{(z-1)}{(z+1)^{2}(z-2)} d z$ whrere $\quad \mathrm{c}$ is $|Z-i|=2$
Solution: the singularities of $\frac{(z-1)}{(z+1)^{2}(z-2)}$ are given by

$$
\begin{aligned}
& (\mathrm{z}+1)^{2}(\mathrm{z}-2)=0 \\
& \quad \Rightarrow \mathrm{Z}=-1 \text { and } \mathrm{z}=2
\end{aligned}
$$

$\mathrm{Z}=-1$ lies inside the circle since $|-1-i|-2<0$
$\mathrm{Z}=2$ lies outside the circle sinceI2-iI- $2>0|2-i|-2>0$
The given line integral can be written as

$$
\begin{equation*}
\int_{c} \frac{(z-1)}{(z+1)^{2}(z-2)} d z=\int_{c} \frac{\frac{(z-1)}{(z-2)}}{(z+1)^{2}} \tag{1}
\end{equation*}
$$

The derivative of analytic function is given by

$$
\begin{equation*}
\int_{c} \frac{f(z)}{(z-a)^{n+1}} d z=\frac{2 \pi i f^{n}(a)}{n!} \tag{2}
\end{equation*}
$$

From (1) and (2) $f(z)=\frac{(z-1)}{(z-2)}, a=-1, n=1$

$$
\begin{aligned}
& f^{1}(z)=\frac{1(z-2)-1(z-1)}{(z-2)^{2}}=\frac{1}{(z-2)^{2}} \\
& f^{1}(-1)=\frac{1}{-9}
\end{aligned}
$$

Substituting in (2), we get

$$
\begin{gathered}
\int_{c} \frac{(z-1)}{(z+1)^{2}(z-2)} d z=\frac{2 \pi i}{1}\left(-\frac{1}{9}\right) \\
=\frac{-2}{9} \Pi \mathrm{i}
\end{gathered}
$$

Problem: Evaluate $\int_{c} \frac{e^{2 z}}{(z+1)^{4}} d z$ where $\mathrm{c}:|z-1|=1$
Solution: the singular points of $\frac{e^{2 z}}{(z+1)^{4}} d z$ are givenby 7

$$
(\mathrm{z}+1)^{4}=0 \Rightarrow z=-1
$$

The singular point $\mathrm{z}=-1$ lies insidethecirclec: $|z-1|=3$
Applying cauchy's integral formula for derivatives

$$
\begin{aligned}
& \quad \int_{c} \frac{f(z)}{(z-a)^{n+1}} d z=\int_{c} \frac{2 \pi i f^{n}(-1)}{n!} d z \\
& \mathrm{f}(\mathrm{z})=\mathrm{e}^{2 z}, \mathrm{n}=3, \mathrm{a}=-1 \\
& \mathrm{f}(\mathrm{z})=2 \mathrm{e}^{2 \mathrm{z}} \\
& \mathrm{f}^{1}(\mathrm{z})=4 \mathrm{e}^{2 \mathrm{z}} \\
& \mathrm{f}^{11}(\mathrm{z})=8 \mathrm{e}^{2 \mathrm{z}} \\
& \left.\left.\mathrm{f}^{111}\right) \mathrm{z}\right)=16 \mathrm{e}^{2 \mathrm{z}}
\end{aligned}
$$

$$
f^{111}(-1)=16 e^{-2}
$$

substituting in(1)

$$
\begin{aligned}
\int_{c} \frac{e^{2 z}}{(z+1)^{4}} d z & =\int_{c} \frac{2 \pi i f^{111}(-1)}{n!} \\
& =\frac{2 \pi i 16 e^{-2}}{2!} \\
& =16 \Pi \mathrm{ie}^{-2}
\end{aligned}
$$

Problem: Use cauchy's integral formula to evaluate $\int_{c} \frac{e^{-2 z}}{(z+1)^{3}} d z \quad$ with $\mathrm{c}:|z|=2$
Solution: $\int_{\text {Given }} \frac{e^{-2 z}}{(z+1)^{3}} d z$

$$
\mathrm{f}(\mathrm{z})=\mathrm{e}^{-2 z}
$$

the singular poin $\mathrm{z}=-1$ lies inside the given circle $|z|=2$
apply Cauchy's integral formula for derivatives

$$
\int_{c} \frac{e^{-2 z}}{(z+1)^{3}} d z=\frac{2 \pi i f^{1}(-1)}{2!} \quad\left[\because \int_{c} \frac{f(z)}{(z-a)^{3}}=\frac{2 \pi i f^{1}(a)}{2!}\right]
$$

Where $f(z)=e^{-2 z}$

$$
\begin{gathered}
\mathrm{f}^{1}(\mathrm{z})=-2 \mathrm{e}^{-2 z} \\
\mathrm{f}^{11}(\mathrm{z})=4 \mathrm{e}^{-2 z} \\
\mathrm{f}^{11}(-1)=4 \mathrm{e}^{2} \\
\therefore \int_{c} \frac{e^{-2 z}}{(z+1)^{3}} d z=\frac{2 \pi i 4 e^{2}}{2}=4 \pi i e^{2}
\end{gathered}
$$

Problem: Evaluate $\int_{c} \frac{d z}{z^{8}(z+4)} d z \quad$ withc: $|z|=2$
Solution: The singularities of $\int_{c} \frac{d z}{z^{8}(z+4)} d z$ are given by
$Z^{8}(\mathrm{z}+4)=0 \Rightarrow \mathrm{z}=0, \mathrm{z}=-4$
The point $\mathrm{z}=0$ lie inside and the $\mathrm{z}=-4$ lies outside the circle $|z|=2$

By the derivative of analytic function.
Problem: Evaluate using integral formula $\oint_{c} \frac{e^{2 z} d z}{(z-1)(z-2)}$ where c is the circle $|z|=3$
Solution: Let $(\mathrm{z})=\mathrm{e}^{\mathrm{z}}$ which is analytic within the circle $\mathrm{c}:|z|=3$ and the two singular points $a=1, a=2$ lie inside $c$.
$\oint_{c} \frac{e^{2 z} d z}{(z-1)(z-2)}==\begin{aligned} & \oint e^{2 z}\left(\frac{1}{z-2}-\frac{1}{z-1}\right) d z \\ & \oint_{c} \frac{e^{2 z}}{z-2} d z-\oint_{c} \frac{e^{2 z}}{z-1} d z\end{aligned}$
Now using cauchy's integral formula, we obtain
$\oint_{c} \frac{e^{2 z} d z}{(z-1)(z-2)}=2 \pi i e^{4}-2 \pi i e^{2}$

$$
=2 \pi i\left(e^{4}-e^{2}\right)
$$

$\oint_{c} \frac{e^{2 z} d z}{(z-1)(z-2)}=2 \pi i\left(e^{4}-e^{2}\right)$
Problem : Evaluate $\oint_{c} \frac{3 z^{2}+z}{z^{2}-1} d z$ where cisthe circle $|z-1|=1$
Solution: Given $f(z)=3 z^{2}+z$
$\mathrm{Z}=\mathrm{a}=+1$ or -1

The circle $|z-1|=1$ has centre at $\mathrm{z}=1$ and radius 1 and includes the point $\mathrm{z}=1, \mathrm{f}(\mathrm{z})=3 \mathrm{z}^{2}+\mathrm{z}$ is an analytic function

$$
\begin{align*}
& \text { Also } \frac{1}{z^{2}-1}=\frac{1}{(z-1)(z+1)}=\frac{1}{2}\left(\frac{1}{z-1}-\frac{1}{z+1}\right) \\
& \oint \frac{3 z^{2}+z}{z^{2}-1}=\frac{1}{2}\left[\int_{c} \frac{3 z^{2}+z}{z-1} d z\right]-\frac{1}{2}\left[\int_{c} \frac{3 z^{2}+z}{z+1} d z\right] \tag{1}
\end{align*}
$$

Since $\mathrm{z}=1$ lies inside c , we have by cauchy's integralformula

$$
\begin{aligned}
\oint_{c} \frac{3 z^{2}+z}{z^{2}-1} d z & =2 \pi i f(i) \\
= & 2 \pi i * 4
\end{aligned}
$$

Bycauchy'sintegral theorem ,since $z=-1$ lies out side c , we have

$$
\oint_{c} \frac{3 z^{2}+z}{z-1} d z=0
$$

From equation(1) we have

$$
\begin{aligned}
\oint_{c} \frac{3 z^{2}+z}{z^{2}-1} d z= & \frac{1}{2}(8 \pi i)-0 \\
& =4 \pi i
\end{aligned}
$$

## EXCERCISE PROBLEMS:

1) Evaluate $\int \frac{d z}{z-z_{0}}$ where $\mathrm{c}:\left|z-z_{0}\right|=\mathrm{r}$
2) Evaluate $\int_{(1,1)}^{(2,2)}(x+y) d x+(y-x) d y$ along the parabola $y^{2}=x$
3)Evaluate $\int_{c} \frac{z^{2}+4}{z^{2}-1} d z$ where $\mathrm{C}:|z|=2$ using Cauchy's Integral formula
4)Evaluate $\int_{c} \frac{e^{2 z}}{(z-1)(z-2)} d z$ where $\mathrm{C}:|z|=4$ using Cauchy's integral formula
3) Evaluate $\int_{c} \frac{z^{3}-z}{(z-2)^{3}}$ where $C:|z|=3$ using Cauchy's integral formula
4) $\operatorname{Expand} \mathrm{f}(\mathrm{z})=\int_{c} \frac{e^{2 z}}{(z-1)^{3}}$ at a point $\mathrm{z}=1$
5) Expand $\mathrm{f}(\mathrm{z})=\int_{c} \frac{1}{z^{2}-4 z+3}$ for $1<|z|<3$
8)Evaluate $\int\left(y^{2}+z^{2}\right) d x+\left(z^{2}+x^{2}\right) d y+\left(x^{2}+y^{2}\right) d z$ from $(0,0,0)$ to $(1,1,1)$, where C is the curve $x=t, x=t^{2}, x=t^{3}$
6) Evaluate $\int_{(0,0)}^{(1,1)}\left(3 x^{2}+4 x y+i x^{2}\right) d z$ along $y=x^{2}$
7) Evaluate $\int_{0}^{1+i}\left(x-y+i x^{2}\right) d z$
(i) along the straight from $\mathrm{z}=0$ to $\mathrm{z}=1+\mathrm{i}$.
(ii) along the real axis from $\mathrm{z}=0$ to $\mathrm{z}=1$ and then along a line parallel to real axis from $\mathrm{z}=$ 1 to $\mathrm{z}=1+\mathrm{i}$
(iii) along the imaginary axis from $\mathrm{z}=0$ to $\mathrm{z}=\mathrm{I}$ and then along a line parallel to real axis $\mathrm{z}=$ i to $\mathrm{z}=1+\mathrm{i}$.
8) Evaluate $\int_{1-i}^{2+i}(2 x+1+i y) \mathrm{dz}$ along (1-i) to $(2+\mathrm{i})$
9) Evaluate $\int_{c}\left(y^{2}+2 x y\right) d x+\left(x^{2}-2 x y\right) d y$ where c is boundary of the region $\mathrm{y}=\mathrm{x}^{2}$ and $\mathrm{x}=\mathrm{y}^{2}$

## The Radus of Convergence of a Power Series

Recall fom the Power Series page that we saw that a power series will converge a ti's center of convergencece, and that it is possible thata p pover series can converge for all $x \in \mathbb{R}$ or on some inteval centereda a the center of convergence. If power series converges on some inteval centered at the center of convergence, then the distance fom the center of convergence to eithe endpoint of that interal is shown as the radus of convergence which we more precisely define below.

Definition: The Radius of Convergence, Ris a non-negative number or o s such that the inteval of convergence for the power series $\sum_{n=0}^{\infty} a_{n}(x-c)^{n}$ is $\left(c-R_{1} c+R\right),\left(c-R_{1} c+R\right),\left(c-R_{1} c+R\right),\left(c-R_{1} c+R\right)$.

For xxample, in the case thata power series $\sum_{n=0}^{\infty} a_{n}(x-c)^{n}$ is convergent only at $x=c$, then the radus of convergence for this pover series is $R=0$ since the interval of convergence is $(c-0, c+0)=[c, c)$. Similaty, if it pe power series is convergent for all $x \in \mathbb{R}$ then the radius of convergence of the pover series is $R=\infty$ since the inteval of convergence is $(-\infty, \infty)$.

## Determining the Radius of Convergence of a Power Seres

We will now look at a technique for determining the radius of convergence of a power series using The Ratio Test tor Positive Seies

$$
\begin{aligned}
& \text { Theorem } 1:: \left.\left|f l i m_{n-\infty}\right| \frac{a_{n+1}}{a_{n}} \right\rvert\,=L \text { where } L \text { is a positive real number or } L=0 \text { or } L=\infty \text {, then the pover series } \sum_{n=0}^{\infty} a_{n}(x-c)^{n} \text { has a addulus of } \\
& \text { convergence } R=\frac{1}{L} \text { wherei if } L=0 \text { then } R=\infty \text { and it } L=\infty \text { then } R=0 \text {. }
\end{aligned}
$$

## An Example.

It's time to exploit this for power series. Consider the series

$$
1-2 x+3 x^{2}-4 x^{3}+5 x^{4}-\ldots
$$

We want to find out for what values of $x$ the series converges. If we view this power series as a series of the form

$$
A_{0}+A_{1}+A_{2}+A_{3}+\ldots
$$

then $\left|A_{0}\right|=1,\left|A_{1}\right|=2|x|,\left|A_{2}\right|=3\left|x^{2}\right|$ and so forth. The general term will have the form

$$
\left|A_{n}\right|=(n+1)\left|x^{n}\right|
$$

(Plug in $n=0,1,2, \ldots$ to see that this formula works!) Consequently the ratios are given by

$$
\frac{\left|A_{n+1}\right|}{\left|A_{n}\right|}=\frac{((n+1)+1)\left|x^{n+1}\right|}{(n+1)\left|x^{n}\right|}=\frac{n+2}{n+1} \cdot|x|
$$

Since

$$
\lim _{n \rightarrow \infty} \frac{n+2}{n+1}=1
$$

we obtain

$$
\lim _{n \rightarrow \infty} \frac{\left|A_{n+1}\right|}{\left|A_{n}\right|}=\lim _{n \rightarrow \infty} \frac{((n+1)+1)\left|x^{n+1}\right|}{(n+1)\left|x^{n}\right|}=|x|
$$

What's next? Do you remember the question we are trying to answer? For what values of $x$ does the power series converge! The ratio test tells us now that the series will converge as long as $|x|<1$. It also tells us that the series will diverge for $|x|>1$. That gives us a pretty complete picture about what's going on:

The biggest interval (it is always an interval!) where a power series is converging is called interval of convergence of the power series. The interval of convergence is always centered at the center of the power series. It is customary to call half the length of the interval of convergence the radius of convergence of the power series. In our example, the center of the power series is 0 , the interval of convergence is the interval from -1 to 1 (note the vagueness about the end points of the interval), its length is 2 , so the radius of convergence equals 1 .

