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74 Unit 4 □ –Contents : Banach Algebra, Invertible | Non-invertible Elements, their Proper- ties and Representations, Continuity of Inverse Mapping, Topological Divisor of Zero, Resolvant Set, Spectrum, Spectral Radius, its formula) 4.1 In a Banach Algebra two apparently diverse trains of disciplines—topological and Algebraic are in conjunction to make a single mathematical system. Definition 4.1.1. An algebra X over a real / Complex field is a system of two compositions, namely, a Vector-space in which multiplication is defined subject to :— (1) $(xy)z = x(yz)$ for any three members x, y, z

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of X (2a) $x(y + z) = xy + xz$, and (2b) $(x + y)z = xz + yz$ for any three members x, y, z

of X , (3) $\alpha(xy) = (\alpha x)y = x(\alpha y)$

for any scalar α and any two element $x, y \in X$. We shall generally deal with complex scalar field \mathbb{C} and term X as an algebra (over \mathbb{C}). Algebra X is said to be commutative if multiplication operation in X is commutative ; That is to say that $xy = yx$ for all members x and y in X . An algebra X is said to possess an identity if there is a member e called the identity in X such that $xe = ex = x$ for all $x \in X$. It is a routine business to see that identity element in a Banach Algebra is unique. Example 4.1.1. The space \mathbb{R} of all reals as a real Vector-space becomes a commutative (real Banach Algebra where multiplication is taken as the usual arithmetic multiplication. Here we see that this multiplication operation in \mathbb{R} is indistinguishable from one of the Principal operation, namely scalar multiplication in

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Vector-space. 75 Example 4.1.2. Let X be a Vector space and

$L(X, X)$ be the collection of all linear operators $T : X \rightarrow X$. Then under usual addition and multiplication (Composition) one makes a routine exercise to check that $L(X, X)$ is an Algebra with identity element as the Identity operator $I : X \rightarrow X$ where $I(x) = x$ for $x \in X$. Notice that I is not the same as additive identity, namely the zero operator O sending every member $x \in X$ to the zero vector in X . In general, $L(X, X)$ is not a commutative Algebra. Neither it has divisor of zero. Definition 4.1.2. An algebra X is said to be a Banach Algebra if X is a Banach space (over \mathbb{C}) with respect to a norm $\| \cdot \|$ such that for $x, y \in X$, $\| xy \| \leq \| x \| \| y \|$. If X possesses the identity element e , then $\| e \| = 1$. Example 4.1.3. Consider the Banach space $C[a, b]$ of all real-valued continuous functions over the closed interval $[a, b]$ of reals with sup norm. Then $C[a, b]$ is a commutative Banach Algebra with identity $e =$ constant function equal to 1 throughout $[a, b]$, and with usual multiplication, namely $(xy)(t) = x(t)y(t)$ in $a \leq t \leq b$ and $x, y \in C[a, b]$. Example 4.1.4. Let \mathbb{C}_n denote the Vector space of all complex polynomials of degree $\leq n$. Since this is a finite dimensional vector space it becomes a Banach space with respect to the norm of \mathbb{C}_n defined as $\| \sum_{i=0}^n x_i t^i \| \leq 1$ where $x(t) = a_0 + a_1 t + a_2 t^2 + \dots + a_n t^n \in \mathbb{C}_n$, and where product xy is defined like $(\sum_{i=0}^n x_i t^i)(\sum_{j=0}^n y_j t^j) = \sum_{k=0}^{2n} y_k t^k$. Then \mathbb{C}_n becomes a Banach Algebra. Example 4.1.5. The collection $BdL(X, X)$ of all bounded Linear operators $T : X \rightarrow X$ becomes a Normed Linear space when X is a Normed Linear space with usual operator norm $\| T \|$ as $T \in BdL(X, X)$. If X is a Banach space, then $BdL(X, X)$ becomes a Banach space. Taking multiplication of two members of $BdL(X, X)$ as their usual

composition it is now a routine exercise to check that $BdL(X, X)$ is a Banach Algebra, where the identity member equals to the Identity operator $I : X \rightarrow X$. As observed earlier Banach Algebra $BdL(X, X)$ may not be commutative. Take the case when $X =$ Euclidean n -space \mathbb{R}^n which is Banach space with usual norm. By Matrix representation Theorem every member of $BdL(\mathbb{R}^n, \mathbb{R}^n)$ is represented by a square matrix of size n over reals. As we know that matrix multiplication is not commutative, so $BdL(\mathbb{R}^n, \mathbb{R}^n)$ is not commutative. Theorem 4.1.1. Multiplication operation in a Banach Algebra X is a continuous operation. Proof : Let $\{x_n\}$ and $\{y_n\}$ be two sequences of elements in X such that $\lim_{n \rightarrow \infty} \| x_n - x \| = 0$ and $\lim_{n \rightarrow \infty} \| y_n - y \| = 0$ in norm of X . So, $\lim_{n \rightarrow \infty} \| x_n y_n - xy \| = 0$. Now,

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$x_n y_n - xy = (x_n - x) y_n + x(y_n - y)$ gives $\| x_n y_n - xy \| \leq \| x_n - x \| \| y_n \| + \| x \| \| y_n - y \|$

$\lim_{n \rightarrow \infty} \| x_n y_n - xy \| = 0$ and in norm of X . Since $\{y_n\}$

is a convergent sequence in X we know that it is bounded and let $\|y_n\| \leq M$ for all n for some +ve real M . Therefore above reads as $\sum_{n=1}^{\infty} \|y_n\| < \infty$ as So, $\lim_{n \rightarrow \infty} \sum_{j=1}^n y_j$ exists. Theorem is proved. Definition 4.1.3. An element x in a Banach Algebra X with identity e is said to be invertible if x^{-1} (inverse of x) exists in X i.e. $x^{-1} \in X$ satisfying $x^{-1}x = xx^{-1} = e$. Otherwise, x is said to be a non-invertible element in X . Explanation : (I) If inverse of x exists in X ($x \in X$), x^{-1} is unique. Because suppose $yx = e = xz$, then we have, $y = ye = y(xz) = (yx)z = ez = z$. (II) If x and y are both invertible, then xy is invertible and $(xy)^{-1} = y^{-1}x^{-1}$. 77 Because, $(xy)(x^{-1}y^{-1}x^{-1}) = x(yy^{-1})x^{-1} = xex^{-1} = xx^{-1} = e$. and similarly, $(y^{-1}x^{-1})(xy) = e$. Theorem 4.1.2. The set G of all invertible elements in X forms a Group. The proof readily follows from Explanations (I) and (II) 4.2. Suppose X is a Banach algebra with identity e . Then, ofcourse, e is an invertible element in X ; There are non-invertible elements in X like 0 (zero vector in X). Below we like to derive a few facts about X where we know that $ex = xe = x$ for every $x \in X$. It will be shown that invertible elements are many in X in the sense that set of all invertible elements of X forms an open set in X . Theorem as under presented demonstrates that even members of X close to e are invertible. Theorem 4.2.1. If $x \in X$ satisfies $\|x - e\| < 1$, then $e - x$ is invertible and $(e - x)^{-1} = \sum_{j=0}^{\infty} x^j$.

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$\ x - e\ < 1$, then $e - x$ is invertible and $(e - x)^{-1} = \sum_{j=0}^{\infty} x^j$			

Proof : By Induction we have $\|x^j\| < 1$ for all +ve integers j . Therefore the series $\sum_{j=0}^{\infty} \|x^j\| < \infty$ is convergent, because $\|x\| < 1$. By completeness in X the infinite series $\sum_{j=0}^{\infty} x^j$ is convergent with sum $s \in X$. Put $s = \sum_{j=0}^{\infty} x^j$. We now verify that Inverse of $e - x$ equal to s i.e. $(e - x)s = s(e - x) = e$. For any natural number n we have

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$(e - x)(e + x + x^2 + \dots + x^n) = (e + x + x^2 + \dots + x^n)(e - x) = e - x^{n+1}$			

Because $\|x\| < 1$, we have $\lim_{n \rightarrow \infty} \|x^{n+1}\| = 0$ (zero vector in X). So one can pass on $\lim_{n \rightarrow \infty}$ in (1) and since multiplication operation is continuous we have, $(e - x)s = s(e - x) = e$. That gives, $(e - x)^{-1}$ exists and it is equal to s . i.e. $(e - x)^{-1} = \sum_{j=0}^{\infty} x^j$.

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Corollary 1. If $x \in X$, satisfies $\ e - x\ < 1$, then x^{-1} exists, and $x^{-1} = \sum_{j=0}^{\infty} x^j$. For proof replaces x in Theorem 4.2.1 by $e - x$ and therefore we get $(e - (e - x))^{-1} = \sum_{j=0}^{\infty} (e - x)^j$ or, $x^{-1} = \sum_{j=0}^{\infty} (e - x)^j$. Corollary 2. Suppose $x \in X$ and a scalar λ satisfies $\ x\ < \lambda $. Then $(\lambda e - x)$ is invertible and $(\lambda e - x)^{-1} = \sum_{j=0}^{\infty} \lambda^{-j-1} x^j$.			

Proof. Write Apply Corollary 1 as above and get $(\lambda e - x)^{-1} = \sum_{j=0}^{\infty} \lambda^{-j-1} x^j$; Therefore, $(\lambda e - x)^{-1}$ exists, and therefore $(\lambda e - x)^{-1}$ exists. Then $(\lambda e - x)^{-1}(\lambda e - x) = (\lambda e - x)(\lambda e - x)^{-1} = \lambda e$.

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$\lambda e - x$ is invertible and $(\lambda e - x)^{-1} = \sum_{j=0}^{\infty} \lambda^{-j-1} x^j$			

Theorem 4.2.2. The set G of all invertible elements of X forms an open set in X . Proof : Take $x_0 \in G$; Take an open ball $B_r(x_0)$ with radius $r > 0$. Then $x \in B_r(x_0)$ if and only if $\|x - x_0\| < r$. Put $y = x - x_0$ and $z = e - y$; Then we have, $z \in G$.

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X and x^{-1} satisfies $\|x\| \geq 1$, then $f(x) \neq 0$. Proof : Let $\|x\| \geq 1$, then $e - x$ is invertible, and $f(e - x) \neq 0$ or, $f(e) - f(x) \neq 0$ or, $f(x) \neq f(e) = 1$, giving $f(x) \neq 1$

Next let $\|x\| < 1$; choose a scalar α with $0 < \alpha < 1$ such that $\alpha \|x\| < 1$ or, $\alpha \|x\| < 1$; As above $f(\alpha x) \neq 1$ or, $\alpha f(x) \neq 1$ or, $f(x) \neq 1/\alpha$, where α satisfies $0 < \alpha < 1$ and therefore $f(x) \neq 1$. The proof is now complete. The converse part of Theorem 4.3.3 is also true. The proof of converse part is rather long and involved; We present a special case for simplicity.

83 Theorem 4.3.5. If ϕ is a linear functional over X with $\phi(e) = 1$ and $\phi(x) \neq 0$ for every invertible element x in X and Null space of ϕ is a sub-algebra of X , then ϕ is a complex homomorphism. Proof : Let Null-space of ϕ be denoted by N . Take $x \notin N$. Put $\phi(x) = \alpha$. Now $\phi(x - \alpha e) = \phi(x) - \alpha \phi(e) = 0$; So $(x - \alpha e) \in N$. Put $x - \alpha e = a$; so that we may represent $x = a + \alpha e = a + \phi(x)e$, where $a \in N$. Similarly, write $y \in X$ as $y = b + \phi(y)e$, where $b \in N$. Therefore, $xy = ab + \phi(y)a + \phi(ab) = 0$; because N is Null-space of ϕ which is also an Algebra ($ab \in N$). Therefore we have from above $\phi(xy) = \phi(x)\phi(y)$ for all $x, y \in X$ and proof is complete. 4.4 Resolvent set; Spectrum As before X is taken as complex Banach Algebra with identity e . Take $x \in X$. Definition 4.4.1 (a) The resolvent set $\rho(x)$ of x is equal to the set of the scalars λ such that $x - \lambda e$ is invertible. i.e. $\rho(x) = \{\lambda \in \mathbb{C} : (x - \lambda e)^{-1} \text{ exists in } X\}$ (b) The Complement $\sigma(x) = \{\lambda \in \mathbb{C} : (x - \lambda e)^{-1} \text{ does not exist in } X\}$ is called spectrum of x , denoted by $\sigma(x)$. Explanation : Any scalar $\lambda \in \sigma(x)$ is called a spectral value of x . Thus we have $\rho(x) \cup \sigma(x) = \mathbb{C}$ with $\rho(x) \cap \sigma(x) = \emptyset$. Take $x \in X$ fixed. Now consider the mapping $\rho : \rho(x) \rightarrow X$ given by $\rho(\lambda) = (x - \lambda e)^{-1}$. We may write $x(\lambda) = (x - \lambda e)^{-1}$

84 This mapping is known as the resolvent function associated with $x \in X$. So a resolvent function is a Vector-valued function over $\rho(x)$ with range in a Banach algebra. Remark : Let us take $\lambda_1, \lambda_2 \in \rho(x)$, Then

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$x(\lambda_1) = (x - \lambda_1 e)^{-1}$ and $x(\lambda_2) = (x - \lambda_2 e)^{-1}$; and $x(\lambda_1)^{-1} x(\lambda_2) = (x - \lambda_1 e) x(\lambda_2) = (x - \lambda_2 e + \lambda_2 e - \lambda_1 e) x(\lambda_2) = (x - \lambda_2 e) x(\lambda_2) + (\lambda_2 e - \lambda_1 e) x(\lambda_2) = e + (\lambda_2 - \lambda_1) x(\lambda_2)$ That gives $x(\lambda_2) = x(\lambda_1) + (\lambda_2 - \lambda_1) x(\lambda_2) x(\lambda_1)$

or, (*) Theorem 4.4.1. The resolvent function $x(\lambda)$ is an analytic function. Proof : Take $\lambda_0 \in \rho(x)$ with $\lambda_0 \in \rho(x)$. From (*) above

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we have $x(\lambda) = (x - \lambda e)^{-1}$. Now $\lim_{\lambda \rightarrow \lambda_0} x(\lambda) = \lim_{\lambda \rightarrow \lambda_0} (x - \lambda e)^{-1} = (x - \lambda_0 e)^{-1} = x(\lambda_0)$. Therefore, $\lim_{\lambda \rightarrow \lambda_0} x(\lambda) = x(\lambda_0)$. $\lambda_0 \in \rho(x)$ whenever $\lambda_0 \in \rho(x)$

So derivative exists at λ_0 , and hence $x(\lambda)$ is an analytic function. This is what was wanted. Definition 4.4.2. For $x \in X$, spectral, radius of x , denoted by $r(x)$ is given as the real number. $r(x) = \sup\{|\lambda| : \lambda \in \sigma(x)\}$ Explanations : If $\lambda \in \sigma(x)$ is such that $\lambda \in \sigma(x)$ we have; Therefore, always.

85 Theorem 4.4.2. $\sigma(x)$ is a compact set of scalars. Proof : For $x \in X$ we have $r(x) < \infty$ and $\lambda \in \sigma(x)$ whenever $|\lambda| \leq r(x)$, it follows that spectrum $\sigma(x)$ is a bounded set of scalars. We show that it is closed in \mathbb{C} , by showing that its complement $\rho(x) = \rho(x)$ is open in \mathbb{C} . Let us look at the function $f : \rho(x) \rightarrow X$ given by $f(\lambda) = x - \lambda e$ as before. This is a continuous function of scalar λ . Take $\lambda_0 \in \rho(x)$. That means $x - \lambda_0 e$ is invertible. So $(x - \lambda_0 e) \in G$. Since G is open, we find an open ball $B(x - \lambda_0 e, r)$ centred at $(x - \lambda_0 e)$ with a +Ve radius r such that $B(x - \lambda_0 e, r) \subset G$. Since f is continuous at λ_0 , choose a +Ve δ such that $\|f(\lambda) - f(\lambda_0)\| < r$ whenever $|\lambda - \lambda_0| < \delta$ i.e., $f(\lambda) \in B(x - \lambda_0 e, r) \subset G$ whenever, $|\lambda - \lambda_0| < \delta$ i.e., $\lambda \in \rho(x)$ whenever $|\lambda - \lambda_0| < \delta$. Hence λ_0 is an Interior point of $\rho(x)$, and $\rho(x)$ is shown to be an open set. Theorem 4.4.3. For $x \in X$, spectrum $\sigma(x)$ is compact. Proof : Consider the Dual

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X^* of X and take $f \in X^*$. For $\epsilon > 0$, Let $f_\epsilon(x) = f(x - \epsilon e) - 1 = f(x)$. Since f is continuous it follows that f_ϵ is a continuous function of x over X . We have already had $x \in X$ for $\epsilon > 0$. i.e., $f_\epsilon(x) = f(x)$.

Now, $\lim_{\epsilon \rightarrow 0} f_\epsilon(x) = f(x)$. This shows that f is analytic on X . Further, $f_\epsilon(x) = f(x)$.
 a f
 e j 1 1
 86 For large value of $|\epsilon|$ we, have and therefore,

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$e \in X$ and $\|e\| = 1$, and therefore, $\|x - \epsilon e\| \geq \|x\| - |\epsilon|$. So, $\|x - \epsilon e\| \geq 1 - |\epsilon|$.

Now from as in above, We pass on the $|\epsilon|$ and obtain $\lim_{\epsilon \rightarrow 0} f_\epsilon(x) = f(x)$. If $f(x) = 0$, we shall have $f(x) = 0$ and f becomes an entire function. So by Liouville's Theorem f must be a constant function, and from limit above we see that this constant = 0. i.e. $f(x) = 0$ for all $x \in X$. This is true for every member f coming from X^* , and therefore it follows that $X^* = \{0\}$. But this is not the case ; because $\|e\| = \|x - \epsilon e\| = \|x\|$ — a contradiction. Therefore conclusion is that $X^* = X$. Theorem 4.4.4 If a Banach Algebra X with identity e has every non-zero member invertible then X is isometrically isomorphic to scalar field \mathbb{C} . (This extraordinary important result is due to Gelfand and Mazur who had left memorable marks in Advanced Functional Analysis) Proof. Take $x \in X$. Then Theorem 4.4.3 says that x is invertible. So there is a scalar λ such that $x - \lambda e$ is not invertible. By assumption every non-zero member of X is invertible. Therefore $x - \lambda e = 0$ or $x = \lambda e$.
 87 Now if λ_1 and λ_2 are two scalars with $x = \lambda_1 e = \lambda_2 e$, Then $\lambda_1 = \lambda_2$. x is a unique multiple of e . Consider the mapping $\phi: X \rightarrow \mathbb{C}$ given by $\phi(x) = \lambda$. This mapping ϕ is $1 - 1$, Linear plus ϕ is onto. Then ϕ is a desired isomorphism as wanted. Theorem 4.4.5 If zero is the only Topological divisor of zero in X , then X is isometrically Isomorphic to the scalarfield \mathbb{C} . Proof. Take $x \in X$. Then $\phi(x) = \lambda$; $\phi(x)$ is also bounded. Let λ be a boundary point of $\phi(x)$. Then $x - \lambda e$ is a topological divisor of zero. By assumption $x - \lambda e = 0$ gives $x = \lambda e$. Now one can copy rest of the proof as in proof of Theorem 4.4.4 to conclude that X is Isomorphic to \mathbb{C} as desired. 4.5 Spectral radius formula Let $x \in X$ and $\sigma(x)$ is spectrum. We know that $\|x^n\| \leq \|x\|^n$. Theorem 4.5.1 If $p(t)$ is a polynomial with complex coefficients and $x \in X$, then $\|p(x)\| \leq p(\|x\|)$. Proof. The proof proceeds by stages. First suppose $p(t)$ is a constant polynomial. Say $p(t) = c \neq 0$, and we have $\|p(x)\| = \|c e\| = |c| \|e\| = |c|$. Now $p(x) = \{p(t) : t \in \sigma(x)\} = \{c\}$. So in this case $\|p(x)\| = p(\|x\|)$. For any member $z \in X$ and any scalar λ we show that $\|p(z)\| = p(\|z\|)$. This is ok when $\lambda = 0$. Suppose $\lambda \neq 0$. Then take $p(z) = \lambda z - \lambda e$ is not invertible $\lambda z - \lambda e$ is not invertible
 88 $\|p(z)\| = p(\|z\|)$. Let us now consider polynomial with leading coefficient equal to 1, and let $p(t) = t^n + a_{n-1}t^{n-1} + \dots + a_1t + a_0$ ($n \geq 1$), and take $p(t) - \lambda$. Since scalarfield \mathbb{C} is algebraically closed we know that $p(t) - \lambda$ is completely factorisable like, $p(t) - \lambda = (t - \lambda_1)(t - \lambda_2) \dots (t - \lambda_n)$. (1) write x for t and set $p(x) - \lambda e = (x - \lambda_1 e)(x - \lambda_2 e) \dots (x - \lambda_n e)$. (2) If $p(x) - \lambda e$ is not invertible, then one of factors $(x - \lambda_j e)$ must be non-invertible and in that case $\lambda_j \in \sigma(x)$. That implies $p(\lambda_j) = \lambda$. (3) Taking λ_j for t in (1) above we see that $p(\lambda_j) = \lambda$ and (3) becomes $\lambda = p(\lambda_j)$, thus we have shown $\|p(x)\| \geq p(\|x\|)$. To obtain opposite inclusion relation let $\lambda \in \sigma(p(x))$; by Definiton of $\sigma(p(x))$, we find $\lambda_j \in \sigma(x)$ such that $\lambda = p(\lambda_j)$. Now from $p(t) - \lambda = (t - \lambda_1)(t - \lambda_2) \dots (t - \lambda_j) \dots (t - \lambda_n)$, it is clear that λ_j is a root of $p(t) - \lambda$. Taking x for t we obtain. $p(x) - \lambda e = (x - \lambda_1 e) \dots (x - \lambda_j e) \dots (x - \lambda_n e)$. (4) If $p(x) - \lambda e$ were invertible, we could have multiplied both sides of (4) on left by $(p(x) - \lambda e)^{-1}$ and move $(x - \lambda_j e)$ all the way to the right to get $e = (p(x) - \lambda e)^{-1} [(x - \lambda_1 e) \dots (x - \lambda_n e)] (x - \lambda_j e)$. (5) to conclude that $(x - \lambda_j e)$ has left inverse. similarly $(x - \lambda_j e)$ has right inverse—a contradiction that $\lambda_j \in \sigma(x)$. Therefore we conclude that $\lambda \in \sigma(p(x))$, and that implies $p(\|x\|) \geq \|p(x)\|$. The proof is now complete. Corollary : $\|p(x)\| = p(\|x\|)$ for any +ve integer n . Theorem 4.5.2. (Spectral radius formula) :

89 Proof : We have $\sup \{ \dots \} \geq \limsup x_n$

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<p>$\limsup x_n = \inf \{ \dots \}$ We always have $\limsup x_n \geq \liminf x_n$ or, $\limsup x_n \geq \liminf x_n$ This gives, $\lim x_n$</p>			

$\limsup x_n = \liminf x_n$ (*)

Since inferior limit of a sequence \leq its superior limit, if it is shown that $\limsup x_n = \liminf x_n$ then $\lim x_n$ exists and $\lim x_n = \limsup x_n = \liminf x_n$. (**)

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<p>$\limsup x_n = \liminf x_n$ (from *) and this implies $\lim x_n$ exists and $\lim x_n = \limsup x_n = \liminf x_n$.</p>			

Now (**) is obtained by computing the radius of convergence of a power series via Cauchy-Hadamard formula. Example 4.5.1. Let X be a Banach algebra with identity e . If $x \in X$ and there are $y, z \in X$ show that x is invertible and $y = z = x^{-1}$. Solution : Here $y = ye = yxz = ez = z$

90 Therefore $yx = e = xz = xy$, showing that x has an inverse y and hence $z = y$ i.e. $y = z = x^{-1}$. EXERCISE-A Short-Answer type questions 1. If x is an invertible element in a Banach Algebra X with identity e such that x commutes with $y \in X$, show that x^{-1} commutes with y . (Here $xy = yx$: So $x^{-1}xy = x^{-1}yx$ or $ey = x^{-1}yx$ or $y = x^{-1}yx$ or $yx^{-1} = x^{-1}yx$ or $yx^{-1} = x^{-1}y$ Here x^{-1} and y commute.) 2. If $\{x_n\}$ and $\{y_n\}$ are two Cauchy sequences in a Banach Algebra X , show that $\{x_n y_n\}$ is a Cauchy sequence in X . 3. For a Banach Algebra X , and for Identity operator $I : X \rightarrow X$, find $\|I\|$. 4. If $A \in B(X, X)$ in a Banach Algebra X with identity e , then show that $\|A^{-1}\| \geq \|A\|^{-1}$. 5. If $e - xy$ is invertible in a Banach Algebra X , then show that $e - yx$ is also invertible in X where $e =$ identity element in X , and $x, y \in X$. 6. Let X be a Banach Algebra and G is the set of all invertible members of X . Show that mapping $\phi : G \rightarrow G$ given by $\phi(x) = x^{-1}$ in G is a Homeomorphism. EXERCISE-B 1. Let X be a complex-Banach space and $B(X, X)$ denote the Banach Space of all bounded linear transformations $T : X \rightarrow X$. If $A \in B(X, X)$ and λ is a scalar satisfying $|\lambda| < \|A\|^{-1}$ show that $\lambda I - A$ is invertible and $\|(\lambda I - A)^{-1}\| \leq \frac{1}{|\lambda| - \|A\|}$ where I is the identity operator. 2. Let X be a commutative Banach Algebra with identity, then for any $x \in X$, show that $\|x^n\| = (\|x\|)^n$.

91 3. If X is a commutative Banach Algebra, and if $x, y \in X$, show that $\|xy\| \leq \|x\| \|y\|$. 4. Let X be a commutative Banach Algebra with identity e with $\|e\| = 1$, and let $f : X \rightarrow \mathbb{C}$ be a non-zero homomorphism; show that $\|f\| = 1$. 5. Let χ be a continuous character of topological Group G , Prove that χ is uniformly continuous. 6. Let H be a closed sub-group of a topological Group G , Prove that dual of G/H is isomorphic and homeomorphic to the sub-group of G^* comprising of all characters that are constants on H and its cosets. 7. Suppose X is a Banach Algebra. If there is a constant $m > 0$ such that $\|xy\| \leq m \|x\| \|y\|$ for all $x, y \in X$, then show that X is isomorphic to \mathbb{C} . EXERCISE-C 1. Show that following statements are equivalent in a Banach Algebra

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<p>X. (i) $\ x^2\ = \ x\ ^2$ for all $x \in X$ and (ii) $\ x\ = \ x\$ for all $x \in X$</p>			

X. 2. In a Banach Algebra X with identity e if $x \in X$ satisfies $\|x\| > 1$, show that $\|(e - x)^{-1} - e^{-1}\| \leq \|x\|^{-2}$. 3. In a Banach Algebra X with identity e if x is invertible and y satisfies $\|yx^{-1}\| > 1$, show that $x - y$ is invertible and $(x - y)^{-1} = x^{-1}y^{-1}x$. 4. If X is a commutative Banach Algebra and $x, y \in X$, show that $\|x + y\| \leq \|x\| + \|y\|$. 5. Let $M_n(\mathbb{C})$ denote the algebra of all complex matrices $n \times n$, show that $\|A\|_1 + \|A\|_2$ is a norm in $M_n(\mathbb{C})$ with respect to which $M_n(\mathbb{C})$ is a Banach Algebra.

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<p>of X (2a) $x(y + z) = xy + xz$, and (2b) $(x + y)z = xz + yz$ for any three members x, y, z</p> <p>SA RMSHEENA_FINAL.pdf (D143097593)</p>				
2/17	SUBMITTED TEXT	11 WORDS	100% MATCHING TEXT	11 WORDS
<p>Vector-space. 75 Example 4.1.2. Let X be a Vector space and</p> <p>SA SHAHEEBA_P.pdf_FINAL.pdf (D143096904)</p>				
3/17	SUBMITTED TEXT	38 WORDS	56% MATCHING TEXT	38 WORDS
<p>$x^n y^n - xy = (x^n - x) y^n + x(y^n - y)$ gives $x y x y x x y x y y y x x x y y n$</p> <p>SA PARASAKTHI_MATHEMATICS_GAC_Krish.pdf (D110268816)</p>				
4/17	SUBMITTED TEXT	15 WORDS	90% MATCHING TEXT	15 WORDS
<p>$X \parallel \&gt; 1$, then $e-x$ is invertible and $(e - x)^{-1} = e x$</p> <p>SA Representations of Locally Compact Groups (Rautio)_PDFA.pdf (D9189229)</p>				
5/17	SUBMITTED TEXT	44 WORDS	68% MATCHING TEXT	44 WORDS
<p>$(e - x)(e + x + x^2 + \dots + x^n) = (e + x + x^2 + \dots + x^n)(e - x) = e - x$</p> <p>SA Representations of Locally Compact Groups (Rautio)_PDFA.pdf (D9189229)</p>				

15/17	SUBMITTED TEXT	50 WORDS	20% MATCHING TEXT	50 WORDS
<p> $\sup_n x_n$ and $\inf_n x_n$ are finite. If $\sup_n x_n < \infty$ then $\lim_{n \rightarrow \infty} x_n$ exists and $\lim_{n \rightarrow \infty} x_n = 0$. We always have $\lim_{n \rightarrow \infty} x_n = 0$ or, $\lim_{n \rightarrow \infty} x_n = b$ or, $\lim_{n \rightarrow \infty} x_n = 1$. This gives, $\lim_{n \rightarrow \infty} x_n = 0$. </p> <p>SA PARASAKTHI_MATHEMATICS_GAC Krish.pdf (D110268816)</p>				
16/17	SUBMITTED TEXT	18 WORDS	46% MATCHING TEXT	18 WORDS
<p> $\lim_{n \rightarrow \infty} x_n = 1$ (from *) and this implies $\lim_{n \rightarrow \infty} x_n = 1$ exists and $\lim_{n \rightarrow \infty} x_n = 1$. </p> <p>SA Fuctional Analysis.pdf (D142230889)</p>				
17/17	SUBMITTED TEXT	21 WORDS	71% MATCHING TEXT	21 WORDS
<p> X. (i) $\ x\ ^2 = \ x\ ^2$ for all $x \in X$ and (ii) $\ x\ = \ x\$ for all x. </p> <p>SA PARASAKTHI_MATHEMATICS_GAC Krish.pdf (D110268816)</p>				

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1 Unit-1 □ Contents : Topological Group : Definition : Examples ; Self- homeomorphisms ; Neighbour- hoods of Identity e, Closure of a Set, Separa- tion Axioms ; Separation Theorams and Consequences. 1.1 Let G (????????????????) be a Group and let G be also a Topological space. If not stated otherwise, group operation is taken as multiplication. Definition 1.1.1.

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K.Rajakumari,M.Phil Dissertation,Mathematics,2 ... (D61273517)

G is said to be a Topological Group if mappings (i) $G \times G \rightarrow G$

given by $(x, y) \mapsto xy$ ($x + y$, in case Group operation is additive), $x, y \in G$ and (ii) $G \rightarrow G$ given by $x \mapsto x^{-1}$ (taking inverse) as $x \in G$ are both continuous. Explanation : The multiplication mapping (i) $(x, y) \mapsto xy$ in G and Inverse mapping (ii) $x \mapsto x^{-1}$ in G are continuous with respect to the given Topology in G and the induced product topology in $G \times G$. If $g_1 : (x, y) \mapsto xy$ in G ; and $g_2 : x \mapsto x^{-1}$ in G , by continuity of multiplication mapping g_1 we mean : Given any neighbourhood W of xy in G , there is a neighbourhood U of x , and there is a neighbourhood V of y in G such that $UV \subset W$. ($U + V \subset W$, in case Group composition is additive). Similarly, by continuity of Inverse mapping g_2 we mean : Given any neighbourhood W of x^{-1} in G ,

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there is a neighbourhood U of x such that $U^{-1} \subset W$ (-

$U \subset W$, in case G is additive).

2 For example, the set \mathbb{R} of all reals is an addition Group (with respect to arithmetic addition $+$ (additive inverse being $-$ ve sign) and \mathbb{R} is also a Topological space with respect to the usual metric topology whose basic open sets are open intervals like (a, b) where a, b are reals with $a < b$. Then \mathbb{R} is a Topological Group. Because, if $x, y \in \mathbb{R}$, and $W = ((x + y) - \epsilon, (x + y) + \epsilon)$, $\epsilon > 0$ is any neighbourhood of $x + y$ in \mathbb{R} , there is a neighbourhood U of x , say, $(x - \delta, x + \delta)$, $\delta > 0$ and there is a neighbourhood V of y , say, $(y - \epsilon, y + \epsilon)$, such that if $u \in U$ and $v \in V$, we have $u + v \in W$; So that $|u + v - (x + y)| = |(u - x) + (v - y)| \leq |u - x| + |v - y| < \delta + \epsilon < \epsilon$; i.e. $(u + v) \in ((x + y) - \epsilon, (x + y) + \epsilon)$. Similarly, if $x \in \mathbb{R}$ and $W = (-x - \epsilon, -x + \epsilon)$, $\epsilon > 0$ be any neighbourhood of $-x$ in \mathbb{R} , We find a neighbourhood $U = (x - \delta, x + \delta)$, $\delta > 0$ of x such that if $u \in U$ i.e. $u \in (x - \delta, x + \delta)$ i.e. $-u \in (-x - \delta, -x + \delta)$ i.e. $-u \in (-x - \epsilon, -x + \epsilon)$ i.e. $-u \in W$ i.e. $u \in U^{-1}$. Therefore we have checked that both group operations, namely addition and its inverse (subtraction) are continuous with respect to the concerned Topology in \mathbb{R} . There \mathbb{R} is a Topological Group. Example 1.1. Let $M_n(\mathbb{R})$ denote the collection of all square matrices with real entries (n is a + ve integer). Then we know that $M_n(\mathbb{R})$ forms an additive Group with respect to usual matrix addition wherein the null matrix becomes the Identity member of this Group. $M_n(\mathbb{R})$ is also a metric space with respect to a metric d given by $d(A, B) = \sum_{i,j} |a_{ij} - b_{ij}|$ and are any two members of $M_n(\mathbb{R})$.

3 Then $M_n(\mathbb{R})$ forms a Topological Group. It is a routine work to verify that $M_n(\mathbb{R})$ is a metric space with respect to the metric d as given above ; There we know that open balls constitute a base for the metric Topology and with respect to this metric topology it is now another exercise to check that group operations are rendered continuous here, and $M_n(\mathbb{R})$ is a Topological Group. Remark 1.1 We may take entries in matrix as complex scalars from \mathbb{C} , then similarly we get the collection $M_n(\mathbb{C})$ of $n \times n$ matrices with complex entries to form a Topological Group. Remark 1.2 Statements (i) and (ii) regarding continuity of mappings g_1 and g_2 may be coupled as under. Theorem 1.1.1 Continuities of g_1 and g_2 are equivalent to the following : For any $x, y \in G$ if W is any neighbourhood of xy^{-1} in G there is a neighbourhood U of x and there is a neighbourhood V of y such that $UV^{-1} \subset W$. Proof : Let us assume continuities of g_1 and g_2 . Take $x, y \in G$ and W any neighbourhood of xy^{-1} in Topological Group G . Then we find a neighbourhood U of x and H a neighbourhood of y^{-1} such that $UH \subset W$. (Applying continuity of g_1)(i) Since $y^{-1} = y^{-1} \cdot e = y^{-1}$; corresponding to neighbourhood H of y^{-1} continuity of g_2 gives a neighbourhood V of y such that $V^{-1} \subset H$ (2) Combining (1) and (2) we have $UV^{-1} \subset UH \subset W$, which was wanted. Conversely assume the opposite. That is, assume the continuity of $(x, y) \mapsto xy^{-1}$ in G . First we deduce that g_2 is continuous i.e. $x \mapsto x^{-1}$ is continuous in G where $x \in G$. Write $ey^{-1} = y^{-1}$ taking $x = e =$ the Identity element e of G . By assumed condition corresponding to a neighbourhood W of y^{-1} , there is a neighbourhood V of e and a neighbourhood U of y such that $VU^{-1} \subset W$. We have $e \in V$, and hence $U^{-1} = eU^{-1} \subset VU^{-1} \subset W$.

4 That means mapping g_2 of taking inverse is continuous. Now write and take W to be any neighbourhood of xy^{-1} ; by assumed condition we find a neighbourhood U of x and a neighbourhood H of y^{-1} in G respectively such that $UH \subset W$. Since H is a neighbourhood of y^{-1} in G , by established continuity of taking inverse (as done above), We find a neighbourhood V of y such that $V^{-1} \subset H$. This gives $V^{-1} \subset H$, and hence from above we deduce $UV^{-1} \subset UH \subset W$. Thus continuity of g_1 of Group composition is established. Corollary 1.1 Composition of any three members of G is a continuous operation. 1.2 If $x, y \in G$, $(x, y) \mapsto x^{-2}y$ is a continuous operation in Topological Group G . 1.3

If x_1, y_2, \dots, x_n are n elements of Topological Group G , and $\alpha_1, \alpha_2, \dots, \alpha_n$ are +ve indices. Then $(x_1, x_2, \dots, x_n) \mapsto x_1^{\alpha_1} x_2^{\alpha_2} \dots x_n^{\alpha_n}$ is a continuous operation in G . We have seen that if G is a Topological Group then G is a Group and it is a Topological space ; but converse is false. Following Example supports this contention. Example 1.1.2 Consider the additive Group \mathbb{R} of all reals and let \mathbb{R} be equipped with the upper limit Topology whose basic open sets look like all left open (and right closed) intervals $\{(a, b] : a, b \in \mathbb{R}; a < b\}$. This topology is strictly stronger than the usual topology of \mathbb{R} . We verify that taking inverse i.e. $x \mapsto -x$ ($x \in \mathbb{R}$) in \mathbb{R} is not a continuous operation. Take a neighbourhood like $(0, \epsilon]$ of 0 in \mathbb{R} with upper limit Topolog. Then there is no neighbourhood V of 0 in \mathbb{R} in this Topology such that $-V \subset (0, \epsilon]$. Therefore \mathbb{R} is not a Topological Group. Theorem 1.1.2 In a Topological Group G if $x_0 \in G$ is a fixed member, then (i) Mapping $\gamma : G \rightarrow G$ given by $x \mapsto x x_0^{-1}$ as $x \in G$, and (ii) Mapping $\delta : G \rightarrow G$ given by $x \mapsto x_0^{-1} x$ as $x \in G$ are homeomorphisms. Proof : (i) The mapping $\gamma : x \mapsto x x_0^{-1}$ as $x \in G$ is 1-1 ; Because let

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$x_1 x_0^{-1} = x_2 x_0^{-1}$ for $x_1, x_2 \in G$: Then $x_1 x_0^{-1} x_0 = x_2 x_0^{-1} x_0$ (by multiplying x_0 from right) $\implies x_1 = x_2$		

e ($e =$ the identity element of G).
 or, $x_1 = x_2$ Hence this mapping is 1-1 (one-one). This mapping is also onto. For any $u \in G$, then $u x_0^{-1} = v \in G$, such that under the mapping $V \mapsto v x_0^{-1} = u$. Thus this mapping is invertible. We now check that the mapping is continuous. Take W to be any neighbourhood of $x x_0^{-1}$; By continuity of Group composition we find

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a neighbourhood U of x and a neighbourhood V of x_0 such that		

UV τ W This gives $Ux \in \tau$ W since $x \in \tau$ V. So the mapping is continuous at $x \in G$. Now its inverse mapping is given by $x \mapsto x^{-1}$ as $x \in G$. Which is essentially of the same type as given one, and hence becomes continuous. Therefore the concerned mapping is bi-continuous, and it is a Homeomorphism. By a similar argument the mapping under (ii) is shown to be a Homeomorphism— and it is a self homeomorphism like (i). Corollary 1.1 Let P be an open set in Topological Group G and Let A be any subset of G , then (i) Pu, uP are open sets in G for any member $u \in G$ (ii) PA and AP are open sets in G . Because (i) the mapping $T_u : G \rightarrow G$ given by $T_u(x) = xu$ for $x \in G$ is a homeomorphism, and further $T_u^{-1} = T_{u^{-1}}$ and by continuity of T_u^{-1} we find $T_u^{-1}(P) = \tau$ W ($P =$ an open set) = an open set in G i.e. $T_u(P) =$ an open set in G i.e. $Pu =$ an open set in G . Similarly employing other multiplying operator we have uP as an open set in G . (ii) Writing $PA = \bigcup_a Pa$ a union of some open sets in $G =$ an open set in G ; and similarly, AP is again an open set in G .

6 Corollary 1.2 Let Q be a closed set in G and $u \in G$, then Qu and uQ are closed sets. Corollary 1.3 if $u, v \in G$, then is a self homeomorphism τ of G such that $\tau(u) = v$. Here put $a = u^{-1}v$; There $a \in G$ and Look at $\tau : G \rightarrow G$ given by $\tau(x) = xa$ as $x \in G$. Then τ is a self homeomorphism of G such that $\tau(u) = ua = uu^{-1}v = v$. Corollary 1.4. In a Topological Group G if $a \in G$, then mapping $\tau : G \rightarrow G$ given by $x \mapsto xax^{-1}$ as $x \in G$ is a self homeomorphism, called an inner automorphism of G . Because Given mapping $\tau : G \rightarrow G$ defined by $x \mapsto xax^{-1}$ as $x \in G$ is a composite mapping out of two self homeomorphisms : $x \mapsto xa^{-1}$ and $x \mapsto ax$ as $x \in G$, and therefor is again a self homeomorphism. Theorem 1.1.3 In a Topological Group G the inverse mapping $f : x \mapsto x^{-1}$ as $x \in G$ is a self-homeomorphism. Proof : This mapping f is 1-1 and onto : and it is continuous ; Further, its inverse f^{-1} is given by $f^{-1} = f$ (i.e. f is self-inverse) and hence is continuous ; So f is a bicontinuous bijective mapping making it a self-homeomorphism of G . Corollary : If P is an open set in G , then P^{-1} is an open set in G ; because $f^{-1}(P) =$ an open set in G , by continuity of f . i.e. $f(P) =$ an open set in G , because $f^{-1} = f$. i.e. $P^{-1} =$ an open set in G . Remarks : We have seen that in a Topological Group G products (Addition) PQ and QP of any two sets P and Q are always open sets. There is a caution! Products of two closed sets may not be a closed set. This would be demonstrated later on. 1.2 Neighbourhood systems of Identity member e of a Topological Group G . Let τ denote the collection of all neighbourhoods of the identity element e of G . Definition 1.2.1 A Sub-collection B of τ is called a fundamental system of

7 neighbourhoods of e if for any member $N \in \tau$, there is a member $B \in B$ such that $B \subseteq N$. For examples, the sub-family comprising of all open intervals like $(-1/n, 1/n)$, e_j , $n = 1, 2, \dots$ constitutes a fundamental system of neighbourhood of $0 =$ the identity element of the additive Topological Group R of the reals with usual Topology. Before we proceed further we recall following Theorem. Theorem 1.2.1 If V is a neighbourhood of e , there is a symmetric neighbourhood U (i.e. $U = U^{-1}$) of e such that $U \subseteq V$. Proof : Put $U = V \cap V^{-1}$. So U is again a neighbourhood of e such that $U \subseteq V$. It remains to check that U is symmetric. Now there is an open set, say O is G with $O \subseteq V$, and therefore, $O^{-1} \subseteq V^{-1}$. Then $(O \cap O^{-1}) \subseteq (V \cap V^{-1})$. If $x \in U$, we have $x \in V$ and $x \in V^{-1}$ as well. Now $x^{-1} \in V^{-1}$ and $x \in V^{-1}$ implies $x^{-1} \in V$; therefore $x^{-1} \in (V \cap V^{-1}) = U$. What we have shown above is when $x \in U$, then $x^{-1} \in U$. Thus $U = U^{-1}$. Theorem 1.2.2 If

V is a

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neighbourhood of e in G , there is a neighbourhood W of e such that $W^2 \subseteq V$

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Proof : We have $e \in V$ and using continuity of group operation corresponding to a neighbourhood V if $e \in G$ we find neighbourhoods V_1 and V_2 of e such that $V_1 V_2 \subseteq V$. Put $U = V_1 \cap V_2$. Then U is a neighbourhood of e in G such that $U^2 = UU \subseteq V_1 V_2 \subseteq V$, and the proof is complete. Remarks 1.2. Without loss of generality one may take U to be symmetric. Remark 2.2. For any integer n there is a neighbourhood U of e such that $U^n \subseteq V$ in G by Induction. Corollary to Theorem 1.2.1 In a Topological Group G there is a fundamental system $\{U\}$ of symmetric neighbourhoods of e in G . 8 In view of Theorem 1.1.2 where it is revealed that translation like homeomorphisms are responsible to send

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a fundamental system of neighbourhoods of e in G to another fundamental system of neighbourhoods of

any other point of G . Following Theorem 1.2.3, let $\{U_\alpha\}$ be a fundamental system of open neighbourhoods of e in G . Then the family $\{xU_\alpha\}$ and $\{x^{-1}U_\alpha\}$ constitutes a base for the Topology on G . Proof: Suppose $a \in G$ and W is an open neighbourhood of a in G . Now the mapping $T_a : G \rightarrow G$ given by $x \mapsto ax$ is a self-homeomorphism of G , we have $T_a^{-1}(W) = a^{-1}W$ as an open set containing e ; it contains a member, say, U_α of the fundamental system of open neighbourhoods of e in G such that $U_\alpha \subset a^{-1}W$; or, $aU_\alpha \subset W$. That shows that $\{xU_\alpha\}$ and $\{x^{-1}U_\alpha\}$ forms a base for the Topology of G . Corollary: Under assumption of Theorem 1.2.3 the family $\{xU_\alpha\}$ and $\{x^{-1}U_\alpha\}$ forms a base for the Topology of G .

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Theorem 1.2.4 Let A be a subset of a Topological Group G . Then $\text{closure}(A) = \bigcap \{U_\alpha \mid A \cap U_\alpha \neq \emptyset\}$; where $\{U_\alpha\}$ denote the system of all

neighbourhoods of the identity e in G . Proof: Take $x \in A$ and $U_\alpha \ni e$; then xU_α^{-1} is a neighbourhood of x in G ; and therefore $(xU_\alpha^{-1}) \cap A \neq \emptyset$. That means $x \in \text{closure}(A)$. Since x is any member of A we have $A \subset \text{closure}(A)$. (1) Conversely, Take any $y \in \text{closure}(A)$ and $U_\alpha \ni e$ and so $yU_\alpha^{-1} \cap A \neq \emptyset$ for each $U_\alpha \ni e$. Then if P is any open neighbourhood of y , we have $P \cap y^{-1}U_\alpha^{-1}$ is a neighbourhood of e in G , and $y \in P \cap y^{-1}U_\alpha^{-1}$ because $P \cap y^{-1}U_\alpha^{-1} \neq \emptyset$. That means $y = ay^{-1}$ for some $a \in A$ and some $p \in P$. Now $y = ay^{-1}$ gives $yy^{-1} = ay^{-1}y$ or, $e = ay^{-1}y$. Hence $y \in A$ showing that $\text{closure}(A) \subset A$. (2) Combining (1) and (2) we have $\text{closure}(A) = A$. Remark: $A = \bigcap \{U_\alpha \mid A \cap U_\alpha \neq \emptyset\}$. The proof is a copy of that of Theorem 1.2.4. Corollary: The closed neighbourhoods of e form a fundamental system of neighbourhoods of e in G . Because given any

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neighbourhood U of the identity element e in G , there is a neighbourhood V of e such that

$V \cap V^{-1} = U$. (1) Now by Theorem 1.2.4 we have $V \cap V^{-1} = U$, and taking $U = V \cap V^{-1}$ we find $V \cap V^{-1} = U$ and from (1) it follows that $V \cap V^{-1} = U$; V being a closed neighbourhood of e in G , the conclusion stands ok. Theorem 1.2.5 In a Topological Group G there is a fundamental system $\{U_\alpha\}$ of closed neighbourhoods of the identity e such that (i) each member U_α is symmetric (ii) for each U_α in the system there is another member U_β satisfying $U_\beta \subset U_\alpha^{-1}U_\alpha$; and (iii) for each member U_α in the system, and the each $a \in G$ there is a member U_β in the system such that $U_\beta \subset a^{-1}U_\alpha a$ or $aU_\beta a^{-1} \subset U_\alpha$. Conversely, given a group G with a filter base $\{U_\alpha\}$ to satisfy (i) – (iii), then there is a unique Topology on G to make G a Topological Group where $\{U_\alpha\}$ forms a fundamental system of neighbourhoods of e in G . Proof: (i) and (ii) are consequences of Theorem 1.2.1 and 1.2.2. And corollary following Theorem 1.1.2 says that mapping $\phi : x \mapsto axa^{-1}$ is a self-homeomorphism in G , and $a^{-1}U_\alpha a$ becomes a neighbourhood of e and hence (iii) follows. Converse part: Let $\{U_\alpha\}$ be a filter base satisfying (i) – (iii). Take any member U_α in the family. By (i) and (ii) we find a member U_β of this family to satisfy $U_\beta \subset U_\alpha^{-1}U_\alpha$ (By symmetry, $U_\beta \subset U_\alpha^{-1}$) If $x \in U_\beta$, then the Identity element in $G = e = xx^{-1} \in U_\beta \subset U_\alpha^{-1}U_\alpha$. Therefore each member U_α of the family contains e . And each member of the family $\{xU_\alpha\}$ and $\{U_\alpha x\}$ contains x for every $x \in G$. Further, $\{xU_\alpha\}$ and $\{U_\alpha x\}$ each forms a filter base at x because so is the family $\{U_\alpha\}$. We now construct a Topology on G . Let \mathcal{T} consist of \emptyset (empty set) and $\{xU_\alpha\}$ as $x \in X$. Since $xU_\alpha \subset X$ by filter property $X \cap \{xU_\alpha\} \neq \emptyset$. Thus $X \in \mathcal{T}$. Suppose U_1, U_2 are two members of \mathcal{T} , and $x \in (U_1 \cap U_2)$, then both U_1, U_2 are members of $\{xU_\alpha\}$ and Filter base property ($U_1 \cap U_2$) is a member of this family implying $(U_1 \cap U_2) \in \mathcal{T}$. Finally, $\{U_r\}_{r \in \mathbb{R}}$ be a family of members of \mathcal{T} . Say $x \in U_r$ for some r . They by choice for some of \mathcal{T} , and $U_r = (xU_\alpha)$. As, by filter-base property the Union is a member of $\{xU_\alpha\}$; That means the Union $\in \mathcal{T}$. Now equipped with this Topology \mathcal{T} , G is a Topological Group if continuity of Group operation $(x, y) \mapsto xy^{-1}; x, y \in G$ is verified with respect to the Topology \mathcal{T} and that we do presently as under: Take $x, y \in G$ and put $xp = u$ and $yq = v$ where $p, q \in G$. Now $(xy^{-1})(uv^{-1}) = yx^{-1}uv^{-1} = yx^{-1}x(p(yq)^{-1}) = ypq^{-1}y^{-1}$

11 Let N_e be a neighbourhood of e (relative to τ); so we find a member U in \mathcal{U}_e satisfying $U \cap N_e = \{e\}$. Now $U \cap N_e = \{e\}$ implies $U \cap N_e = \{e\}$ if $U \cap N_e = \{e\}$ (1) Using (iii) we find a member $U' \in \mathcal{U}_e$ satisfying $U' \cap N_e = \{e\}$. Again from (i) and (ii) there is a member $W \in \mathcal{U}_e$ such that $W \cap N_e = \{e\}$. So that $W \cap N_e = \{e\}$. So $W \cap N_e = \{e\}$. Let $p, q \in W$. Then we have $pq^{-1} \in W$ and $q^{-1} \in W$ (since W is symmetric) i.e. $pq^{-1} \in W$ and $q^{-1} \in W$. From (1) we conclude that $U \cap N_e = \{e\}$ or, $(xy^{-1}) \in U$ or, $(uv^{-1}) \in U$ whenever $p, q \in W$. That confirms that G is a Topological Group. The proof is complete. Example 1.2 Let E_1 and E_2 be compact subsets of a Topological Group G . Then $E_1 \cap E_2$ is compact. Consider the mapping $h : G \times G \rightarrow G$ where $h(x, y) = xy$ as $x, y \in G$. Since E_1 and E_2 are compact, the product sub-space $E_1 \times E_2$ is compact. The mapping h is a continuous mapping and since continuous image of a compact space is compact, $E_1 \cap E_2 = \text{image of } E_1 \times E_2 \text{ under } h$ becomes compact. 1.3 Separation Axioms : First the recall Definitions of separation axioms like T_0, T_1, T_2, \dots in a Topological space (X, τ) as under : Definition 1.3.1 (X, τ) is called a T_0 -space if given two distinct points in X , there is an open set containing any one without containing the other. For example, real number space \mathbb{R} with usual topology is a T_0 -space; because if $x, y \in \mathbb{R}$ and $x < y$, there is an open interval containing x keeping y outside. On the other hand there are topological spaces that are not T_0 . Example 1.3.1 Let $X = \{a, b, c\}$ and let τ be a family of subset of X consisting of $\emptyset, X, \{a\}$ and $\{b, c\}$. Then (X, τ) is a Topological space which is not T_0 ; because distinct elements b and c in X have no T_0 -separation. Definition 1.3.2 (X, τ) is called a T_1 -space if given any two distinct elements in X , there is an open set to contain each one of them without containing the other. Explanation : A very common example of a T_1 -space is real number space \mathbb{R} with usual topology. On the other hand if $X = \{a, b, c\}$ where a, b, c are all distinct, and if $\tau = \{\emptyset, X, \{a, b\}\}$, Then (X, τ) is Topological space where T_1 -stipulation is missing. Because pair (a, b) of distinct elements in X has no attracting open sets as demanded by T_1 -condition. Thus (X, τ) is not T_1 . Remark : Definitions 1.3.1. and 1.3.2 are so framed that a T_1 -space is always T_0 ; but opposite implication is, however, false. For example, taking $X = \{a, b\}$, a $\tau = \{\emptyset, X, \{a\}\}$ is a T_0 -space without being T_1 . Because only open set to take b inside is $\{a, b\}$ that does not leave a . Definition 1.3.3 A topological space (X, τ) is called a T_2 -space or a Hausdorff space if given any two distinct members x and y in X , there are open sets U and V in X such that $x \in U$ and $y \in V$ with $U \cap V = \emptyset$. As per Definitions we at once see that $T_2 \implies T_1$ Example 1.3.2 Let $X = \{a, b\}$, a $\tau = \{\emptyset, X, \{a, b\}\}$. Then (X, τ) is a Topological space where there is T_2 separation. And there are topological spaces that are T_1 without being T_2 . Example 1.3.3 shall bear it out. Example 1.3.3 (Cofinite Topology) : Let X be an infinite set and Let $\tau = \{G \subseteq X : (X \setminus G) \text{ is a finite set (may be empty)}\} \cup \{\emptyset\}$; Then τ becomes a Topology in X ; very often this Topology is named as Co-finite Topology in X . This Topological space (X, τ) is T_1 without being T_2 . Take two members x, y in X without $x = y$; Put $U = X \setminus \{y\}$ and $V = X \setminus \{x\}$. Then U and V are members of τ such that U contains x leaving y outside and V contains y leaving x outside. Therefore (X, τ) is T_1 . If possible, let any two distinct elements u and v in X have T_2 separation. Then there are two open sets, say, H and K in X such that $u \in H, v \in K$ with $H \cap K = \emptyset$. So $(X \setminus H) \cap (X \setminus K) = X \setminus (H \cap K) = X$; since $H \cap K = \emptyset$. — a contradiction; because X is not a finite set. Thus (X, τ) is not T_2 . Here we quote some important Theorems whose proofs may be found in any text of General Topology. Theorem 1.3.1 If (X, τ) is T_0 , then closures of distinct points in X are distinct. Theorem 1.3.2 (X, τ) is T_1 if and only if each singleton in X is closed. Theorem 1.3.3 (X, τ) is T_2 (

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Hausdorff) if and only if every net in X converges to atmost one point in

X. Theorem 1.3.4 A product of T_2 -spaces is a T_2 -

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space. Definition 1.3.4.(a) (X, τ) is called a regular space if given any closed set F

in X , and an outside point $x \in X \setminus F$ there open sets U and V in X such that $x \in U$ and $F \subset V$ with $U \cap V = \emptyset$. (b) A regular space that is also a T_1 -space is called a T_3 -space. Explanation : If $X = \{x, y, z\}$, and $\tau = \{\emptyset, X, \{x\}, \{y, z\}\}$, Then (X, τ) is a Topological space whose only closed sets are $X, \emptyset, \{y, z\}$ and $\{x\}$. We easily check that (X, τ) is a regular space ; (X, τ) is not T_1 -space; because singleton $\{z\}$ is not a closed set in X . Further we have $T_3 \iff T_2 \iff T_1 \iff T_0$. Definition 1.3.5 (a) (X, τ) is called a Normal space if given any pair of disjoint closed sets F and G is

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X , there are disjoint open sets, U and V satisfying $F \subset U$ and $G \subset V$. (b) A normal space that is

also a T_1 is called a T_4 -space. Example 1.3.4 Take $X = \{$

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$a, b, c, d, e, f\}$ and $\tau = \{\emptyset, X, \{e\}, \{f\}, \{e, f\}, \{a, b, c\}, \{c, d, f\}, \{a, b, e, f\}, \{c, d, e, f\}\}$.

Then we can verify that (X, τ) is a Normal space where we find four pairs of disjoint non-empty closed sets only : $\{\{$

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$a, b\}, \{c, d\}, \{a, b\}, \{c, d, f\}, \{a, b, e\}, \{c, d\}\}$ and $\{a, b, e\}, \{c, d,$

$f\}$. Here each pair is separated by disjoint pair of open sets $\{a, b, e\}, \{c, d, f\}$. Here we observe that this Normal space is not regular; because $\{a, b\}$ is a closed in X with an outside element $e \notin \{a, b\}$; and there is no disjoint pair of open sets in X to separate them. Further we note that $T_4 \iff T_3$; because if F is a closed set in a T_4 -space X with $x \in X \setminus F$ as an outside point in X ; Then singleton $\{x\}$ is a closed set ; So normality in X attracts desired separation. So X is T_3 . Definition 1.3.6(a)

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A topological space (X, τ) is called completely regular if given any closed set F and an outside point $x \in X \setminus F$ there is a

continuous function $f : X \rightarrow [0, 1]$ (

closed unit interval of reals) such that $f(x) = 0$ and $f(u) = 1$ for $u \in F$. (b) A completely regular space which is also T_1 is called a Tychonoff space, often designated as -

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space. Theorem 1.3.4 A topological space (X, τ) is a Normal space if and only if

given any pair of sets (F, H) where F is closed and H is open with $F \cap H = \emptyset$, there is another open Set G in X such that $F \subset G \subset \overline{G} \subset H$, \overline{G} denoting the closure. Proof : The condition is necessary : Let (X, τ) be a Normal space where (F, H) is a pair of closed and open sets such that $F \cap H = \emptyset$ ($F =$ a closed set ; $H =$ an open set). The complement of $H = H^c$ is a closed set in X with $F \cap H^c = \emptyset$. By normality in X we find a pair of disjoint open sets, say, G and M satisfying. $F \subset G$ and $H^c \subset M$ with $G \cap M = \emptyset$. Thus $G \cap M^c$ and $H^c \cap M$ gives $M^c \cap (H^c)^c = H$ As M^c is a closed set, we obtain $F \subset G \subset M^c \subset H$ That is, $F \subset G \subset H$. The condition is sufficient : Let the condition hold in (X, τ) . Suppose F_1 and F_2 are a pair of disjoint closed sets in X .