## Another example.

Consider the power series
$5+(x-1)+\frac{(x-1)^{2}}{2!}+\frac{(x-1)^{3}}{3!}+\frac{(x-1)^{4}}{4!}+\frac{(x-1)^{5}}{5!}+\ldots$

- Step 1. Find the general term of the power series. In our case

$$
A_{n}=\frac{(x-1)^{n}}{n!}
$$

will do the job for $n \geq 1$. Since we will be taking the limit as $n$ goes to infinity, the odd " 5 " at the beginning is of no consequence!

- Step 2. Compute the ratios $\frac{\left|A_{n+1}\right|}{\left|A_{n}\right|}$. Don't forget the absolute values!

$$
\begin{aligned}
\frac{\left|A_{n+1}\right|}{\left|A_{n}\right|} & =\frac{\left|x^{n+1}\right| /(n+1)!}{\left|x^{n}\right| / n!} \\
& =\frac{1 \cdot 2 \cdot 3 \cdots n}{1 \cdot 2 \cdot 3 \cdots n \cdot(n+1)} \cdot|x| \\
& =\frac{1}{n+1} \cdot|x|
\end{aligned}
$$

- Step 3. Compute the iminit of fhe ratios. Since $\lim _{n \rightarrow \infty} \frac{1}{n+1}=0$, in our case

$$
\lim _{n \rightarrow \infty} \frac{\left|A_{n+1}\right|}{\left|A_{n}\right|}=0 \cdot|x|=0
$$

for fll $x$.

- Step 4.Apply the ratio test. Since © $<1$ (in thisis exmmple the iminit dees not depend on the aluve ofx), the serese conterges for all $x$.

Thus the interal of conreggence is the interal $(-\infty, \infty)$. The radius of conregegence in this case is sad to be $\infty$.

## MODULE-III

## POWER SERIES EXPANSIONS OF COMPLEX FUNCTION

## Power series:

A series expansion is a representation of a particular function as a sum of powers in one of its variables, or by a sum of powers of another (usually elementary) function $f(z)$.

A power series in a variable zis an infinite sum of the form

$$
\sum a_{i} z^{i}
$$

A series of the form $\sum a_{n} z^{n}$ is called as power series.
That is $\sum a_{n} z^{n}=a_{1} z+a_{2} z^{2}+\ldots \ldots \ldots .+a_{n} z^{n}+\ldots .$.

## Taylor's series:

Taylor's theorem states that any function satisfying certain conditions may be represented by a Taylor series.

The Taylor series is an infinite series, whereas a Taylor polynomial is a polynomial of degree n and has a finite number of terms. The form of a Taylor polynomial of degree $n$ for a function
$f(z)$ at $x=a$ is

$$
\begin{aligned}
f(z) & =f(a)+f^{\prime}(a)(z-a)+f^{\prime \prime}(a) \frac{(z-a)^{2}}{2!}+f^{\prime \prime \prime}(a) \frac{(z-a)^{3}}{3!}+\ldots \ldots \ldots+f^{n}(a) \frac{(z-a)^{n}}{n!}+\ldots \ldots . . \\
|z-a| & <r
\end{aligned}
$$

## Maclaurin series:

A Maclaurin series is a Taylor series expansion of a function about $\mathrm{x}=0$,

$$
f(z)=f(0)+f^{\prime}(0)(z)+f^{\prime \prime}(0) \frac{(z)^{2}}{2!}+f^{\prime \prime \prime}(0) \frac{(z)^{3}}{3!}+\ldots \ldots \ldots+f^{n}(0) \frac{(z)^{n}}{n!}+\ldots \ldots .
$$

This series is called as maclurins series expansion of $f(z)$.
Some important result:

$$
\begin{array}{ll}
\ln (1+x)=x-\frac{1}{2} x^{2}+\frac{1}{3} x^{3}-\frac{1}{4} x^{4}+\ldots & \text { for }-1<x \leq 1 \\
e^{x}=1+x+\frac{1}{2} x^{2}+\frac{1}{6} x^{3}+\frac{1}{24} x^{4}+\ldots & \text { for }-\infty<x<\infty \\
\cos x=1-\frac{1}{2} x^{2}+\frac{1}{24} x^{4}-\frac{1}{720} x^{6}+\ldots & \text { for }-\infty<x<\infty \\
\frac{1}{1-x}=1+x+x^{2}+x^{3}+x^{4}+x^{5}+\ldots & \text { for }-1<x<1
\end{array} \quad \begin{array}{ll}
\sin x=\sum_{n=0}^{\infty} \frac{(-1)^{n}}{(2 n+1)!} x^{2 n+1} &
\end{array}
$$

## Problems

Problem: Determine the first four terms of the power series for $\sin 2 \mathrm{x}$ using Maclaurin's series.

## Solution:

Let

$$
\begin{array}{ll}
\mathrm{f}(\mathrm{x})=\sin 2 \mathrm{x} & \mathrm{f}(0)=\sin 0=0 \\
\mathrm{f}^{\prime}(\mathrm{x})=2 \cos 2 \mathrm{x} & \mathrm{f}^{\prime}(0)=2 \cos 0=2 \\
\mathrm{f}^{\prime \prime}(\mathrm{x})=-4 \sin 2 \mathrm{x} & \mathrm{f}^{\prime \prime}(0)=-4 \sin 0=0 \\
\mathrm{f}^{\prime \prime \prime}(\mathrm{x})=-8 \cos 2 \mathrm{x} & \mathrm{f}^{\prime \prime \prime}(0)=-8 \cos 0=-8 \\
\mathrm{f}^{\mathrm{iv}}(\mathrm{x})=16 \sin 2 \mathrm{x} & \mathrm{f}^{\text {iv }}(0)=16 \sin 0=0 \\
\mathrm{f}^{\mathrm{v}}(\mathrm{x})=32 \cos 2 \mathrm{x}(0) & \mathrm{f}^{\mathrm{v}}(0)=32 \cos 0=32 \\
\mathrm{f}^{\text {vi }}(\mathrm{x})=-64 \sin 2 \mathrm{x} & \mathrm{f}^{\text {vi }}(0)=-64 \sin 0=0 \\
\mathrm{f}^{\text {vii }}(\mathrm{x})=-128 \cos 2 \mathrm{x} & \mathrm{f}^{\text {vii }}(0)=-128 \cos 0=-128 \\
\mathrm{f}(\mathrm{x})=\sin 2 \mathrm{x}=0+2 \mathrm{x}+0 \mathrm{x}^{2}+(-8) \frac{x^{3}}{3!}+0 \cdot \mathrm{x}^{4}+32 \frac{x^{5}}{5!} \\
=2 \mathrm{x}-\frac{4 x^{3}}{3}+\frac{4 x^{5}}{15}
\end{array}
$$

Problem : Find the Taylor series about $\mathrm{z}=-1$ for $\mathrm{f}(\mathrm{x})=1 / \mathrm{z}$. Express your answer in sigma notation.

## Solution:

$$
\begin{array}{cl}
\text { Let } \mathrm{f}(\mathrm{z})=\mathrm{z}^{-1} & \mathrm{f}(-1)=-1 \\
\mathrm{f}^{\prime}=-\mathrm{z}^{-2} & \mathrm{f}^{( }(-1)=-1 \\
\mathrm{f}^{\prime \prime}=2 \mathrm{z}^{-3} & \mathrm{f} "(-1)=-2 \\
\mathrm{f}^{\prime \prime \prime}=-6 \mathrm{z}^{-4} & \mathrm{f}^{\prime \prime \prime}(-1)=-6 \\
\mathrm{f}^{\prime \prime \prime}=24 \mathrm{z}^{-5} & \mathrm{f}^{\prime \prime \prime}(-1)=-24 \\
\mathrm{f}(\mathrm{z})=-1-1(\mathrm{z}+1)-\frac{2}{2!}(z+1)^{2}-\frac{6}{3!}(z+1)^{3}-\frac{24}{4!}(z+1)^{4}-\ldots \ldots \ldots . . \\
=\sum_{n=0}^{\infty}-1(z+1)^{n} &
\end{array}
$$

Problem : Find the Maclaurin series for $f(z)=z^{2}{ }^{z}$ Express your answer in sigma notation. Solution:
Let $f(z)=\mathrm{ze}^{\mathrm{z}}$
$\mathrm{f}(0)=0$
$f^{\prime}=e^{z}+z e^{z}$
$\mathrm{f}^{\prime}(0)=1+0=1$
$f^{\prime \prime}=e^{z}+e^{z}+e^{z}$
$\mathrm{f}^{\prime \prime}(0)=1+1+0=2$

$$
\begin{array}{rlrl}
\mathrm{f} " ' & =\mathrm{e}^{\mathrm{z}}+\mathrm{e}^{\mathrm{z}}+\mathrm{e}^{\mathrm{z}}+\mathrm{z} \mathrm{e}^{\mathrm{z}} & \mathrm{f} & \text { '"(0) }=1+1+1+0=3 \\
\mathrm{f}^{\prime \prime \prime} & =\mathrm{e}^{\mathrm{z}}+\mathrm{e}^{\mathrm{z}}+\mathrm{e}^{\mathrm{z}}+\mathrm{e}^{\mathrm{z}}+\mathrm{z} \mathrm{e}^{\mathrm{z}} & \mathrm{f} \text { '"' }(0)=1+1+1+1+0=4 \\
f(z) & =0+1 z+\frac{2}{2!} z^{2}+\frac{3}{3!} z^{3}+\frac{4}{4!} z^{4}+\ldots \ldots \ldots \ldots . . \\
& =z+z^{2}+\frac{1}{2} z^{3}+\frac{1}{6} z^{4}+\ldots \ldots \ldots \ldots . . \\
& =\sum_{n=1}^{\infty} \frac{z^{n}}{(n-1)!} &
\end{array}
$$

Problem: Expand $\log \mathrm{z}$ by taylor's series about $\mathrm{z}=1$.

## Solution:

$$
\begin{aligned}
& \text { Let } \mathrm{f}(\mathrm{z})=\log \mathrm{z} \\
& \mathrm{z}-\mathrm{l}=\mathrm{w} \\
& \mathrm{z}=1+\mathrm{w} \\
& \log \mathrm{z}
\end{aligned}=\log (1+\mathrm{w}) .
$$

## Laurent series:

In mathematics, the Laurent series of a complex function $f(z)$ is a representation of that function as a power series which includes terms of negative degree. It may be used to express complex functions in cases where a Taylor series expansion cannot be applied.

The Laurent series for a complex function $f(z)$ about a point $c$ is given by:

$$
\begin{aligned}
& f(z)=\sum_{n=-\infty}^{\infty} a_{n}(z-a)^{n} \\
& f(z)=\sum_{n=0}^{\infty} a_{n}(z-a)^{n}+\sum_{n=1}^{\infty} b_{n} \frac{1}{(z-a)^{n}}
\end{aligned}
$$

where the $a_{n}$ and $a$ are constants.

## Laurent polynomials:

A Laurent polynomial is a Laurent series in which only finitely many coefficients are non-zero. Laurent polynomials differ from ordinary polynomials in that they may have terms of negative degree.

## Principal part:

The principal part of a Laurent series is the series of terms with negative degree, that is

$$
f(z)=\sum_{K=-\infty}^{-1} a_{K}(z-a)^{K}
$$

If the principal part of $f$ is a finite sum, then $f$ has a pole at $c$ of order equal to (negative) the degree of the highest term; on the other hand, if $f$ has an essential singularity at $c$, the principal part is an infinite sum (meaning it has infinitely many non-zero terms).

Two Laurent series with only finitely many negative terms can be multiplied: algebraically, the sums are all finite; geometrically, these have poles at $c$, and inner radius of convergence 0 , so they both converge on an overlapping annulus.
Thus when defining formal Laurent series, one requires Laurent series with only finitely many negative terms.
Similarly, the sum of two convergent Laurent series need not converge, though it is always defined formally, but the sum of two bounded below Laurent series (or any Laurent series on a punctured disk) has a non-empty annulus of convergence.

## Zero's of an analytic function:

A zero of an analytic function $f(z)$ is a value of $z$ such that $f(z)=0$. Particularly a point a is called a zero of an analytic function $f(z)$ if $f(a)=0$.
Eg: $f(z)=\frac{(z+1)^{2}}{\left(z^{2}+1\right)^{2}}$
Now, $(z+1)^{2}=0$
$Z=-1, z=-1$ are zero's of an analytic function.
Zero's of $\mathbf{m}^{\text {th }}$ order:
If an analytic function $\mathrm{f}(\mathrm{z})$ can be expressed in the form $f(z)=(z-a)^{m} \Phi(z)$ where $\Phi(z)$ is analytic function and $\Phi(a) \neq 0$ then $\mathrm{z}=\mathrm{a}$ is called zero of $\mathrm{m}^{\text {th }}$ order of the function $\mathrm{f}(\mathrm{z})$.

- A simple zero is a zero of order 1 .

Eg: 1. $f(z)=(z-1)^{3}$
$\Rightarrow(z-1)^{3}=0$
$\mathrm{z}=1$ is a zero of order 3 of the function $\mathrm{f}(\mathrm{z})$.
2. $f(z)=\frac{1}{1-z}$
i.e $z=\infty$ is a simple zero of $\mathrm{f}(\mathrm{z})$.
3. $f(z)=\sin z$

$$
\text { i.e } z=n \pi \quad \forall n=0,1,2,3, \ldots \ldots \text {. are simple zero's of } \mathrm{f}(\mathrm{z}) \text {. }
$$

## Problems

Problem: Find the first four terms of the Taylor's series expansion of the complex function
$f(z)=\frac{z+1}{(z-3)(z-4)}$ about $\mathrm{z}=2$. Find the region of convergence.

## Solution:

The singularities of the function $f(z)=\frac{z+1}{(z-3)(z-4)}$ are $\mathrm{z}=3$ and $\mathrm{z}=4$
Draw a circle with centre at $\mathrm{z}=2$ and radius 1 .Then the distance of singularities from the centre are 1 and 2 .
Hence within the circle $|z-2|=1$, the given function is analytic .Hence ,it can be extended in Taylor's series within the circle $|z-2|=1$.
Hence $|z-2|=1$ is the circle of convergence.
Now $f(z)=\frac{5}{z-4}-\frac{4}{z-3}$ (partial fraction), $\mathrm{f}(2)=3 / 2$

$$
f^{\prime}(z)=-\frac{5}{(z-4)^{2}}+\frac{4}{(z-3)^{2}}, \quad f^{\prime}(2)=\frac{11}{4}
$$

$$
f^{\prime \prime}(z)=-\frac{8}{(z-3)^{3}}+\frac{10}{(z-4)^{3}}, f^{\prime \prime}(2)=\frac{27}{4}
$$

$$
f^{\prime \prime \prime}(z)=\frac{24}{(z-3)^{4}}-\frac{30}{(z-4)^{4}}, f^{\prime \prime \prime}(2)=\frac{177}{8}
$$

Taylor's series expansion for $\mathrm{f}(\mathrm{z})$ at $\mathrm{z}=\mathrm{a}$ is

$$
\begin{aligned}
& f(z)=f(a)+f^{\prime}(a)(z-a)+f^{\prime \prime}(a) \frac{(z-a)^{2}}{2!}+f^{\prime \prime \prime}(a) \frac{(z-a)^{3}}{3!}+\ldots \ldots \ldots+f^{n}(a) \frac{(z-a)^{n}}{n!}+\ldots \ldots . \\
& \frac{z+1}{(z-3)(z-4)}=\frac{3}{2}+(z-2) \frac{11}{4}+\frac{(z-2)^{2}}{2!}\left(\frac{27}{4}\right)+\frac{(z-2)^{3}}{3!}\left(\frac{177}{8}\right) \\
& f(z)=\frac{3}{2}+(z-2) \frac{11}{4}+(z-2)^{2}\left(\frac{27}{8}\right)+(z-2)^{3}\left(\frac{59}{16}\right) .
\end{aligned}
$$

Problem: ObtainLaurent series for $f(z)=\frac{e^{2 z}}{(z-1)^{3}}$ about $\mathrm{z}=1$.

## Solution:

$$
\text { Given } f(z)=\frac{e^{2 z}}{(z-1)^{3}}
$$

Put $\mathrm{z}-1=\mathrm{w}$ so that $\mathrm{z}=\mathrm{w}+1$

$$
\begin{aligned}
& f(z)=\frac{e^{2(1+w)}}{w^{3}} \\
& f(z)=\frac{e^{2} e^{2 w}}{w^{3}} \\
& =\frac{e^{2}}{w^{3}}\left[1+2 w+\frac{(2 w)^{2}}{2!}+\frac{(2 w)^{3}}{3!}+\ldots \ldots . .\right] \text { if } w \neq 0 \\
& =e^{2} \sum_{n=0}^{\infty} \frac{2^{n}}{n!} w^{n-3} \\
& =e^{2} \sum_{n=0}^{\infty} \frac{2^{n}}{n!}(z-1)^{n-3} \quad \text {, if } z-1 \neq 0 \\
& =e^{2} \sum_{n=0}^{\infty} \frac{2^{n}}{n!}(z-1)^{n-3} \quad, \text { if }|z-1| \neq 0 \\
& f(z)=e^{2} \sum_{n=0}^{\infty} \frac{2^{n}}{n!}(z-1)^{n-3} \quad \text {, if }|z-1|>0
\end{aligned}
$$

Since points $|z-1| \leq 0$ will be singular points.

Singular point of an analytic function: A point at which an analytic function $f(z)$ is not analytic, i.e. at which $\mathrm{f}^{\prime}(\mathrm{z})$ fails to exist, is called a singular point or singularity of the function.

There are different types of singular points:
Isolated and non-isolated singular points: A singular point $\mathrm{z}_{0}$ is called an isolated singular point of an analytic function $f(z)$ if there exists a deleted $\varepsilon$-spherical neighborhood of $z_{0}$ that contains no singularity. If no such neighborhood can be found, $\mathrm{z}_{0}$ is called a non-isolated singular point. Thus an isolated singular point is a singular point that stands completely by itself, embedded in regular points. In fig 1a where $\mathrm{z}_{1}, \mathrm{z}_{2}$ and $\mathrm{z}_{3}$ are isolated singular points. Most singular points are isolated singular points. A non-isolated singular point is a singular point such that every deleted $\varepsilon$-spherical neighborhood of it contains singular points. See Fig. $1 b$ where $z_{0}$ is the limit point of a set of singular points. Isolated singular points include poles, removable singularities, essential singularities and branch points.


Fig. 1 a


Fig. 1b

## Types of isolated singular points:

Pole: An isolated singular point $\mathrm{z}_{0}$ such that $\mathrm{f}(\mathrm{z})$ can be represented by an expression that is of the form

$$
f(z)=\frac{\phi(z)}{\left(z-z_{0}\right)^{n}}
$$

Where n is a positive integer, $\phi(\mathrm{z})$ is analytic at $\mathrm{z}_{0}$, and $\phi\left(\mathrm{z}_{0}\right) \neq 0$. The integer n is called the order of the pole. If $\mathrm{n}=1, \mathrm{z}_{0}$ is called a simple pole.

Example: 1.The function

$$
f(z)=\frac{5 z+1}{(z-2)^{3}(z+3)(z-2)}
$$

has a pole of order 3 at $\mathrm{z}=2$ and simple poles at $\mathrm{z}=-3$ and $\mathrm{z}=2$.

1. A point z is a pole for f if f blows up at z (f goes to infinity as you approach z ). An example of a pole is $\mathrm{z}=0$ for $\mathrm{f}(\mathrm{z})=1 / \mathrm{z}$.

Simple pole : A pole of order 1 is called a simple pole whilst a pole of order 2 is called a double pole.

If the principal part of the Laurent series has an infinite number of terms then $\mathrm{z}=\mathrm{z}_{0}$ is called an isolated essential singularity of $f(z)$. The function $f(z)=i / z(z-i) \equiv 1 /(z-i)-(1 / z)$ has a simple pole at $\mathrm{z}=0$ and another simple pole at $\mathrm{z}=\mathrm{i}$.

The function $e^{\frac{1}{z-2}}$ has an isolated essential singularity at $\mathrm{z}=2$. Some complex functions have non-isolated singularities called branch points. An example of such a function is $\sqrt{z}$.

Removable singular point: An isolated singular point $\mathrm{z}_{0}$ such that f can be defined, or redefined, at $\mathrm{z}_{0}$ in such a way as to be analytic at $\mathrm{z}_{0}$. A singular point $\mathrm{z}_{0}$ is removable if $\lim _{z \rightarrow z_{0}} f(z)$ exist.

Example: 1.The singular point $z=0$ is a removable singularity of $f(z)=(\sin z) / z$ since $\lim _{z \rightarrow 0} \frac{\sin z}{z}=1$

A point z is a removable singularity for f if f is defined in a neighborhood of the point z , but not at z , but f can be defined at z so that f is a continuous function which includes z . Here is an example of this: if $f(z)=z$ is defined in the punctured disk, the disk minus 0 , then $f$ is not defined at $\mathrm{z}=0$, but it can certainly be extended continuously to 0 by defining $\mathrm{f}(0)=0$. This means at $\mathrm{z}=0$ is a removable singularity.

Essential singular point: A singular point that is not a pole or removable singularity is called an essential singular point.

Example: 1. $f(z)=e^{1 /(z-3)}$ has an essential singularity at $\mathrm{z}=3$.
2. A point z is an essential singularity if the limit as f approaches z takes on different values as you approach $z$ from different directions. An example of this is $\exp (1 / z)$ at $z=0$. As $z$ approaches 0 from the right, $\exp (1 / \mathrm{z})$ blows up and as z approaches 0 from the left, $\exp (1 / z)$ goes to 0 .

Singular points at infinity: The type of singularity of $f(z)$ at $z=\infty$ is the same as that of $f(1 / w)$ at $\mathrm{w}=0$. Consult the following example.

Example: The function $f(z)=z^{2}$ has a pole of order 2 at $z=\infty$, since $f(1 / w)$ has a pole of order 2 at $\mathrm{w}=0$.

Using the transformation $\mathrm{w}=1 / \mathrm{z}$ the point $\mathrm{z}=0$ (i.e. the origin) is mapped into $\mathrm{w}=\infty$, called the point at infinity in the $w$ plane. Similarly, we call $z=\infty$ the point at infinity in the $z$ plane. To consider the behavior of $f(z)$ at $z=\infty$, we let $z=1 / w$ and examine the behavior of $f(1 / w)$ at $w=$ 0 .

## Residues:

The constant $\mathrm{a}_{-1}$ in the Laurent series

$$
\begin{equation*}
f(z)=\sum_{n=-\infty}^{\infty} a_{n}\left(z-z_{0}\right)^{n} \tag{1}
\end{equation*}
$$

of about a point $\mathrm{z}_{0}$ is called the residue of $\mathrm{f}(\mathrm{z})$. If $f$ is analytic at $\mathrm{z}_{0}$, its residue is zero, but the converse is not always true (for example, $\frac{1}{z^{2}}$ has residue of 0 at $\mathrm{z}=0$ but is not analytic at $\mathrm{z}=0$. The residue of a function f at a point $\mathrm{z}_{0}$ may be denoted $\underset{Z \rightarrow Z_{0}}{\operatorname{Re}} f(z)$.

Residue: Let $f(z)$ have a nonremovable isolated singularity at the point $z_{0}$. Then $f(z)$ has the Laurent series representation for all z in some disk $\mathrm{D}_{\mathrm{R}}^{*}\left(\mathrm{z}_{0}\right)$ given by

$$
\begin{equation*}
\text { - } f(z)=\sum_{n=-\infty}^{\infty} a_{n}\left(z-z_{0}\right)^{n} \tag{1}
\end{equation*}
$$

The coefficient $a_{-1}$ of $\frac{1}{z-z_{0}}$ is called the residue of $f(z)$ at $z_{0}$ and we use the notation

$$
\operatorname{Res}\left[f, z_{0}\right]=a_{-1}
$$

Example : If $f(z)=e^{\frac{2}{z}}$, then the Laurent series of f about the point $\mathrm{z}_{0}=0$ has the form

$$
\begin{aligned}
& f(z)=1+2 \frac{1}{z}+\frac{2^{2}}{2!z^{2}}+\frac{2^{3}}{3!z^{3}}+\frac{2^{4}}{4!z^{4}}+\frac{2^{5}}{5!z^{5}}+\cdots \text {, and } \\
& \operatorname{Res}[f, 0]=a_{-1}=2
\end{aligned}
$$

The residue of a function $f$ around a point $z_{0}$ is also defined by
$\operatorname{Re} s f=\frac{1}{2 \pi i} \int_{c} f(z) d z$
where C is counterclockwise simple closed contour, small enough to avoid any other poles of $f$. In fact, any counterclockwise path with contour winding number 1 which does not contain any other poles gives the same result by the Cauchy integral formula. The above diagram shows a suitable contour for which to define the residue of function, where the poles are indicated as black dots.

It is more natural to consider the residue of a meromorphic one-form because it is independent of the choice of coordinate. On a Riemann surface, the residue is defined for a meromorphic oneform $\alpha$ at a point $p$ by writing $\alpha=f d z$ in a coordinate $z$ around $p$. Then
$\underset{p}{\operatorname{Res} \alpha}=\underset{z=p}{\operatorname{Res}} f$.

The sum of the residues of $\int f d z$ is zero on the Riemann sphere. More generally, the sum of the residues of a meromorphic one-form on a compact Riemann surface must be zero.

The residues of a function $f(z)$ may be found without explicitly expanding into a Laurent series as follows. If $f(z)$ has a pole of order $m$ at $z_{0}$, then $a_{n}=0$ for $n<-m$ and $a_{-m} \neq 0$. Therefore,

$$
\begin{align*}
f(z)=\sum_{n=-m}^{\infty} a_{n}\left(z-z_{0}\right)^{n} & =\sum_{n=0}^{\infty} a_{-m+n}\left(z-z_{0}\right)^{-m+n}  \tag{4}\\
\left(z-z_{0}\right)^{m} f(z) & =\sum_{n=0}^{\infty} a_{-m+n}\left(z-z_{0}\right)^{n}  \tag{5}\\
\frac{d}{d z}\left[\left(z-z_{0}\right)^{m} f(z)\right] & =\sum_{n=0}^{\infty} n a_{-m+n}\left(z-z_{0}\right)^{n-1}  \tag{6}\\
& =\sum_{n=1}^{\infty} n a_{-m+n}\left(z-z_{0}\right)^{n-1} \tag{7}
\end{align*}
$$

$$
\begin{align*}
& =\sum_{n=0}^{\infty}(n+1) a_{-m+n+1}\left(z-z_{0}\right)^{n}  \tag{8}\\
\frac{d^{2}}{d z^{2}}\left[\left(z-z_{0}\right)^{m} f(z)\right] & =\sum_{n=0}^{\infty} n(n+1) a_{-m+n+1}\left(z-z_{0}\right)^{n-1}  \tag{9}\\
& =\sum_{n=1}^{\infty} n(n+1) a_{-m+n+1}\left(z-z_{0}\right)^{n-1}  \tag{10}\\
& =\sum_{n=0}^{\infty}(n+1)(n+2) a_{-m+n+2}\left(z-z_{0}\right)^{n} . \tag{11}
\end{align*}
$$

Iterating,

$$
\begin{align*}
& \frac{d^{m-1}}{d z^{m-1}}\left[\left(z-z_{0}\right)^{m} f(z)\right]=\sum_{n=0}^{\infty}(n+1)(n+2) \cdots(n+m-1) a_{n-1}\left(z-z_{0}\right)^{n} \\
& =(m-1)!a_{-1}+\sum_{n=1}^{\infty}(n+1)(n+2) \cdots(n+m-1) a_{n-1}\left(z-z_{0}\right)^{n-1} . \tag{12}
\end{align*}
$$

So

$$
\begin{align*}
\lim _{z \rightarrow z_{0}} \frac{d^{m-1}}{d z^{m-1}}\left[\left(z-z_{0}\right)^{m} f(z)\right] & =\lim _{z \rightarrow z_{0}}(m-1)!a_{-1}+0  \tag{13}\\
& =(m-1)!a_{-1}, \tag{14}
\end{align*}
$$

and the residue is
$a_{-1}=\frac{1}{(m-1)!} \frac{d^{m-1}}{d z^{m-1}}\left[\left(z-z_{0}\right)^{m} f(z)\right]_{z=z_{0}}$.

The residues of a holomorphic function at its poles characterize a great deal of the structure of a function, appearing for example in the amazing residue theorem of contour integration.

If $f(z)$ has a removable singularity at $z_{0}$ then $a_{-1}$ $=0$ for $n=1,2, \ldots \ldots$. Therefore, $\operatorname{Res}\left[f, z_{0}\right]=0$.