15 Then we have $F \cap C^c$ (complement of $F \cap C$), which is open. Hence by assumed condition we find an open set G in X such that $F \subset G \subset F \cap C^c$. Now $G \cap C^c$ gives $F \subset G \subset C^c$, and of course, $G \cap C = \emptyset$. So $G \cap C = \emptyset$. Thus, $F \subset G$ and $F \cap C = \emptyset$ where G and $G \cap C$ form a pair of disjoint open sets to bring the desired separation. Hence (X, τ) is Normal. Theorem 1.3.5 (Separation Theorem in Topological Group G) In a Topological Group G let F be a closed set and C a compact set such that $F \cap C = \emptyset$. Then there is a neighbourhood W of the identity e in G such that (i) $F \cap CW = \emptyset$ (ii) $WF \cap WC = \emptyset$. Proof : To establish (i) it suffices to look for a neighbourhood U of the identity in G such that $(FU \cup -1) \cap C = \emptyset$. If U is a neighbourhood of e , put $F \cup U = \bar{FU}$, bar denoting the closure. So $F \cup U$ is closed and we have $F \cup FU \cup U \cup V = e \cap \tau$, denoting the neighbourhood system at e . $W = U \cup -1 \cup V = F$ (closure of F) = F , because F is closed ; and $F \cup U$ is closed. Thus as per assumption, $F \cap C = \emptyset$ we have $F \cup U \cap C = \emptyset$. This is true for all open neighbourhood U of e . Therefore the family $\{G \setminus F \cup U\}$ is

16 an open cover for C . By compactness of C there is a finite sub-family, say,

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$F \cup U_1, F \cup U_2, \dots, F \cup U_n$ such that $G \setminus F \cup U_1, G \setminus F \cup U_2, \dots, G \setminus F \cup U_n$			

forms an open cover for C . Therefore, $F \cap C = \emptyset$. (1) Put $W = U \cup -1$. Then W is neighbourhood of e in G . Now $W \cap -1 = U \cup U \cap -1 = U \cup -1$ and hence $F \cap W \cap -1 = F \cap U \cup -1 = F \cap U \cup -1$. So taking closure $F \cap W \cup -1 = F \cap U \cup -1$. That means, $F \cap W \cap -1 = F \cap W \cup -1$. From (1) it is clear that $F \cap W \cap -1 \cap C = \emptyset$. Therefore $F \cap CW = \emptyset$. This is exactly what has been wanted in (i). Similarly, one can establish (ii) i.e., $V \cap VC = \emptyset$ for some $V \in \tau$. Remark : If one takes $U = W \cap V$, this neighbourhood U of e works in (i) as well as in (ii). Theorem 1.3.6 Let F be a closed set and C a compact set in a Topological Group G . Then FC (CF) is closed. Proof : Take $x \in G \setminus FC$; So, $(F \cap -1) \cap C = \emptyset$. $F \cap -1$ is closed ($F \cap -1$ is homeomorphic image of F under homeomorphism : $u \mapsto u^{-1}$ in G ; and therefore $F \cap -1$ is a closed set in G . Thus $F \cap -1$ is closed and

17 C is Compact in G and we apply Theorem 1.3.5 (separation Theorem) to obtain a neighbourhood U of the identity e in G such that $(F \cap -1 \cup U) \cap C = \emptyset$. That means $(x \cup U \cup -1) \cap C = \emptyset$. Now $x \cup U \cup -1$ is a neighbourhood of x because $U \cup -1$ is a neighbourhood of e in G . And as x is any member outside FC , it follows that FC is closed. Similarly we show that CF is closed, and Theorem is proved. Remarks 1.3.1 Under hypothesis of the Theorem 1.3.5 $F \cap CW = \emptyset$, bar denoting the closure. Because, if $p \in F \cap CW$; p becomes a limit point of $F \cap CW$ and therefore any neighbourhood of p shall meet $F \cap CW$. Without loss of generality taking W to be open we find CW to be an open set with p as an inside point and therefore CW acts as a neighbourhood of p . That calls for $F \cap CW = \emptyset$ — a contradiction. Therefore $F \cap CW = \emptyset$. Theorem 1.3.7 In

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a Topological Group G following statements are equivalent. (i) G is a T_0 -space (ii) G is a T_1 -space. (iii) G is a T_2 -space or a Hausdorff space. (iv) $\tau \cup \mathcal{F} = \tau$ where $\mathcal{F} = \{U \cup -1 \mid U \in \tau\}$ denoting a fundamental system of neighbourhood of e .			

Proof : Suppose statement (i) is true. Take $x, y \in G$ with $x \neq y$. Because of T_0 -separation in G , say, x has an open neighbourhood N_x such that $y \notin N_x$. Now $x \cup -1 \cap N_x = V$ (say) is an open neighbourhood of identity e in G . Therefore

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$V \cup -1 = W$ (say) is an open symmetric neighbourhood of e ; and hence $y \in W$ is neighbourhood of y .			

We claim that $x \in y \cup W$. Otherwise, $x \cup -1 \cap W \cap -1 = W \cap -1$ (W symmetric) $\cap y \cup -1 \cap x \cup -1$. So $e = x \cup -1 \cap x \cup -1 \cap N_x \cup -1 = N_x \cup -1$

18 giving $y \in N_x$ which is not the case. Therefore $x \notin yW$. Thus T_1 , separation holds in G . So statement (ii) stands OK. Now we check that (ii) \implies (iii). Suppose $x, y \in G$ with $x \neq y$. Since T_1 separation holds in G we know that each singleton is closed; Thus $\{x\}$ is closed. Put $P = G \setminus \{x\}$. Then P

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is an open neighbourhood of y and therefore $y^{-1}P$ is an open neighbourhood of the identity e in G . Choose an open neighbourhood V of e such that $VV^{-1} \subseteq y^{-1}P$. Thus yV is an open neighbourhood of y . Put $Q = G \setminus yV$;

So Q is open set. Here $x \notin Q$; otherwise, $x \in yV \implies x \in yV \cap y^{-1}P$. That means $x \in yV \cap y^{-1}P = P \cap y^{-1}P = P - a$ contradiction. Further, $Q \cap yV = \emptyset$ and $x \in Q$ where yV and Q are open sets. Hence T_2 -separation is established i.e. statement (iii) is true. Now let statement (iii) be true. We show that statement (iv) remains true. Suppose \mathcal{F} denote a fundamental system of neighbourhood of e in G . Let $x \in U \in \mathcal{F}$. Assume that $x \neq e$. Then by T_2 separation property, we find

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a neighbourhood P of e such that $x \notin P$. Let $U \in \mathcal{F}$ such that $U \subseteq P$. Then $x \in U$ (because $x \in U \in \mathcal{F}$)—a contradiction that $x \notin P$. Hence

we have shown that $x = e$ and (iv) is established. Finally the proof shall be completed by showing that statement (iv) \implies (i). Take $x, y \in G$ with $x \neq y$; Then $xy^{-1} \neq e$, and therefore from (iv) we find a member $U \in \mathcal{F}$ such that $xy^{-1} \notin U$; Now Uy is a neighbourhood of y such that $x \notin Uy$ —

19 confirming T_0 -separation in G . Thus statement (i) holds. The cycle of implication being complete, we have proved Theorem. Example 1.3.5 Let E be a compact set and O an open set in a Topological Group G . If $E \cap O \neq \emptyset$, show that there is a neighbourhood V of the identity e in G such that $VE \subseteq O$. Solution: Take $x \in E \cap O$; write $x = ex$ and using continuity of group operation find

a neighbourhood V_x of the identity e in G , such that $V_x x \subseteq O$ ($O =$ open set containing x). Find

an open neighbourhood W_x of e such that $W_x V_x x \subseteq O$. So one writes $E = \bigcup_{x \in E} W_x x$

$E = \bigcup_{x \in E} W_x x$ i.e. $\{W_x x\}$ is an open cover of E which is compact in G . So we pick up a finite number of members like $W_{x_1} x_1, W_{x_2} x_2, \dots, W_{x_n} x_n$ such that $E \subseteq \bigcup_{i=1}^n W_{x_i} x_i$. Construct a neighbourhood V of the identity e where $V W_{x_i} x_i \subseteq O$ for $i = 1, 2, \dots, n$. It is now clear that $V W_{x_i} x_i \subseteq O$ for $i = 1, 2, \dots, n$; that means $V \bigcup_{i=1}^n W_{x_i} x_i \subseteq O$; This gives $Vx \subseteq O$ and this is true for any $x \in E$, and hence $VE \subseteq O$. Theorem 1.3.8 A Topological Group that is Hausdorff (T_2) is completely regular. Proof: Let G be a Topological Group which is Hausdorff. Let \mathcal{F}_e denote

a fundamental system of neighbourhoods of the identity e in G satisfying (i) each member of \mathcal{F}_e is symmetric (ii) for each member $U \in \mathcal{F}_e$ there is member $V \in \mathcal{F}_e$ such that $V^2 \subseteq U$ and (iii) for each member $U \in \mathcal{F}_e$ and $a \in G$, there is a member $V \in \mathcal{F}_e$ to satisfy $V \subseteq a^{-1}Ua$ or $aVa^{-1} \subseteq U$.

20 Take C be a closed subset of G such that $e \notin C$. Put $U_0 = G \setminus C$. Then U_0 is an open neighbourhood of e in G . For each natural number n there is a member $U_n \in \mathcal{F}_e$ such that (i) $U_n \subseteq U_{n-1}$ and (ii) $U_n^2 \subseteq U_{n-1}$. If $D =$ set of all dyadic rationals of form $\frac{k}{2^n} = \frac{K}{2^n}$, $K \in \mathbb{N}$, $0 \leq K < 2^n$, then for each $\frac{k}{2^n} \in D$, by Induction, let us define (ii) $V = U_n$, $n \geq 0$. Suppose $V(\frac{k}{2^n})$ has been defined for all $\frac{k}{2^n} \in D$, then define (iii) if K is even, $V(\frac{K}{2^n}) = V(\frac{K/2}{2^{n-1}})$ and (iv) if K is odd, $V(\frac{K}{2^n}) = V(\frac{K-1}{2^n}) \cup V(\frac{K+1}{2^n})$. If $0 < \frac{k}{2^n} = \frac{2m}{2^{n+1}} < \frac{1}{2}$ we have $V(\frac{k}{2^n}) = V(\frac{2m}{2^{n+1}}) = V(\frac{m}{2^n})$ by (iii) since $e \in U_n$ by (i) by (iv) Therefore, (v) for all $0 < \frac{k}{2^n} < \frac{1}{2}$, $K = 2m$. Similarly, one can prove (v) when $K = 2m + 1$. So, (v) is true for all integers K such that $0 < \frac{k}{2^n} < \frac{1}{2}$. We now check that for $\frac{1}{2} < \frac{k}{2^n} < 1$ we have $V(\frac{k}{2^n}) = V(\frac{2k-2m}{2^n})$

21 Suppose $\frac{k}{2^n} = \frac{2m+1}{2^n}$ and $\frac{k}{2^n} = \frac{2m+2}{2^n}$. Then $V(\frac{k}{2^n}) = V(\frac{2m+1}{2^n}) \cup V(\frac{2m+2}{2^n})$ and hence $V(\frac{k}{2^n}) = V(\frac{2m+1}{2^n}) \cup V(\frac{2m+2}{2^n})$. Clearly, if $m+1 \leq 2^n$ then $V(\frac{2m+1}{2^n}) = V(\frac{m}{2^{n-1}})$ by (v). And we have $V(\frac{k}{2^n}) = V(\frac{m}{2^{n-1}}) \cup V(\frac{2m+2}{2^n})$

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$b g b g b g \dots$ in p steps where $K p K n n 1 2 2 2 1 ? ? ?$. But $? 1 = K n n n 1 2 2 2 1 2 ?$ and $? 2 = K n n n 2 2 2 1 1 2 ?$, we see that $V(? 1) ??V (? 2)$. Now we define f over G as under : $f x D x V V r V x V r ? ? ? ? ? ? R S T$ Inf if $? ? ? ; () n s 1 1$ Since $e?V ? ?$ for all $? ?D$ and $\text{Inf } D = 0$ we see $f(e) = 0$. Further more $V 1 = V 1 2 ? e j = U 0 = G \setminus C$, and, hence $f(C) = 1$. By definition of f we have $0 ??f(x) ??1$ for all $x?G$. We know show that f is continuous. Take $x?G$, such that $f(x) = 1$. If $y?V 1 2 n e j x$ then $y?G \setminus (K/2 n)$, $K \> 2 n - 2$. Otherwise, $y?V K n 2 e j x$ and symmetry of V 's shows that $x?V K n 2 e j y ??V 1 2 n e j V K n 2 e j ??V K n ? 1 2 e j$ by (v). Hence $f(x) \> 1$, contradicting assumption that $f(x) = 1$. Thus it follows that $1 1 2 2 1 2 ? ? ? n n n ??f(y) ??1$. Hence $| f(y) - f(x) | ??? 1 2 1 n ? ?$. If for a given $??\epsilon < 0$, appropriately large n satisfies $1 2 1 n ? \> ?$. Continuity of f at x follows. It is more easy to establish continuity of f when $f(x) = 0$. Now let $0 \> f(x) \> 1$ for some $x?G$. Then there are integers m, K with $K \> 2 m, m \< n + 1$ such that $x?V (K/2 m) \setminus K m ? 1 2 e j 2 2$ because $f(x) = \text{Inf } \{??D : x?V ?\}$ and D is dense in $[0, 1]$. Using (v) as before, for each $y?V x, y?V$. But $x?V$ implies $y?V$ by (v). Hence by Definition of $f, ??f(y) ??$. Since $(K - 1)/2 m ??f(x) ?$, We have $| f(x) - f(y) | ??$ Hence employing same argument as above f is shown to be continuous in all cases that arise. As we know translations have homeomorphism effect, above construction may be carried out at any point $x?G$ instead of the identity e in G . The proof of Theorem is now complete. Example 1.3.6 In a Topological Group G if U is

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any neighbourhood of the identity e in G and F any compact subset of G . Then there is a neighbourhood V of e such that $xVx^{-1} ??U$

for all $x?F$. Solution : Let $S e$ denote family of all symmetric neighbourhoods of e in G . First we check that for a fixed y in G , there is a member $V?S e$ such that $x?Vy$ implies $xVx^{-1} ??U$ Take a member $V 1 ?S e$ such that $V 1 3 ??U$ and take a member $V 2 ?S e$ such that $yV 2 y^{-1} ??V 1$. (see Theorem 1.2.5) Put $V = V 1 ??V 2$. Let $x?Vy$, i.e. $xy^{-1} ??V ??V 1$ and hence $yx^{-1} ??V 1^{-1} = V 1$ ($V 1$ symmetric) Hence $xVx^{-1} ??xV 2 x^{-1} = xy^{-1} y V 2 y^{-1} yx^{-1} ??V 1 3$, (because $xy^{-1} ??V 1, yx^{-1} ??V 1$ and $yV 2 y^{-1} ??V 1$ see above) ??U (see above) Therefore (1) holds. Now for each $y?F$, there is a $V y ??S e$ such that $x?V y y$ implies $xV y x^{-1} ??U$. Since F and F is compact, we find a finite number of members, say, $y 1, y 2, \dots, y n ?F$ such that $F V y 1 V y 2 V y n ? ? ? ? 1 2 1 2 \dots d i$ Put $V = \bigcap_{i=1}^n V y i$. If $x?F$, then $x?V y K y K$ for some $K (= 1, 2, \dots, n)$, and hence $xVx^{-1} ??xV y K x^{-1} ??U$.

23 EXERCISE A Short answer type questions 1. If $X = [0, 1)$ with a Topology $\tau = \{?, [0, ?) : 0 \> ?; ??\> 1\}$. Show that (X, τ) is not $T 1$. 2. Show that any sub-space of a Hausdorff space is Hausdorff. 3. Let G be an algebraic Group with discrete Topology. Examine if G is a Topological Group. 4. Show that an algebraic Group G with indiscrete Topology is a Topological Group, and examine if it is $T 0$. 5. Let G be an infinite Group with co-finite Topology. Examine if G is a Topological Group. 6. Show that every Topological vector space when treated as an additive group is a Topological Group. 7. Show that additive Group Z of all integers with usual Topology of reals is a discrete Topological group that satisfies second axiom of countability. 8. If R is the set of all reals, Show that $R \setminus \{0\}$ with arithmetic multiplication and with usual Topology of reals forms a multiplicative commutative Topological Group. EXERCISE B 1. Let X be a Hausdorff space and let C and D are disjoint compact sets in X . Show that there are open sets H and K in X such that $C \subset H$ and $D \subset K$ with $H \cap K = \emptyset$. 2. In a Topological Group G if $x?G$, and V is any

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neighbourhood of x , Show that there is a neighbourhood W of x such that $W \subset \bar{V}$, bar denoting the closure. 3. If a Topological Group

G is T_1 show that G is Hausdorff. 4. Let \mathcal{N}_e be the system of all neighbourhoods of the identity e of a Topological Group G , show that for any subset A of G , $\overline{A} = \bigcap \{N \in \mathcal{N}_e : A \cap N \neq \emptyset\}$. 5. If \mathbb{R} is the set of all reals, show that $\mathbb{R} \setminus \{0\}$ with arithmetic multiplication as Group composition and with usual Topology of reals forms a multiplicative commutative Topological Group.

24.6. In a Topological Group G if A and B are closed subsets, show that AB need not be closed. (Solution : Consider the additive Group \mathbb{R} of reals equipped with usual Topology. Then \mathbb{R} is a Topological Group. Here the set \mathbb{Z} of all integers is a closed subset ; If α is any irrational number, then $\alpha\mathbb{Z}$ is a closed set. The set $\mathbb{Z} + \alpha\mathbb{Z}$ consisting of all numbers $m + n\alpha$ where m and n are integers is not closed. This set is a dense subset of \mathbb{R} .) 7. Let A and B be subsets of Topological Group G . Then show that (a) $\overline{AB} \subseteq \overline{A} \overline{B}$, $\overline{gB} \subseteq \overline{g} \overline{B}$, $\overline{Bg} \subseteq \overline{B} \overline{g}$, $\bar{}$ denoting the closure, (b) $\overline{A} \overline{B} \subseteq \overline{AB}$, (c) $x \overline{A} y \subseteq \overline{xAy}$, for all $x, y \in G$, ..

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	G is said to be a Topological Group if mappings (i) $G \times G \rightarrow G$			
SA	K.Rajakumari,M.Phil Dissertation,Mathematics,2019.pdf (D61273517)			
2/23	SUBMITTED TEXT	11 WORDS	100% MATCHING TEXT	11 WORDS
	there is a neighbourhood U of x such that $U^{-1} \subseteq W$ (W is a neighbourhood of e)			
SA	MS - 334.docx (D110841764)			
3/23	SUBMITTED TEXT	53 WORDS	47% MATCHING TEXT	53 WORDS
	$x^1 \times x^0 = x^2 \times x^0$ for $x \in G$, $x^2 \in G$: Then $x^1 \times x^0 \times x^0^{-1} = x^2 \times x^0 \times x^0^{-1}$ (by multiplying x^0^{-1} from right) 5 or, $x^1 e = x^2$			
W	https://socratic.org/questions/how-do-you-use-the-binomial-series-to-expand-1-x-1-2-2			
4/23	SUBMITTED TEXT	14 WORDS	88% MATCHING TEXT	14 WORDS
	a neighbourhood U of x and a neighbourhood V of x^{-1} such that $UV \subseteq W$ (W is a neighbourhood of e)			
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5/23	SUBMITTED TEXT	17 WORDS	73% MATCHING TEXT	17 WORDS
<p>neighbourhood of e in G, there is a neighbourhood W of e such that $W \cap W = \{e\}$</p> <p>SA MS - 334.docx (D110841764)</p>				
6/23	SUBMITTED TEXT	16 WORDS	65% MATCHING TEXT	16 WORDS
<p>a fundamental system of neighbourhoods of e in G to another fundamental system of neighbourhoods of e in G</p> <p>SA MS - 334.docx (D110841764)</p>				
7/23	SUBMITTED TEXT	29 WORDS	45% MATCHING TEXT	29 WORDS
<p>Theorem 1.2.4 Let A be a subset of a Topological Group G. Then $(\text{closure of } A) = \overline{A}$; where \mathcal{B} denote the system of all</p> <p>SA 120004039-Project-1982444.pdf (D19454576)</p>				
8/23	SUBMITTED TEXT	17 WORDS	63% MATCHING TEXT	17 WORDS
<p>neighbourhood U of the identity element e in G, there is a neighbourhood V of e such that $V \cap V = \{e\}$</p> <p>SA Selvi C Chapter3.docx (D35106187)</p>				
9/23	SUBMITTED TEXT	15 WORDS	66% MATCHING TEXT	15 WORDS
<p>Hausdorff) if and only if every net in X converges to atmost one point in X</p> <p>SA suriyaprakasam REG.NO P17CAK8118.pdf (D58411288)</p>				
10/23	SUBMITTED TEXT	16 WORDS	70% MATCHING TEXT	16 WORDS
<p>space. Definition 1.3.4.(a) (X, \mathcal{B}) is called a regular space if given any closed set F and any point $x \notin F$, there exist disjoint open sets U and V such that $x \in U$ and $F \subset V$</p> <p>SA 120004039-Project-1982444.pdf (D19454576)</p>				

11/23	SUBMITTED TEXT	20 WORDS	47% MATCHING TEXT	20 WORDS
<p>X, there are disjoint open sets, U and V satisfying $F \cap U$ and $G \cap V$. (b) A normal space that is</p> <p>SA Plag_Rama pathak_33.pdf (D15260422)</p>				
12/23	SUBMITTED TEXT	28 WORDS	40% MATCHING TEXT	28 WORDS
<p>a, b, c, d, e, f) and $\{?, X, (e), (f), (e, f), (a, b, c), (c, d, f), (a, b, e, f), (c, d, e, f)\}$.</p> <p>a.1, b.1, c.1, d.1, e.1 3.7.11 b.1 3.6.6a a.1, c.1, d.1, e.1 3.8.1 a.3, b.4, c.2, d.1, d.2, e.1, f.1, g.1 3.6.6b a.1, b.1, c.1, d.1, d.2, e.1, f.1 3.8.4</p> <p>W https://www.nrc.gov/docs/ML1034/ML103470148.pdf</p>				
13/23	SUBMITTED TEXT	21 WORDS	52% MATCHING TEXT	21 WORDS
<p>a, b), (c, d)}, $\{(a, b), (c, 14 d, f)\}$, $\{(a, b, e), (c, d)\}$ and $\{(a, b, e), (c, d)$,</p> <p>a, b), (c, -d:)}$\{(1, 0), (0, 1):)$ then find a, b,c and d. If $\{(1, 0),(0, 1):)\{(a, b+c),(b-c, d):) = \{(4, -5),(3, 2):)$, then $(a-b) + (c-$</p> <p>W https://www.doubtnut.com/question-answer/simplify-1-1-1-1-1-1-a-0-b-1-c-2-d-3-3639419</p>				
14/23	SUBMITTED TEXT	25 WORDS	54% MATCHING TEXT	25 WORDS
<p>A topological space $(X, ?)$ is called completely regular if given any closed set F and an outside point x (i.e., $x \notin F$) there is a</p> <p>SA 120004039-Project-1982444.pdf (D19454576)</p>				
15/23	SUBMITTED TEXT	16 WORDS	71% MATCHING TEXT	16 WORDS
<p>space. Theorem 1.3.4 A topological space (X, τ) is a Normal space if and only if</p> <p>SA Totally na-Feebly regular continuous Function and its various structure.doc (D22998329)</p>				
16/23	SUBMITTED TEXT	28 WORDS	55% MATCHING TEXT	28 WORDS
<p>$F \cup U_1, F \cup U_2, \dots, F \cup U_n$ such that $G \cap F \cup U_1, G \cap F \cup U_2, \dots, G \cap F \cup U_n$</p> <p>SA M Asha Merlin, Reg.No.182311720920010, Chapters 2-5..pdf (D113514283)</p>				

17/23	SUBMITTED TEXT	48 WORDS	53% MATCHING TEXT	48 WORDS
<p>a Topological Group G following statements are equivalent. (i) G is a T_0-space (ii) G is a T_1-space. (iii) G is a T_2-space or a Hausdorff space. (iv) \mathcal{U}_e $\neq \emptyset$ $U = \{e\}$, \mathcal{U}_e denoting a fundamental system of neighbourhood of e.</p> <p>SA MS - 334.docx (D110841764)</p>				
18/23	SUBMITTED TEXT	22 WORDS	55% MATCHING TEXT	22 WORDS
<p>$V^{-1} = W$ (say) is an open symmetric neighbourhood of e; and hence yW is neighbourhood of y.</p> <p>SA SITHEESWARI (16PMAVO31).docx (D38133619)</p>				
19/23	SUBMITTED TEXT	49 WORDS	69% MATCHING TEXT	49 WORDS
<p>is an open neighbourhood of y and therefore $y^{-1}P$ is an open neighbourhood of the identity e in G. Choose an open neighbourhood V of e such that $VV^{-1} \subset y^{-1}P$. Thus yV is an open neighbourhood of y. Put $Q = yV \setminus \{y\}$;</p> <p>SA SITHEESWARI (16PMAVO31).docx (D38133619)</p>				
20/23	SUBMITTED TEXT	24 WORDS	36% MATCHING TEXT	24 WORDS
<p>a neighbourhood P of e such that $x \in P$. Let $U \subset F^2$ such that $U \cap P \neq \emptyset$. Then $x \in U$ (because $x \in U \cap F^2$)—a contradiction that $x \in P$. Hence</p> <p>SA SITHEESWARI (16PMAVO31).docx (D38133619)</p>				
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<p>V K V K n n n n n n n n 1 1 2 2 2 2 1 2 2 2 1 2 2 1 2 1 1 2 ? ? ? ? ? ? ?</p> <p>SA M Asha Merlin, Reg.No.182311720920010, Chapters 2-5..pdf (D113514283)</p>				

22/23	SUBMITTED TEXT	27 WORDS	55% MATCHING TEXT	27 WORDS
<p>any neighbourhood of the identity e in G and F any compact subset of G. Then there is a neighbourhood V of e such that $xVx^{-1} \cap U \neq \emptyset$</p> <p>SA MS - 334.docx (D110841764)</p>				

23/23	SUBMITTED TEXT	28 WORDS	43% MATCHING TEXT	28 WORDS
<p>neighbourhood of x, Show that there is a neighbourhood W of x such that $W \subseteq V \bar{V}$, bar denoting the closure. 3. If a Topological Group</p> <p>SA MS - 334.docx (D110841764)</p>				

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Vice-Chancellor

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25 Unit 2 □ -Sub-group, Normal Sub-group, Locally Compact Group, Topological Group Involving Connectedness, Locally Euclidean Group, Homomorphisms between Topological Groups, Lie Group. Structure 2.1 Introduction 2.2 Given A Topological Group G and Closed sub-group H in G 2.3 Locally compact Groups 2.4 Topological Groups Involving Connectedness 2.5 Linear Groups, Locally Eudidean Groups and lie Groups 2.6 Lie Groups 2.1 Introduction

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Let G be a Topological Group and H be a subgroup of G. Then H

inherits topology in G. Now Group operation : $(x, y) \mapsto xy^{-1}$ from $G \times G$ to G is continuous as $(x, y) \in G \times G$. Its restriction from $H \times H \subset (G \times G)$ to $H \subset G$ therefore remains continuous. Therefore H forms a Topological Group in its own right. H is called a Topological sub-group or simply a sub-group of G. There are always two sub-groups in a group ; namely G itself and singleton $\{e\}$ where e is the Identity member of G. These two sub-groups are so called trivial sub-groups. If H is a sub-group of G. Then we have $H^2 = H$ becomes coincident with H and we write $H^2 = H$ and similarly, $H^{-1} = H$. Let $a \in G$, it is a routine exercise to see that $a^{-1}Ha$ is also a sub-group of G. By chance, if a is also a member of H, then, ofcourse, $a^{-1}Ha = H$. In case $a \notin H$, then $a^{-1}Ha$ need not coincide with H.

26 Definition 2.1.1. For a sub-group H of G, if $a^{-1}Ha = H$ for every member $a \in G$, then H is said to be a normal sub-group or an invariant sub-group of G. Explanation : If Z is the additive group of all integers and is endowed with usual topology of reals. Then Z is a topological group of which $2Z$ forms a subgroup. It is a normal sub-group of Z. Trivially, the singleton $\{e\}$ of any topological group G whose identity equals to e forms a normal sub-group of G. In this connection following Theorem is an additional information. Theorem 2.1.1. If H is a sub-group of a topological group G, then its closure \bar{H} is so. Proof : A subset P of an algebraic group G is again a sub-group if $PP^{-1} \subset P$ i.e. $uv^{-1} \in P$ for all $u, v \in P$. In a topological group G we have seen that for any subsets A, B of G we have (i) $\bar{A} \bar{B} \subset \overline{AB}$, bar denoting the closure. (ii) $\overline{AB} \subset \bar{A} \bar{B}$ (iii) $\overline{xAx^{-1}} \subset \bar{A}$ for any $x \in G$; bar denoting closure. Here H is a given subgroup of G ; so $\bar{H} \bar{H}^{-1} = \bar{H}$ Now, $\bar{H} \bar{H} \subset \overline{HH} \subset \bar{H}$ from (i) $\bar{H} \bar{H} \subset \bar{H}$ and $\bar{H} \subset \bar{H}$ from (ii) = H, because H is a subgroup ; $\bar{H} \bar{H}^{-1} = \bar{H}$. This confirms that H is an algebraic subgroup of G ; Finally, continuity of group operation : $(x, y) \mapsto xy^{-1}$ in G works in respect of H to make H a Topological sub-group of G. Corollary : If H is normal sub-group of G, then \bar{H} is so. because if $x \in G$, we have $x \bar{H} x^{-1} = \overline{xHx^{-1}} \subset \bar{H}$ since $H = xHx^{-1}$, H is normal.

27 So H is a normal sub-group (algebraic) of G ; Also as above, group operation : $(x, y) \mapsto xy^{-1}$ in H is continuous. That makes H a topological normal sub Group. Remarks : In Topological Group G with identity e, the closure of $e = \{e\}$ is a closed normal sub-group of G and it is the smallest closed sub-group of G. Further, closure of a singleton $\{a\}$ ($a \in G$) i.e. $\bar{\{a\}} = \overline{a\{e\}}$. Theorem 2.1.2. (a) A sub-group H of a topological group G is open if and only if its interior $(\text{int } H) \neq \emptyset$ (b) Every open sub-group of G is closed. Proof : (a) Let $\text{Int } H \neq \emptyset$; and $x \in \text{Int } H$. Then there is an open neighbourhood U of the identity e of G such that $xU \subset H$. Now take any $y \in H$; we have $yU = yx^{-1}xU \subset yx^{-1}H$ (because $x \in H$). Since H is a sub-group and $x, y \in H$ we have $yx^{-1} \in H$; Therefore $yU \subset H$; So H is open, as every number of H is an interior point of H. Conversely, if H is open we have $\text{Int } H \neq \emptyset$. (b) Let H be an open sub-group of G ; then uH is an open set for every member $u \in G$. Now write $H = (G \setminus \{xH\})$ where $x \in G$ such that $\{xH\}$ is the family of all pairwise disjoint left cosets in G other than H. Clearly xH is an open set in G and hence H is its complement ; it follows that H is closed. Corollary : It is open and closed.