## Residues at Poles:

(i) If $\mathrm{f}(\mathrm{z})$ has a simple pole at $\mathrm{z}_{0}$, then $\operatorname{Res}\left[f, z_{0}\right]=\lim _{Z \rightarrow z_{0}}\left(z-z_{0}\right) f(z)$
(ii) If $\mathrm{f}(\mathrm{z})$ has a pole of order 2 at $\mathrm{z}_{0}$, then $\operatorname{Re} s\left[f, z_{0}\right]=\lim _{z \rightarrow z_{0}} \frac{d}{d z}\left(z-z_{0}\right)^{2} f(z)$
(iii) If $\mathrm{f}(\mathrm{z})$ has a pole of order 3 at $\mathrm{z}_{0}$, then $\operatorname{Re} s\left[f, z_{0}\right]=\frac{1}{2!\lim _{z \rightarrow z_{0}} \frac{d^{2}}{d z^{2}}\left(\left(z-z_{0}\right)^{3} f(z)\right), ~(z)}$
(v) If $\mathrm{f}(\mathrm{z})$ has a pole of order k at $\mathrm{z}_{0} \mathrm{z}_{0}$, then $\operatorname{Re} s\left[f, z_{0}\right]=\frac{1}{(k-1)!} \lim _{z \rightarrow z_{0}} \frac{d^{k-1}}{d z^{k-1}}\left(\left(z-z_{0}\right)^{k} f(z)\right)$

## Cauch's Residue Theorem:

An analytic function $\mathrm{f}(\mathrm{z})$ whose Laurent series is given by $f(z)=\lim _{z \rightarrow Z_{0}}\left(z-z_{0}\right) f(z)$ can be integrated term by term using a closed contour C encircling $\mathrm{z}_{0}$,

$$
\begin{align*}
& \int_{c} f(z) d z=\sum_{n=-\infty}^{\infty} a_{n} \int_{c}\left(z-z_{0}\right)^{n} d z \\
& =\sum_{n=-\infty}^{-2} a_{n} \int_{c}\left(z-z_{0}\right)^{n} d z+a_{-1} \int_{c} \frac{d z}{\left(z-z_{0}\right)}+\sum_{n=0}^{\infty} a_{n} \int_{c}\left(z-z_{0}\right)^{n} d z \tag{2}
\end{align*}
$$

The Cauchy integral theorem requires that the first and last terms vanish, so we have

$$
\begin{equation*}
\int_{c} f(z) d z=a_{-1} \int_{c} \frac{d z}{z-z_{0}} \tag{3}
\end{equation*}
$$

where $\mathrm{a}_{-1}$ is the complex residue. Using the contour $\mathrm{z}=\mathrm{c}(\mathrm{t})=\mathrm{e}^{\mathrm{it}}+\mathrm{Z}_{0}$ gives

$$
\begin{equation*}
\int_{c} \frac{d z}{z-z_{0}}=\int_{0}^{2 \pi} \frac{i e^{i t}}{e^{i t}} d t=2 \pi i \tag{4}
\end{equation*}
$$

so we have

$$
\begin{equation*}
\int_{c} f(z) d z=a_{-1} 2 \pi i \tag{5}
\end{equation*}
$$

If the contour C encloses multiple poles, then the theorem gives the general result

$$
\begin{equation*}
\int_{c} f(z) d z=2 \pi i \sum_{a \in A} \operatorname{Re}_{z=a_{i}} s f(z) \tag{6}
\end{equation*}
$$

Where A is the set of poles contained inside the contour. This amazing theorem therefore says that the value of a contour integral for any contour in the complex plane depends only on the properties of a few very special points inside the contour.

## Residue at infinity:

The residue at infinity is given by:
$\operatorname{Re} s[f(z)]_{Z=\infty}=-\frac{1}{2 \pi i_{C}} \int_{C} f(z) d z$
Where f is an analytic function except at finite number of singular points and C is a closed countour so all singular points lie inside it.
Problem: Determine the poles of the function $f(z)=\frac{z+2}{(z+1)^{2}(z-2)}$ and the residue at each pole.

Solution: The poles of $f(z)$ are given by $(z+1)^{2}(z-2)=0$
Here $\mathrm{z}=2$ is a simple pole and $\mathrm{z}=-1$ is a pole of order 2 .
Residue at $\mathrm{z}=2$ is

$$
\lim _{z \rightarrow 2}(z-2) f(z)=\lim _{z \rightarrow 2}(z-2) \frac{z+2}{(z+1)^{2}(z-2)}=\frac{4}{9}
$$

Residue at $\mathrm{z}=-1$ is

$$
\begin{aligned}
\lim _{z \rightarrow-1} \frac{d}{d z}(z+1)^{2} f(z) & =\lim _{z \rightarrow-1} \frac{d}{d z}(z+1)^{2} \frac{z+2}{(z+1)^{2}(z-2)} \\
\lim _{z \rightarrow-1} \frac{d}{d z} \frac{(z+2)}{(z-2)} & =\lim _{z \rightarrow-1} \frac{-4}{(z-2)^{2}}=\frac{-4}{9}
\end{aligned}
$$

Problem: Find the residue of the function $f(z)=\frac{1-e^{2 z}}{z^{4}}$ at the poles.
Solution: Let $f(z)=\frac{1-e^{2 z}}{z^{4}}$

$$
\mathrm{z}=0 \text { is a pole of order } 4
$$

Residue of $f(z)$ at $z=0$ is

$$
\begin{aligned}
& =\frac{1}{3!} \lim _{z \rightarrow 0} \frac{d^{3}}{d z^{3}}(z-0)^{4} \frac{\left(1-e^{2 z}\right)}{z^{4}} \\
& =\frac{1}{3!} \lim _{z \rightarrow 0} \frac{d^{3}}{d z^{3}}\left(1-e^{2 z}\right) \\
& =\frac{1}{3!} \lim _{z \rightarrow 0} \frac{d^{2}}{d z^{2}}\left(-2 e^{2 z}\right) \\
& =\frac{1}{3!} \lim _{z \rightarrow 0} \frac{d}{d z}\left(-4 e^{2 z}\right) \\
& =\frac{1}{3!} \lim _{z \rightarrow 0}\left(-8 e^{2 z}\right) \\
& =\frac{-8}{3!}=\frac{-4}{3} .
\end{aligned}
$$

Problem: Find the residue of the function $f(z)=z^{3} \cos \left(\frac{1}{z}\right)$ at $z=\infty$.
Solution: Let $f(z)=z^{3} \cos \left(\frac{1}{z}\right)$

$$
\begin{aligned}
& g(t)=f\left(\frac{1}{t}\right)=\frac{1}{t^{3}} \cos t \\
& \quad=\frac{1}{t^{3}}\left[1-\frac{t^{2}}{2!}+\frac{t^{4}}{4!}-\ldots \ldots \ldots .\right] \\
& =\left[\frac{1}{t^{3}}-\frac{1}{2 t}+\frac{t}{24}-\ldots \ldots . . . .\right]
\end{aligned}
$$

There fore $\operatorname{Re} \operatorname{sf}(z)=-$ coefficient ot t in the eapansion of $\mathrm{g}(\mathrm{t})$ about $\mathrm{t}=0$

$$
=-1 / 24
$$

Problem: Evaluate $\int_{c} \frac{4-3 z}{z(z-1)(z-2)} d z$ where c is the circle $|z|=\frac{3}{2}$. Using Residue theorem.
Solution: Let $\int_{c} \frac{4-3 z}{z(z-1)(z-2)} d z$
The poles of $f(z)$ are $z(z-1)(z-2)=0$

$$
\mathrm{z}=0, \mathrm{z}=1, \mathrm{z}=2
$$

These poles are simple poles.
The poles $\mathrm{z}=0$ and $\mathrm{z}=1$ lie within the circle $\mathrm{c}:|z|=\frac{3}{2}$
Residue of $\mathrm{f}(\mathrm{z})$ at $\mathrm{z}=0$ is $R_{1}=\lim _{z \rightarrow 0}(z-0) f(z)=\lim _{z \rightarrow 0} z \frac{4-3 z}{z(z-1)(z-2)}=\frac{4}{2}=2$
Residue of $\mathrm{f}(\mathrm{z})$ at $\mathrm{z}=1$ is $R_{2}=\lim _{z \rightarrow 1}(z-1) f(z)=\lim _{z \rightarrow 1}(z-1) \frac{4-3 z}{z(z-1)(z-2)}=\frac{4-3}{1-2}=-1$
By Residue theorem, $\int_{c} \frac{4-3 z}{z(z-1)(z-2)} d z=2 \pi i\left(R_{1}+R_{2}\right)=2 \pi i(2-1)=2 \pi i$.

Problem: Evaluate $\int_{0}^{2 \pi} \frac{d \theta}{5+4 \sin \theta}$ Using Residue theorem.
Solution: Let $I=\int_{-\pi}^{\pi} \frac{d \theta}{5+4 \sin \theta}$
Put $z=e^{i \theta}, d \theta=\frac{d z}{i z}$ and $\sin \theta=\frac{1}{2 i}\left(z-\frac{1}{z}\right)$
$z=e^{i \theta}$ unit circle $\mathrm{c}:|z|=1$

$$
\begin{aligned}
& I=\int_{c} \frac{1}{5+4 \frac{1}{2 i}\left(z-\frac{1}{z}\right)} \frac{d z}{i z} \\
& \begin{aligned}
=\int_{c} \frac{d z}{z^{2}+5 i z+2 i^{2}} & =\int_{c} \frac{d z}{(2 z+i)(z+2 i)} \\
& =\int_{c} \frac{d z}{\left(z+\frac{i}{2}\right)(z+2 i)} \\
& =\frac{1}{2} \int_{c} f(z) d z
\end{aligned}
\end{aligned}
$$

Where $f(z)=\left(z+\frac{i}{2}\right)(z+2 i)$
The poles of $f(z)$ are $z=-i / 2$ and $z=-2 i$
The pole $\mathrm{z}=-\mathrm{i} / 2$ lies inside the unit circle .
Residue of $f(z)$ at $z=-i / 2$ is

$$
\begin{aligned}
& =\lim _{z \rightarrow-i / 2}\left(z+\frac{i}{2}\right) f(z) \\
& =\lim _{z \rightarrow-i / 2}\left(z+\frac{i}{2}\right) \frac{1}{\left(z+\frac{i}{2}\right)(z+2 i)} \\
& =\lim _{z \rightarrow-i / 2} \frac{1}{(z+2 i)} \\
& =\frac{1}{\frac{-i}{2}+2 i}=\frac{2}{3 i}
\end{aligned}
$$

By cauchy's residue theorem

$$
\begin{aligned}
I= & \frac{1}{2} \int_{c} f(z) d z=\frac{1}{2} 2 \pi i\left(\frac{2}{3 i}\right)=\frac{2 \pi}{3} \\
& \int_{0}^{2 \pi} \frac{d \theta}{5+4 \sin \theta}=\frac{2 \pi}{3} .
\end{aligned}
$$

Problem: Prove that $\int_{-\infty}^{\infty} \frac{x^{2} d x}{\left(x^{2}+a^{2}\right)\left(x^{2}+b^{2}\right)}=\frac{\pi}{a+b} \quad(a>0, b>0, a \neq b)$

Solution: To evaluate the given integral, consider $\int_{c} \frac{z^{2} d z}{\left(z^{2}+a^{2}\right)\left(z^{2}+b^{2}\right)}=\int_{c} f(z) d z$
Where c is the contour consisting of the semi circle $C_{R}$ of radius R together with the real part of the real axis from -R to R .

The poles of $f(z)=\frac{z^{2}}{\left(z^{2}+a^{2}\right)\left(z^{2}+b^{2}\right)}$ are $z= \pm a i \quad ; z= \pm b i$
But $\mathrm{z}=\mathrm{ia}$ and $\mathrm{z}=\mathrm{ib}$ are the only two poles lie in the upper half of the plane .
$[\operatorname{Re} s f(z)]_{z=i a}=\underset{z \rightarrow a i}{\operatorname{Lt}}(z-i a) f(z)$

$$
=\underset{z \rightarrow a i}{ } \frac{z^{2}}{(z+i a)\left(z^{2}+b^{2}\right)}=\frac{-a^{2}}{2 i a\left(-a^{2}+b^{2}\right)}=\frac{a}{2 i\left(a^{2}-b^{2}\right)}
$$

Also $[\operatorname{Re} s f(z)]_{z=i b}=\underset{z \rightarrow b i}{\operatorname{Lt}}(z-i b) f(z)$

$$
=\underset{z \rightarrow b i}{\operatorname{Lt}} \frac{z^{2}}{(z+i b)\left(z^{2}+a^{2}\right)}=\frac{-b^{2}}{2 i b\left(a^{2}-b^{2}\right)}=\frac{-b}{2 i\left(a^{2}-b^{2}\right)}
$$

By Cauchy's Residue theorem, we have $\int_{c} f(z) d z=2 \pi i$ (sum of the residues with in C)

$$
\int_{c} f(z) d z=2 \pi i\left[\frac{a}{2 i\left(a^{2}-b^{2}\right)}-\frac{b}{2 i\left(a^{2}-b^{2}\right)}\right]=\pi\left[\frac{a-b}{\left(a^{2}-b^{2}\right)}\right]=\frac{\pi}{a+b}
$$

We have

$$
\int_{-R}^{R} f(x) d x+\int_{C_{R}} f(z) d z=\frac{\pi}{a+b}
$$

But $\int_{C_{R}} f(z) d z \rightarrow 0 \quad$ as $\quad z=\operatorname{Re}^{i \theta}$ and $\quad R \rightarrow \infty$
Hence $\int_{-\infty}^{\infty} f(x) d x=\int_{-\infty}^{\infty} \frac{x^{2} d x}{\left(x^{2}+a^{2}\right)\left(x^{2}+b^{2}\right)}=\frac{\pi}{a+b}$

## EXCERCISE PROBLEMS:

1)Evaluate $\int_{0}^{2 \pi} \frac{d \theta}{a+b \cos \theta}$ where $\mathrm{C}:|z|=1$
2)Prove that $\int_{-\infty}^{\infty} \frac{d x}{a^{2}+x^{2}}=\frac{\pi}{a}$
3)Show that $\int_{-\infty}^{\infty} \frac{\mathrm{dx}}{(\mathrm{x}+1)^{3}}=\frac{3 \pi}{8}$
4)Prove that $\int_{-\infty}^{\infty} \frac{d x}{\left(x^{2}+a^{2}\right)\left(x^{2}+b^{2}\right)}=\frac{\pi}{a b(a+b)}$
5)Evaluate $\int\left(y^{2}+z^{2}\right) d x+\left(z^{2}+x^{2}\right) d y+\left(x^{2}+y^{2}\right) d z$ from $(0,0,0)$ to $(1,1,1)$, where C is the curve $x=t, x=t^{2}, x=t^{3}$
6) Evaluate $\int_{(0,0)}^{(1,1)}\left(3 x^{2}+4 x y+i x^{2}\right) d z$ along $y=x^{2}$
7)Obtain the Taylor series expansion of $\mathrm{f}(\mathrm{z})=\frac{1}{z}$ about the point $\mathrm{z}=1$
8)Obtain the Taylor series expansion of $\mathrm{f}(\mathrm{z})=e^{z}$ about the point $\mathrm{z}=1$
9)Expand $\mathrm{f}(\mathrm{z})=\frac{z-1}{z+1}$ in Taylor's series about the point (i) $\mathrm{z}=0$ (ii) $\mathrm{z}=1$
10) Expand $f(z)=\frac{1}{z^{2}}$ in Taylor's series in powers of $z+1$
11)Obtain Laurent's series expansion of $\mathrm{f}(\mathrm{z})=\frac{z^{2}-4}{z^{2}+5 z+4}$ valid in $1<\mathrm{z}<2$
12)Give two Laurent's series expansions in powers of $Z$ for $f(z)=\frac{1}{z^{2}(1-z)}$
13) Expand $f(z)=\frac{1}{(1-z)(z-2)}$
14)Maclaurin's series expansion of $f(z)$

$$
1<|z|<2
$$

15)Laurent's series expansion in the annulus region in
16)Find the residue of the function $f(z)=\frac{z^{3}}{\left(z^{2}-1\right)}$ at $z=\infty$
17) Find the residue of $\frac{z^{2}}{z^{4}+1}$ at these singular points which lie inside the circle $|z|=2$
18) Find the residue of the function $f(z)=\frac{z^{2}-2 z}{\left(z^{2}+1\right)(z+1)^{2}}$ at each pole

## MODULE-IV

## SINGLE RANDOM VARIABLES

## Probability:

Probability is a branch of mathematics that deals with calculating the likelihood of a given event's occurrence, which is expressed as a number between 1 and 0 . An event with a probability of 1 can be considered a certainty: for example, the probability of a coin toss resulting in either "heads" or "tails" is 1 , because there are no other options, assuming the coin lands flat. An event with a probability of .5 can be considered to have equal odds of occurring or not occurring: for example, the probability of a coin toss resulting in "heads" is .5 , because the toss is equally as likely to result in "tails." An event with a probability of 0 can be considered an impossibility: for example, the probability that the coin will land (flat) without either side facing up is 0 , because either "heads" or "tails" must be facing up. A little paradoxical, probability theory applies precise calculations to quantify uncertain measures of random events.

In its simplest form, probability can be expressed mathematically as: the number of occurrences of a targeted event divided by the number of occurrences plus the number of failures of occurrences (this adds up to the total of possible outcomes):

$$
\mathrm{p}(\mathrm{a})=\mathrm{p}(\mathrm{a}) /[\mathrm{p}(\mathrm{a})+\mathrm{p}(\mathrm{~b})]
$$

Calculating probabilities in a situation like a coin toss is straightforward, because the outcomes are mutually exclusive: either one event or the other must occur. Each coin toss is an independent event; the outcome of one trial has no effect on subsequent ones. No matter how many consecutive times one side lands facing up, the probability that it will do so at the next toss is always $.5(50-50)$. The mistaken idea that a number of consecutive results (six "heads" for example) makes it more likely that the next toss will result in a "tails" is known as the gambler's fallacy, one that has led to the downfall of many a bettor.
Probability theory had its start in the 17th century, when two French mathematicians, Blaise Pascal and Pierre de Fermat carried on a correspondence discussing mathematical problems
dealing with games of chance. Contemporary applications of probability theory run the gamut of human inquiry, and include aspects of computer programming, astrophysics, music, weather prediction, and medicine.

Trial and Event: Consider an experiment, which though repeated under essential and identical conditions, does not give a unique result but may result in any one of the several possible outcomes. The experiment is known as Trial and the outcome is called Event
E.g. (1) Throwing a dice experiment getting the no's $1,2,3,4,5,6$ (event)
(2) Tossing a coin experiment and getting head or tail (event)

## Exhaustive Events:

The total no. of possible outcomes in any trial is called exhaustive event.
E.g.: (1) In tossing of a coin experiment there are two exhaustive events.
(2) In throwing an n-dice experiment, there are $6^{n}$ exhaustive events.

## Favorable event:

The no of cases favorable to an event in a trial is the no of outcomes which entities the happening of the event.
E.g. (1) In tossing a coin, there is one and only one favorable case to get either head or tail.

Mutually exclusive Event: If two or more of them cannot happen simultaneously in the same trial then the event are called mutually exclusive event.
E.g. In throwing a dice experiment, the events $1,2,3,-----6$ are M.E. events

Equally likely Events: Outcomes of events are said to be equally likely if there is no reason for one to be preferred over other. E.g. tossing a coin. Chance of getting $1,2,3,4,5,6$ is equally likely.

## Independent Event:

Several events are said to be independent if the happening or the non-happening of the event is not affected by the concerning of the occurrence of any one of the remaining events.

An event that always happen is called Certain event, it is denoted by ' S '.
An event that never happens is called Impossible event, it is denoted by ' $\phi$ '.
Eg: In tossing a coin and throwing a die, getting head or tail is independent of getting no's 1 or 2 or 3 or 4 or 5 or 6 .

## Definition: probability (Mathematical Definition)

If a trial results in n-exhaustive mutually exclusive, and equally likely cases and $m$ of them are favorable to the happening of an event $E$ then the probability of an event $E$ is denoted by $P(E)$ and is defined as
$\mathrm{P}(\mathrm{E})=\frac{\text { no of favourable cases to event }}{\text { Total no of exaustive cases }}=\frac{m}{n}$

## Sample Space:

The set of all possible outcomes of a random experiment is called Sample Space .The elements of this set are called sample points. Sample Space is denoted by S.

Eg. (1) In throwing two dies experiment, Sample $S$ contains 36 Sample points.
$S=\{(1,1),(1,2)$,----------(1,6), --------(6,1),(6,2),--------(6,6) $\}$
Eg. (2) In tossing two coins experiment , $\mathrm{S}=\{\mathrm{HH}, \mathrm{HT}, \mathrm{TH}, \mathrm{TT}\}$
A sample space is called discrete if it contains only finitely or infinitely many points which can be arranged into a simple sequence $\mathrm{w}_{1}, \mathrm{w}_{2}, \ldots \ldots$. while a sample space containing non denumerable no. of points is called a continuous sample space.

## Statistical or Empirical Probability:

If a trial is repeated a no. of times under essential homogenous and identical conditions, then the limiting value of the ratio of the no. of times the event happens to the total no. of trials, as the number of trials become indefinitely large, is called the probability of happening of the event. ( It is assumed the limit is finite and unique)

Symbolically, if in ' $n$ ' trials and events E happens ' $m$ ' times, then the probability ' $p$ ' of the happening of E is given by $\mathrm{p}=\mathrm{P}(\mathrm{E})=\lim _{n \rightarrow \infty} \frac{m}{n}$.

An event $E$ is called elementary event if it consists only one element.
An event, which is not elementary, is called compound event.
Example 1: What is the probability of getting a 2 or a 5 when a die is rolled?
Solution:
Taking the individual probabilities of each number, getting a 2 is $1 / 6$ and so is getting a 5 .
Applying the formula of compound probability,

Probability of getting a 2 or a 5,
$\mathrm{P}(2$ or 5$)=\mathrm{P}(2)+\mathrm{P}(5)-\mathrm{P}(2$ and 5$)$

$$
\begin{array}{ll}
=\Rightarrow & 1 / 6+1 / 6-0 \\
=\Rightarrow & 2 / 6=1 / 3 .
\end{array}
$$

Example 2: Consider the example of finding the probability of selecting a black card or a 6 from a deck of 52 cards.

## Solution:

We need to find out $\mathrm{P}(\mathrm{B}$ or 6$)$
Probability of selecting a black card $=26 / 52$
Probability of selecting a $6=4 / 52$
Probability of selecting both a black card and a $6=2 / 52$
$\mathrm{P}(\mathrm{B}$ or 6$) \quad=\mathrm{P}(\mathrm{B})+\mathrm{P}(6)-\mathrm{P}(\mathrm{B}$ and 6$)$
$=26 / 52+4 / 52-2 / 52$
$=28 / 52$
$=7 / 13$.

## Conditional probability:

Conditional probability is calculating the probability of an event given that another event has already occured.

The formula for conditional probability $\mathrm{P}(\mathrm{A} \mid \mathrm{B})$, read as $\mathrm{P}(\mathrm{A}$ given B$)$ is

$$
\mathbf{P}(\mathbf{A} \mid \mathrm{B})=\mathbf{P}(\mathbf{A} \text { and } \mathrm{B}) / \mathbf{P}(\mathbf{B})
$$

Consider the following example:
Example: In a class, $40 \%$ of the students study math and science. $60 \%$ of the students study math. What is the probability of a student studying science given he/she is already studying math?

## Solution

$\mathrm{P}(\mathrm{M}$ and S$)=0.40$
$\mathrm{P}(\mathrm{M})=0.60$
$\mathrm{P}(\mathrm{S} \mid \mathrm{M})=\mathrm{P}(\mathrm{M}$ and S$) / \mathrm{P}(\mathrm{S})=0.40 / 0.60=2 / 3=0.67$
Complement of an event
A complement of an event A can be stated as that which does NOT contain the occurrence of A.
A complement of an event is denoted as $\mathrm{P}\left(\mathrm{A}^{\mathrm{c}}\right)$ or $\mathrm{P}\left(\mathrm{A}^{\prime}\right)$.

$$
\mathrm{P}\left(\mathrm{~A}^{\mathrm{c}}\right)=1-\mathrm{P}(\mathrm{~A})
$$

or it can be stated, $\mathrm{P}(\mathrm{A})+\mathrm{P}\left(\mathrm{A}^{\mathrm{c}}\right)=1$
For example,
if $A$ is the event of getting a head in coin toss, $A^{c}$ is not getting a head i.e., getting a tail.
if $A$ is the event of getting an even number in a die roll, $A^{c}$ is the event of NOT getting an even number i.e., getting an odd number.
if A is the event of randomly choosing a number in the range of -3 to $3, \mathrm{~A}^{\mathrm{c}}$ is the event of choosing every number that is NOT negative i.e., $0,1,2 \& 3$ ( 0 is neither positive or negative).

Consider the following example:
Example: A single coin is tossed 5 times. What is the probability of getting at least one head?

## Solution:

Consider solving this using complement.
Probability of getting no head $=\mathrm{P}($ all tails $)=1 / 32$
$\mathrm{P}($ at least one head $)=1-\mathrm{P}($ all tails $)=1-1 / 32=31 / 32$.

Example 1: A dice is thrown 3 times .what is the probability that atleast one head is obtained?
Sol: Sample space $=[$ HHH, HHT, HTH, THH, TTH, THT, HTT, TTT $]$
Total number of ways $=2 \times 2 \times 2=8$. Fav. Cases $=7$
$P(A)=7 / 8$
OR
$\mathrm{P}($ of getting at least one head $)=1-\mathrm{P}($ no head $) \Rightarrow 1-(1 / 8)=7 / 8$
Example 2: Find the probability of getting a numbered card when a card is drawn from the pack of 52 cards.

Sol: Total Cards $=52$. Numbered Cards $=(2,3,4,5,6,7,8,9,10) 9$ from each suit $4 \times 9=36$
$\mathrm{P}(\mathrm{E})=36 / 52=9 / 13$

Example 3: There are 5 green 7 red balls. Two balls are selected one by one without replacement. Find the probability that first is green and second is red.
Sol: $\mathrm{P}(\mathrm{G}) \times \mathrm{P}(\mathrm{R})=(5 / 12) \times(7 / 11)=35 / 132$
Example 4: What is the probability of getting a sum of 7 when two dice are thrown?
Sol: Probability math - Total number of ways $=6 \times 6=36$ ways. Favorable cases $=(1,6)(6,1)$ $(2,5)(5,2)(3,4)(4,3)---6$ ways. $P(A)=6 / 36=1 / 6$

Example 5: 1 card is drawn at random from the pack of 52 cards.
(i) Find the Probability that it is an honor card.
(ii) It is a face card.

Sol: (i) honor cards $=(A, J, Q, K) 4$ cards from each suits $=4 \times 4=16$
$P($ honor card $)=16 / 52=4 / 13$
(ii) face cards $=(\mathrm{J}, \mathrm{Q}, \mathrm{K}) 3$ cards from each suit $=3 \times 4=12$ Cards.
$\mathrm{P}($ face Card $)=12 / 52=3 / 13$
Example 6: Two cards are drawn from the pack of 52 cards. Find the probability that both are diamonds or both are kings.
Sol: Total no. of ways $={ }^{52} \mathrm{C}_{2}$
Case I: Both are diamonds $={ }^{13} \mathrm{C}_{2}$
Case II: Both are kings $={ }^{4} \mathrm{C}_{2}$
P (both are diamonds or both are kings) $=\left({ }^{13} \mathrm{C}_{2}+{ }^{4} \mathrm{C}_{2}\right) /{ }^{52} \mathrm{C}_{2}$
Example 7: Three dice are rolled together. What is the probability as getting at least one '4'?
Sol: Total number of ways $=6 \times 6 \times 6=216$. Probability of getting number ' 4 ' at least one time $=1-($ Probability of getting no number 4$)=1-(5 / 6) \times(5 / 6) \times(5 / 6)=91 / 216$

Example 8: A problem is given to three persons P, Q, R whose respective chances of solving it are $2 / 7,4 / 7,4 / 9$ respectively. What is the probability that the problem is solved?
Sol: Probability of the problem getting solved $=1-$ (Probability of none of them solving the problem)

$$
P(P)=\frac{2}{7} \Rightarrow P(\bar{P})=1-\frac{2}{7}=\frac{5}{7}, P(Q)=\frac{4}{7} \Rightarrow P(\bar{Q})=1-\frac{4}{7}=\frac{3}{7}, P(R)=\frac{4}{9} \Rightarrow P(\bar{R})=1-\frac{4}{9}=\frac{5}{9}
$$

Probability of problem getting solved $=1-(5 / 7) \times(3 / 7) \times(5 / 9)=(122 / 147)$
Example 9: Find the probability of getting two heads when five coins are tossed.
Sol: Number of ways of getting two heads $={ }^{5} \mathrm{C}_{2}=10$. Total Number of ways $=2^{5}=32$
$P($ two heads $)=10 / 32=5 / 16$
Example 10: What is the probability of getting a sum of 22 or more when four dice are thrown?
Sol: Total number of ways $=6^{4}=1296$. Number of ways of getting a sum 22 are $6,6,6,4=4$ ! / 3 !
$=4$
$6,6,5,5=4!/ 2!2!=6$. Number of ways of getting a sum 23 is $6,6,6,5=4!/ 3!=4$.
Number of ways of getting a sum 24 is $6,6,6,6=1$.
Fav. Number of cases $=4+6+4+1=15$ ways. $P($ getting a sum of 22 or more $)=15 / 1296=$ 5/432

Example 11: Two dice are thrown together. What is the probability that the number obtained on one of the dice is multiple of number obtained on the other dice?