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a symmetric open neighbourhood of the identity e in topological group G and $L = \{x \in G : xU \subset H\}$, then L is an open and closed (clo-open) sub-group of G. Because we have

the following reasons. Take a and $y \in L$. Then let $x = y^k$ and $y = x^l$ for some indices k and l (say). Then $xy = y^{k+l}$ and $x^{-1} = (y^{-1})^k$ which is the same as y^{-k} because L is symmetric. That means L is a subgroup of G . We appeal to Theorem 2.12. to conclude that L is closed because L is open. Definition 2.1.2. Given a group G the set $C = \{x \in G : xa = ax \text{ for all } a \in G\}$ is called the centre of the Group G . Explanation : Centre C of the group G comprises of those member of G that commute with every element of G . Then C becomes a sub-group of G ; because let $p, q \in C$. So $pa = ap$ and $qa = aq$ for all $a \in G$. Now $(pq)a = p(qa) = p(aq)$ since $aq = qa = (pa)q$ by associativity $= (ap)q = a(pq)$ and this is o.k. for every member $a \in G$. Therefore $pq \in C$. Again for $a \in G$ we have $pa = ap$ So, $p^{-1}pap^{-1} = p^{-1}app^{-1}$ or, $ap^{-1} = p^{-1}a$ Thus p^{-1} commutes with every member $a \in G$ making sure that $p^{-1} \in C$. Hence C forms a sub-group of G . Finally, take any member $a \in G$, if $x \in C$ we have, ofcourse $ax = xa$ or, $axa^{-1} = x \in C$. That means, $aCa^{-1} \subset C$ and C is normal subgroup of G . Theorem 2.1.1. has corollary to tell us that its closure i.e. \bar{C} is a normal subgroup if C is the centre of a topological group G . Theorem 2.1.3. The centre C of a Hausdorff Topological Group G is a closed Normal sub-group. Proof : Now \bar{C} (= closure of C) is a normal subgroup of G . We now show that $\bar{C} \subset C$. Take $x \in \bar{C}$, let there be a member a in G such that $a^{-1}xa \neq x$. Since G is Hausdorff, and G is regular, Therefore we find open sets, U and V in G such that $x \in U$ and $(a^{-1}xa) \in V$ with $\bar{U} \cap \bar{V} = \emptyset$, $\bar{}$ denoting the closure. As $x \in \bar{C}$, it is easy to see that $x \in C$; So $(a^{-1}xa) \in C$ and $a^{-1}xa = a^{-1}xa$ because C is the centre of Group G —This is a contradiction and proof is complete. Example 2.1.1. In a Topological Group G if H is a sub-group of G such that \bar{H} is closed in G for some neighbourhood U of e in G , then H is closed.

29 Solution : Suppose U is a neighbourhood of the identity e of G such that $\bar{H} \cap U = H \cap U$ (bar denoting the closure) is closed. Take a symmetric neighbourhood V of e satisfying $V^2 = V$. Let x be a limit point of H ; we show that $x \in H$. take $x \in D$ (directed set), $\{x_\alpha\}$ be a net in H converging to x . Clearly, $x_\alpha \in H$ and since H is also a sub-group we find $x_\alpha^{-1} \in H$. So, the neighbourhood Vx_α^{-1} of x_α^{-1} shall cut H i.e. $(Vx_\alpha^{-1}) \cap H \neq \emptyset$. Take $y_\alpha \in (Vx_\alpha^{-1}) \cap H$. Since $x_\alpha \rightarrow x$, we see that $x_\alpha^{-1}y_\alpha \rightarrow x^{-1}x = e$ for $\alpha \rightarrow 0$ for some $\alpha \in D$. Thus for $\alpha \rightarrow 0$ we find $(y_\alpha x_\alpha^{-1}) \in (Vx_\alpha^{-1}) \cap H = Vx_\alpha^{-1} \cap H$. Therefore $(y_\alpha x_\alpha^{-1}) \in H$. Now the net $\{y_\alpha x_\alpha^{-1} : \alpha \in D, \alpha \neq 0\}$ converges to yx^{-1} , and \bar{H} being closed, we have $(yx^{-1}) \in \bar{H}$. Hence, $x = (yx^{-1})x \in H$ i.e. $H = \bar{H}$; that makes H to be closed. 2.2 Given a Topological Group G and Closed sub-group H in G . Given a Topological Group G and closed sub-group H in G . Suppose G/H denotes

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the family of all (Left) cosets of H in G

i.e. $G/H = \{aH : a \in G\}$. If H is a normal sub-group we need not make any distinction between left and right cosets of H in G . Thus G/H consists of all distinct cosets of H in G . We now take H to be a normal closed sub-group of G . Now G/H forms a group with respect to binary composition $aH \cdot bH = abH$ for $a, b \in G$, where in it is well known that H itself serves as the identity element in group G/H , and inverse member of aH ($a \in G$) in G/H is $a^{-1}H$. Definition 2.2.1. If H is a normal closed sub-group of Topological Group G , then the group G/H of all cosets of H in G is called the Quotient Group (also known as factor group) of G by H .

30 Example 2.2.1. Algebraically if G is the additive Group of all integers and $H = \{2n : n \in \mathbb{Z}\}$, then H is a normal sub-group of G and the Quotient Group G/H consists of two members H and $1 + H$. Example 2.2.2. Algebraically if G denotes the additive group of all rationals and $H = \mathbb{Z}$ the set of all integers in G , then H is a normal sub-group of G and a typical member of the quotient group G/H looks like m/n where m is an integer $\>$; n , and prime to n (n is a natural number). Therefore the Quotient Group G/H is an infinite group. We are now after an appropriate topology for the Quotient Group G/H in order to make G/H a topological Group, called very often, Quotient Topological Group or simply Quotient Group. Let $f : G \rightarrow G/H$ be the canonical mapping where $f(a) = aH$ as $a \in G$. Desired Topology in G/H shall make f continuous. We call a subset W of G/H to be 'open' if and only if $f^{-1}(W)$ is an open set in Topological Group G . We verify that the collection \mathcal{W} of such open sets W in Quotient group G/H forms a Topology in G/H . (2.1.1.) Since $f^{-1}(?) = ?$ and $f^{-1}(G/H) = G$ we see that \emptyset and G/H are members of \mathcal{W} . (2.1.2) Let W_1, W_2 be any two members of \mathcal{W} , then we have $f^{-1}(W_1)$ and $f^{-1}(W_2)$ are open sets in G , and so is $f^{-1}(W_1) \cap f^{-1}(W_2)$ which equals to $f^{-1}(W_1 \cap W_2)$. That means $W_1 \cap W_2 \in \mathcal{W}$. (2.1.3) Finally take $\{w_\alpha\}$ as a collection of member $w_\alpha \in \mathcal{W}$, then we know that $f^{-1}(w_\alpha)$ is an open set in G for each α , and $f^{-1}(\bigcup_\alpha w_\alpha)$ is also open set in G . i.e. $f^{-1}(\bigcup_\alpha w_\alpha) = \bigcup_\alpha f^{-1}(w_\alpha)$ is also open set in G ; that means, $\bigcup_\alpha w_\alpha \in \mathcal{W}$ and $\bigcup_\alpha w_\alpha$ is an open set in G/H .

31 So, (2.1.1) (2.1.3) verify that W is a topology in G/H ; This topology is called the Quotient topology in G/H . The Quotient Topology in G/H is one that makes canonical mapping f (see above) to be continuous. Theorem 2.2.1. With respect to Quotient Topology in G/H the Canonical mapping $f : G \rightarrow G/H$ is an open mapping. Proof : Take O be an open set in Topological Group G . We check that $f(O)$ is an open set in G/H . We need showing $f^{-1}(f(O))$ is open in G . Now, $f(O) = \{aH : a \in O\} = OH$. Take $x \in f^{-1}(f(O))$; so, $f(x) \in f(O) = OH$; there we find a member $y \in O$ such that $f(x) = yH$ or, $xH = yH$ or, $\{xh : h \in H\} = \{yh : h \in H\}$. Since H is a sub-group, $e \in H$ and we see $x = xe \in \{xh : h \in H\} = \{yh : h \in H\}$. Therefore, $x = yh$ for some $h \in H$. That means $x \in OH$ or we have $f^{-1}(f(O)) \subseteq OH$. Reversing the argument we deduce $OH \subseteq f^{-1}(f(O))$; and therefore $f^{-1}(f(O)) = OH$ which is, of course, an open set in G . Theorem 2.2.2. (H a closed sub-group) In the quotient Group G/H Quotient Topology is Hausdorff. Proof : Let $y, x \in G$ with $xH \neq yH$. So, $x \notin yH$. As H is closed, we see that yH is closed with x as an outside point; and x is not a limit point of yH ; so we find a neighbourhood U of the identity e in G such that $(xU) \cap (yH) = \emptyset$. We now find a symmetric open neighbourhood of e satisfying $W \subseteq U$. We assert that $(WxH) \cap (WyH) = \emptyset$ (1)
32 Otherwise, we find some $w \in U, w \in W$ and $h \in H$ such that $w \in xhH = wyH$. Thus $w \in W \cap xhH = wyH$. Therefore, $w \in W \cap wyH$.
Now

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$w \in W \cap wyH$ (since W is symmetric, $W^{-1} = W$) So, $w \in W \cap wyH$ (since W is symmetric). Therefore, $w \in W \cap wyH$

$x \in xH \cap yH$ while $y \in yH \cap xH$, because H is subgroup. i.e. $(Ux) \cap (yH) \neq \emptyset$ — a contradiction. Thus our assertion (1) stands. i.e. $(WxH) \cap (WyH) = \emptyset$ and that means, $(Wx) \cap (Wy) = \emptyset$; (taking $e \in H$) $(Wx) \cap (WxH)$, and similarly $(Wy) \cap (WyH)$. Put $W' = f(Wx) = WxH$ and $W'' = f(Wy) = WyH$ showing $W' \cap W'' = \emptyset$. To complete the proof we now show that $(xH) \cap W'$ and $(yH) \cap W''$ (here W', W'' are open in G/H ; f sending open sets to open sets). To that end we recall $f(x) = f(Wx)$ because $e \in W$. or $(xH) \cap W'$, and similarly, $(yH) \cap W''$ and therefore W' and W'' are respectively disjoint open covers for xH and yH in G/H .
Theorem 2.2.3.

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Let G be a Topological Group and H a closed normal subgroup of G , then the quotient Group G/H is a Topological Group

with quotient Topology. Proof : Consider the canonical mapping $f : G \rightarrow G/H$. In preceding Theorems we have seen that f is a continuous and open mapping. Now we check that f is a group Homomorphism. Take $x, y \in G$. Then $f(xy) = (xy)H = xHyH = f(x)f(y)$. Thus f is a Homomorphism.
33 We now show that Group operation in Quotient Group G/H shall be continuous with respect to underlying topologies. i.e. one must show that the mapping $(xH, yH) \mapsto xH(yH) = (xy)H : (G/H) \times (G/H) \rightarrow G/H$ is continuous. Suppose W be an open neighbourhood of $xH(yH) = (xy)H$ ($x, y \in G$), then $f^{-1}(W)$ is open in G with $(xy) \in f^{-1}(W)$. By continuity of group operation in G (a Topological Group), we find open sets U and V in G such that $x \in U$ and $y \in V$ with $UV \subseteq f^{-1}(W)$, or $f(UV) \subseteq W$. Since f is also a Group homomorphism, we have from above $f(U)f(V) = f(UV) \subseteq W$. Let $u \in U$ and $v \in V$; Then $uH \cdot vH = f(u)f(v) \in f(UV) \subseteq W$ i.e. $uH \cdot vH \in W$. This shows that group operation in G/H is continuous to make the quotient group G/H a Topological Group with quotient topology. Definition 2.2.2. A

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Topological space X is said to a Homogeneous space if for any two member x_1, x_2 in X there is a self Homeomorphism f in X such that $f(x_1) =$

x_2 . For example, every topological group G is always a Homogeneous space ; because if $x_1, x_2 \in G$, let us take $x_1^{-1}x_2 = u \in G$ and defining the mapping $f : G \rightarrow G$ where $f(x) = xu$ for $x \in G$, we see atonce that f is a self-homeomorphism of G such that $f(x_1) = x_2$. Theorem 2.2.3. (a). If H is a sub Group of a Topological Group G , then G/H , the quotient Topological Group is homogeneous. Proof : Take two members x_1H and x_2H in G/H with $x_1, x_2 \in G$. Taking $x_1^{-1}x_2 = u \in G$ consider a homeomorphism $f : G/H \rightarrow G/H$ given by $f(xH) = (xu)H (= xHuH)$ for all $(xH) \in G/H$. Then we have $f(x_1H) = (x_1u)H = (x_1^{-1}x_2)H = x_2H$. Hence G/H

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is homogeneous. 34 Theorem 2.2.3.(b). Let G be a Topological Group and H a sub-group of G . Then G/H is T_1 if and only if H

is closed. Proof : Suppose G/H is T_1 . Then every singleton in G/H is closed. Therefore $\{eH\} = \{eH\}$ is closed in G/H ; Under conomical mapping $f : G \rightarrow G/H$ which is continuous we have $f^{-1}(eH) = H$. Therefore H is closed in G . Conversely let the sub-group H be closed in G . Take any member xH in G/H . consider the singleton $\{xH\}$ in G/H . Since H is closed we know that xH is closed making $G \setminus \{xH\}$ to be open in G . Therefore under cononical mapping $f : G \rightarrow G/H$, we have $f(G \setminus \{xH\})$ is open in G/H . Now $f(G \setminus \{xH\}) = (G/H) \setminus \{xH\}$ we conclude that $\{xH\}$ is closed in G/H . Therefore every singleton in G/H is closed and that makes G/H T_1 .

The proof is complete. Theorem 2.2.3(c)
Let G be a Topological Group and

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H a sub-group of G . Then G/H is a discrete space if and only of H is open. Proof :

Suppose G/H is a discrete space. Therefore each singleton of G/H is open. In particular, $eH = H$ (e being the identity of G) is open. Under cononical mapping $f : G \rightarrow G/H$ which is continuous, we have $f^{-1}(eH) = H$ becomes open in G . Conversely let sub-group H be open. If $x \in G$, we have xH is open. That means every singleton in G/H is open in G/H and this is why G/H is a discrete space. Theorem 2.2.3(d) : Let H be a sub-group of a Topological Group G , and $f : G \rightarrow G/H$ be the cononical mapping. If $\{V_i\}$ be a fundamental system of neighbourhoods of the identity e in G , then the family $\{f(V_i)\}$ is a fundamental system of neighbourhoods of the identity $eH = H$ of G/H . Proof : Let $f : G \rightarrow G/H$ be the canonical mapping. By property of f we see that if V_i is any member of

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a fundamental system $\{V_i\}$ of neighbourhoods of the identity e in G , then $f(V_i)$ is a neighbourhood of

the identity eH in G/H . Suppose V is any neighbourhood of eH in G/H . Then $f^{-1}(V)$ by continuity of f , is a

35 neighbourhood of the identity e in G . So we find a member, say, γ in the family $\{U_\alpha\}_{\alpha \in I}$ such that $\gamma \in \bigcap_{\alpha \in I} U_\alpha$ or, $\gamma \in \bigcap_{\alpha \in I} \gamma^{-1} U_\alpha \gamma$. This shows that the family $\{\gamma^{-1} U_\alpha \gamma\}_{\alpha \in I}$ is a fundamental system of neighbourhoods of the identity e . $H = H$ in G/H . Definition 2.2.3. A Topological Group G is said to be totally disconnected if the component of the identity e in G equals to $\{e\}$. Theorem 2.2.3. (e) : Let C be the component of the identity e in a Topological Group G . Then the quotient topological Group G/C becomes a totally disconnected T_2 space. Proof : First we show that C is a closed normal sub-group of G . Since C is the component ; by maximality C becomes closed. Now take $a \in C$. Then $a^{-1} C \subset C$, because $a^{-1} C$ is the image of C under the homeomorphism $x \mapsto a^{-1} x$ becomes connected with e $a^{-1} C \subset C$; therefore, $a C \subset C$ $a C \subset C$ $a^{-1} C \subset C$ So, C is a closed sub-group of G . Further, by continuity of the mapping : $x \mapsto a^{-1} x a$ we have for $a \in G$, $a^{-1} C a$ is also connected ; thus $a^{-1} C a \subset C$ for each $a \in G$ because C is the component. Therefore C is a Normal sub-group of G . As C is closed it follows that quotient G/C is T_1 -space and hence it is T_2 . We have now to show that G/C is totally disconnected. Let U be the component of the identity member ($eC = C$) in G/C . If π is the natural homomorphism of $G \rightarrow G/C$, we have $\pi^{-1}(U) = C$ and $C \cap \pi^{-1}(U) = C$. If G/C is not totally disconnected there is a member $(x.C) \neq (e.C)$ such that $(x.C) \in U$. That

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means C is a proper subset of $\pi^{-1}(U)$. Since C is a maximal connected set containing e , $\pi^{-1}(U)$

is not connected. Let a disconnection of $\pi^{-1}(U)$ be like : $\pi^{-1}(U) = [P \cap \pi^{-1}(U)] \cup [Q \cap \pi^{-1}(U)] \dots (1)$ where

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P and Q are open sets in G , such that $[P \cap \pi^{-1}(U)] \cap [Q \cap \pi^{-1}(U)] = \emptyset$ and neither is empty. So $U = [\pi(P) \cup \pi(Q)] \cup U$. Taking $U = UC$ Let $x \in U$ such that $x \in C$; Hence from (1) we have $x \in C = (P \cap xC) \cup (Q \cap xC)$ Since xC is connected, either $x \in (P \cap xC)$ or, $x \in (Q \cap xC)$. Consequently, images $P \cap UC$ and $Q \cap UC$ under π are disjoint, since they are unions of cosets of C . $\pi(P) \cap \pi(Q) = \emptyset$

$U \cap \pi(Q) \cap U = \emptyset$ Now π is an open mapping, so $\pi(P)$ and $\pi(Q)$ are open sets, and hence we have shown that U is not connected—a contradiction what U is the component of eC . Hence we have proved that G/C is totally disconnected. Remark : Given a topological Group G and a closed normal sub-group H in G , we have seen that cononical mapping $f : G \rightarrow G/H$, where G/H is topological group with quotient topology, becomes a continuous mapping which is also an open mapping. This mapping may not be a closed mapping. Example 2.2.3. Let R be the topological Group with addition as Group Composition and with usual topology of reals ; If Z is the sub-group of R consisting of all integers, then we see that Z is closed and a Normal sub-group of R . Here canonical mapping $f : R \rightarrow R/Z$ is not closed. Solution : Consider the set $E = \{x \in \mathbb{R} : 0 \leq x < 1\}$. Then E is a closed set in topological Group R . Every coset $x + Z$ in R contains the number $x - [x]$, $[x]$ denoting the largest integer not larger than real x and no other real number in $[0, 1)$. Therefore, $[0, 1)$ may be treated as the quotient space R/Z . The Topology imposed in $[0, 1)$ as a model of the space R/Z has basic open sets like (γ, δ) , and $[0, \gamma) \cap [0, 1)$ where $0 < \gamma < \delta < 1$. Now canonical mapping f sends E into a non-closed set (having 0 as a limit point outside the image set $f(E)$). Hence the conclusion stands OK. However we have following Theorem in this connection. Theorem 2.2.4. If H is a compact normal sub-group of a Topological Group G , then the cononical mapping : $G \rightarrow G/H$ is a closed mapping where G/H is the quotient topological Group. Proof : Suppose C is a closed set in G ; and the canonical mapping $f : G \rightarrow G/H$ is in action to send $x \in G$ to xH if $f(x) = xH$ as $x \in G$. Take $xH \in (G/H) \setminus f(C)$, and $x \notin CH$. As C is closed and H is compact we know that CH is closed. Therefore x is an outside point of the closed set CH , and we find an open set U in G such that $x \notin (U \cap CH)$. Cononical mapping f being an open mapping $f(U \cap CH)$ is an open set containing $f(x) = xH$ i.e. $f(U)$ is an open neighbourhood of xH such that $f(U \cap CH) \cap f(C) = \emptyset$, showing that $(G/H) \setminus f(C)$ is open and hence $f(C)$ is closed. The proof is complete. 2.3 Locally compact Groups : We recall following Definition : Definition 2.3.1. A topological space X is called locally compact if each point x in X has an open neighbourhood U whose closure \bar{U} is compact. Then it is true that a Hausdorff topological space is locally compact if and only if, each point has a compact neighbourhood. Also we remember that every Hausdorff locally compact topological space is completely regular (and hence regular). Theorem 2.3.1.

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A Locally compact Hausdorff topological space X is normal if it is the union of an increasing

sequence $\{U_n\}$ of open sets such that each U_n is compact. Proof : We have by assumption

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$n \in \mathbb{N}$, and we write $X_n = \overline{U_n} \setminus U_{n-1}$ and X_n is compact. Also $X = \bigcup_{n \in \mathbb{N}} X_n$.

a f , where $f(0) = 0$. Suppose $\mathcal{P} = \{U_\alpha\}$ is an open cover for X . Since each X_n is compact, there shall be a finite sub-cover of \mathcal{P} for X_n . This is true for each n . As countable union of finite families of sets constitute a countable family, one has a countable sub-family of \mathcal{P} to cover X -making X a Lindeloff space. Since every Lindeloff regular space is normal, the conclusion is arrived at as desired. Theorem 2.3.2. Every compact Hausdorff space is normal. For proof see any text book on general topology.

38 Theorem 2.3.3. A Topological Group is a locally compact topological group if and only if its identity e has a compact neighbourhood. Proof : Suppose G is a locally compact topological group. So its identity e has a neighbourhood V whose closure \overline{V} is compact. Conversely, suppose G is a Topological Group where identity e has a compact neighbourhood V . Choose

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a neighbourhood V of e such that $VV = V$. Now, $V \cap \overline{V} = V$; Hence V is a closed subset of compact set \overline{V} and therefore V is a compact neighbourhood of e .

Let x be any element in G . Then xV is a neighbourhood of x and we have $xV \cap xV = x(V \cap V) = xV$ becomes compact, because translation operator is a homeomorphism in G . The proof is now complete. Theorem 2.3.4. A locally compact Hausdorff topological Group is normal. Proof : First we establish that in a general topological Group G if H is a symmetric neighbourhood of its identity e , then H is a clo-open (closed and open) sub-group of G . Because if H is a neighbourhood of e and $x, y \in H$, say $x \in H$ and $y \in H$. Then $xy \in H$. Further, $x^{-1} \in H$ and $(x^{-1})^{-1} = x \in H$ (H being symmetric). Therefore H is a sub-group of G . If $y \in H$, we have $y \in \text{int}(H)$ showing every member of H is an interior point of H and H is open, and every open sub-group of G is also closed. Hence the assertion follows. Now it is known that

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H is a Hausdorff locally compact sub-group of G . Also $H = \bigcup_{n \in \mathbb{N}} H_n$ for $n \in \mathbb{N}$. Hence H_n is compact. Thus H is a union (countable union) of increasing sequence of compact sets. H is normal (see theorem 2.3.1). Consider the collection $\{aH\}$ of pairwise disjoint cosets in G . Since

translation (left or right) is always a homeomorphism in G , each member aH ($a \in G$) is homeomorphic to H and becomes normal. Therefore $G = \bigcup aH$ becomes normal. The proof is complete.

39 Corollary : If G is a locally compact Hausdorff Topological Group and C is a closed subset in G and U an open set with $C \cap U = \emptyset$, then there is a real-valued continuous function f over G such that $f(x) = 1$ if $x \in C$ and $f(x) = 0$ if $x \in U$. Because G is normal by Theorem above and C and U are a pair of disjoint closed sets, by Urysohn's Lemma we find a continuous function $f : G \rightarrow [0, 1]$ satisfying $f(x) = 0$ if $x \in U$ and $f(x) = 1$ if $x \in C$.

Theorem 2.3.5. Let G be a locally compact Topological Group, and

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Let C be a compact subset and V an open subset of G such that $C \subseteq V$. Then there

is a neighbourhood V of e such that

$C \cap V$ is compact with $C \subseteq V$. Proof: As C is compact, if $x \in C$, we find an open neighbourhood V_x

of the identity e in G such that $xV_x \subseteq V$. Also choose an open neighbourhood W_x of e such that $W_x \cap W_x = W_x^2 \subseteq V_x$. Now the family $\{xW_x \mid x \in C\}$ becomes an open

cover

for C . By compactness of C , there is a finite sub-cover, say

$x_1V_{x_1}, x_2V_{x_2}, \dots, x_nV_{x_n}$ to cover C . Now put $W = \bigcap_{i=1}^n W_{x_i}$, then W is an open neighbourhood of e in G .

Clearly $C \subseteq W \cap W = W^2 \subseteq V$ (since $W_{x_i} \cap W_{x_i} = W_{x_i}^2 \subseteq V_{x_i}$). By a similar argument we produce an open neighbourhood W' of e in G such that $W' \subseteq C$. Since $W \cap W' = W$

is an

open neighbourhood of e

in G , we choose a neighbourhood V of e in G such that

its closure \bar{V} is compact and $V \cap \bar{V} = V$. Therefore $C \cap \bar{V} \subseteq V$. As C is compact and \bar{V} is compact we know that $C \cap \bar{V}$ and $V \cap \bar{V}$ are each closed set; Also $C \cap \bar{V} \subseteq V \cap \bar{V}$ as a Union of two compact sets becomes a compact set.

40 Further $C \cap \bar{V} = C \cap V$ and $V \cap \bar{V} = V \cap C$. Therefore $(C \cap \bar{V}) \cap (V \cap \bar{V}) = C \cap V \cap \bar{V} = C \cap V$ as compact with $(C \cap \bar{V}) \cap (V \cap \bar{V}) \subseteq C \cap V$. Theorem 2.3.6. Let V be an open neighbourhood of the Identity e

in a Topological Group G and C be a compact set in G . Then there is an

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open neighbourhood V of e such that $C \cap V = -1 \cap C \cap V$. Proof: Choose a symmetric open neighbourhood W_1 of the identity e in G such that $W_1 \cap W_1 = W_1^2 \subseteq V$, and for a fixed $a \in G$ take a symmetric open neighbourhood W_2 of e such that $aW_2a^{-1} \subseteq W_1$. Put $W = W_1 \cap W_2$. Now $x \in W$ gives $(xa^{-1}) \in W_2 \cap W_1$; and $ax^{-1} \in W_1 \cap W_2$ (

W_1 is symmetric). Therefore, $xWx^{-1} \cap W_2 \cap W_1 = (xa^{-1})aW_2a^{-1} \cap (ax^{-1})W_1 \cap W_1 = W_1 \cap W_2 \subseteq V$. Since W is dependent on $a \in G$, we designate W by W_a . Now the family $W_a \mid a \in C$ is

an open cover for C ; by compactness of C , there is a finite sub-cover, say,

$W_{a_1}, W_{a_2}, \dots, W_{a_n}$ to cover C . Let us put $V = \bigcap_{i=1}^n W_{a_i}$. Then V is an open symmetric neighbourhood of e in G . If $x \in C$ we see that $x \in W_{a_k}$ for some k , and this implies $xW_{a_k} \subseteq V$. Therefore, $xVx^{-1} \subseteq V$.

This completes

the proof. Example 2.3.1.

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Let G be a Topological Group and N is a closed Normal sub-group. (i) if G is compact, then G/N is a compact quotient Topological Group; and (ii) if G is locally compact, then G/N is a Locally compact

quotient Topological Group. Solution. given that N is a closed normal sub-group. Then the quotient group G/N becomes a Topological Group (See Theorem 2.2.3). (i) Suppose G is compact. Now the canonical mapping $f: G \rightarrow G/N$ where

41 $f(x) = xN \text{ } G/N$ as $x \text{ } G$ is continuous, and therefore $f(G)$ is compact since G is compact. Here $f(G) = G/N$. So, G/N becomes compact. (ii) Suppose G is locally compact. So there is an open neighbourhood O of the Identity e in G such that O (closure of O) is compact. Now $f(e) = eN = N$; Therefore $N = f(e) \text{ } f(O) \text{ } f(O)$ as $O \text{ } O$ By continuity of f we also have $f(O)$, is compact. So $f(O)$ is a compact subset of a Hausdorff space, and therefore $f(O)$ is closed. Also $f(O)$ is an open neighbourhood of N (f is an open mapping) and $f(O) \text{ } f(O) \text{ } f(O) \text{ } f(O) \text{ } f(O) \text{ } f(O)$, because $f(O)$ is closed. Thus $f(O)$ is closed subset of $f(O)$ which is compact. Therefore $f(O)$ is compact. Hence G/N is locally compact. 2.4 Topological Groups Involving

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Connectedness : Definition 2.4.1. A topological space X is said to be connected if

X does not admit of a decomposition like $X = P \cup Q$ Where P and Q are non-empty disjoint open sets in X . Explanation : A connected Topological space X is thus such a strong piece of objects that it does not allow its partition in the manner as above. Definition 2.4.1 shows that a Topological space X is connected if and only if in the space X there are no clopen (Closed and open) sets other than \emptyset and X . A subset E of X shall be taken as a connected set if it is a connected space in respect of relative topology of E . In the real number space R with usual topology it is known that a subset of R is connected if and only if it is an interval. Definition 2.4.2. Given a point in X , the maximal connected subset in X containing the point is said to be the component of that point. In consequence, we recall that given a connected set A in X , its closure \bar{A} is also a connected set, and thus every component in X is a closed set. Further, if $\{E_i\}$ is a family of connected sets in X , with $\bigcap E_i \neq \emptyset$, then $\bigcup E_i$ is

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is connected. 42 A topological space X is said to be Locally Connected if each

open neighbourhood of every point in X contains a connected open neighbourhood. We also recall that continuous image of a connected space becomes connected, and this gives as a special case that every real-valued continuous function over an interval enjoys Intermediate value property. In the following we present some basic properties of Topological groups depending upon connectedness of the Group when taken as a topological space. Theorem 2.4.1.

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Let G be a Topological Group and H be the component of the Identity e of G . Then H is a closed Normal sub-group of G .

Proof :

We know that in G group operation inversion : $x \mapsto x^{-1}$ as $x \mapsto G$ is a homeomorphism, and therefore H^{-1} being continuous image of connected set H becomes connected, and it is a connected set containing e . This shows that by maximal property of H as the component, $H^{-1} \subseteq H$. Take $x \in H$, then $x^{-1} \in H^{-1} \subseteq H$. So, $x^{-1} \in H$. then xH is a connected set such that $e = xx^{-1} \in xH$; Therefore xH is a connected set containing e ; By maximal property of H as the component containing e we have $xH \subseteq H$ That means, $H \text{ } H$ (x is any member of H) (2) From (1) and (2) it follows that H is a sub-group in G . Now take $a \in G$; then the mapping : $x \mapsto axa^{-1}$ as $x \mapsto G$ is a homeomorphism; thus by its continuity we have as a continuous image aHa^{-1} is a connected set containing e . Again by maximality of H as a maximal connected set containing e , we produce. $aHa^{-1} \subseteq H$. That means H is a normal sub-group in G . Since every component in a topological space is always closed we have H as a closed set. Thus

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H is a closed Normal sub-group in G. Example 2.4.1. Let G be a Topological Group and H be the component of

the identity

e in G ; If $a \in G$, aH (= Ha) is the component of a. Solution : Here H is a Normal sub-group of G (see Theorem 2.4.1) If $a \in G$, we have $aHa^{-1} = H$; giving $aH = Ha$.

43 Since translation is a homeomorphism we have aH as a connected set containing a (H is the component e, hence a connected set with $e \in H$). Let C be the component of a. Then we have $aH \cap C \neq \emptyset$ (1) Let L as a connected set with $a \in L$. Then $a^{-1}L$ is a connected set containing e. because $e = a^{-1}a \in a^{-1}L$. H being the component of e we have $a^{-1}L \cap H \neq \emptyset$ or, $L \cap aH \neq \emptyset$ This being true for any connected set containing a we have, the component C of a satisfies. $C \cap aH \neq \emptyset$ (2) From (1) and (2) we have $C = aH$ i.e. aH is the component of a. Theorem 2.4.2. Let G be a connected Topological Group and H is the component of this identity e. If N is any neighbourhood of e, then $G = \bigcup_{n \in \mathbb{N}} H^n$. Proof : Choose a symmetric neighbourhood V of the identity e in G such that $V \cap N \neq \emptyset$. Then we have (see corollary following Theorem 2.1.2.) $V \cap N$ is open and closed. Since G is connected, G is the only non-empty open and closed (clopen) set in G. Hence we have $G = V \cap N$, giving $G = N$. Now Suppose $\{G_i\}_{i \in I}$ be a family of Topological Groups. Put $G = \prod_{i \in I} G_i$; G is called product of G_i 's. Let G have the product topology. Let $x = \{x_i\}_{i \in I}$, $y = \{y_i\}_{i \in I}$ be two elements of G. Then xy is defined as : $xy = \{x_i y_i\}_{i \in I}$; with this definition of composition of two members of G we easily verify that G forms a Group where the identity element e of G is given by $e = \{e_i\}_{i \in I}$ where e_i is the identity element of G_i for $i \in I$. This Group G is called the Direct product of $\{G_i\}_{i \in I}$, where individual members G_i are called factors.

44 Theorem 2.4.3.

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Let $\{G_i\}_{i \in I}$ be a family of Topological Groups. If $G = \prod_{i \in I} G_i$ is

the Direct product of G_i endowed with product topology, then G is a Topological Group. Proof : Actually we need showing that the mapping $(x, y) \mapsto xy^{-1}$ of $G \times G$ onto G is continuous. To that end take W as a neighbourhood of xy^{-1} in G as $(x, y) \in G \times G$. Then there is a finite number of indices, say i_1, i_2, \dots, i_n such that $W \cap \prod_{i \in I} G_i \neq \emptyset$, with $W \cap \prod_{i \in I} G_i = G$ for $i \in I \setminus \{i_1, i_2, \dots, i_n\}$ and V_i as open neighbourhoods of $x_{i_j} y_{i_j}^{-1} \in W$ ($i_j \in \{i_1, i_2, \dots, i_n\}$), and $V = \prod_{i \in I} V_i$. Since $(x, y) \mapsto xy^{-1}$ is a continuous operation in topological Group G for each i_j , we obtain neighbourhood V_i of $x_{i_j} y_{i_j}^{-1}$ in G_{i_j} ($i_j \in \{i_1, i_2, \dots, i_n\}$) such that $V_i \cap G_{i_j} \neq \emptyset$; $V_i \cap G_{i_j} \neq \emptyset$. Put $V = \prod_{i \in I} V_i$ where $V_i = G_{i_j}$ for $i_j \in \{i_1, i_2, \dots, i_n\}$, and $V_i = G_i$ for $i \in I \setminus \{i_1, i_2, \dots, i_n\}$. Similarly construct V' ; Then V and V' respectively form neighbourhoods of x and y, and we have $VV' \cap W \neq \emptyset$. Therefore we have checked that direct product

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G is a Topological Group with Group composition and product Topology. Theorem 2.4.4. Let $G = \prod_{i \in I} G_i$ be the direct product of Topological Groups $\{G_i\}_{i \in I}$

$\{G_i\}_{i \in I}$, and let G have the product Topology. Then following statements are true. (i) G is a compact Topological Group if and only if each G_i is a compact Topological Group. (ii) G is a T₂ Topological Group if and only if each G_i is so. (iii) G is a locally compact Topological Group if all G_i are compact Topological Group except for a finite number of them that are each a locally compact Topological Group.

47 Theorem 2.5.2. G_n forms an open set in M_n . Proof : Consider the mapping $\det : M_n \rightarrow \mathbb{R}$, defined by $\det(A) = \det A$ as $A \in M_n$. Now $G_n = \{A \in M_n : \det A \neq 0\}$. Since $\det^{-1}(0) = \{A \in M_n : \det(A) = 0\} = \{A \in M_n : \det A = 0\}$ we have. $G_n = M_n \setminus \det^{-1}(0)$ Since \det is continuous we see that $\det^{-1}(0)$ is a closed set in M_n and therefore G_n is open in M_n . Theorem 2.5.3. G_n is a T_2 -multiplicative Topological Group with respect to relative Topology induced by M_n . Proof : We know that product of two non-singular square matrices of order n is again a non-singular matrix of the same size. Further if $A \in G_n$, then $(A^{-1})^{-1} = A$, and we see that $A^{-1} \in G_n$. Thus with matrix multiplication G_n forms a Group whose identity element is the identity matrix $I = 1 \ 1 \ 0 \ 0 \dots \ 0 \ n \times n$ with upper and lower blocks comprise of zeros since M_n is T_2 , one sees that G_n with respect to relative topology inherited from M_n is also T_2 . We now examine continuity of group composition of G_n in this topology. Let $A, B \in G_n$, and $A = (a_{ij})_{n \times n}$, $B = (b_{ij})_{n \times n}$. If $AB = C$ where $C = (c_{ij})_{n \times n}$ and $c_{ij} = \sum_{k=1}^n a_{ik} b_{kj}$. Now mappings $A \rightarrow a_{ij}$ are continuous, because they are projections of \mathbb{R}^{n^2} onto co-ordinate spaces. Similarly $B \rightarrow b_{ij}$ are continuous, therefore $AB \rightarrow C_{ij}$ is also continuous. So mapping $(A, B) \rightarrow AB$ is continuous. Finally, if $A \in G_n$, we have $A^{-1} = \frac{1}{\det A} \text{adj } A$ where $\text{adj } A$'s are minors in and are poly nomials in coefficients in A . As $\det A \neq 0$, mapping $A \rightarrow A^{-1}$ is also continuous.

48 Therefore $A \rightarrow A^{-1}$ is continuous. Thus G_n is a T_2 -topological multiplication Group. Definition 2.5.2. A topological space X is called locally Euclidean if there is a +ve integer n such that every $x \in X$ has a neighbourhood U_x which is homeomorphic to the open unit ball of the Euclidean n -space \mathbb{R}^n , namely $U_x = \{(x_1, x_2, \dots, x_n) \in \mathbb{R}^n : \sum_{i=1}^n x_i^2 < 1\}$. Explanation : A Topological Group is locally Euclidean if and only if for some +ve integer n , its identity e has a neighbourhood homeomorphic to the open unit ball of the Euclidean- n space \mathbb{R}^n . Theorem 2.5.4. M_n is a locally Euclidean Group. Proof : By identification technique and defining Topology in M_n , we see that M_n is homeomorphic to \mathbb{R}^{n^2} which in turn is homeomorphic to the Euclidean space \mathbb{R}^{2n^2} . Thus conclusion stands O.K. Remark : One can prove a similar Theorem saying that G_n is locally Euclidean. Thus examples of locally Euclidean topological group are not scarce. However, we note that there are topological Groups that are not locally Euclidean. For example, take $G = \mathbb{R}^{\mathbb{R}}$ where $\mathbb{R} =$ the space \mathbb{R} of all reals for each $x \in \mathbb{R}$ and G is the direct product of an infinite number of copies of \mathbb{R} . G is equipped with the product topology. Then G is a topological Group that is not locally Euclidean. 2.6. Lie Groups : Consider a real-valued function f over an open set $S \subset \mathbb{R}^n$ (Euclidean n -space). f is said to belong to the class C^∞ if all partial derivatives including mixed derivatives of all orders of f exist and they are continuous in S . Now X is a T_2 (Hausdorff) space. We now explain what is meant by an atlas A of class C^∞ on X . Definition 2.6.1. (i) A family of ordered pairs like $\{(U_\alpha, \phi_\alpha) : \alpha \in I\}$ where $\{U_\alpha\}$ forms an open cover of X and $\phi_\alpha : U_\alpha \rightarrow \mathbb{R}^n$ is a homeomorphism for each $\alpha \in I$ is called an atlas, denoted by A if following conditions are satisfied :

49 (a) For $U_\alpha, U_\beta \in A$, $U_\alpha \cap U_\beta \neq \emptyset$, $\phi_\alpha \circ \phi_\beta^{-1} : \phi_\beta(U_\alpha \cap U_\beta) \rightarrow \phi_\alpha(U_\alpha \cap U_\beta)$ is of class C^∞ . (b) Let (U, ϕ) be a pair of an open set U in X and a homeomorphism ϕ of U onto an open set of \mathbb{R}^n . If for each pair $(U_\alpha, \phi_\alpha) \in A$ for which $U \cap U_\alpha \neq \emptyset$, $\phi \circ \phi_\alpha^{-1} : \phi_\alpha(U \cap U_\alpha) \rightarrow \phi(U \cap U_\alpha)$ is of class C^∞ , then $(U, \phi) \in A$. (ii) A Hausdorff topological space X with an atlas A is called a manifold. Explanation : Consequence of Definition 2.6.1. is that every manifold is locally Euclidean and therefore it is locally compact. We recall that $M_n(\mathbb{R})$ may be identified with the Euclidean space \mathbb{R}^{n^2} and that M_n may be identified with the Euclidean space \mathbb{R}^{2n^2} ; therefore they are each a Manifold. Definition 2.6.2. A manifold G which is also a Group is called a Lie Group if mappings (i) $(x, y) \rightarrow xy$ of $G \times G$ onto G and (ii) $x \rightarrow x^{-1}$ of G onto G are analytic functions. For example, the Euclidean n -space \mathbb{R}^n is a lie Group, because, for $x = (x_1, x_2, \dots, x_n) \in \mathbb{R}^n$; x_i are reals, taking the identity mapping $I(x) = (x_1, x_2, \dots, x_n) = x \in \mathbb{R}^n$, we verify that I belongs to C^∞ , and all requirements are O.K. for \mathbb{R}^n to be a manifold. So \mathbb{R}^n is a manifold. Further \mathbb{R}^n is additively a commutative Group such that $(x, y) \rightarrow x + y$ as $x, y \in \mathbb{R}^n$ is analytic, and similarly $x \rightarrow -x \in \mathbb{R}^n$ is also analytic. Therefore, \mathbb{R}^n is a lie group. Example 6.2.1. Topological Group M_n is a lie Group. Every lie Group is locally Euclidean and hence locally compact. The famously well known fifth problem of Hilbert says that every locally Euclidean topological Group is a lie Group. For compact and Abelian Topological Group Problem had been solved long before the general solution was found. one may see pontrjagin "Topological Group". Homomorphism between Topological Groups : Theorem 2.6.2. If

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G and H are two Topological Groups and $f : G \rightarrow H$ is a homomorphism then (

- a) For any two subsets A and B in G, $f(AB) = f(A)f(B)$ (b) For any two subsets C and D in H, $f^{-1}(C)f^{-1}(D) = f^{-1}(CD)$
 50 (c) If S is a symmetric set in G, then $f(S)$ is symmetric in H (d) If T is a symmetric set in H, then $f^{-1}(T)$ is symmetric in G. Proof : (a) Since f is a homomorphism, we have $f(AB) = f(A)f(B)$ whenever $A \subseteq G$ and $B \subseteq G$. (b) Let f be a homomorphism : $G \rightarrow H$ and take $x \in f^{-1}(C)$ and $y \in f^{-1}(D)$; So we have $f(x) \in C$ and $f(y) \in D$ Now $f(xy) = f(x)f(y) \in CD$; since f is a homomorphism. Therefore, $xy \in f^{-1}(CD)$. So, we write, $f^{-1}(C)f^{-1}(D) \subseteq f^{-1}(CD)$. (c) Let S be a symmetric set in G. We show that $f(S)$ is a symmetric set by showing $f(S) = (f(S))^{-1}$. Take $y \in f(S)$; So $f(x) = y$ for some $x \in S$. Since S is symmetric, we have $x^{-1} \in S$ Hence $y^{-1} = (f(x))^{-1} = f(x^{-1})$, Since f is a homomorphism. So, $y^{-1} \in f(S)$ or $y \in (f(S))^{-1}$ This gives $f(S) \subseteq (f(S))^{-1}$ (1) Conversely, take $x \in (f(S))^{-1}$ Then $x^{-1} \in f(S)$ So $x^{-1} = f(u)$ for some $u \in S$. Thus $(f(u))^{-1} = f(u^{-1})$ (f is a homomorphism) $\in f(S)$ because S is symmetric. This gives $x = (f(u))^{-1} \in f(S)$ or, $(f(S))^{-1} \subseteq f(S)$ (2) (1) and (2) give $f(S) = (f(S))^{-1}$, showing that $f(S)$ is symmetric. (d) proof shall be similar to that of (c).

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Theorem 2.6.3. Let G and H be two topological Groups and $f : G \rightarrow H$ be a Homomorphism. Then (

- a) For any two subsets A and B of G, $f(A)f(B) \subseteq f(AB)$ (b) For any two subsets C and D of H, $f(C)f(D) \subseteq f(CD)$
 51 (d) For any symmetric subset T in H, $f^{-1}(T)$ is symmetric in G and $f(B)f(B) \subseteq f(B)$, bar denoting the closure. Proof : First we observe that for any two subsets A and B in G, using continuity of group operation in G we have $AB \subseteq \overline{A}\overline{B}$, bar denoting the closure. Taking note of this inclusion relation proof of (a) and (b) shall follow from (a) and (b) parts of Theorem 2.6.2. above. (c) Inverse mapping in a Topological Group is a homomorphism, Therefore for any subset E in G we have $E \subseteq \overline{E}$. Let S be a symmetric set in G. Then Theorem 2.6.2. Says that $f(S)$ is symmetric. Consider $f(S)$ in topological Group H. By the remark above we have $f(S) \subseteq \overline{f(S)}$, because $f(S)$ is symmetric. That means $f(S)$ is symmetric. (d) The proof is similar to that in part (c). Theorem 2.6.4. If G and H are two Topological Groups, and $f : G \rightarrow H$

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is a Homomorphism. Then f is continuous if and only if f is continuous

at the identity e in G. Proof : Let $f : G \rightarrow H$ be a Homomorphism, and let f be continuous. Then of course

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f is continuous at the identity element e of G. Conversely, suppose f is continuous at e, and $x \in G$ ($x \neq e$). Let W be a neighbourhood of $f(x)$ in H.

Choose a neighbourhood V of the identity e_H in H such that $W = f(x)V$. Now f being a Homomorphism we know that $f(e) = e_H$, and using continuity of f at e, we find a neighbourhood U of e in G such that $f(U) \subseteq V$. Clearly U is a neighbourhood of x in G such that $f(xU) = f(x)f(U) \subseteq f(x)V = W$ That shows, f is continuous at x. The proof is complete.

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Theorem 2.6.5. Let G and H be two topological Groups and $f : G \rightarrow H$ be a Homomorphism. Then

f sends any open set in G to an open set in H iff $f(O)$ is open in H for every open set O containing the identity e in G .

52 Proof : Suppose f is an open mapping i.e. f sends any open set in G to an open set in H . Then, of course, $f(O)$ is open in H whenever O is an open set in G containing the identity e of G . Conversely, suppose the condition holds and take any open set U in G . If $x \in U$, then U is a neighbourhood of x in G and choose an open set V containing the identity e in G such that $U = xV$; Now $f(V)$ is open in H . Then $f(U) = f(xV) = f(x)f(V)$, because f is a Homomorphism. Now $f(V)$ being open we have $f(U)$ is an open set in H . i.e. $f(U)$ is open in H . So, f sends an open set in G to an open set in H .

Example 2.6.2. A continuous Homomorphism between two topological Groups may not be an open mapping. Solution. Take R as the set of all reals. Treat R as an additive commutative Topological Group with discrete Topology. Also treat R as an additive commutative Topological Group with usual Euclidean Topology and call it R_u . Then consider the identity mapping $I : R \rightarrow R_u$ as a Homomorphism which is, in this case, 1-1 and onto. Since discrete topology is strictly finer than the usual topology of reals in R we see at once that I is not an open mapping; nevertheless, I is continuous.

EXERCISE-A Short Answer type Questions

1. Define a sub-group of a Topological Group with an example with justification.
2. When is a sub-group of a Topological Group called discrete? Find a discrete sub-group of the Topological additive Group R of all reals.
3. If H is a sub-group of a Topological Group G , show that its closure \bar{H} is a sub-group of G .
4. In a Topological Group G is $x \rightarrow x^{-1}$ a self-homeomorphism of G such that $f(x^{-1}) = x$.
5. Let G be a locally compact Topological Group and $f : G \rightarrow F$ is an open continuous homeomorphism where F is another Topological Group. Show that F is locally compact.

53 6. If a topological Group G is connected and H is a sub-group of G , show that G/H is connected. 7. If H is a sub-group (a normal sub-group) of a Topological Group G , show that its closure \bar{H} is a subgroup (a normal sub-group) of G . 8. Example if the set Q of all rationals forms a Topological sub-group of Topological additive Group R of all reals. 9. Find a discrete sub-group of Topological additive Group R of all reals with reasons.

EXERCISE-B

1. Let G be a topological Group and H a sub-group of G . If H is a neighbourhood of the identity e in G such that $H^2 = H$, show that H is closed in G .

Solution : Take a symmetric neighbourhood V of the identity e in G such that $V^2 = V$. Let $\{x_\alpha : \alpha \in D, -\}$ is a net in H such that $\{x_\alpha\}$ converges to x in G . Now $x^{-1} \in H$ (H , a sub-group). So $(Vx^{-1}) \cap H \neq \emptyset$. Take $y \in (Vx^{-1}) \cap H$. Let $x_\alpha \in V$ for $\alpha \in D$ (say), $\alpha \in D$; then we have $yx_\alpha \in (Vx^{-1})(xV) = V^2 = V$, and hence $(yx_\alpha) \in V \cap H$. As the net $\{yx_\alpha\}$ converges to yx , and H is closed, we have $yx \in H$. Therefore $x = y^{-1}yx \in H$, showing H is closed. 2. If H is a normal sub-group of a topological group G , show that quotient Group G/H is homogenous. 3.

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Let G be a Topological Group and H a sub-group of G . If H

and G/H are locally compact, show that G is so. 4. Let G be a locally compact topological group, and C be the component of the identity e in G . Show that $C = \{H : H \text{ is any open sub-group of } G\}$. 5. Let $\{U_\alpha\}$ be the neighbourhood system of the identity e in a Topological Group G and $A \subseteq G$. Prove that $\bar{A} = \bar{\{U_\alpha A\}}$, bar denoting the closure.

54 6. Let G be a Topological Group with the identity e . Show that $e \in G$ is a normal closed sub-group of G , and hence, G/e is a Hausdorff Topological Group. 7. Prove that the component of the identity of a Topological Group is a closed Normal

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a Hausdorff Topological Group. 7. Prove that the component of the identity of a Topological Group is a closed Normal

sub-group. 8. Let R^2 be an additive topological Group and H be the st. line $y = \alpha x$ in R^2 which is a sub-group of R^2 . If $f: R^2 \rightarrow R^2/N$ be the cononical mapping where $N = \{(m, n) : m, n \text{ an integers}\}$ is a sub-group of R^2 , examine if $f(H)$ is a closed sub-group of Topological Group G/N for α to be (i) a rational number and (ii) an irrational number. 9. Let G be the additive Topological Group of all reals, and Z be the sub-group of G . show that Z is a discrete sub-group of G and factor Group G/Z is homeomorphic to a Circle. 10. Prove that topological product of two Eudidean spaces R^n and R^m is homeomorphic to the Euclidean space R^{n+m} .

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1/33	SUBMITTED TEXT	16 WORDS	96% MATCHING TEXT	16 WORDS
	Let G be a Topological Group and H be a subgroup of G . Then H			
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2/33	SUBMITTED TEXT	29 WORDS	32% MATCHING TEXT	29 WORDS
	a symmetric open neighbourhood of the identity e in topological group G and $L = \{x \in G : x^{-1} \in U\}$, then L is an open and closed (clo-open) sub-group of G . Because we have			
SA	Selvi C Chapter 6.docx (D35106226)			
3/33	SUBMITTED TEXT	11 WORDS	100% MATCHING TEXT	11 WORDS
	the family of all (Left) cosets of H in G			
SA	T-Group.pdf (D48987336)			
4/33	SUBMITTED TEXT	26 WORDS	58% MATCHING TEXT	26 WORDS
	$w \in W$ $W^{-1} = W$ (since W is symmetric, $W^{-1} = W$) So, $w \in W^{-1}$ (since $W^{-1} = W$). Therefore, $w \in W$		$w^4(w+1)^2(-w^2-1)(w^4+w^3+w^2+w+1)^{-1}(1-w)^4(w+1)^2(-w^2-1)(w^4+w^3+w^2+w+1)^{-1}$	
W	https://www.tiger-algebra.com/drill/(1-w)(1-w2)(1-w4)(1-w5)/			

5/33	SUBMITTED TEXT	23 WORDS	56% MATCHING TEXT	23 WORDS
<p>Let G be a Topological Group and H a closed normal subgroup of G, then the quotient Group G/H is a Topological Group</p>		<p>Let G be a topological group and let H a normal subgroup of G. From we know that G/H is a group.</p>		
<p>SA B.Viba Nandhini-205207145.pdf (D136277979)</p>				
6/33	SUBMITTED TEXT	36 WORDS	38% MATCHING TEXT	36 WORDS
<p>Topological space X is said to a Homogeneous space if for any two member x_1, x_2 in X there is a self Homeomorphism f in X such that $f(x_1) =$</p>				
<p>SA T-Group.pdf (D48987336)</p>				
7/33	SUBMITTED TEXT	27 WORDS	46% MATCHING TEXT	27 WORDS
<p>is homogeneous. 34 Theorem 2.2.3.(b). Let G be a Topological Group and H a sub-group of G. Then G/H is T_1 if and only if H</p>				
<p>SA MS - 334.docx (D110841764)</p>				
8/33	SUBMITTED TEXT	19 WORDS	61% MATCHING TEXT	19 WORDS
<p>H a sub-group of G. Then G/H is a discrete space if and only of H is open. Proof :</p>				
<p>SA 120004039-Project-1982444.pdf (D19454576)</p>				
9/33	SUBMITTED TEXT	25 WORDS	64% MATCHING TEXT	25 WORDS
<p>a fundamental system $\{ \mathcal{U}_\alpha \}$ of neighbourhoods of the identity e in G, then $f(\mathcal{U}_\alpha)$ is a neighbourhood of</p>				
<p>SA MS - 334.docx (D110841764)</p>				
10/33	SUBMITTED TEXT	21 WORDS	76% MATCHING TEXT	21 WORDS
<p>means C is a proper subset of $\mathbb{R}^{-1}(U)$. Since C is a maximal connected set containing $e, \mathbb{R}^{-1}($</p>				
<p>SA MS - 334.docx (D110841764)</p>				

11/33	SUBMITTED TEXT	79 WORDS	42% MATCHING TEXT	79 WORDS
<p>P and Q are open sets in G, such that $[P \cap U] \cap [Q \cap U] = \emptyset$ and neither is empty. So $U = [P \cap U] \cup [Q \cap U]$. Taking $U = UC$ Let $x \in U$ such that $x \in C$; Hence from (1) we have $x \in C = (P \cap C) \cup (Q \cap C)$ Since $x \in C$ is connected, either $x \in (P \cap C)$ or, $x \in (Q \cap C)$. Consequently, images $P \cap UC$ and $Q \cap UC$ under π are disjoint, since they are unions of cosets of C. $\pi(P \cap UC)$</p> <p>SA MS - 334.docx (D110841764)</p>				
12/33	SUBMITTED TEXT	17 WORDS	91% MATCHING TEXT	17 WORDS
<p>A Locally compact Hausdorff topological space X is normal if it is the union of an increasing</p> <p>SA MS - 334.docx (D110841764)</p>				
13/33	SUBMITTED TEXT	40 WORDS	45% MATCHING TEXT	40 WORDS
<p>$n \in \mathbb{N}$, and we write $\{x \in X \mid x \in U_n\}$ and $\{x \in X \mid x \in V_n\}$ is compact Also $X = \bigcup_{n \in \mathbb{N}} U_n = \bigcup_{n \in \mathbb{N}} V_n$</p> <p>W https://brainly.in/question/18854841</p>				
14/33	SUBMITTED TEXT	37 WORDS	53% MATCHING TEXT	37 WORDS
<p>a neighbourhood V of e such that $VV = V$. Now, $V \cap V = V$; Hence V is a closed subset of compact set and therefore V is a compact neighbourhood of e.</p> <p>SA MS - 334.docx (D110841764)</p>				
15/33	SUBMITTED TEXT	67 WORDS	33% MATCHING TEXT	67 WORDS
<p>H is a Hausdorff locally compact sub-group of G. Also $\{x \in H \mid x \in U_n\} = U_n$ for $n \in \mathbb{N}$. Hence $H = \bigcup_{n \in \mathbb{N}} U_n$. But U_n is compact, because U_n is compact. Thus H is a union (countable union) of increasing sequence of compact sets. H is normal (see theorem 2.3.1). Consider the collection $\{aH\}$ of pairwise disjoint cosets in G. Since</p> <p>SA MS - 334.docx (D110841764)</p>				

16/33	SUBMITTED TEXT	19 WORDS	94% MATCHING TEXT	19 WORDS
<p>let C be a compact subset and U an open subset of G such that $C \cap U \neq \emptyset$. Then there</p> <p>SA MS - 334.docx (D110841764)</p>				

17/33	SUBMITTED TEXT	39 WORDS	35% MATCHING TEXT	39 WORDS		
<table style="width: 100%; border: none;"> <tr> <td style="width: 50%; border: none; vertical-align: top;"> <p>Let G be a Topological Group and N is a closed Normal sub-group. (i) if G is compact, then G/N is a compact quotient Topological Group ; and (ii) if G is locally compact, then G/N is a Locally compact</p> <p>SA B.Viba Nandhini-205207145.pdf (D136277979)</p> </td> <td style="width: 50%; border: none; vertical-align: top;"> <p>Let G be a topological group and let H be a closed normal in G. The following statements hold: 1. If G is compact, then G/H is compact. 2. If G is locally compact, then G/H is locally compact. Proof: Assume that G is a compact</p> </td> </tr> </table>					<p>Let G be a Topological Group and N is a closed Normal sub-group. (i) if G is compact, then G/N is a compact quotient Topological Group ; and (ii) if G is locally compact, then G/N is a Locally compact</p> <p>SA B.Viba Nandhini-205207145.pdf (D136277979)</p>	<p>Let G be a topological group and let H be a closed normal in G. The following statements hold: 1. If G is compact, then G/H is compact. 2. If G is locally compact, then G/H is locally compact. Proof: Assume that G is a compact</p>
<p>Let G be a Topological Group and N is a closed Normal sub-group. (i) if G is compact, then G/N is a compact quotient Topological Group ; and (ii) if G is locally compact, then G/N is a Locally compact</p> <p>SA B.Viba Nandhini-205207145.pdf (D136277979)</p>	<p>Let G be a topological group and let H be a closed normal in G. The following statements hold: 1. If G is compact, then G/H is compact. 2. If G is locally compact, then G/H is locally compact. Proof: Assume that G is a compact</p>					

18/33	SUBMITTED TEXT	73 WORDS	31% MATCHING TEXT	73 WORDS
<p>open neighbourhood V of e such that $C \cap V \neq \emptyset$. Proof : Choose a symmetric open neighbourhood W_1 of the identity e in G such that $W_1^3 \subseteq V$, and for a fixed $a \in G$ take a symmetric open neighbourhood W_2 of e such that $aW_2a^{-1} \subseteq W_1$. Put $W = W_1 \cap W_2$. Now $x \in W$ gives $(xa^{-1}) \in W$; and $ax^{-1} \in W$.</p> <p>SA MS - 334.docx (D110841764)</p>				

19/33	SUBMITTED TEXT	14 WORDS	87% MATCHING TEXT	14 WORDS
<p>Connectedness : Definition 2.4.1. A topological space X is said to be connected if</p> <p>SA chapter01.pdf (D95760313)</p>				

20/33	SUBMITTED TEXT	14 WORDS	85% MATCHING TEXT	14 WORDS
<p>is connected. 4.2 A topological space X is said to be Locally Connected if each</p> <p>SA chapter01.pdf (D95760313)</p>				

21/33	SUBMITTED TEXT	25 WORDS	46% MATCHING TEXT	25 WORDS
<p>Let G be a Topological Group and H be the component of the Identity e of G. Then H is a closed Normal subgroup of G.</p> <p>SA 120004039-Project-1982444.pdf (D19454576)</p>				
22/33	SUBMITTED TEXT	21 WORDS	83% MATCHING TEXT	21 WORDS
<p>Let $\{G_i\}_{i \in I}$ be a family of Topological Groups. If $G = \prod_{i \in I} G_i$ is</p> <p>let G_i be a family of topological groups. Then, $G = \prod_{i \in I} G_i$ is</p> <p>SA B.Viba Nandhini-205207145.pdf (D136277979)</p>				
23/33	SUBMITTED TEXT	21 WORDS	70% MATCHING TEXT	21 WORDS
<p>G is a Topological Group with Group composition and product Topology. Theorem 2.4.4. Let $G = \prod_{i \in I} G_i$ be the direct product of Topological Groups $\{G_i\}_{i \in I}$</p> <p>G_i is a topological group with the product topology. Theorem 4.0.9 Let $G = \prod_{i \in I} G_i$ be a direct product of topological groups,</p> <p>SA B.Viba Nandhini-205207145.pdf (D136277979)</p>				
24/33	SUBMITTED TEXT	21 WORDS	52% MATCHING TEXT	21 WORDS
<p>H is a closed Normal sub-group in G. Example 2.4.1. Let G be a Topological Group and H be the component of</p> <p>SA Selvi C Chapter 4.docx (D35106194)</p>				
25/33	SUBMITTED TEXT	21 WORDS	57% MATCHING TEXT	21 WORDS
<p>$f(x + y) = f(x) + f(y)$, and (ii) $f(\lambda x) = \lambda f(x)$ for all $x, y \in F$ and $\lambda \in F$</p> <p>$f(x + y) = f(x) + f(y)$ for all $x, y \in F$ and $f(\lambda x) = \lambda f(x)$ for all $x \in F$ and $\lambda \in F$</p> <p>SA B.Viba Nandhini-205207145.pdf (D136277979)</p>				
26/33	SUBMITTED TEXT	28 WORDS	68% MATCHING TEXT	28 WORDS
<p>an open neighbourhood U of e such that $x \in U \implies x^{-1} \in U$. Now $P_r = \{x \in G \mid x^{-1} \in U\}$ is an open neighbourhood of</p> <p>SA SITHEESWARI (16PMAVO31).docx (D38133619)</p>				