Sol:Total number of cases $=6^{2}=36$
Since the number on a die should be multiple of the other, the possibilities are
$(1,1)(2,2)(3,3)-----(6,6)---6$ ways
$(2,1)(1,2)(1,4)(4,1)(1,3)(3,1)(1,5)(5,1)(6,1)(1,6)---10$ ways
$(2,4)(4,2)(2,6)(6,2)(3,6)(6,3)--6$ ways
Favorable cases are $=6+10+6=22$. So, $P(A)=22 / 36=11 / 18$
Example 12: From a pack of cards, three cards are drawn at random. Find the probability that each card is from different suit.
Sol: Total number of cases $={ }^{52} \mathrm{C}_{3}$
One card each should be selected from a different suit. The three suits can be chosen in ${ }^{4} \mathrm{C}_{3}$ was The cards can be selected in a total of $\left({ }^{4} \mathrm{C}_{3}\right) \times\left({ }^{13} \mathrm{C}_{1}\right) \times\left({ }^{13} \mathrm{C}_{1}\right) \times\left({ }^{13} \mathrm{C}_{1}\right)$
Probability $={ }^{4} \mathrm{C}_{3} \times\left({ }^{13} \mathrm{C}_{1}\right){ }^{3} /{ }^{52} \mathrm{C}_{3}$
$=4 \mathrm{x}(13)^{3} /{ }^{52} \mathrm{C}_{3}$
Example 13: Find the probability that a leap year has 52 Sundays.
Sol: A leap year can have 52 Sundays or 53 Sundays. In a leap year, there are 366 days out of which there are 52 complete weeks \& remaining 2 days. Now, these two days can be (Sat, Sun)
(Sun, Mon) (Mon, Tue) (Tue, Wed) (Wed, Thur) (Thur, Friday) (Friday, Sat).
So there are total 7 cases out of which (Sat, Sun) (Sun, Mon) are two favorable cases. So, P (53
Sundays) $=2 / 7$
Now, $\mathrm{P}(52$ Sundays $)+\mathrm{P}(53$ Sundays $)=1$
So, $P(52$ Sundays $)=1-P(53$ Sundays $)=1-(2 / 7)=(5 / 7)$
Example 14: Fifteen people sit around a circular table. What are odds against two particular people sitting together?
Sol: 15 persons can be seated in 14! Ways. No. of ways in which two particular people sit together is $13!\times 2$ !

The probability of two particular persons sitting together $13!2!/ 14!=1 / 7$
Odds against the event $=6: 1$
Example 15: Three bags contain 3 red, 7 black; 8 red, 2 black, and 4 red \& 6 black balls respectively. 1 of the bags is selected at random and a ball is drawn from it. If the ball drawn is red, find the probability that it is drawn from the third bag.
Sol: Let E1, E2, E3 and A are the events defined as follows.
$\mathrm{E} 1=$ First bag is chosen
$\mathrm{E} 2=$ Second bag is chosen
E3 $=$ Third bag is chosen
$\mathrm{A}=\mathrm{Ball}$ drawn is red
Since there are three bags and one of the bags is chosen at random, so $P(E 1)=P(E 2)=P(E 3)=$ $1 / 3$
If E1 has already occurred, then first bag has been chosen which contains 3 red and 7 black balls. The probability of drawing 1 red ball from it is $3 / 10$. So, $\mathrm{P}\left(\mathrm{A} / \mathrm{E}_{1}\right)=3 / 10$, similarly $\mathrm{P}\left(\mathrm{A} / \mathrm{E}_{2}\right)=$ $8 / 10$, and $P\left(A / E_{3}\right)=4 / 10$. We are required to find $P\left(E_{3} / A\right)$ i.e. given that the ball drawn is red, what is the probability that the ball is drawn from the third bag by Baye's rule

$$
=\frac{\frac{1}{3} \times \frac{4}{10}}{\frac{1}{3} \times \frac{3}{10}+\frac{1}{3} \times \frac{8}{10}+\frac{1}{3} \times \frac{4}{10}}=\frac{4}{15} .
$$

## Random Variables

- A random variable $X$ on a sample space $S$ is a function $X: S \rightarrow R$ from $S$ onto the set of real numbers $R$, which assigns a real number $X$ (s) to each sample point ' $s$ ' of $S$.
- Random variables (r.v.) bare denoted by the capital letters X,Y,Z,etc..
- Random variable is a single valued function.
- Sum, difference, product of two random variables is also a random variable .Finite linear combination of r.v is also a r.v .Scalar multiple of a random variable is also random variable.
- A random variable, which takes at most a countable number of values, it is called a discrete r.v. In other words, a real valued function defined on a discrete sample space is called discrete r.v.
- A random variable $X$ is said to be continuous if it can take all possible values between certain limits .In other words, a r.v is said to be continuous when it's different values cannot be put in 1-1 correspondence with a set of positive integers.
- A continuous r.v is a r.v that can be measured to any desired degree of accuracy. Ex : age , height, weight etc..
- Discrete Probability distribution: Each event in a sample has a certain probability of occurrence. A formula representing all these probabilities which a discrete r.v. assumes is known as the discrete probability distribution.
- The probability function or probability mass function (p.m.f) of a discrete random variable $X$ is the function $f(x)$ satisfying the following conditions.
i) $f(x) \geq 0$
ii) $\sum_{x} f(x)=1$
iii) $P(X=x)=f(x)$
- Cumulative distribution or simply distribution of a discrete r.v. X is $\mathrm{F}(\mathrm{x})$ defined by $\mathrm{F}(\mathrm{x})$ $=\mathrm{P}(\mathrm{X} \leq \mathrm{x})=\sum_{t \leq x} f(t)$ for $-\infty<x<\infty$
- If X takes on only a finite no. of values $\mathrm{x}_{1,} \mathrm{X}_{2, \ldots \ldots}, \mathrm{x}_{\mathrm{n}}$ then the distribution function is given by

$$
F(x)=
$$



$$
\begin{aligned}
& \mathrm{F}(-\infty)=0, \mathrm{~F}(\infty)=1,0 \leq \mathrm{F}(\mathrm{x}) \leq 1, \mathrm{~F}(\mathrm{x}) \leq \mathrm{F}(\mathrm{y}) \text { if } \mathrm{x}<\mathrm{y} \\
& \mathrm{P}\left(\mathrm{x}_{\mathrm{k}}\right)=\mathrm{P}\left(\mathrm{X}=\mathrm{x}_{\mathrm{k}}\right)=\mathrm{F}\left(\mathrm{x}_{\mathrm{k}}\right)-\mathrm{F}\left(\mathrm{x}_{\mathrm{k}-1}\right)
\end{aligned}
$$

- For a continuous r.v. X, the function $f(x)$ satisfying the following is known as the probability density function(p.d.f.) or simply density function:
i) $\quad \mathrm{f}(\mathrm{x}) \geq 0,-\infty<\mathrm{x}<\infty$
ii)

$$
\int_{-\infty}^{\infty} f(x) d x=1
$$

iii) $\quad \mathrm{P}(\mathrm{a}<\mathrm{X}<\mathrm{b})=\int_{a}^{b} f(x) d x=$ Area under $\mathrm{f}(\mathrm{x})$ between ordinates $\mathrm{x}=\mathrm{a}$ and $\mathrm{x}=\mathrm{b}$

$$
\mathrm{P}(\mathrm{a}<\mathrm{X}<\mathrm{b})=\mathrm{P}(\mathrm{a} \leq \mathrm{x}<\mathrm{b})=\mathrm{P}(\mathrm{a}<\mathrm{X} \leq \mathrm{b})=\mathrm{P}(\mathrm{a} \leq \mathrm{X} \leq \mathrm{b})
$$

(i.e) In case of continuous it does not matter weather we include the end points of the interval from a to $b$.This result in general is not true for discrete r.v.

- Probability at a point $\mathrm{P}(\mathrm{X}=\mathrm{a})=\int_{a-\Delta x}^{a+\Delta x} f(x) d x$
- Cumulative distribution for a continuous r.v. X with p.d.f. $\mathrm{f}(\mathrm{x})$, the cumulative distribution $\mathrm{F}(\mathrm{x})$ is defined as

$$
\mathrm{F}(\mathrm{x})=\mathrm{P}(\mathrm{X} \leq \mathrm{x})=\int_{-\infty}^{\infty} f(t) d t \quad-\infty<\mathrm{x}<\infty
$$

It follows that $\mathrm{F}(-\infty)=0, \mathrm{~F}(\infty)=1,0 \leq \mathrm{F}(\mathrm{x}) \leq 1$ for $-\infty<\mathrm{x}<\infty$

$$
f(x)=d / d x(F(x))=F^{1}(x) \geq 0 \text { and } P(a<x<b)=F(b)-F(a)
$$

- In case of discrete r.v. the probability at a point i.e., $\mathrm{P}(\mathrm{x}=\mathrm{c})$ is not zero for some fixed c however in case of continuous random variables the probability at appoint is always zero. I.e., $\mathrm{P}(\mathrm{x}=\mathrm{c})=0$ for all possible values of c .
- $\quad \mathrm{P}(\mathrm{E})=0$ does not imply that the event E is null or impossible event.
- If $X$ and $Y$ are two discrete random variables the joint probability function of $X$ and $Y$ is given by $P(X=x, Y=y)=f(x, y)$ and satisfies
(i) $\mathrm{f}(\mathrm{x}, \mathrm{y}) \geq 0$
(ii) $\sum_{x} \sum_{y} f(x, y)=1$

The joint probability function for X and Y can be reperesented by a joint probability table.

Table


The probability of $\mathrm{X}=\mathrm{x}_{\mathrm{j}}$ is obtained by adding all entries in arrow corresponding to $\mathrm{X}=\mathrm{x}_{\mathrm{j}}$
Similarly the probability of $\mathrm{Y}=\mathrm{y}_{\mathrm{k}}$ is obtained by all entries in the column corresponding to Y $=y_{k}$
$f_{1}(x)$ and $f_{2}(y)$ are called marginal probability functions of $X$ and $Y$ respectively.
The joint distribution function of X and Y is defined by $\mathrm{F}(\mathrm{x}, \mathrm{y})=\mathrm{P}(\mathrm{X} \leq \mathrm{x}, \mathrm{Y} \leq \mathrm{y})=\sum_{u \leq x v \leq y} \sum_{v} f(u, v)$

- If X and Y are two continuous r.v.'s the joint probability function for the r.v.'s X and Y is defined by
(i) $f(x, y) \geq 0$
(ii) $\int_{-\infty}^{\infty} \int_{-\infty}^{\infty} f(x, y) d x d y=1$
- $\mathrm{P}(\mathrm{a}<\mathrm{X}<\mathrm{b}, \mathrm{c}<\mathrm{Y}<\mathrm{d})=\int_{x=a}^{b} \int_{y=c}^{d} f(x, y) d x d y$
- The joint distribution function of $X$ and $Y$ is $F(x, y)=P(X \leq x, Y \leq y)=$ $\int_{u=-\infty}^{\infty} \int_{v=-\infty}^{\infty} f(u, v) d u d v$
- $\frac{\partial^{2} F}{\partial x \partial y}=f(x, y)$

The Marginal distribution function of X and Y are given by $\mathrm{P}(\mathrm{X} \leq \mathrm{x})=\mathrm{F}_{1}(\mathrm{x})=\mathrm{E}(\mathrm{X})=$ $\left\{\sum_{i} x_{i} f\left(x_{i}\right) \quad X \quad\right.$ is discrete
$\left\{\int_{-\infty}^{\infty} x f(x) d x \quad X \quad\right.$ is Continuous

- $\int_{u=-\infty}^{\infty} \int_{v=-\infty}^{\infty} f(u, v) d u d v$ and $\mathrm{P}(\mathrm{Y} \leq \mathrm{y})=\mathrm{F}_{2}(\mathrm{y})=\int_{u=-\infty}^{\infty} \int_{v=-\infty}^{\infty} f(u, v) d u d v$
- The marginal density function of $X$ and $Y$ are given by

$$
\mathrm{f}_{1}(\mathrm{x})=\int_{v=-\infty}^{\infty} f(x, v) d v \text { and } \mathrm{f}_{2}(\mathrm{y})=\int_{u=-\infty}^{\infty} f(u, y) d u
$$

- Two discrete random variables X and Y are independent iff

$$
\begin{aligned}
& \mathrm{P}(\mathrm{X}=\mathrm{x}, \mathrm{Y}=\mathrm{y})=\mathrm{P}(\mathrm{X}=\mathrm{x}) \mathrm{P}(\mathrm{Y}=\mathrm{y}) \forall \mathrm{x}, \mathrm{y} \quad \text { (or) } \\
& \mathrm{f}(\mathrm{x}, \mathrm{y})=\mathrm{f}_{1}(\mathrm{x}) \mathrm{f}_{2}(\mathrm{y}) \quad \forall \mathrm{x}, \mathrm{y}
\end{aligned}
$$

- Two continuous random variables X and Y are independent iff

$$
\begin{aligned}
& \mathrm{P}(\mathrm{X} \leq \mathrm{x}, \mathrm{Y} \leq \mathrm{y})=\mathrm{P}(\mathrm{X} \leq \mathrm{x}) \mathrm{P}(\mathrm{Y} \leq \mathrm{y}) \forall \mathrm{x}, \mathrm{y} \quad \text { (or }) \\
& \mathrm{f}(\mathrm{x}, \mathrm{y})=\mathrm{f}_{1}(\mathrm{x}) \mathrm{f}_{2}(\mathrm{y}) \quad \forall \mathrm{x}, \mathrm{y}
\end{aligned}
$$

If $X$ and $Y$ are two discrete r.v. with joint probability function $f(x, y)$ then

$$
\mathrm{P}(\mathrm{Y}=\mathrm{y} \mid \mathrm{X}=\mathrm{x})=\frac{f(x, y)}{f_{1}(x)}=\mathrm{f}(\mathrm{y} \mid \mathrm{x})
$$

Similarly, $\mathrm{P}(\mathrm{X}=\mathrm{x} \mid \mathrm{Y}=\mathrm{y})=\frac{f(x, y)}{f_{2}(y)}=\mathrm{f}(\mathrm{x} \mid \mathrm{y})$

If X and Y are continuous r.v. with joint density function $\mathrm{f}(\mathrm{x}, \mathrm{y})$ then $\frac{f(x, y)}{f_{1}(x)}=\mathrm{f}(\mathrm{y} \mid \mathrm{x})$ and

$$
\frac{f(x, y)}{f_{2}(y)} \quad=\mathrm{f}(\mathrm{x} \mid \mathrm{y})
$$

Expectation or mean or Expected value : The mathematical expectation or expected value of r.v. X is denoted by $\mathrm{E}(\mathrm{x})$ or $\mu$ and is defined as

- If X is a r.v. then
 for Discrete For Continuous
- If X, Y are r.v.'s with joint probability function $f(x, y)$ then

$$
\begin{aligned}
& \mathrm{E}[\mathrm{~g}(\mathrm{X}, \mathrm{Y})]=\sum_{x} \sum_{y} g(x, y) f(x, y) \\
& \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} g(x, y) f(x, y) d x d y
\end{aligned} \quad \text { for discrete r.v.'s } \quad \begin{aligned}
& \\
& \text { for continuous r.v.'s }
\end{aligned}
$$

If $X$ and $Y$ are two continuous r.v.'s the joint density function $f(x, y)$ the conditional expectation or the conditional mean of Y given X is $\mathrm{E}(\mathrm{Y} \mid \mathrm{X}=\mathrm{x})=\int_{-\infty}^{\infty} y f(y \mid x) d y$

Similarly, conditional mean of X given Y is $\mathrm{E}(\mathrm{X} \mid \mathrm{Y}=\mathrm{y})=\int_{-\infty}^{\infty} x f(x \mid y) d x$

- Median is the point, which divides the entire distribution into two equal parts. In case of continuous distribution median is the point, which divides the total area into two equal parts. Thus, if M is the median then $\int_{-\infty}^{M} f(x) d x=\int_{M}^{\infty} f(x) d x=1 / 2$. Thus, solving any one of the equations for $M$ we get the value of median. Median is unique
- Mode: Mode is the value for $\mathrm{f}(\mathrm{x})$ or $\mathrm{P}\left(\mathrm{x}_{\mathrm{i}}\right)$ at attains its maximum

For continuous r.v. $X$ mode is the solution of $\mathrm{f}^{1}(\mathrm{x})=0$ and $\mathrm{f}^{11}(\mathrm{x})<0$ provided it lies in the given interval. Mode may or may not be unique.