27/33	SUBMITTED TEXT	16 WORDS	66% MATCHING TEXT	16 WORDS
<p>G and H are two Topological Groups and $f : G \rightarrow H$ is a homomorphism then (</p> <p>SA Selvi C Chapter3.docx (D35106187)</p>				
28/33	SUBMITTED TEXT	19 WORDS	67% MATCHING TEXT	19 WORDS
<p>Theorem 2.6.3. Let G and H be two topological Groups and $f : G \rightarrow H$ be a Homomorphism. Then (</p> <p>SA Selvi C Chapter3.docx (D35106187)</p>				
29/33	SUBMITTED TEXT	15 WORDS	78% MATCHING TEXT	15 WORDS
<p>is a Homomorphism. Then f is continuous if and only if f is continuous</p> <p>SA MS - 334.docx (D110841764)</p>				
30/33	SUBMITTED TEXT	30 WORDS	39% MATCHING TEXT	30 WORDS
<p>f is continuous at the identity element e of G. Conversely, suppose f is continuous at e, and $x \in G$ ($x \neq e$). Let W be a neighbourhood of $f(x)$ in H.</p> <p>SA Selvi C Chapter 6.docx (D35106226)</p>				
31/33	SUBMITTED TEXT	20 WORDS	67% MATCHING TEXT	20 WORDS
<p>Theorem 2.6.5. Let G and H be two topological Groups and $f : G \rightarrow H$ be a Homomorphism. Then</p> <p>SA Selvi C Chapter3.docx (D35106187)</p>				
32/33	SUBMITTED TEXT	19 WORDS	73% MATCHING TEXT	19 WORDS
<p>a Hausdorff Topological Group. 7. Prove that the component of the identity of a Topological Group is a closed Normal</p> <p>a Hausdorff topological group is a closed normal subgroup. The component of the identity of a topological group is a closed normal</p> <p>SA B.Viba Nandhini-205207145.pdf (D136277979)</p>				

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Let G be a Topological Group and H a sub-group of G . If
 H

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PREFACE In the auricular structure introduced by this University for students of Post- Graduate degree programme, the opportunity to pursue Post-Graduate course in Subject introduced by this University is equally available to all learners. Instead of being guided by any presumption about ability level, it would perhaps stand to reason if receptivity of a learner is judged in the course of the learning process. That would be entirely in keeping with the objectives of open education which does not believe in artificial differentiation. Keeping this in view, study materials of the Post-Graduate level in different subjects are being prepared on the basis of a well laid-out syllabus. The course structure combines the best elements in the approved syllabi of Central and State Universities in respective subjects. It has been so designed as to be upgradable with the addition of new information as well as results of fresh thinking and analysis. The accepted methodology of distance education has been followed in the preparation of these study materials. Co-operation in every form of experienced scholars is indispensable for a work of this kind. We, therefore, owe an enormous debt of gratitude to everyone whose tireless efforts went into the writing, editing and devising of a proper lay-out of the materials. Practically speaking, their role amounts to an involvement in invisible teaching. For, whoever makes use of these study materials would virtually derive the benefit of learning under their collective care without each being seen by the other. The more a learner would seriously pursue these study materials the easier it will be for him or her to reach out to larger horizons of a subject. Care has also been taken to make the language lucid and presentation attractive so mat they may be rated as quality self- learning materials. If anything remains still obscure or difficult to follow, arrangements are there to come to terms with them through the counselling sessions regularly available at the network of study centres set up by the University. Needless to add, a great deal of these efforts is still experimental-in fact, pioneering in certain areas. Naturally, there is every possibility of some lapse or deficiency here and there. However, these do admit of rectification and further improvement in due course. On the whole, therefore, these study materials are expected to evoke wider appreciation the more they receive serious attention of all concerned. Professor (Dr.) Subha Sankar Sarkar Vice-Chancellor

NETAJI SUBHAS OPEN UNIVERSITY PG (MT)–IX A(I) Unit 1 ? Analytic Continuation 7-25 Unit 2 ? Harmonic Functions 26-40 Unit 3 ? Conformal Mappings 41-49 Unit 4 ? Multi-valued Functions and Riemann Surface 50-82 Unit 5 ? Conformal Equivalence 83-104 Unit 6 ? Entire and Meromorphic Functions 105-155 N E T A J I S UBHAS OPEN UN I V E R S I T Y

7 Unit 1 Analytic Continuation Structure 1.0 Objectives of this chapter 1.1 The idea of analytic continuation 1.2 Direct analytic continuation 1.3 Analytic continuation of elementary functions 1.4 Analytic continuation by power series 1.5 Analytic continuation along a curve 1.6 Multi-valued Functions and Analytic continuation 1.0 Objectives of this Chapter

In this chapter we shall introduce the idea of direct analytic continuation of an analytic function. The concepts of analytic continuation by means of power series, complete analytic function, natural boundary, analytic continuation along a curve will be explained with the help of examples. Homotopic curves, analytic continuation of multi-valued function and Monodromy theorem will also be discussed. 1.1 The idea of analytic continuation The idea of analytic continuation rests on the notion of analytic function. A function $f(z)$ is analytic at $z = z_0$ if it is differentiable in some ϵ -neighbourhood of z_0 or, equivalently if it can be expressed in the form of a Taylor series in a neighbourhood of that point. The domain of convergence of this power series will be the region of analyticity of the function $f(z)$. Following Uniqueness Theorem : "If two functions $f(z)$ and $g(z)$, analytic on a region D , are such that $f(z) = g(z)$ on a set $A \subset D$ having a limit point in D , then $f(z) = g(z) \forall z \in D$," we know that if two analytic functions agree in some small neighbourhood of a point situated in their common region of analyticity D , they

8 coincide everywhere in D . We first introduce the idea of analytic continuation by the following examples. The geometric series $1 + z + z^2 + \dots$ converges for $|z| < 1$ and its sum function $g(z) = \frac{1}{1-z}$ is an analytic function for $|z| < 1$. The geometric series diverges for $|z| \geq 1$. However, the function $h(z) = \frac{1}{1-z}$ is analytic for all z except $z = 1$. But we observe that $h(z) = g(z)$ for $|z| < 1$. Thus, we may regard $h(z)$ as determining an analytic continuation of $g(z)$ from the domain $|z| < 1$ into the domain $\mathbb{C} \setminus \{1\}$. Example 1.1 Consider the Laplace transform of 1 in the z -plane, $F(z) = \int_0^\infty e^{-zt} dt = \frac{1}{z}$ for $\text{Re } z > 0$. We introduce a function $\phi(z) = \frac{1}{z}$ which is analytic in the complex plane \mathbb{C} except the origin. Here $\phi(z) = F(z)$ for $\text{Re } z > 0$ and we consider $\phi(z)$ as analytic continuation of $F(z)$ from the domain $\text{Re } z > 0$ into the complex plane with the point $z = 0$ deleted. We put these ideas more precisely in the following discussion. 1.2 Direct analytic continuation Let (i) $f(z)$ and $g(z)$ be analytic functions on domains D_1 and D_2 respectively. (ii) $D_1 \cap D_2 \neq \emptyset$ (iii) $f(z) = g(z)$ for all z belonging to $D_1 \cap D_2$. Then $g(z)$ is called a direct analytic continuation of $f(z)$ to D_2 , and vice versa.

9 Theorem 1.1. A direct analytic continuation, if it exists, is unique. Proof. Let $f(z)$ be an analytic function with domain of definition D_1 and let $g(z)$, another analytic function with domain of definition D_2 , be its direct analytic continuation. We shall show that $g(z)$ is unique. On the contrary suppose $\phi(z)$ be another analytic continuation of $f(z)$ into D_2 . Then $f(z) = g(z) = \phi(z)$ for all $z \in D_1 \cap D_2$. Also, $f(z) = \phi(z)$ for all $z \in D_1$

and so $\phi(z)$ coincides with $g(z)$ in $D_1 \cap D_2$. Thus we have, by the Uniqueness theorem, $\phi(z) = g(z)$ in D_2 . 1.3 Analytic continuation of elementary functions The functions $e^z, \sin z, \cos z, \sinh z$ etc are already known to us. These functions are regular in the entire complex plane. Let us assume, by definition, that $e^z = \sum_{n=0}^{\infty} \frac{z^n}{n!}$ and observe that it coincides with e^x , known earlier, for real values of z . Thus we can take e^z as the analytic continuation of e^x from real axis into the entire complex plane. Likewise introducing

$\sin z, \cos z, \sinh z, \cosh z$ in the form of power series— $\sin z = \sum_{n=0}^{\infty} \frac{(-1)^n z^{2n+1}}{(2n+1)!}, \cos z = \sum_{n=0}^{\infty} \frac{(-1)^n z^{2n}}{(2n)!}$

$$\begin{aligned} \sin z &= z - \frac{z^3}{3!} + \frac{z^5}{5!} - \frac{z^7}{7!} + \dots \\ \cos z &= 1 - \frac{z^2}{2!} + \frac{z^4}{4!} - \frac{z^6}{6!} + \dots \\ \sinh z &= z + \frac{z^3}{3!} + \frac{z^5}{5!} + \frac{z^7}{7!} + \dots \\ \cosh z &= 1 + \frac{z^2}{2!} + \frac{z^4}{4!} + \frac{z^6}{6!} + \dots \end{aligned}$$

and We

can treat them as the analytic continuation of the functions $\sin x, \cos x, \sinh x$ and $\cosh x$ respectively from the real axis into the entire complex plane. D 1 D 2 Fig. 1

denote this expansion by $h(z)$, which converges in the right-hand circle $|z - 3/4| < 5/4$; and coincides with $g(z)$ in the shaded region. We see that $h(z)$ is clearly a direct analytic continuation of $g(z)$. Let us construct another analytic continuation of $g(z)$. Now we consider a neighbourhood of the point $z = 1$ (though it is a boundary point of the unit circle the function $f(z)$ is analytic there) and obtain an expansion represented by $\phi(z) = (z-1)^{-1} - (z-1)^{-2} + \dots$ for $|z-1| < 1$. In this way we can determine all possible direct analytic continuations of $g(z)$ and then continuations of these continuations and so on. A complete analytic function is defined as consisting of the original function and the collection of all the continuations so achieved. Here the complete analytic function is $1/(z^2 + 1)$, defined in the whole complex plane barring the points $z = \pm i$.

Example 1.3 Consider the function $f(z) = z + z^2 + z^3 + \dots$ (4) Clearly this function is analytic everywhere except at $z = -1$. We take a point $z = -1/4$ inside the region of convergence of $\phi(z)$ and in a neighbourhood of this point we determine $\psi(z) = z + z^2 + z^3 + \dots$

(5) It can be checked easily that $\phi(z)$ and $\psi(z)$ are direct analytic continuation of each other. Again in the neighbourhood of $z = i/2$ we obtain an expansion $k(z) = z + z^2 + z^3 + \dots$

(6) In performing analytic continuations we notice that there are certain points which always lie on the boundary of domains in which expansions are not valid. These points are nothing but the singularities of the complete analytic function. In example 1.2 these are $z = \pm i$ whereas it is $z = -1$ for example 1.3.

Regular and Singular points Let $f(z)$ be an analytic function defined in the domain D , bounded by a simple closed curve Γ . A point $\zeta \in \Gamma$ is called a regular point of the function $f(z)$ if there exist a neighbourhood $|z - \zeta| < \epsilon$ of the point ζ and an analytic function $\phi_\zeta(z)$ such that $\phi_\zeta(z) = f(z)$ in $D \cap \{z : |z - \zeta| < \epsilon\}$. The boundary point ζ which is not a regular point is called a singular point of $f(z)$.

In any neighbourhood of the point ζ , there cannot be any analytic function coinciding with $f(z)$ in the part common to the neighbourhood of ζ and the domain D . Natural boundary In examples 1.2 and 1.3 we have encountered with finite number of singular points situated on the boundary of the region of analyticity of the given function. It might happen that the boundary is dense with singular points. In this case analytic continuation across the boundary of the region is not possible. Such a boundary is called a natural boundary. Example 1.4 Test whether analytic continuation of the function $f(z) = \sum_{n=0}^{\infty} z^{2^n}$ is possible outside its circle of convergence. Solution : Applying the ratio test we find that the given series

$f(z) = z + z^2 + z^4 + z^8 + \dots$ (7) converges for $|z| < 1$. The point $z = 1$ is a singular point of $f(z)$ as it is seen for real z that the sum $\sum_{n=0}^{\infty} z^{2^n}$ increases indefinitely as $x \rightarrow 1$. Now to test whether the circle of convergence, the unit circle, is a natural boundary we examine the behaviour of the given function at the points $z = e^{2\pi i k/s}$, $k = 0, 1, 2, \dots, s-1$ (where s is any natural number). For this sake we consider the points $z = e^{2\pi i k/s}$ and evaluate $f(z)$ at these points. Then $f(z) = \sum_{n=0}^{\infty} z^{2^n} = \sum_{n=0}^{\infty} e^{2\pi i k 2^n / s}$. The first term consists of a finite number of terms and hence bounded in absolute value, whereas the second term is absolute value reduces to $\sum_{n=2}^{\infty} r^{2^n}$. Clearly this sum increases indefinitely as $r \rightarrow 1$. This shows that the points $z = e^{2\pi i k/s}$ (as $k \rightarrow \infty$) are singular points of the given function $f(z)$. Now as $k \rightarrow \infty$ these points form an everywhere dense set of points on the boundary of the unit circle. Thus analytic continuation outside the circle of convergence of the given function is not possible. Example 1.5 Show that the function $f(z) = \sum_{n=0}^{\infty} z^{n!}$

has unit circle as its natural boundary. Theorem 1.2 Every power series has at least one singular point on its circle of convergence. Proof. Let $f(z) = a_0 + a_1(z - z_0) + a_2(z - z_0)^2 + \dots$ be any power series with region of convergence $K: |z - z_0| < R$. We shall have to prove there lies at least one singular point on the circle of convergence $\Gamma: |z - z_0| = R$ of the function. Suppose, on the contrary, that every point on Γ are regular points. Let $\zeta_1, \zeta_2, \dots, \zeta_i, \dots$ be certain number of regular points belonging to Γ and $N(\zeta_1), N(\zeta_2), \dots, N(\zeta_i), \dots$ be their neighbourhoods respectively. The points ζ_i 's are chosen in such a way that $N(\zeta_i)$ has non null intersection with $N(\zeta_{i-1})$ and $N(\zeta_{i+1})$ and the union of these neighbourhoods completely cover the boundary Γ . Let D be the union of K and all these neighbourhoods $N(\zeta_i)$. D is open since K and every $N(\zeta_i)$ are open. D is also connected since. (i) any two points lying in $K \subset D$ can be connected by a straight line segment lying in K , since K is connected. (ii) one point $z_1 \in N(\zeta_1)$ and the other $z_2 \in K$ can be connected by two straight line segments $z_1 z_1 \zeta_1$ and $\zeta_1 z_2$ lying within $N(\zeta_1) \cup K \subset D$. (iii) one point $z_m \in N(\zeta_m)$ and $z_n \in N(\zeta_n)$ can be connected by a curve consisting of $z_m z_n$ since $z_m z_n \subset D$. (Fig. 8)

16 and finally if two points lie in the same neighbourhood $N(\zeta_i)$ it is always connected by a curve $\gamma \subset N(\zeta_i) \subset D$. Now we introduce an analytic function $\psi(z)$ on the open connected set D which satisfies $\psi(z) = \varphi(\zeta_i(z))$, $z \in N(\zeta_i)$ $f(z)$, $z \in K$ where $\varphi(\zeta_i(z))$ is a direct analytic continuation of $f(z)$ in the neighbourhood $N(\zeta_i)$ of the regular point ζ_i . We now prove that $\psi(z)$ is well-defined on D . Let α, β be any two points on Γ such that $H \cap N = \emptyset$ and since α, β are regular points there exist functions $\varphi_\alpha(z)$ and $\varphi_\beta(z)$ as direct analytic continuations of $f(z)$ in $N(\alpha)$ and $N(\beta)$ respectively i.e. $\varphi_\alpha(z) = \varphi_\beta(z)$ in $H \cap G$. Now since $\varphi_\alpha(z), \varphi_\beta(z)$ are analytic in H and G is a part of H , by the uniqueness theorem $\varphi_\alpha(z) \equiv \varphi_\beta(z) \forall z \in H$. As α and β are arbitrary points of Γ we conclude that $\psi(z)$ is a well-defined analytic function on D . Let C be the boundary of D and let $\rho = \min_{z \in C} |z - z_0|$, C be the minimum distance from z_0 to the boundary C of D . Then clearly $\rho < R$ as ζ lies outside the circle Γ . Thus we observe that $\psi(z)$ coincides with $f(z)$ on the disc $|z - z_0| < \rho$. Then it is obvious to conclude that the radius of convergence of the given power series $\sum a_n(z - z_0)^n$ is ρ , not R , which is a contradiction. Hence every point on Γ cannot be regular points, i.e., there must be at least one singular point on Γ .

1.5 Analytic continuation along a curve Earlier, analytic continuation by power series method, we have extended $f(z)$ to a

17 larger domain considering its power series expansion about a point a from its original circle of convergence with centre at z_0 ($-a \neq z_0$) and radius r . We know, this power series converges in the disc $D_1: |z - a| < R$, where $R \geq r - |z_0 - a|$ [(see Fig. 9), for example 1.2]. Then it converges to an analytic function $g(z)$ defined on D_1 , which is equal to $f(z)$ on $D \cap D_1$. Analytic continuation along a curve is an extension of this idea to the situation where a curve is covered by an overlapping sequence of discs and an analytic function defined on the first disc, can be extended successively to each disc in the sequence (see figure 10). We will make this idea more precise after introducing the definition of function element. Definition 1. An ordered pair (f, D) , where D is a region and f is an analytic function on D is called a function element. We say that it is a function element at z_0 if z_0 belongs to D . Two function elements (φ, G) and (ψ, H) are equal if and only if $\varphi(z) \equiv \psi(z)$, $G = H$. Clearly a function element (f_1, D_1) is a direct analytic continuation of another function element (f_2, D_2) when $D_1 \cap D_2 \neq \emptyset$ and $f_1 = f_2$ in $D_1 \cap D_2$. In this case the two function elements (f_1, D_1) and (f_2, D_2) are said to be equivalent. Definition 2. Let $\gamma: [0, 1] \rightarrow \mathbb{C}$ be a curve and (f_0, D_0) be a function element at $z_0 = \gamma(0)$. Suppose there exists (i) a partition $0 = t_0 < t_1 < \dots < t_n = 1$ of $[0, 1]$ and (ii) a finite sequence of function elements $(f_0, D_0), (f_1, D_1), \dots, (f_n, D_n)$ with $\gamma([t_j, t_{j+1}]) \subset D_j$ and (iii) $f_j(z) = f_{j+1}(z)$ on $D_j \cap D_{j+1}$ for $j = 0, 1, \dots, n-1$. Then (f_n, D_n) is called an analytic continuation of (f_0, D_0) along γ . Apparently, it seems that the function element (f_n, D_n) of the above definition, depends on the choice of partition $0 = t_0 < t_1 < \dots < t_n = 1$ of $[0, 1]$ and the finite sequence $(f_0, D_0), (f_1, D_1), \dots, (f_n, D_n)$ of function elements. It turns out that up to equivalence, it is actually independent of these choices. Fig. 9 Fig. 10

at z_0 and R be a simply connected region containing D_0 , ζ be any point lying in R . We suppose (i) (f_0, D_0) can be analytically continued along every curve in R . (ii) γ_0 and γ_1 are homotopic curves from z_0 to ζ . Then the continuations of the function element (f_0, D_0) along γ_0 and γ_1 at ζ are equivalent. Fig. 12 Fig. 13

24 Proof. A homotopy from γ_0 to γ_1 determines a continuous one parameter family of curves $\{\gamma_s\}$, $0 \leq s \leq 1$ from z_0 to ζ given by the equations $z = \sigma_s(t)$, $0 \leq t \leq 1$. By hypothesis, the function element (f_0, D_0) has an analytic continuation along each of the curves, γ_s . Denote the terminal function element at ζ for the continuation along γ_s by ϕ_s . We claim that, for each $k \in [0, 1]$, there is a $\delta > 0$ such that ϕ_s is equivalent to ϕ_k whenever $|s-k| < \delta$. Let $0 = t_0 < t_1 < \dots < t_n = 1$ be a partition and $(f_0, D_0), (f_1, D_1), \dots, (f_n, D_n)$ be a finite sequence of function elements defining $\phi_k = (f_n, D_n)$ as the terminal function element at ζ for the analytic continuation of (f_0, D_0) along γ_k . Then $E_j = \sigma_k([t_j, t_{j+1}]) \subset D_j$ for $j = 0, 1, \dots, n-1$. For each $j = 0, 1, \dots, n-1$, let ϵ_j be the minimum distance from the compact set E_j to the boundary of the D_j . If $|\sigma_s(t) - \sigma_k(t)| < \epsilon_j$, $t \in [0, 1]$, then it will also be true that $\sigma_s([t_j, t_{j+1}]) \subset D_j$. Thus, if $\epsilon = \min\{\epsilon_0, \epsilon_1, \dots, \epsilon_{n-1}\}$ and we choose $\delta > 0$ such that $|\sigma_s(t) - \sigma_k(t)| < \epsilon$ whenever $|s-k| < \delta$, then for each s with $|s-k| < \delta$, the partition $0 = t_0 < t_1 < \dots < t_n = 1$ and sequence of function elements $(f_0, D_0), (f_1, D_1), \dots, (f_n, D_n)$ also defines (f_n, D_n) as the terminal function element at ζ for the analytic continuation of (f_0, D_0) along γ_s . Since, by the previous theorem 1.3, any other continuation of (f_0, D_0) along γ_s results function element equivalent to this one, we conclude that ϕ_k is equivalent to ϕ_s . This proves that ϕ_s is equivalent to ϕ_k whenever $|s-k| < \delta$. This means that for every $s \in I = [0, 1]$ there is a positive $\delta(s)$ such that if s lies in the interval $I_s = (s-\delta(s), s+\delta(s))$, then the analytic continuation of $f_0(z)$ along all such curves γ_s , result equivalent function elements at the point ζ . Now by the Heine-Borel theorem, we can always choose a finite number of intervals I_{s_j} , $0 = s_0 < s_1 < \dots < s_n = 1$ that cover the segment I and are such that the intervals I_{s_j} and $I_{s_{j+1}}$, $0 \leq j \leq n-1$ have a non-empty intersection. Then, if $s \in I_{s_0} \cap I_{s_1}$, the analytic continuation of $f_0(z)$ result equivalent function elements at the point ζ . The same is true for $s \in I_{s_1} \cap I_{s_2}$ and so on. Continuing in this way we observe that the analytic continuation of the function element (f_0, D_0) along all the curves γ_s , $0 \leq s \leq 1$ produce equivalent function elements at the point ζ . This completes the proof of the theorem. The above theorem leads us to the following most important corollary. Corollary. Let R be a simply connected region and (i) (f_0, D_0) be a function element at z_0 belonging to R (ii) (f_0, D_0) admit analytic continuation along every curve in R . Then there is a function F which is analytic on R and coincides with f_0 on D_0 . Proof. Let z_1 be a point in R . Then, since R is simply connected any two curves from z_0 , to z_1 are homotopic in R . The Monodromy theorem implies that any two terminal function elements of analytic continuations of (f_0, D_0) along curves from z_0 to z_1 in R will be equivalent and hence, will determine a function F_1 analytic in some neighbourhood of z_1 , say Q_1 . Clearly, $F_1(z) = f_0(z)$ on D_0 , $F_1(z) = f_1(z)$ on D_1 , ..., etc for the continuation along the curve γ_1 from z_0 to z_1 . Again let z_2 be a point in R , and γ_2 be a curve in R joining z_0 to z_2 and let (g_n, E_n) be the function element at z_2 continuing along the curve γ_2 with $f_0 = g_0$ on $D_0 = E_0$. We simply join z_2 to z_1 by a curve γ and claim that continuation of (F_1, Q_1) , along the curve γ to z_2 , will be equivalent to (g_n, E_n) (since the curves $\gamma_1 \cup \gamma$ and γ_2 are homotopic), which gives rise to the fact that there is a function F_2 analytic in some neighbourhood of z_2 , say Q_2 , which coincides with F_1 on Q_1 . Clearly, $F_2(z)$ possesses larger domain of analyticity than $F_1(z)$. Proceeding in this way finite number of times we can achieve a function F analytic throughout the region R .

26 Unit 2 Harmonic Functions Structure 2.0 Objectives 2.1 Harmonic Function 2.2 Gauss' Mean Value Theorem for harmonic 2.3

Inverse point of a given point with respect to a

circle 2.4 The Dirichlet Problem 2.5 Subharmonic & Superharmonic Functions 2.0 Objectives In this chapter we shall

mainly study harmonic functions and their basic properties. Gauss' mean value theorem, Poisson's integral formula, Dirichlet's problem for a disc and Harnack inequality for harmonic functions will be discussed. Subharmonic and super harmonic functions will be explained through examples. 2.1 Harmonic Function A function $u(x, y)$ of two real variables x and y defined in an open set D is said to be harmonic in D if it has continuous derivatives of the second order and satisfies the equation $\Delta u = 0$ known as Laplace's equation. The differential operator $\Delta = \frac{\partial^2}{\partial x^2} + \frac{\partial^2}{\partial y^2}$ is called the Laplacian and is denoted by ∇^2 . We introduce the differential operators $\frac{\partial}{\partial x} = \frac{\partial}{\partial x} + i \frac{\partial}{\partial y}$ and $\frac{\partial}{\partial y} = -i \frac{\partial}{\partial x} + \frac{\partial}{\partial y}$ and (17)

27 in order to achieve a condition equivalent to (16) for

f(z). If

we write x

z

$z = x + iy$ and $y = -ix + z$ (18) then $\frac{\partial^2 f}{\partial x^2} + \frac{\partial^2 f}{\partial y^2} = 0$

$f(x, y)$

z

$\frac{\partial f}{\partial x} + i \frac{\partial f}{\partial y} = 2 \frac{\partial f}{\partial z}$

(19a)

b) $\frac{\partial^2 f}{\partial x^2} + \frac{\partial^2 f}{\partial y^2} = 0$

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

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

x

y

x

$\frac{\partial f}{\partial x} + i \frac{\partial f}{\partial y} = 2 \frac{\partial f}{\partial z}$

$- \frac{\partial f}{\partial x} + i \frac{\partial f}{\partial y} = 2 \frac{\partial f}{\partial \bar{z}}$

$\frac{\partial f}{\partial x} + i \frac{\partial f}{\partial y} = 2 \frac{\partial f}{\partial z}$

f

if $f = u + iv$ and consequently the condition equivalent to (16) is $\nabla^2 f = 0$ (20) A function $f(z)$ is said to be harmonic in D if f has continuous second derivatives in D and satisfies $\nabla^2 f = 0$, $z \in D$ (21) Result 1 : If $f = u + iv$ is analytic in a domain D , then $\frac{\partial^2 u}{\partial x^2} + \frac{\partial^2 u}{\partial y^2} = 0$ and $\frac{\partial^2 v}{\partial x^2} + \frac{\partial^2 v}{\partial y^2} = 0$. Proof : u and v satisfy the Cauchy-Riemann equations and using (19b) we have, $\frac{\partial^2 f}{\partial x^2} + \frac{\partial^2 f}{\partial y^2} = 0$, u and v satisfy the Cauchy-Riemann equations and using C-R equations = 0

The

real and imaginary parts of an analytic function are harmonic. Proof :

Let $f = u + iv$ be analytic in a domain D .

By

Cauchy-Riemann equations $u_x = v_y$ and $u_y = -v_x$

i.e. u

$u_{xx} = v_{yy}$

$u_{xy} = -v_{yx}$

$u_{xy} = -v_{yx}$ [since $v_{xy} = v_{yx}$, partial derivatives being continuous] and on addition it proves that u is harmonic in D . Likewise v is also harmonic in D . Harmonic conjugates :

Let $u(x, y)$ and $v(x, y)$ be two harmonic functions in a domain $D \subset \mathbb{C}$.

28 If they

satisfy

the Cauchy-Riemann equations : $u_x = v_y$ and $u_y = -v_x$

$u_x = v_y$ and $u_y = -v_x$

in D , then we say that v is a harmonic conjugate of u .

It follows

that

$f(z) = u(x, y) + i v(x, y)$ is analytic in

a domain

D if and only if

$v(x, y)$

$v(x, y)$

is a harmonic conjugate of $u(x, y)$

in D. Remark : We know that the real part as well as the imaginary part of an analytic function are harmonic. Now the questions arise : 1. Can any real harmonic function be the real part of an analytic function? 2. Whether every real harmonic function has a harmonic conjugate? Existence of Harmonic conjugates Theorem 2.1 Let $u(x, y)$ be a real-valued harmonic function in a simply connected domain $D \subset \mathbb{C}$. Then there is an analytic function f in D such that $u = \operatorname{Re} f$ (or, equivalently there is a function v , a harmonic conjugate of u) which is unique to within addition of an arbitrary real constant. Proof. Since the function $u(x, y)$ is harmonic in a simply connected domain D ,

we have $\frac{\partial^2 u}{\partial x^2} + \frac{\partial^2 u}{\partial y^2} = 0$ which can be rewritten as $\frac{\partial}{\partial x} \left(\frac{\partial u}{\partial y} \right) - \frac{\partial}{\partial y} \left(\frac{\partial u}{\partial x} \right) = 0$, where $\frac{\partial u}{\partial x}$ and $\frac{\partial u}{\partial y}$ are given functions with continuous first partial derivatives. This implies that $-\frac{\partial u}{\partial x} dx + \frac{\partial u}{\partial y} dy$ is exact. So there is a single-valued function $v(x, y)$ which is unique to within an additive arbitrary constant, i.e. $\int_{(x_0, y_0)}^{(x, y)} -\frac{\partial u}{\partial x} dx + \frac{\partial u}{\partial y} dy = v(x, y) - v(x_0, y_0)$ (22) $K \equiv$ real constant, where (x_0, y_0) is an initial point and (x, y) is any variable point lying in D and the integral on the curve connecting (x_0, y_0) to (x, y) is path independent. From (22) we find that $\frac{\partial v}{\partial x} = -\frac{\partial u}{\partial x}$ and $\frac{\partial v}{\partial y} = \frac{\partial u}{\partial y}$,

29 which in turn ensures that $v(x, y)$ is harmonic in D and harmonic conjugate to $u(x, y)$ i.e. $f = u + iv$ forms an analytic function in D . Observation : If D is multiply connected then the integral in (22) may take different values for different paths connecting (x_0, y_0) , to (x, y) giving $v(x, y)$ as a multi-valued function, unless the paths are restricted to a simply connected sub domain contained in D . Example 1. Let D be the whole plane cut along the negative real axis including the origin ($y = 0, x \leq 0$).

Show

that $u(x, y) = \sin x \cosh y$ is harmonic in D , and find its harmonic conjugate.

Also find the corresponding analytic function. Solution : Here $u(x,$

$y)$ possesses continuous second order partial derivatives in D and also satisfies the Laplace equation : $u_{xx} + u_{yy} = 0$.

Hence $u(x, y)$ is harmonic in D . Let $v(x, y)$ be its harmonic conjugate. Then according to the formula (22), we have $\int_{(1, 0)}^{(x, y)} -\frac{\partial u}{\partial x} dx + \frac{\partial u}{\partial y} dy = v(x, y) - v(1, 0)$ real constant, where $M(1, 0)$ is the initial point. Here,

$u(x,$

$y) = \sin x \cosh y$ $u_x = \cos x \cosh y$ $u_y = \sin x \sinh y$

Now let the point $Q(x, y)$ lie in the 1st quadrant of the right-half plane. Then integrating along MNQ , we find that $\int_{(1, 0)}^{(x, y)} -\frac{\partial u}{\partial x} dx + \frac{\partial u}{\partial y} dy = -\int_1^x \cos t \cosh y dt + \int_0^y \sin x \sinh t dt = -\cos x \cosh y + \sin x \sinh y + K$ Again, if the point (x, y) lies in the 2nd quadrant of the left-half plane, then we obtain

$\int_{(1, 0)}^{(x, y)} -\frac{\partial u}{\partial x} dx + \frac{\partial u}{\partial y} dy = -\int_1^x \cos t \cosh y dt + \int_0^y \sin x \sinh t dt = 0$

y

$\cos 1 \cosh$

$y dy + 1 \times -\sin x \sinh y dx + K = \cos 1 \sinh y + \cos x \sinh y - \cos 1 \sinh y + K = \cos x$

$\sinh y +$

K The

expression for $v(x, y)$ in both the cases turns out to be the same apart from an additive constant. It results from the fact that the two paths in determining the Fig. 15 $N O M(1, 0) Q(x, y) Q(x, y) N$

30 integral lie in a simply connected domain. Thus, $v(x, y) = \cos x \sinh y + K$ at all points of D . Therefore, an analytic function with the given real part will be of the form

$f(z) = \sin x \cosh y + i \cos x \sinh y +$

$iK, K \equiv$ real constant $= \sin(x + iy) + iK = \sin z + iK$ As for uniqueness, if two analytic functions in D have the same real part, then their difference has derivative zero, by the Cauchy-Riemann equations. In that case the functions differ by a constant.

2.2 Gauss' Mean Value Theorem for harmonic functions Let $u(z) = u(x, y), z = x + iy$, be harmonic in the disk $K : |z - z_0| \leq R$ and continuous on the closed disk K . Then $u(z_0) = \frac{1}{2\pi} \int_0^{2\pi} u(z_0 + re^{i\theta}) d\theta$ (23) Proof. Let $f(z)$ be an analytic function defined in K such that $\operatorname{Re} f(z) = u(z)$. It follows from Cauchy's integral formula that $f(z) = \frac{1}{2\pi i} \int_{|z-z_0|=r} \frac{f(\zeta) d\zeta}{\zeta - z}$ $\int_0^{2\pi} \frac{f(z_0 + re^{i\theta}) d\theta}{e^{i\theta} - \frac{z-z_0}{r}}$ $= \frac{1}{2\pi} \int_0^{2\pi} f(z_0 + re^{i\theta}) \frac{d\theta}{1 - \frac{z-z_0}{r} e^{i\theta}}$ $= \frac{1}{2\pi} \int_0^{2\pi} f(z_0 + re^{i\theta}) d\theta + \frac{z-z_0}{r} \int_0^{2\pi} f(z_0 + re^{i\theta}) e^{i\theta} d\theta + \dots$ $= \frac{1}{2\pi} \int_0^{2\pi} f(z_0 + re^{i\theta}) d\theta + \frac{z-z_0}{r} \int_0^{2\pi} f(z_0 + re^{i\theta}) e^{i\theta} d\theta + \dots$ Equating the real parts, we obtain $u(z_0) = \frac{1}{2\pi} \int_0^{2\pi} u(z_0 + re^{i\theta}) d\theta + \frac{z-z_0}{r} \int_0^{2\pi} u(z_0 + re^{i\theta}) e^{i\theta} d\theta + \dots$ whence taking the limit $r \rightarrow R$, we obtain the desired result (23) 2.3

Inverse point of a given point with respect to a

circle Let $\gamma : |z - \alpha| = R$ and z_0 be a given point. Let z_1 be another point on the radius through z_0 such that $|z_0 - \alpha| |z_1 - \alpha| = R^2$. Then either of the points z_0 and z_1 is called the inverse point of the other with respect to γ . The centre of the circle γ is called the centre of inversion. It follows from the definition that (i) if z_0 lies inside γ , then z_1 must lie outside

γ , (ii) if z_0 lies on γ , then z_1 must also lie on γ and it coincides with z_0 , (iii) if z_0 lies outside γ , then z_1 must lie inside γ . Every point, except the centre of the circle, on the plane has a unique inverse point with respect to the circle. We

associate the point at infinity to the inverse point of the centre. Result : Let $\gamma : |z| = R$ and z_0 be a given point. Then the inverse point of z_0 with respect to γ is given by R^2 / \bar{z}_0 . Proof : Let $z_0 = re^{i\theta}$. Then its inverse point with respect to γ is given by $z_1 = r_1 e^{i\theta}$, where $rr_1 = R^2$. Hence $r_1 = R^2 / r$ and so $z_1 = R^2 / \bar{z}_0$. Poisson's

integral formula : Theorem : Let $u(x, y)$ be a harmonic function in a simply connected region D and $\gamma : |\zeta| = R$ be a circle contained in D . Then for any $z = re^{i\theta}$, $r < R$, u can be written as $u(r, \theta) = \frac{1}{2\pi} \int_0^{2\pi} u(R, \phi) \frac{R^2 - r^2}{R^2 - 2Rr \cos(\theta - \phi) + r^2} d\phi$. Proof : Since $u(x, y)$ is harmonic in D , there exists a conjugate harmonic function

$v(x, y)$ in D so

that

$f(z) = u(x, y) + iv(x, y)$ is analytic

in

D .

Then

$f(z)$ is analytic

within and on γ and so for any z within γ , by Cauchy's integral formula, $f(z) = \frac{1}{2\pi i} \int_{\gamma} \frac{f(\zeta) d\zeta}{\zeta - z}$. The inverse point of z with respect to γ lies outside γ and is given by R^2 / \bar{z} . Hence by Cauchy-Goursat theorem, $0 = \frac{1}{2\pi i} \int_{\gamma} \frac{f(\zeta) d\zeta}{\zeta - z} - \frac{1}{2\pi i} \int_{\gamma} \frac{f(\zeta) d\zeta}{\zeta - R^2/\bar{z}}$

Subtracting (25) from (24) we get, $f(z) - \frac{R^2}{\bar{z}} f\left(\frac{R^2}{\bar{z}}\right) = -\frac{1}{2\pi i} \int_{\gamma} f(\zeta) \left[\frac{1}{\zeta - z} - \frac{1}{\zeta - R^2/\bar{z}} \right] d\zeta$

$= -\frac{1}{2\pi i} \int_{\gamma} f(\zeta) \frac{R^2 - z\bar{\zeta}}{(\zeta - z)(\zeta - R^2/\bar{z})} d\zeta$

Let $\zeta = Re^{i\phi}$. Also, $z = re^{i\theta}$. Then (26) becomes $f(r, \theta) - \frac{R^2}{r} f\left(\frac{R^2}{r}, \theta\right) = -\frac{1}{2\pi} \int_0^{2\pi} f(R, \phi) \frac{R^2 - z\bar{\zeta}}{R^2 - 2Rr \cos(\theta - \phi) + r^2} d\phi$

$= -\frac{1}{2\pi} \int_0^{2\pi} f(R, \phi) \frac{R^2 - z\bar{\zeta}}{R^2 - 2Rr \cos(\theta - \phi) + r^2} d\phi$

$= -\frac{1}{2\pi} \int_0^{2\pi} f(R, \phi) \frac{R^2 - z\bar{\zeta}}{R^2 - 2Rr \cos(\theta - \phi) + r^2} d\phi$

Let $f(re^{i\theta}) = u(r, \theta) + iv(r, \theta)$. Then (27) becomes $u(r, \theta) - \frac{R^2}{r} u\left(\frac{R^2}{r}, \theta\right) = -\frac{1}{2\pi} \int_0^{2\pi} u(R, \phi) \frac{R^2 - z\bar{\zeta}}{R^2 - 2Rr \cos(\theta - \phi) + r^2} d\phi$

$u(r, \theta) - \frac{R^2}{r} u\left(\frac{R^2}{r}, \theta\right) = -\frac{1}{2\pi} \int_0^{2\pi} u(R, \phi) \frac{R^2 - z\bar{\zeta}}{R^2 - 2Rr \cos(\theta - \phi) + r^2} d\phi$

Equating real parts in (28) we get, $u(r, \theta) - \frac{R^2}{r} u\left(\frac{R^2}{r}, \theta\right) = -\frac{1}{2\pi} \int_0^{2\pi} u(R, \phi) \frac{R^2 - z\bar{\zeta}}{R^2 - 2Rr \cos(\theta - \phi) + r^2} d\phi$

Formula (29) is known as Poisson's integral formula. Note : Let $Rr < R^2$. Then, the function $P(R, r, \theta - \phi)$ is called the Poisson Kernel.

Hence we can write (29) in the form $u(r, \theta) - \frac{R^2}{r} u\left(\frac{R^2}{r}, \theta\right) = -\frac{1}{2\pi} \int_0^{2\pi} P(R, r, \theta - \phi) u(R, \phi) d\phi$

$u(r, \theta) - \frac{R^2}{r} u\left(\frac{R^2}{r}, \theta\right) = -\frac{1}{2\pi} \int_0^{2\pi} P(R, r, \theta - \phi) u(R, \phi) d\phi$

We can also get a formula similar to (29) for the imaginary part of $f(z)$ by equating the imaginary part in (28). The

corresponding formula is $v(r, \theta) - \frac{R^2}{r} v\left(\frac{R^2}{r}, \theta\right) = -\frac{1}{2\pi} \int_0^{2\pi} P(R, r, \theta - \phi) v(R, \phi) d\phi$

$v(r, \theta) - \frac{R^2}{r} v\left(\frac{R^2}{r}, \theta\right) = -\frac{1}{2\pi} \int_0^{2\pi} P(R, r, \theta - \phi) v(R, \phi) d\phi$

Remark :

Cauchy's integral formula expresses the values of an analytic function inside a circle

in terms of its values on the

boundary of

the

circle whereas Poisson's

integral formula expresses the values of a harmonic function inside a circle

in terms of its values on the

boundary of the

circle. Result 3. $\frac{1}{2\pi} \int_0^{2\pi} P(R, r, \theta - \phi) d\phi = 1$. Proof : By Poisson's integral formula we have, $u(r, \theta) - \frac{R^2}{r} u\left(\frac{R^2}{r}, \theta\right) = -\frac{1}{2\pi} \int_0^{2\pi} P(R, r, \theta - \phi) u(R, \phi) d\phi$

Taking $u(r, \theta) \equiv 1$ we get, $1 - \frac{R^2}{r} = -\frac{1}{2\pi} \int_0^{2\pi} P(R, r, \theta - \phi) d\phi$

Result 4. $P(R, r, \theta) = \frac{R^2 - r^2}{R^2 - 2Rr \cos(\theta - \phi) + r^2}$

Proof : Let $\zeta = Re^{i\phi}$, $z = re^{i\theta}$, $r < R$. Then, $\zeta + z = R e^{i\phi} + r e^{i\theta} = R \left(\cos \phi + i \sin \phi \right) + r \left(\cos \theta + i \sin \theta \right) = (R \cos \phi + r \cos \theta) + i (R \sin \phi + r \sin \theta)$

$\zeta - z = R e^{i\phi} - r e^{i\theta} = R \left(\cos \phi + i \sin \phi \right) - r \left(\cos \theta + i \sin \theta \right) = (R \cos \phi - r \cos \theta) + i (R \sin \phi - r \sin \theta)$

$\zeta \bar{z} = R r e^{i(\phi - \theta)} = R r (\cos(\phi - \theta) + i \sin(\phi - \theta))$

$\zeta \bar{\zeta} = R^2$

$r \cos \theta$

Simplifying we get, $\operatorname{Re} \cos(z) = \frac{e^{z} + e^{-z}}{2}$. Result 5. Poisson Kernel $P(R, r, \varphi - \theta)$

is harmonic in $|z| < R$. Proof: Let $f(z) = \frac{e^{z} + e^{-z}}{2}$. Then $f(z)$ is analytic in $|z| < R$. By result 4, $P(R, r, \varphi - \theta) = \operatorname{Re} f(z)$. Hence the Poisson Kernel is the real part of an analytic function. Hence $P(R, r, \varphi - \theta)$ is harmonic in $|z| < R$. Note: We can easily show that $\operatorname{Re} \frac{e^{z} + e^{-z}}{2} = \frac{e^{\operatorname{Re} z} \cos \operatorname{Im} z + e^{-\operatorname{Re} z} \cos \operatorname{Im} z}{2} = \cos(\operatorname{Im} z)$

where $z = re^{i\theta}$, $r < R$. Hence $\operatorname{Re} \frac{e^{z} + e^{-z}}{2} = \cos(\theta)$ and Poisson's integral formula (29) can be written as $u(r, \theta) = \frac{1}{2\pi} \int_0^{2\pi} P(R, r, \varphi - \theta) u(R, \varphi) d\varphi$ (32) The function $P(R, r, \varphi - \theta)$ is the Poisson Kernel. Theorem 2.2 Let $u(x, y) \neq \text{constant}$ be harmonic on a simply connected domain D .

Then $u(x, y)$ has neither a maximum nor a minimum at any point of D . Proof. Let $z_0 = x_0 + iy_0$ be an arbitrary point of D . Then following theorem 2.1 there is an analytic function $f(z)$ in a neighbourhood $N(z_0)$ of z_0 such that $\operatorname{Re} f = u$. Then $g(z) = e^{f(z)}$ is analytic on $N(z_0)$ and not equal to constant since $u(x, y) \neq \text{constant}$ and $|g(z)| = e^{u(x,y)}$. Again exponential function is strictly increasing, so a maximum for u at (x_0, y_0) is also a maximum for e^u , and hence also a maximum of $|g(z)|$ i.e. of $|g(z)|$ at z_0 .

The function $u(x, y)$ cannot have a maximum at (x_0, y_0) , since otherwise $|g(z)|$ would have a maximum at z_0 , thereby contradicting the maximum modulus principle. Likewise, following the minimum modulus principle $|g(z)|$ cannot have a minimum value at z_0 since $|g(z)| \neq 0$ on D . Therefore $u(x, y)$ cannot possess minimum value at (x_0, y_0) . Corollary. Let $u(x, y)$ be harmonic on a domain D and continuous on \bar{D} . Then $u(x, y)$ attains its maximum and

its minimum on the boundary of D . Proof. Since $u(x, y)$ is continuous on the compact set \bar{D} , it attains both its maximum and its minimum on \bar{D} , but $u(x, y)$ cannot possess a maximum or a minimum at a point of D . Therefore the corollary follows. Example 2. Given $u(x, y)$ harmonic in the disk $|z| < R$ and $A(r, j)$ its maximum value on the circle $|z| = r$, $r, j < R$, $j = 1, 2, 3$. Prove that $A(r, 1) \log r + \alpha \log r \leq A(r, 2) \log r + \alpha \log r \leq A(r, 3) \log r + \alpha \log r$ for $0 < r < R$. Solution. Since $u(x, y)$ is harmonic in $|z| < R$, $u(x, y) + \alpha \log r$ is also harmonic in the annulus $r_1 < |z| < r_2$. Hence its

maximum is attained on the boundary of the annulus i.e. on $|z| = r_1$ or $|z| = r_2$ or, on both. Either $A(r_1) + \alpha \log r_1$ or $A(r_2) + \alpha \log r_2$ is maximum. We define α so that $A(r_1) + \alpha \log r_1 = A(r_2) + \alpha \log r_2$ or, $\alpha = \frac{A(r_2) - A(r_1)}{\log r_2 - \log r_1}$. The circle $|z| = r$ lies inside the annulus $r_1 < |z| < r_2$ and according to corollary of the theorem 2.2 regarding maximum value of the harmonic function $u(x, y) + \alpha \log r$ we have $A(r) + \alpha \log r \leq A(r_2) + \alpha \log r_2$ or, $A(r) \leq A(r_2) + \alpha(\log r_2 - \log r)$.

2.4 The Dirichlet Problem Let D be a domain with boundary Γ and let (x, y) be a continuous real function defined on Γ . The Dirichlet problem is to find a function $u(x, y)$, harmonic on D and continuous on \bar{D} , which coincides with (x, y) at every point of Γ . Existence of a solution of Dirichlet's problem for a disc Theorem 2.3 Let D be the disc $|z| < R$ with boundary $\Gamma : |z| = R$ and let $U(\varphi)$ be a continuous real function on the interval $[0, 2\pi]$ such that $U(0) = U(2\pi)$. Then the function $u(r, \theta)$ defined by the integral $u(r, \theta) = \frac{1}{2\pi} \int_0^{2\pi} P(R, r, \varphi - \theta) U(\varphi) d\varphi$ (33) for any point (r, θ) on D and $u(R, \varphi) = U(\varphi)$ (34) for any point (R, φ) on Γ , solves the Dirichlet problem for the disc D . In other words, (i) u is harmonic on D and continuous on \bar{D} and (ii) $\lim_{(x,y) \rightarrow (x_0,y_0)} u(x,y) = u(x_0,y_0)$ where (x_0,y_0) is any fixed point on Γ . Proof: To prove that $u(r, \theta)$ defined by (33) on D is harmonic on D we observe that

$\operatorname{Re} \frac{e^{z} + e^{-z}}{2} = \frac{e^{\operatorname{Re} z} \cos \operatorname{Im} z + e^{-\operatorname{Re} z} \cos \operatorname{Im} z}{2} = \cos(\operatorname{Im} z)$ where $P(R, r, \varphi - \theta)$ is the Poisson Kernel

and $\zeta = Re^{i\varphi}$, $z = re^{i\theta}$, $r < R$. The r.h.s. is the real part of the function $\frac{e^{\zeta} + e^{-z}}{2}$ which is analytic in D . Hence the Poisson Kernel $P(R, r, \varphi - \theta)$ is harmonic in D . So, differentiation under the sign of integration is valid. Applying the Laplacian ∇^2 in (r, θ) to both sides of (33) we get, $\nabla^2 u = \nabla^2 \frac{1}{2\pi} \int_0^{2\pi} P(R, r, \varphi - \theta) U(\varphi) d\varphi$

$P(R, r, \varphi - \theta) U(\varphi)$ [Since $P(R, r, \varphi - \theta)$ is harmonic in $D \Rightarrow \nabla^2 P(R, r, \varphi - \theta) = 0$]. \Rightarrow

u is harmonic on D . Next we prove that the function $u(r, \theta)$ defined by the integral (33) approaches $U(\varphi_0)$ as the point (r, θ) in D tends to any fixed point (R, φ_0) on Γ . Let (r_n, θ_n) be an arbitrary sequence of points in D converging to the boundary point (R, φ_0) . We now consider the difference $u(r_n, \theta_n) - U(\varphi_0)$. (35) Since $U(\varphi)$ is continuous on Γ , for given $\epsilon > 0$ there exists a $\delta(\epsilon) > 0$ such that $|U(\varphi) - U(\varphi_0)| < \epsilon/2$ (36) whenever $|\varphi - \varphi_0| < \delta$; (37) we choose δ so small that (36) is satisfied and $\varphi_0 - 2\delta < \varphi_0 + 2\delta < 2\pi$. We break the integral on r.h.s. of (35) as $u(r_n, \theta_n) - U(\varphi_0) = \dots$ (38)

37 Now, $| \dots | \leq \dots + \pi \varphi \theta \varphi \varphi \varphi \delta \varphi \delta \epsilon \pi |, \dots$ (39) To estimate the other two terms we choose n so large that $|\varphi_0 - \theta_n| < \delta$. Then, $|\varphi - \theta_n| = |\varphi - \varphi_0 + \varphi_0 - \theta_n| \geq |\varphi - \varphi_0| - |\varphi_0 - \theta_n| < 2\delta - \delta = \delta$ since $|\varphi - \varphi_0| < 2\delta$ whenever φ belongs to either of the intervals $[0, \varphi_0 - 2\delta]$ or $[\varphi_0 + 2\delta, 2\pi]$. Then, $|\dots| \leq \dots \cos \dots$ (40) Using (39) and (40) in (38) we get, $| \dots | \leq \dots$ for sufficiently large n ; i.e. $\lim_{n \rightarrow \infty} u(r_n, \theta_n) = U(\varphi_0)$ (41) where (r_n, θ_n) is an arbitrary sequence of points in D approaching (R, φ_0) . Equation (41) still holds if some or all the points (r_n, θ_n) lie on Γ since in that case we can directly use the fact that $U(\varphi)$ is continuous on Γ . This implies $u(r, \theta)$ is continuous on D . This completes the proof. Uniqueness of the solution to the Dirichlet problem for a disc. Let u_1 and u_2 be two solutions of the Dirichlet problem. Then their difference $u_1 - u_2 = h$ is harmonic in D and continuous in the closed disk and takes the value zero on the boundary. Hence h attains its upper bounds at some points of the closed disk. If $l > 0$, the upper bound will occur in the open disk, since on the boundary Γ h is zero. This contradicts the conclusions of theorem 2.2. So then $l = 0$. In the same way we can show that the lower bound of h on D is zero. Thus there is no alternative but h to be zero on D .

38 Theorem 2.4 Any continuous function $u(z)$ possessing the mean-value property in a domain D is harmonic in D . Proof. Let K be a closed disk contained in D . By hypothesis of the theorem u satisfies the mean value property in K . We shall prove that u is harmonic in K . By the theorem 2.3 on the Dirichlet problem for a disk there exists a continuous function $\sim u$ in K , which is harmonic in the interior of K and coincides with $u(z)$ on the boundary of K . The difference $u - \sim u$ is continuous and satisfies the mean-value property in K . By the corollary to the theorem 3.7 [(14) page-58] $u - \sim u$ satisfies the maximum modulus principle in K . Now as $u - \sim u$ is zero on the boundary of K , it will be identically zero in K .

Therefore u coincides with the harmonic function $\sim u$ in the interior of K and since K is arbitrary, u is harmonic in the domain D . The Harnack Inequality : Let u be a non-negative Harmonic function on a closed disk $D(0, R)$. Then, for any point $z \in D(0, R)$ $R - z \leq u(z) \leq \dots$ (42) where $D(0, R)$ denotes a disk with centre 0 and radius R . Proof. From the Poisson's integral formula for u on $D(0, R)$: $u(z) = \dots$

Now, $R - z \leq \dots \leq \dots = \dots + \dots$ Combining these two, we see that $u(z) \leq \dots$ where we make use of the mean value theorem. Similarly, the other inequality in (42) will follow from $R - z \geq \dots \geq \dots = \dots - \dots$

Corollary Let u be a non-negative harmonic function on a closed disk $D(\zeta, \rho)$. Then for any $z \in D(\zeta, \rho)$ $R - z \leq u(z) \leq \dots$ (43)

39 2.5 Subharmonic & Superharmonic Functions Definition : A real-valued continuous function $u(x, y)$ in an open set D of the complex plane $C/$ is said to be (i) subharmonic if, for any $\zeta \in D$ $u(\zeta) \leq \dots$ hold for sufficiently small $r > 0$. (ii) superharmonic if, for any $a \in D$ $u(a) \geq \dots$ hold for sufficiently small $r > 0$. From the definition it follows that every harmonic function is subharmonic as well as superharmonic. Example 3. If $f(z)$ is analytic on a domain D , then $|f(z)|$

is subharmonic but not harmonic in D unless $f(z) \equiv \text{constant}$. Solution : Using the Cauchy's integral formula $f(a) = \frac{1}{2\pi i} \int_{\gamma} \frac{f(z)}{z-a} dz$ (44) for every $a \in D$ and $r > 0$ is small enough. Here equality holds only if $f(z) \equiv \text{constant}$. We now show that the integral $I(r) = \frac{1}{2\pi} \int_0^{2\pi} |f(re^{i\theta})|^2 d\theta$ is a strictly increasing function of r , if $f(z) \neq \text{constant}$. Let $0 < r_1 < r_2 < \infty$; $k(\theta)$ and $g(\theta)$ be continuous on $[0, 2\pi]$ and $F(z)$ be defined by (i) $f(z) = \sum_{n=0}^{\infty} k_n z^n$, $\theta \in [0, 2\pi]$ (ii) $F(z) = \sum_{n=0}^{\infty} g_n z^n$, $\theta \in [0, 2\pi]$ (iii) $k_n = \min_{|z|=r} |f(z)|$ distance between a and the boundary of D . $F(z)$ is regular for $|z| \leq r_2$ and attains its maximum on the boundary of the disc, say at $z = r_2 e^{i\phi}$. Then $I(r_2) - I(r_1) = \frac{1}{2\pi} \int_0^{2\pi} (|f(r_2 e^{i\theta})|^2 - |f(r_1 e^{i\theta})|^2) d\theta = \frac{1}{2\pi} \int_0^{2\pi} (|k(\theta)|^2 r_2^{2n} - |g(\theta)|^2 r_1^{2n}) d\theta$, taking $\phi = \theta = \psi = -\theta + 2\pi$ in the third integral, we find that it cancels the second term) $= I(r_2)$. Hence equality in (44) is possible if and only if $f(z) \equiv \text{constant}$. Therefore $|f(z)|$ is subharmonic but not harmonic

in D unless $f(z) \equiv \text{constant}$. Example 4.

If $f(z) \neq 0$ is analytic in a domain D , then $\log |f(z)|$

is subharmonic in

D . Solution : Let $\Phi(z) = \log |f(z)|$. Here at the zeros of $f(z)$, $\Phi(z)$ has poles and takes the value $-\infty$ there. In every closed disk contained in D there are at most a finite number of points where $\log |f(z)| = -\infty$. Now let $a \in D$ be any point at which $f(z)$ is distinct from zero. Since $f(z)$ is analytic and not identically zero, there exists a small neighbourhood of a where $f(z)$ is distinct from zero. We find that $\log f(z) = \log |f(z)| + i \arg f(z)$ is analytic in this neighbourhood and hence $\log |f(z)|$ is harmonic there and we have the equality $\Delta \Phi = 0$ (45) for all sufficiently small values of r . On the other hand, if a is a zero of $f(z)$, we have $\Delta \Phi = -\infty$ (46) Combining (45) with (46) we obtain $\Phi(z)$ is subharmonic in D .

41 Unit 3 Conformal Mappings Structure 3.0 Objectives of this Chapter 3.1 Conformal Mappings 3.2 Basic Properties of Conformal Mapping 3.0 Objectives of this Chapter This chapter deals with conformal mappings and their basic properties. Many examples are given to explain different concepts on conformal mappings. The inverse function theorem is also discussed. 3.1 Conformal Mappings Let X be an open set in \mathbb{C} and suppose a function $f : X \rightarrow \mathbb{C}$ is given. We know from functional analysis that if f is

continuous, a compact set of X is mapped onto a compact set in $f(X)$ and a connected set of X onto a connected set of $f(X)$. If moreover, f is single-valued and analytic there occur several interesting results. In this chapter we study mappings which transform different curves and regions from one complex plane to other complex plane with reference to magnitude and orientation. Such type of mappings play an important role in the study of various physical problems defined on domains and curves of arbitrary shape. Level Curves Let $w = f(z)$ with $z = x + iy$ and $w = u + iv$ where $f(z)$ is analytic. $u =$

$u(x, y) = v(x, y)$

satisfy Cauchy-Riemann equations $u_x = v_y, u_y = -v_x$

from which it follows that $u_{xx} + u_{yy} = 0, v_{xx} + v_{yy} = 0$ Also, $\nabla u \cdot \nabla v = 0$, where Fig. 16 $u(x, y) = \text{constant}, v(x, y) = \text{constant}$

42 $\nabla = \partial_x \partial_x + \partial_y \partial_y$, So that the

level curves $u(x, y) = \text{constant}$ and $v(x, y) = \text{constant}$ are orthogonal.

Now

$f'(z) = u_x + iv_x =$

$u_x - iv_y = v_y + iv_x$ so that $f'(z) = u_x + iv_x = v_y + iv_x$

$f'(z) = u_x + iv_x = v_y + iv_x$

$f'(z) = u_x + iv_x = v_y + iv_x$

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Now

$f'(z) = u_x + iv_x =$

$u_x - iv_y = v_y + iv_x$ so that $f'(z) = u_x + iv_x = v_y + iv_x$

$f'(z) = u_x + iv_x = v_y + iv_x$

D
containing z_0 . If $f'(z_0) \neq 0$, then $f(z)$ is conformal at
 z_0 .

Proof. Let $C_1 : z = z_1(t)$ and $C_2 : z = z_2(t)$, $t \equiv$ parameter, be two curves which intersect at some $t = t_0$ where $z_1(t_0) = z_2(t_0) = z_0$, C_1, C_2 , are their images under the mapping $w = f(z)$. Then following the result given in eq. (47) $\text{Arg } w_1 - \text{Arg } z_1 = \text{Arg } f'(z_1)$ and $\text{Arg } w_2 - \text{Arg } z_2 = \text{Arg } f'(z_2)$

Subtracting, $\text{Arg } w_1 - \text{Arg } w_2 = \text{Arg } f'(z_1) - \text{Arg } f'(z_2)$
i.e., $\theta = \phi$, where $\theta =$

angle between the curves C_1 and C_2 at z_0 and $\phi =$ angle between the curves C_1 and C_2 at w_0 . Observation : From the basic results proved earlier we learn that if f is a conformal mapping, then orthogonal curves are mapped onto orthogonal curves.