- Variance: Variance characterizes the variability in the distributions with same mean can still have different dispersion of data about their means

Variance of r.v. X denoted by $\operatorname{Var}(\mathrm{X})$ and is defined as
$\begin{aligned} \operatorname{Var}(\mathrm{X})=\mathrm{E}\left[(\mathrm{X}-\mu)^{2}\right]= & \sum_{x}(x-\mu)^{2} f(x)\left\{\begin{array}{ll}\text { for discrete } \\ & \int_{-\infty}^{\infty}(x-\mu)^{2} f(x) d x\end{array} \quad \begin{array}{l} \\ \\ \\ \\ \\ \end{array} \quad \begin{array}{l}\text { for continuous }\end{array}\right.\end{aligned}$
where $\mu=\mathrm{E}(\mathrm{X})$

- If c is any constant then $\mathrm{E}(\mathrm{cX})=\mathrm{c} \mathrm{E}(\mathrm{X})$
- If $X$ and $Y$ are two r.v.'s then $E(X+Y)=E(X)+E(Y)$
- IF $X, Y$ are two independent r.v.'s then $E(X Y)=E(X) E(Y)$
- If $X_{1}, X_{2},-\cdots----, X_{n}$ are random variables then $E\left(c_{1} X_{1}+c_{2} X_{2}+-----+c_{n} X_{n}\right)=$ $c_{1} E\left(X_{1}\right)+c_{2} E\left(X_{2}\right)+----+c_{n} E\left(X_{n}\right)$ for any scalars $c_{1}, c_{2},-----, c_{n}$ If all expectations exists
- If $\mathrm{X}_{1}, \mathrm{X}_{2},-\cdots----, \mathrm{X}_{\mathrm{n}}$ are independent r.v's then $\mathrm{E}\left(\prod_{i=1}^{n} X_{i}\right)=\prod_{i=1}^{n} E\left(X_{i}\right)$ if all expectations exists.
- $\quad \operatorname{Var}(\mathrm{X})=\mathrm{E}\left(\mathrm{X}^{2}\right)-[\mathrm{E}(\mathrm{X})]^{2}$
- If ' $c$ ' is any constant then $\operatorname{var}(c X)=c^{2} \operatorname{var}(X)$
- The quantity $\mathrm{E}\left[(\mathrm{X}-\mathrm{a})^{2}\right]$ is minimum when $\mathrm{a}=\mu=\mathrm{E}(\mathrm{X})$
- If X and Y are independent r.v.'s then $\operatorname{Var}(\mathrm{X} \pm \mathrm{Y})=\operatorname{Var}(\mathrm{X}) \pm \operatorname{Var}(\mathrm{Y})$


## Module-V

## PROBABILITY DISTRIBUTION

## Binomial Distribution

- A random variable X is said to follow binomial distribution if it assumes only non-negative values and its probability mass function is given by

$$
\mathrm{P}(\mathrm{X}=\mathrm{x})=\mathrm{P}(\mathrm{x})=\binom{n}{x} p^{x} q^{n-x} \quad \text { where } \mathrm{x}=0,1,2,3, \ldots \mathrm{n} \quad \mathrm{q}=1-\mathrm{p}
$$

where $\mathrm{n}, \mathrm{p}$ are known as parameters, n - number of independent trials p - probability of success in each
trial, q- probability of failure.

- Binomial distribution is a discrete distribution.
- The notation $X \sim B(n, p)$ is the random variable $X$ which follows the binomial distribution with parameters n and p
- If n trials constitute an experiment and the experiment is repeated N times the frequency function of the binomial distribution is given by $f(x)=N P(x)$. The expected frequencies of $0,1,2, \ldots . n$ successes are the successive terms of the binomial expansion $N(p+q)^{n}$
- The mean and variance of Binomial distribution are np, npq respectively.
- Mode of the Binomial distribution: Mode of B.D. Depending upon the values of ( $\mathrm{n}+1$ ) p
(i) If $(\mathrm{n}+1) \mathrm{p}$ is not an integer then there exists a unique modal value for binomial distribution and it is ' $m$ ' $=$ integral part of $(n+1) p$
(ii) If $(\mathrm{n}+1) \mathrm{p}$ is an integer say m then the distribution is Bi-Modal and the two modal values are m and $\mathrm{m}-1$
- Moment generating function of Binomial distribution: If $X \sim B(n, p)$ then $M_{X}(t)=\left(q+\mathrm{pe}^{t}\right)^{n}$
- The sum of two independent binomial variates is not a binomial varaite. In other words, Binomial distribution does not posses the additive or reproductive property.
- For B.D. $\gamma_{1}=\sqrt{ } \beta_{1}=\frac{1-2 p}{\sqrt{n p q}} \gamma_{2}=\beta_{2}-3=\frac{1-6 p q}{n p q}$
- If $X_{1} \sim B\left(n_{1}, p\right)$ and $X_{2} \sim B\left(n_{2}, p\right)$ then $X_{1}+X_{2} \sim B\left(n_{1}+n_{2}, p\right)$.Thus the B.D. Possesses the additive or reproductive property if $\mathrm{p}_{1}=\mathrm{p}_{2}$


## Poisson Distribution

- Poisson Distribution is a limiting case of the Binomial distribution under the following conditions:
(i) n , the number of trials is infinitely large.
(ii) P , the constant probability of success for each trial is indefinitely small.
(iii) $n \mathrm{p}=\lambda$, is finite where $\lambda$ is a positive real number.
- A random variable X is said to follow a Poisson distribution if it assumes only non-negative values and
its p.m.f. is given by

$$
\mathrm{P}(\mathrm{x}, \lambda)=\mathrm{P}(\mathrm{X}=\mathrm{x}) \begin{cases}\frac{e^{-\lambda} \lambda^{x}}{x!} & \mathrm{x}=0,1,2,3, \ldots \ldots . \lambda>0 \\ 0 & \text { Other wise }\end{cases}
$$

Here $\lambda$ is known as the parameter of the distribution.

- We shall use the notation $\mathrm{X} \sim \mathrm{P}(\lambda)$ to denote that X is a Poisson variate with parameter $\lambda$
- Mean and variance of Poisson distribution are equal to $\lambda$.
- The coefficient of skewness and kurtosis of the poisson distribution are $\gamma_{1}=\sqrt{ } \beta_{1}=1 / \sqrt{ } \lambda$ and $\gamma_{2}=$ $\beta_{2}-3=1 / \lambda$. Hence the poisson distribution is always a skewed distribution. Proceeding to limit as $\lambda$ tends to infinity we get $\beta_{1}=0$ and $\beta_{2}=3$
- Mode of Poisson Distribution: Mode of P.D. Depending upon the value of $\lambda$
(i) when $\lambda$ is not an integer the distribution is uni- modal and integral part of $\lambda$ is the unique modal value.
(ii) When $\lambda=\mathrm{k}$ is an integer the distribution is bi-modal and the two modals are $\mathrm{k}-1$ and k .
- Sum of independent poisson variates is also poisson variate.
- The difference of two independent poisson variates is not a poisson variate.
- Moment generating function of the P.D.

If $\mathrm{X} \sim \mathrm{P}(\lambda)$ then $\mathrm{M}_{\mathrm{X}}(\mathrm{t})=e^{\lambda\left(e^{t}-1\right)}$

- Recurrence formula for the probabilities of P.D. ( Fitting of P.D.)
$\mathrm{P}(\mathrm{x}+1)=\frac{\lambda}{x+1} p(x)$
- Recurrence relation for the probabilities of B.D. (Fitting of B.D.)

$$
\mathrm{P}(\mathrm{x}+1)=\left\{\frac{n-x}{x+1} \cdot \frac{p}{q}\right\} p(x)
$$

## Normal Distribution

- A random variable $X$ is said to have a normal distribution with parameters $\mu$ called mean and $\sigma^{2}$ called variance if its density function is given by the probability law

$$
\mathrm{f}(\mathrm{x} ; \mu, \sigma)=\frac{1}{\sigma \sqrt{2 \pi}} \exp \left[\frac{-1}{2}\left\{\frac{x-\mu}{\sigma}\right\}^{2}\right], \quad-\infty<\mathrm{x}<\infty,-\infty<\mu<\infty, \sigma>0
$$

- A r.v. X with mean $\mu$ and variance $\sigma^{2}$ follows the normal distribution is denoted by $X \sim N\left(\mu, \sigma^{2}\right)$
- If $\mathrm{X} \sim \mathrm{N}\left(\mu, \sigma^{2}\right)$ then $\mathrm{Z}=\frac{X-\mu}{\sigma}$ is a standard normal variate with $\mathrm{E}(\mathrm{Z})=0$ and $\operatorname{var}(\mathrm{Z})=0$ and we write $\mathrm{Z} \sim \mathrm{N}(0,1)$
- The p.d.f. of standard normal variate Z is given by $\mathrm{f}(\mathrm{Z})=\frac{1}{\sqrt{2 \pi}} e^{-z^{2} / 2},-\infty<\mathrm{Z}<\infty$
- The distribution function $\mathrm{F}(\mathrm{Z})=\mathrm{P}(\mathrm{Z} \leq \mathrm{z})=\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{z} e^{-t^{2} / 2} d t$
- $\mathrm{F}(-\mathrm{z})=1-\mathrm{F}(\mathrm{z})$
- $\mathrm{P}(\mathrm{a}<\mathrm{z} \leq \mathrm{b})=\mathrm{P}(\mathrm{a} \leq \mathrm{z}<\mathrm{b})=\mathrm{P}(\mathrm{a}<\mathrm{z}<\mathrm{b})=\mathrm{P}(\mathrm{a} \leq \mathrm{z} \leq \mathrm{b})=\mathrm{F}(\mathrm{b})-\mathrm{F}(\mathrm{a})$
- If $\mathrm{X} \sim \mathrm{N}\left(\mu, \sigma^{2}\right)$ then $\mathrm{Z}=\frac{X-\mu}{\sigma}$ then $\mathrm{P}(\mathrm{a} \leq \mathrm{X} \leq \mathrm{b})=F\left(\frac{b-\mu}{\sigma}\right)-F\left(\frac{a-\mu}{\sigma}\right)$
- N.D. is another limiting form of the B.D. under the following conditions:
i) $n$, the number of trials is infinitely large.
ii) Neither p nor q is very small
- Chief Characteristics of the normal distribution and normal probability curve:
i) The curve is bell shaped and symmetrical about the line $x=\mu$
ii) Mean median and mode of the distribution coincide.
iii) As $x$ increases numerically $f(x)$ decreases rapidly.
iv) The maximum probability occurring at the point $x=\mu$ and is given by

$$
[\mathrm{P}(\mathrm{x})]_{\max }=1 / \sigma \sqrt{ } 2 \Pi
$$

v) $\beta_{1}=0$ and $\beta_{2}=3$
vi) $\mu_{2 \mathrm{r}+1}=0(\mathrm{r}=0,1,2 \ldots \ldots)$ and $\mu_{2 \mathrm{r}}=1.3 \cdot 5 \ldots .(2 \mathrm{r}-1) \sigma^{2 \mathrm{r}}$
vii) Since $f(x)$ being the probability can never be negative no portion of the curve lies below x - axis.
viii) Linear combination of independent normal variate is also a normal variate.
ix) X - axis is an asymptote to the curve.
x) The points of inflexion of the curve are given by $\mathrm{x}=\mu \pm \sigma, \mathrm{f}(\mathrm{x})=\frac{1}{\sigma \sqrt{2 \pi}} e^{-1 / 2}$
xi) Q.D. : M.D.: S.D. :: $\frac{2}{3} \sigma: \frac{4}{5} \sigma: \sigma:: \frac{2}{3}: \frac{4}{5}: 1$ Or Q.D. : M.D.: S.D. ::10:12:15
xii) Area property: $\quad \mathrm{P}(\mu-\sigma<\mathrm{X}<\mu+\sigma)=0.6826=\mathrm{P}(-1<\mathrm{Z}<1)$

$$
\begin{aligned}
\mathrm{P}(\mu-2 \sigma<\mathrm{X}<\mu+2 \sigma)= & 0.9544=\mathrm{P}(-2<\mathrm{Z}<2) \\
& \mathrm{P}(\mu-3 \sigma<\mathrm{X}<\mu+3 \sigma)=0.9973=\mathrm{P}(-3<\mathrm{Z}<3) \\
& \mathrm{P}(|\mathrm{Z}|>3)=0.0027
\end{aligned}
$$

- m.g.f. of $N . D$. If $X \sim N\left(\mu, \sigma^{2}\right)$ then $M_{X}(t)=e^{\mu t}+t^{2} \sigma^{2} / 2$

IfZ $\sim \mathrm{N}(0,1)$ then $\mathrm{M}_{\mathrm{Z}}(\mathrm{t})=e^{t^{2} / 2}$

## Continuity Correction:

- The N.D. applies to continuous random variables. It is often used to approximate distributions of discrete r.v. Provided that we make the continuity correction.
- If we want to approximate its distribution with a N.D. we must spread its values over a continuous scale. We do this by representing each integer $k$ by the interval from $k-1 / 2$ to $k+1 / 2$ and at least $k$ is represented by the interval to the right of $k-1 / 2$ to at most $k$ is represented by the interval to the left of $k+1 / 2$.
- Normal approximation to the B.D:

$$
\begin{aligned}
& \mathrm{X} \sim \mathrm{~B}(\mathrm{n}, \mathrm{p}) \text { and if } \mathrm{Z}=\frac{X-n p}{\sqrt{n p(1-p)}} \text { then } \mathrm{Z} \sim \mathrm{~N}(0,1) \text { as } \mathrm{n} \text { tends to infinity and } \mathrm{F}(\mathrm{Z})= \\
& \mathrm{F}(\mathrm{Z})=\mathrm{P}(\mathrm{Z} \leq \mathrm{z})=\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{z} e^{-t^{2} / 2} d t-\infty<\mathrm{Z}<\infty
\end{aligned}
$$

- Use the normal approximation to the B.D. only when (i) np and $n(1-p)$ are both greater than 15 (ii) n is small and p is close to $1 / 2$
- Poisson process: Poisson process is a random process in which the number of events (successes) x occurring in atime interval of length T is counted. It is continuous parameter, discrete stable process. By dividing T into n equal parts of length $\Delta \mathrm{t}$ we have $\mathrm{T}=\mathrm{n} . \Delta \mathrm{T}$. Assuming that (i) $\mathrm{P} \propto \Delta \mathrm{T}$ or $\mathrm{P}=\alpha \Delta \mathrm{t}$ (ii) The occurrence of events are independent (iii) The probability of more than one substance during a small time interval $\Delta t$ is negligible.

As $\mathrm{n} \rightarrow \infty$, the probability of x success during a time interval T follows the P.D. with parameter $\lambda=n \mathrm{p}=\alpha \mathbf{T}$ where $\alpha$ is the average(mean) number of successes for unit time.

## PROBLEMS:

1:A random variable x has the following probability function:


Find (i) $k$ (ii) $P(x<6)$ (iii) $P(x>6)$

## Solution:

(i) since the total probability is unity, we have $\sum_{x=0}^{n} p(x)=1$
i.e., $0+\mathrm{k}+2 \mathrm{k}+2 \mathrm{k}+3 \mathrm{k}+\quad k^{2}+7 k^{2}+\mathrm{k}=1$
i.e., $8 k^{2}+9 \mathrm{k}-1=0$
$k=1,-1 / 8$
(ii) $\mathrm{P}(\mathrm{x}<6)=0+\mathrm{k}+2 \mathrm{k}+2 \mathrm{k}+3 \mathrm{k}$

$$
=1+2+2+3=8
$$

iii) $\quad \begin{aligned} & \mathrm{P}(\mathrm{x}>6)=k^{2}+7 k^{2}+\mathrm{k} \\ & =9\end{aligned}$
2. Let $X$ denotes the minimum of the two numbers that appear when a pair of fair dice is thrown once. Determine (i) Discrete probability distribution (ii) Expectation (iii) Variance

## Solution:

When two dice are thrown, total number of outcomes is $6 \times 6-36$

$$
\begin{aligned}
&\{(1,1)(1,2)(1,3)(1,4)(1,5)(1,6) \\
&(2,1)(2,2)(2,3)(2,4)(2,5)(2,6) \\
&=(3,1)(3,2)(3,3)(3,4)(3,5)(3,6) \\
&=(4,1)(4,2)(4,3)(4,4)(4,5)(4,6) \\
&(5,1)(5,2)(5,3)(5,4)(5,5)(5,6) \\
&(6,1)(6,2)(6,3)(6,4)(6,5)(6,6)
\end{aligned}
$$

$$
\text { In this case, sample space } S=\frac{(3,1)(3,2)(3,3)(3,4)(3,5)(3,6)}{(4,1)(4,2)(4,3)(4,4)(4,5)(4,6)}
$$

If the random variable $X$ assigns the minimum of its number in $S$, then the sample space $S=$
$\left\{\begin{array}{llllll}1 & 1 & 1 & 1 & 1 & 1 \\ 1 & 2 & 2 & 2 & 2 & 2 \\ 1 & 2 & 3 & 3 & 3 & 3 \\ 1 & 2 & 3 & 4 & 4 & 4 \\ 1 & 2 & 3 & 4 & 5 & 5 \\ 1 & 2 & 3 & 4 & 5 & 6\end{array}\right\}$

The minimum number could be $1,2,3,4,5,6$
For minimum 1, the favorable cases are 11
Therefore, $\mathrm{P}(\mathrm{x}=1)=11 / 36$
$\mathrm{P}(\mathrm{x}=2)=9 / 36, \mathrm{P}(\mathrm{x}=3)=7 / 36, \mathrm{P}(\mathrm{x}=4)=5 / 36, \mathrm{P}(\mathrm{x}=5)=3 / 36, \mathrm{P}(\mathrm{x}=6)=1 / 36$
The probability distribution is

| X | 1 | 2 | 3 | 4 | 5 | 6 |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| $\mathrm{P}(\mathrm{x})$ | $11 / 36$ | $9 / 36$ | $7 / 36$ | $5 / 36$ | $3 / 36$ | $1 / 36$ |

(ii) Expectation mean $=\sum p_{i} x_{i}$
$E(x)=1 \frac{11}{36}+2 \frac{9}{36}+3 \frac{7}{36}+4 \frac{5}{36}+5 \frac{3}{36}+6 \frac{1}{36}$

Or $\mu=\frac{1}{36}[11+8+21+20+15+6]=\frac{9}{36}=2.5278$
(ii) $\quad$ variance $=\sum p_{i} x^{2}{ }_{i}-\mu^{2}$

$$
\begin{aligned}
& E(x)=\frac{11}{36} 1+\frac{9}{36} 4+\frac{7}{36} 9+\frac{5}{36} 16+\frac{3}{36} 25+\frac{1}{36} 36-(2.5278)^{2} \\
& =1.9713
\end{aligned}
$$

3: A continuous random variable has the probability density function
$f(x)=\left\{\begin{array}{l}k x e^{-\lambda x}, \text { for } x \geq 0, \lambda>0 \\ 0, \text { otherwise }\end{array}\right.$
Determine (i) k (ii) Mean (iii) Variance

## Solution:

(i) since the total probability is unity, we have $\int_{-\infty}^{\infty} f(x) d x=1$

$$
\begin{aligned}
& \int_{-\infty}^{0} 0 d x+\int_{0}^{\infty} k x e^{-\lambda x} d x=1 \\
& \text { i.e., } \int_{0}^{\infty} k x e^{-\lambda x} d x=1 \\
& k\left[x\left(\frac{e^{-\lambda x}}{-\lambda}\right)-1\left(\frac{e^{-\lambda x}}{\lambda^{2}}\right)\right]_{0}^{\infty} \text { or } k=\lambda^{2}
\end{aligned}
$$

(ii) mean of the distribution $\mu=\int_{-\infty}^{\infty} x f(x) d x$

$$
\begin{aligned}
& \int_{-\infty}^{0} 0 d x+\int_{0}^{\infty} k x^{2} e^{-\lambda x} d x \\
& \lambda^{2}\left[x^{2}\left(\frac{e^{-\lambda x}}{-\lambda}\right)-2 x\left(\frac{e^{-\lambda x}}{\lambda^{2}}\right)+2\left(\frac{e^{-\lambda x}}{\lambda^{3}}\right)\right]_{0}^{\infty} \\
& =\frac{2}{\lambda}
\end{aligned}
$$

Variance of the distribution $\sigma^{2}=\int_{-\infty}^{\infty} x^{2} f(x) d x-\mu^{2}$

$$
\begin{aligned}
& \sigma^{2}=\int_{-\infty}^{\infty} x^{2} f(x) d x-\frac{4}{\lambda^{2}} \\
& \lambda^{2}\left[x^{3}\left(\frac{e^{-\lambda x}}{-\lambda}\right)-3 x^{2}\left(\frac{e^{-\lambda x}}{\lambda^{2}}\right)+6 x\left(\frac{e^{-\lambda x}}{\lambda^{3}}\right)-6\left(\frac{e^{-\lambda x}}{\lambda^{4}}\right)\right]_{0}^{\infty}-\frac{4}{\lambda^{2}} \\
& =\frac{2}{\lambda^{2}}
\end{aligned}
$$

## 4:

Out of 800 families with 5 children each, how many would you expect to have (i) 3 boys (ii) 5 girls (iii)either 2 or 3 boys ? Assume equal probabilities for boys and girls

## Solution

$$
\mathrm{P}(3 \text { boys })=\mathrm{P}(\mathrm{r}=3)=\mathrm{P}(3)=\frac{1}{2^{5}}{ }^{5} C_{3}=\frac{5}{16} \text { per family }
$$

Thus for 800 families the probability of number of families having 3 boys $=\frac{5}{16}(800)=250$ families
$\mathrm{P}(5$ girls $)=\mathrm{P}($ no boys $)=\mathrm{P}(\mathrm{r}=0)=\frac{1}{2^{5}}{ }^{5} C_{0}=\frac{1}{32}$ per family
Thus for 800 families the probability of number of families having 5 girls $=\frac{1}{32}(800)=25$ families
$\mathrm{P}($ either 2 or 3 boys $=\mathrm{P}(\mathrm{r}=2)+\mathrm{P}(\mathrm{r}=3)=\mathrm{P}(2)+\mathrm{P}(3)$
$\frac{1}{2^{5}}{ }^{5} C_{2}+\frac{1}{2^{5}}{ }^{5} C_{3}=5 / 8$ per family
Expected number of families with 2 or 3 boys $=\frac{5}{8}(800)=500$ families.

5: Average number of accidents on any day on a national highway is 1.8 . Determine the probability that the number of accidents is (i) at least one (ii) at most one

## Solution:

Mean $=\lambda=1.8$
We have $\mathrm{P}(\mathrm{X}=\mathrm{x})=\mathrm{p}(\mathrm{x}) \frac{e^{-\lambda} \lambda^{x}}{x!}=\frac{e^{-1.8} 1.8^{x}}{x!}$
$P($ at least one $)=P(x \geq 1)=1-P(x=0)$
$=1-0.1653$
$=0.8347$
$\mathrm{P}($ at most one $)=\mathrm{P}(\mathrm{x} \leq 1)$
$=P(x=0)+P(x=1)$
$=0.4628$
6: The mean weight of 800 male students at a certain college is 140 kg and the standard deviation is 10 kg assuming that the weights are normally distributed find how many students weigh I) Between 130 and 148 kg ii) more than 152 kg

## Solution:

Let $\mu$ be the mean and $\sigma$ be the standard deviation. Then $\mu=140 \mathrm{~kg}$ and $\sigma=10$ pounds
(i) When $\mathrm{x}=138, z=\frac{x-\mu}{\sigma}=\frac{138-140}{10}=-0.2=z_{1}$

When $\mathrm{x}=138, z=\frac{x-\mu}{\sigma}=\frac{148-140}{10}=0.8=z_{2}$
$\therefore \mathrm{P}(138 \leq \mathrm{x} \leq 148)=\mathrm{P}(-0.2 \leq \mathrm{z} \leq 0.8)$
$=\mathrm{A}\left(z_{2}\right)+\mathrm{A}\left(z_{1}\right)$
$=\mathrm{A}(0.8)+\mathrm{A}(0.2)=0.2881+0.0793=0.3674$
Hence the number of students whose weights are between 138 kg and 140 kg $=0.3674 \times 800=294$
(ii) When $\mathrm{x}=152, \frac{x-\mu}{\sigma}=\frac{152-140}{10}=1.2=\mathrm{z}_{1}$

Therefore $\mathrm{P}(\mathrm{x}>152)=\mathrm{P}\left(\mathrm{z}>\mathrm{Z}_{1}\right)=0.5-\mathrm{A}\left(\mathrm{z}_{1}\right)$

$$
=0.5-0.3849=0.1151
$$

Therefore number of students whose weights are more than $152 \mathrm{~kg}=800 \times 0.1151=92$.

## Exercise Problems:

1. Two coins are tossed simultaneously. Let $X$ denotes the number of heads then find $i$ ) $\mathrm{E}(\mathrm{X})$ ii) $\mathrm{E}\left(\mathrm{X}^{2}\right)$ iii $) \mathrm{E}\left(\mathrm{X}^{3}\right)$ iv) $\mathrm{V}(\mathrm{X})$
2. If $\mathrm{f}(\mathrm{x})=\mathrm{k} e^{-|x|}$ is probability density function in the interval, $-\infty<x<\infty$, then find i) k ii) Mean iii) Variance iv) $\mathrm{P}(0<\mathrm{x}<4)$
3. Out of 20 tape recorders 5 are defective. Find the standard deviation of defective in the sample of 10 randomly chosen tape recorders. Find (i) $\mathrm{P}(\mathrm{X}=0)$ (ii) $\mathrm{P}(\mathrm{X}=1)$ (iii) $\mathrm{P}(\mathrm{X}=2)$ (iv) $\mathrm{P}(1<\mathrm{X}<4)$.
4. In 1000 sets of trials per an event of small probability the frequencies $f$ of the number of x of successes are

| f | 0 | 1 | 2 | 3 | 4 | 5 | 6 | 7 | Total |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| x | 305 | 365 | 210 | 80 | 28 | 9 | 2 | 1 | 1000 |

Fit the expected frequencies.
5. If X is a normal variate with mean 30 and standard deviation 5. Find the probabilities that
i) $\mathrm{P}(26 \leq \mathrm{X} \leq 40)$
ii) $\mathrm{P}(\mathrm{X} \geq 45)$
6. The marks obtained in Statistics in a certain examination found to be normally distributed. If

15\% of the students greater than or equal to 60 marks, $40 \%$ less than 30 marks. Find the
mean and standard deviation.
7.If a Poisson distribution is such that $P(X=1)=\frac{3}{2} P(X=3)$ then find (i) $P(X \geq 1)$ (ii) $P(X \leq 3)$ (iii) $P(2 \leq X \leq 5)$.

A random variable X has the following probability function:

| X | -2 | -1 | 0 | 1 | 2 | 3 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\mathrm{P}(\mathrm{x})$ | 0. | K | 0.2 | 2 K | 0.3 | K |
|  | 1 |  |  |  |  |  |

Then find (i) k (ii) mean (iii) variance (iv) $\mathrm{P}(0<\mathrm{x}<3)$