3.2 Basic Properties of conformal Mappings

Let $f(z)$ be an analytic function in a domain D , and let z_0 be a point in D .

If $f'(z_0) \neq 0$, then we can express $f(z)$ in the form $f(z) = f(z_0) + (z - z_0)f'(z_0) + \eta(z)$, where $\eta(z) \rightarrow 0$ as $z \rightarrow z_0$.

If z is near z_0 , then the transformation $w = f(z)$ has the linear approximation $G(z) = A + B(z - z_0)$, where $A = f(z_0)$ and $B = f'(z_0)$.

As $\eta(z) \rightarrow 0$ when $z \rightarrow z_0$, for points near z_0 the transformation $w = f(z)$ has an effect much like the linear mapping $w = G(z)$. The effect of the linear mapping G is a rotation of the plane through the angle $\alpha = \text{Arg } (f'(z_0))$, followed by a magnification by the factor $|f'(z_0)|$, followed by a translation by the vector $A + Bz_0$.

Remark : If $f'(z_0) = 0$, the angle may not be preserved. Let us consider, $w = f(z) = z^2$, then we have $f'(0) = 0$ and the angle at $z = 0$ is not preserved but is doubled. Definition : Let $f(z)$ be a nonconstant analytic function. If $f'(z_0) = 0$, the z_0 is called a critical point of $f(z)$, and the mapping $w = f(z)$ is not conformal at z_0 .

We shall see afterwards what happens at a critical point. Fig. 23 Fig. 24 z-plane w-plane

00 47 The Inverse Function theorem 3.2 Let $f(z)$ be

analytic at z_0 and $f'(z_0) \neq 0$. Then there exists a neighbourhood $N(w_0, \epsilon)$ of $w_0 = f(z_0)$ in which the inverse function $z = F(w)$ exists and is analytic. Moreover, $F'(w_0) = 1/f'(z_0)$. (48) Proof : Given $w = f(z)$, ($z = x + iy$, $w = u + iv$) is analytic in a neighbourhood of z_0 , $K : |z - z_0| < \rho$. We shall show that for each $w \in L : |w - w_0| < \epsilon$ there is a unique solution $z = F(w)$, where $z \in K$. We express the mapping $w = f(z)$ in terms of the set of equations $u = u(x, y)$ and $v = v(x, y)$ (49) which represents a transformation from the xy plane to the uv plane, u, v , possess continuous first-order partial derivatives satisfying C-R equations. The Jacobian determinant $J(x, y)$, is defined by $J(x, y) = \begin{vmatrix} u_x & u_y \\ v_x & v_y \end{vmatrix}$ (50) The transformation in equations (49) has a local inverse in L provided $J(x, y) \neq 0$ in K [(3) pp. 358-361]. Expanding r.h.s. of equation (50) and using the C-R equations, we obtain $J(x, y) = 2(u_x^2 + v_x^2)$ (51) $\neq 0$, by the given hypothesis. Utilising the continuity of $J(x, y)$ in a small neighbourhood of (x_0, y_0) , equations (49) and (51) imply that a local inverse $z = F(w)$ exists in a neighbourhood of the point $w_0 = f(z_0)$. The derivative of $F(w)$ is given by the familiar expression

$$F'(w) = \frac{1}{f'(z)}$$

$$F'(w_0) = \frac{1}{f'(z_0)}$$

$$F'(w) = \frac{1}{f'(z)}$$

holds in a neighbourhood of the point w_0 , as $f(z)$ is analytic in K . In particular, $F'(w_0) = 1/f'(z_0)$

Theorem 3.3 Let $f(z)$ be analytic at the point z_0 . If $f'(z_0) = 0, f''(z_0) = 0, \dots, f^{(k-1)}(z_0) = 0$ and $f^{(k)}(z_0) \neq 0$, then the mapping $w = f(z)$ magnifies angles at z_0

by k times. Proof. By the given hypothesis, $f(z)$ has the Taylor expansion in a neighbourhood of z_0 in the form $f(z) = f(z_0) + c_k(z - z_0)^k + c_{k+1}(z - z_0)^{k+1} + \dots$, $c_k \neq 0$ so that we can express

$$f(z) - f(z_0) = (z - z_0)^k + h(z) \quad (52)$$

where $h(z)$ is analytic at z_0 and $h'(z_0) \neq 0$.

Now let $w = f(z)$ and $w_0 = f(z_0)$ and we obtain from (52) $\text{Arg}(w - w_0) = k \text{Arg}(z - z_0) + \text{Arg}(h(z))$ Let $z \rightarrow z_0$ along a curve γ . Then $w \rightarrow w_0$ along the image curve Γ

and the slope of tangent to the curve γ at z_0 and that of the tangent to the curve Γ at w_0 are connected by the relation $\lim_{z \rightarrow z_0} \frac{dw}{dz} = k \lim_{z \rightarrow z_0} \frac{dw}{dz}$ i.e., $\theta_2 - \theta_1 = k(\phi_2 - \phi_1)$ with the sense remain unchanged. Example 2. Show that the mapping $w = f(z) = z^2$ maps the rectangle R in the z -plane of unit area onto the region enclosed by the parabolas $v = u$ and $v = 4 - u$. Solution : Here $f(z) = z^2$ and the mapping $w = z^2$ is conformal for all $z \neq 0$. We note that the right angles at the vertices $z_1 = 1, z_2 = 1 + i/2, z_3 = -1 + i/2$ and $z_4 = -1$ are mapped into right angles at the vertices $w_1 = 1, w_2 = 1 - 1/4, w_3 = 3/4 - i/4$ and $w_4 = -1$ respectively.

49 The parabolas shown in the figure are obtained as follows : Let $w = u + iv$. Then $u = x^2 - y^2, v = 2xy$... (53) The line $x = 1$ corresponds to the curve $u = 1 - y^2, v = 2y$. Eliminating y , we get $v^2 = -4(u - 1)$, which is a parabola with vertex $(1, 0)$ and opens towards the negative side of the u -axis in the w -plane. Also, the part of the line $x = 1$ lying above the real axis corresponds to the part of the parabola lying above the u -axis in the w -plane. The same parabola in the w -plane is the image of the line $x = -1$. In this case, the part of the line $x = -1$ lying above the real axis corresponds to the part of the parabola lying below the u -axis in the w -plane. Again, when $y = 1/2$, from (53) we get $u = x^2 - 1/4$ and $v = x$. Eliminating x we get, $v^2 = 4(u + 1/4)$ which is also a parabola with vertex $-1/4, 0$, and opening towards the positive side of the u -axis in the w -plane. By similar argument as before we can say that the mapping $w = z^2$ maps the rectangle R in the z -plane onto the region enclosed by the parabolas $v = u$ and $v = 4 - u$. Note : It is not hard to prove that the parabolas intersect each other orthogonally at w_2 and w_3 . At the point $z_0 = 0$, we have $f'(z_0) = f'(0) = 0$ and $f''(z_0) = 2 \neq 0$. Hence the angles at the origin $z_0 = 0$ are magnified by the factor $k = 2$. In particular the straight angle at $z_0 = 0$ is mapped onto 2π angle at $w_0 = 0$. Fig. 25 Fig. 26

50 Unit 4 Multi-valued functions and Riemann Surface Structure 4.0 Objectives of this Chapter 4.1 Multi-valued functions 4.2 The logarithm function 4.3 Properties of $\log z$ 4.4 Branch, Branch point and Branch cut 4.5 Integrals of Multi-valued function 4.6 Branch points at infinity 4.7 Detection of branch points 4.8 The Riemann Surface for $w = z^{1/2}$ 4.9 Concept of neighbourhood 4.10 The Riemann Surface for $w = \log z$ 4.11 The Inverse Trigonometric Functions 4.0 Objectives of this Chapter In this chapter we shall study multi-valued functions and their Riemann surfaces. In particular, multi-valued logarithm function, the power function z^α both z, α complex numbers, $z \neq 0$ will be discussed. The ideas of branch, branch point, branch cut, branch point at infinity will be explained by means of different examples. A few contour integrations of multi-valued functions will be performed. Also Riemann surfaces for different multi-valued functions will be constructed. 4.1 Multi-valued functions So far we have considered single-valued functions i.e., one-to-one mapping or, many-to-one mapping. In the later case, under certain restrictions, inverse mappings give rise to multi-valued functions i.e., one-to-many. For example,

51 $z = e^{i\omega}, z = \omega^2, z = \sin \omega, z = \cos \omega$ For each of these functions, a given value of z corresponds to more than one value of ω . $\omega = f^{-1}(z)$ is multi-valued and $z = f(\omega)$ is single-valued, given ω , there is a unique value of z . The aim of this chapter is as follows : (i) To determine all possible values of the inverse function ω and (ii) To construct an inverse function which is single-valued in some region of the complex plane. Let $\omega = f(z)$ be a multi-valued function.

A branch of f is any single-valued function f_0 that is continuous in some domain (except, perhaps, on the boundary). At each point z in the domain, it assigns one of the values of $f(z)$. Example 1 : We consider branches of the two-valued square-root function $f(z) = z^{1/2} (z \neq 0)$. The principal branch of the square root function is $f_1(z) = \sqrt{r} e^{i\theta/2}$ where $r = |z|$ and $-\pi < \theta \leq \pi$. The function f_1 is a branch of f . Using the same notation, we can find other branches of the function f . For example if we let $f_2(z) = \sqrt{r} e^{i(\theta/2 + \pi)}$ then f_2 is another branch of f . So, f_1 and f_2 can be taken as the two branches of the multi-valued square root function. The negative real axis is called a branch cut for the functions f_1 and f_2 . Each point on the branch cut is a point of discontinuity for both functions f_1 and f_2 . Result 1 : Show that the function f_1 is discontinuous on the negative real axis. Fig. 27

52 Solution : Let $z_0 = r_0 e^{i\pi}$ be any point on the negative real axis. We compute the limit as z approaches z_0 through the upper half plane $\text{Im } z < 0$ and the limit as z approaches z_0 through the lower half plane $\text{Im } z > 0$. The limits are $\lim_{z \rightarrow z_0^-} f_1(z) = \sqrt{r_0} e^{i\pi/2}$ and $\lim_{z \rightarrow z_0^+} f_1(z) = \sqrt{r_0} e^{3\pi/2}$. The limits are $\lim_{z \rightarrow z_0^-} f_2(z) = \sqrt{r_0} e^{3\pi/2}$ and $\lim_{z \rightarrow z_0^+} f_2(z) = \sqrt{r_0} e^{5\pi/2}$.

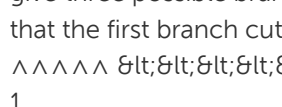
63 we find that like $z^{1/2}$, $\zeta^{-1/2}$ possesses double sheeted Riemann surface. We see that each time we walk around the origin, the argument of $\zeta^{-1/2}$ changes by $-\pi$. This means that the value of the function changes by the factor $e^{-i\pi} = -1$, i.e. the function changes sign. If we walk around the origin twice, the argument changes by -2π , so that the value of the function does not change, $e^{-2\pi i} = 1$. Now, since $\zeta^{-1/2}$ has a branch point at zero, we conclude that $z^{1/2}$ has a branch point at infinity. Example 10 : Again consider the multi-valued logarithm function $f(z) = \log z$. Mapping the point at infinity to the origin, we have $f(\zeta) = \log \zeta$, $\zeta = z^{-1}$. But $\log \zeta$ has a branch point at $\zeta = 0$. Thus $\log z$ has a branch point at infinity. Branch points at infinity : Paths around infinity We can also check for a branch point at infinity by considering a path that encloses the point at infinity and no other singularities. This can be done by drawing a simple closed curve that separates the complex plane into a bounded region that contains all the singularities of the function in the finite plane. Then, depending upon the orientation, the curve is a contour enclosing all the finite singularities, or the point at infinity and no other singularities. Once again consider the function $z^{1/2}$. We know that the function changes value on a curve that goes around the origin. Such a curve can be considered to be either a path around the origin or a path around the point at infinity. In either case the path encloses one branch point. Now consider a curve that does not go around the origin. Such a curve can be considered to be either a path around neither of the branch points or both of them. Thus we see that $z^{1/2}$ does not change value when we follow a path that encloses neither or both of its branch points. Example 11 : Consider the multi-valued function $f(z) = (z^2 - 1)^{1/2}$. Rewriting the function $f(z) = (z - 1)^{1/2} (z + 1)^{1/2}$, we see that there are branch points at $z = \pm 1$. Now consider the point at infinity. $f(\zeta^{-1}) = (\zeta^{-2} - 1)^{1/2} = \pm \zeta^{-1} (1 - \zeta^2)^{1/2}$ which shows that $f(\zeta^{-1})$ does not have a branch point at $\zeta = 0$ and $f(z)$ does not have a branch point at infinity. We might reach the same conclusion by considering a path around the point at infinity. Consider a path that encircles the branch points at $z = \pm 1$ once in the positive direction. Equivalently it encircles the point at infinity once in the negative direction. In traversing this path, the value of $f(z)$ is multiplied by the factor $(e^{2i\pi})^{1/2} (e^{-2i\pi})^{1/2} = e^{2i\pi} = 1$. Thus the value of the function remains unchanged. There is no branch point at infinity.

64 4.7 Detection of branch points We have the definition of a branch point, but we do not have a convenient criterion for determining if a particular function has a branch point. We have noticed that $\log z$ and z^k for non-integer k have branch points at zero and infinity. The inverse trigonometric functions like $\sin^{-1} z$, $\cos^{-1} z$ etc. also have branch points, but they can be written in terms of the logarithm and the square root. In fact all the elementary functions with branch points can be written in terms of the functions $\log z$ and z^k . Furthermore, note that the multi-valuedness of z^k comes from the logarithm, $z^k = e^{k \log z}$. This gives us a way of determining branch points of a function if there is any. Result : Let $f(z)$ be a single-valued function. Then $\log f(z)$ and $(f(z))^k$ may have branch points only where $f(z)$ is zero or singular. Example 12 : Consider the functions 1. $(z^2)^{1/2}$ 2. $(z^{1/2})^2$ 3. $(z^{1/2})^3$ Are they multi-valued? Do they have branch points? Solution 1. $z^2 z^{1/2} z^{1/2} = z^2$ Because of $(z^2)^{1/2}$, the function is multi-valued. The only possible branch points are at zero and point at infinity. If $(e^{i\theta})^2)^{1/2} = 1$, then as $((e^{2\pi i})^2)^{1/2} = (e^{4\pi i})^{1/2} = e^{2\pi i} = 1$ the function does not change value when we walk around the origin. We can also consider this to be a path around infinity. This function is multi-valued, but has no branch points. 2. $z^2 z^{1/2} z^{1/2} = z^2$ This function is single-valued. 3. $z^2 z^{1/2} z^{1/2} z^{1/2} z^{1/2} = z^3$ This function is multi-valued. We consider the possible branch point at $z = 0$. If $(e^{i\theta})^{1/2})^3 = 1$, then as $((e^{2\pi i})^{1/2})^3 = ((e^{i\pi})^{1/2})^3 = (e^{i\pi})^3 = e^{3\pi i} = -1$, the function changes value when we walk around the origin. So it has a branch point at $z = 0$. Since this is also a path around infinity, there is a branch point at the point at infinity. Example 13 : Consider the function $f(z) = \log(1/z - 1)$. Since $1/z - 1$ has only zero at infinity and its only singularity (a pole here) is at $z = 1$, the only, possible branch points are at $z = 1$ and $z = \infty$.

65 Here $f(z) = \log(1/z - 1) = \log(1 - z) - \log z$, $z \neq 0, 1$ say We know that $\log w$ has branch points at zero and infinity, so $f(z)$ has branch points at $z = 1$ and $z = \infty$. Example 14 : Consider the functions 1. $e^{\log z}$ 2. $\log e^z$ Are they multi-valued? Do they have branch points? Solution : 1. $e^{\log z} = e^{\log z + i2\pi k}$, $k = 0, \pm 1, \dots = e^{\log z} e^{i2\pi k} = z$ The function is single-valued. 2. $\log e^z = \text{Log} e^z + i2\pi k = z + i2\pi k$, $k = 0, \pm 1, \dots$ This function is multi-valued. It may have branch points only where e^z is zero or infinite. This occurs only at $z = \infty$. Thus there are no branch points in the finite plane. The function does not change when traversing a simple closed path and since this path can be considered to enclose the point at infinity, there is no branch point at infinity. Note : Let $f(z)$ be single-valued and have either a zero or a singularity at $z = z_0$. Then $\{f(z)\}^k$ may have a branch point at $z = z_0$. If $f(z)$ is not a power of z , then we are not sure whether $\{f(z)\}^k$ changes value when we walk around z_0 . Now if $f(z)$ can be decomposed into factors $f(z) = h(z)g(z)$, where $h(z)$ is finite and non zero at z_0 , then from $g(z)$

we know how fast $f(z)$ vanishes or tends to infinity. Again $\{f(z)\}^k = \{h(z)\}^k \{g(z)\}^k$ and $\{h(z)\}^k$ does not have a branch point at $z = 0$. So that $\{f(z)\}^k$ has a branch point at $z = 0$ if and only if $\{g(z)\}^k$ has a branch point there. Similarly, we can decompose $\log \{f(z)\} = \log \{h(z)g(z)\} = \log \{h(z)\} + \log \{g(z)\}$ to see that $\log \{f(z)\}$ has a branch point at $z = 0$ if and only if $\log \{g(z)\}$ has a branch point there.

Example 15 : Consider the functions : 1. $\sin z^{1/2}$ 2. $(\sin z)^{1/2}$ 3. $z^{1/2} \cos z^{1/2}$ 4. $(\sin z^2)^{1/2}$. Find the branch points and the number of branches. **Solution :** 1. $\sin z^{1/2} = \pm \sin z^{1/2}$ So it is multi-valued. It has two branches and the possible branch points are zero and infinity. Consider the unit circle $|z| = 1$ which is a path around the origin and infinity. If $\sin(e^{i0})^{1/2} = \sin(1)$, then as $66 \sin((e^{i2\pi})^{1/2}) = \sin(e^{i\pi}) = \sin(-1) = -\sin 1$, there are branch points at the origin and infinity 2. $(\sin z)^{1/2} = \pm \sin z^{1/2}$ The function is multi-valued and has two branches. The sine function vanishes at $z = n\pi$ and is singular at infinity. These may be branch points of the function. Consider the point $z = n\pi$. We can express $\sin(z - n\pi) = \sin z \cos n\pi - \cos z \sin n\pi = \sin z (-1)^n$ an integer. But $\lim_{z \rightarrow n\pi} \sin(z - n\pi) = \lim_{z \rightarrow n\pi} \cos(-1)z = \cos(-1)z \rightarrow \cos(-1)z = \cos z$ So, $(\sin z)^{1/2}$ has branch points at $z = n\pi$ since $(z - n\pi)^{1/2}$ has a branch point at $z = n\pi$. Here the branch points are $z = n\pi$, $n = 0, \pm 1, \dots$ and they go to infinity. So it is not possible to make a path that encloses infinity and no other singularities. The point at infinity is a non-isolated singularity. A point can be a branch point only if it is an isolated singularity. 3. $z^{1/2} \cos z^{1/2} = \pm z^{1/2} \cos z^{1/2}$ The function is multi-valued. It may possess branch points at $z = 0$ and $z = \infty$. If $(e^{i0})^{1/2} \cos(e^{i0})^{1/2} = \cos(1)$, then as $(e^{i2\pi})^{1/2} \cos((e^{i2\pi})^{1/2}) = (-1)\cos(e^{i\pi}) = -\cos(-1) = -\cos 1$, there are branch points at the origin and infinity. 4. $(\sin z^2)^{1/2} = \pm \sin z^2$ The function is multi-valued. Now since $\sin z^2 = 0$ at $z = (n\pi)^{1/2}$, there may be branch points there. We consider first the point $z = 0$. We can write $\sin z^2 = \sin z^2$ but $\lim_{z \rightarrow 0} \sin z^2 = \lim_{z \rightarrow 0} \cos z^2 = 1$ So, $(\sin z^2)^{1/2}$ does not have a branch point at $z = 0$ as $(z^2)^{1/2}$ does not have a branch point there. Next consider the point $z = \pi$ $67 \sin(\pi) = \sin(\pi)$ $z = \pi$ $n z^2 = \pi^2$ but $\lim_{z \rightarrow \pi} \sin(z^2) = \lim_{z \rightarrow \pi} \cos(z^2) = \cos(\pi^2)$ Since $z = \pi$ has a branch point at $z = \pi$, $(\sin z^2)^{1/2}$, too as a branch point there. Thus we see that $(\sin z^2)^{1/2}$ has branch points at $z = (n\pi)^{1/2}$ for $n \in \mathbb{Z} \setminus \{0\}$. This is the set of numbers : $\pm \sqrt{\pi}, \pm \sqrt{4\pi}, \pm \sqrt{9\pi}, \dots$. The point at infinity is a non-isolated singularity and hence it is not included in the set of branch points.

Example 16 : Find the branch points of $f(z) = (z^3 - z)^{1/3}$ and introduce the branch cuts. If $f(1) = 3^{1/3}$ find $f(-3)$. **Solution :** Here $f(z) = z^{1/3} (z - 1)^{1/3} (z + 1)^{1/3}$ So the branch points are at $z = -1, 0$ and 1 . We consider the point at infinity $f(1/\zeta) = (1/\zeta^3 - 1/\zeta)^{1/3} = (1/\zeta^3)^{1/3} (1 - \zeta^3)^{1/3} = 1/\zeta (1 - \zeta^3)^{1/3}$ Since $f(1/\zeta)$ does not have a branch point at $\zeta = 0$, $f(z)$ does not have a branch point at infinity. Here we give three possible branch cuts : In the first and third the function is single-valued but in the second it is not. It is clear that the first branch cut does not allow us to walk around any of the branch points.  Fig. 41 \odot Three possible branch cuts for $f(z) = (z^3 - z)^{1/3}$ \odot $-1, -1 - 1 - 1$

68 The second branch cut allows us to walk around the branch points at $z = \pm 1$. If we walk around these two once in the positive direction, the value of the function would change by the factor $e^{i4\pi/3}$. The third branch cut allows us to walk around all the three branch points, the value of the function will not change (since $e^{i6\pi/3} = e^{i2\pi} = 1$). To find $f(-3)$, we consider the third branch cut with $f(z) = z^{1/3} (z - 1)^{1/3} (z + 1)^{1/3}$ $f(-3) = (-3)^{1/3} (-3 - 1)^{1/3} (-3 + 1)^{1/3} = -3^{1/3} (-4)^{1/3} (-2)^{1/3} = -3^{1/3} 4^{1/3} 2^{1/3} = -24^{1/3}$ The value of $f(-3)$ is $-24^{1/3}$ Example 17 : Determine the branch points of the function $f(z) = (z^3 - 1)^{1/2}$. Construct branch cuts and define a branch so that $z = 0$ and $z = -1$ do not lie on a cut, such that $f(0) = -i$; then what is $f(-1/2)$? **Solution :** The roots of the equation $z^3 - 1 = 0$ are $1, \omega, \omega^2, \dots, -1, -1 - i\sqrt{3}, \dots$ so that, $z^3 - 1 = (z - 1)(z - \omega)(z - \omega^2)$ There are branch points at each of the cube roots of unity $z = 1, \omega, \omega^2$ $1, \omega, \omega^2, -1, -1 - i\sqrt{3}, -1 + i\sqrt{3}$ Now we examine the point at infinity. We make the change of variable $z = 1/\zeta$ $f(1/\zeta) = (1/\zeta^3 - 1)^{1/2} = \zeta^{-3/2} (1 - \zeta^3)^{1/2}$ $\zeta^{-3/2}$ has a branch point at $\zeta = 0$, while $(1 - \zeta^3)^{1/2}$ is not singular there. Since $f(1/\zeta)$

72 4.10 The Riemann Surface for $w = \log z$ The Riemann surface for the multivalued function $\omega = \log z$ is similar to the one we presented for the square root function. However, it

requires infinitely many copies of the z -plane cut along the negative x -axis,

which mark S_k for $k = \dots, -n, \dots, -1, 0, 1, \dots, n, \dots$. Now we stack these cut planes directly on each other so that the corresponding points have the same position. We join the sheet S_k to S_{k+1} as follows. For each integer k , the edge of the sheet S_k in the upper half-plane is joined to the edge of the sheet S_{k+1} in the lower-half plane.

The Riemann surface for the domain of $\log z$ looks like a spiral staircase that extends upward on the sheets S_1, S_2, \dots , and downward on the sheets S_{-1}, S_{-2}, \dots

as shown in figure 45. For S_k , we use $z = re^{i\theta} = r(\cos \theta + i \sin \theta)$, where $r = |z|$ and $2\pi k - \pi < \theta \leq \pi + 2\pi k$. Again, for S_k , the correct branch of $\log z$ on each sheet is $\log z = \log r + i\theta$, where $r = |z|$ and $2\pi k - \pi < \theta \leq \pi + 2\pi k$.

Example 19 : Form a Riemann surface for $f(z) = (z - 1)^{1/3}$ taking a branch cut along the line $y = 0, x \geq 1$. Detect the point where the function takes the value $\sqrt{2}(i - 1)$. Solution : Let $r = |z - 1|$ and $\theta = \arg(z - 1)$, where $0 \leq \theta < 2\pi$. Here the Riemann surface consists of three domains D_1, D_2 and D_3 : $f_1(z) = r^{1/3} e^{i\theta/3}, 0 \leq \theta < 2\pi$ (D_1) $f_2(z) = r^{1/3} e^{i\theta/3}, 2\pi \leq \theta < 4\pi$ (D_2) $f_3(z) = r^{1/3} e^{i\theta/3}, 4\pi \leq \theta < 6\pi$ (D_3)

Each function f_1, f_2 and f_3 is single-valued on the domain formed by cutting the z -plane along the line $y = 0, x \geq 1$. We place D_1 on the top, then D_2 and D_3 . The edge of D_1 in the upper-half plane is joined to the edge of D_2 in the lower-half plane and the edge of D_2 in the upper-half plane is joined to the edge of D_3 in the lower-half plane and finally the edge of D_3 in the upper-half plane is joined to the edge of D_1 in the lower-half plane. Say at $z = \zeta, f(\zeta) = \sqrt{2}(i - 1)$ i.e. $f_1(\zeta) = \sqrt{2}(i - 1)$ lying in the domain D_2 .

Example 20 : Form the Riemann surface for the function $f(z) = (z^2 - 1)^{1/2}$. Solution : Here the given function $f(z) = (z^2 - 1)^{1/2}$ has branch points at $z = \pm 1$. To examine the point at infinity, we substitute $z = 1/\zeta$ and examine the point $\zeta = 0$. $f(1/\zeta) = (1/\zeta^2 - 1)^{1/2} = (1 - \zeta^2)^{1/2} / \zeta$. Since there is no branch point at $\zeta = 0$, $f(z)$ has no branch point at infinity. Let $z - 1 = r_1 e^{i\phi_1}$ and $z + 1 = r_2 e^{i\phi_2}$, where $\phi_1 = \arg(z - 1)$ and $\phi_2 = \arg(z + 1)$. Then $\omega = f(z) = (z^2 - 1)^{1/2} = (z - 1)^{1/2} (z + 1)^{1/2} = (r_1 r_2)^{1/2} e^{i(\phi_1 + \phi_2)/2}$.

Case-I $0 \leq \phi_1 < 2\pi, 0 \leq \phi_2 < 2\pi$ on the $\phi_1 + \phi_2 \in [0, 4\pi)$ Continuity segment of $f(z)$ $B = 0, i, 2\pi - i, C = 0, 1$ Yes $C' = 2\pi, 2\pi - 1$ Yes $D' = \pi, \pi - 1$ Fig. 49 Case-II $0 \leq \phi_1 < 2\pi, -\pi \leq \phi_2 < \pi$ on the $\phi_1 + \phi_2 \in [0, \pi)$ Continuity segment of $f(z)$ $B = 0, i$ Yes $B' = \pi, 0$ i $C = 0, 1$ No $C' = 2\pi, 0 - 1$ $D = \pi, \pi - 1$ No $D' = \pi, -\pi$ Two branches of $(z - 1)^{1/2}$ can be taken as $f = z r e^{z r e^{i\theta}}$ and $g = z r e^{z r e^{i\theta}}$ Again two branches of $(z + 1)^{1/2}$ can be taken as $g = z r e^{z r e^{i\theta}}$ and $g = z r e^{z r e^{i\theta}}$ Let us construct a Riemann surface for $\omega = (z^2 - 1)^{1/2}$ considering case I. Here a Riemann surface consists of two sheets S_0 and S_1 . Let S_0 be an extended complex plane cut along the real axis from $z = -1$ to $z = 1$ and S_1 be another extended complex plane cut of similar nature. $S_0 = \{z \in \mathbb{C} : \arg z \in [0, 2\pi), z \neq [-1, 1]\}$ and $S_1 = \{z \in \mathbb{C} : \arg z \in [\pi, 3\pi), z \neq [-1, 1]\}$. The sheets S_0 and S_1 are joined along the segment $[-1, 1]$ in such a way that the lower edge of the slit in S_0 is joined to the upper edge of the slit in S_1 , and the lower edge of the slit in S_1 is joined to the upper edge of the slit in S_0 .

k or k is any integer. Likewise, the set of values for other inverse trigonometric functions can be ascertained. Example 22 : Discuss the mapping $\omega = \sinh z$ that transforms the infinite strip $-\infty < x < \infty; 0 < y < \pi$ into the ω -plane. Find cuts in the ω -plane which make the mapping continuous both ways. What are the images of the lines (a) $y = 1/\pi$ (b) $x = 1$? Solution : First we express $\sinh z$ in cartesian form $\omega = \sinh z = \sinh x \cos y + i \cosh x \sin y = u + iv$ We consider the line segment $x = c, y \in (0, \pi)$. Its image is

$78 \{ \sinh c \cos y + i \cosh c \sin y | y \in (0, \pi) \}$ Clearly, u and v then satisfy the equation for the ellipse $u^2/c^2 + v^2/\cosh^2 c = 1$. The ellipse starts at the point $(\sinh c, 0)$, passes through the point $(0, \cosh c)$ and ends at $(-\sinh c, 0)$. As c varies from zero to ∞ or from zero to $-\infty$, the semi-ellipses cover the upper-half of ω -plane. Thus the mapping is 2-to-1. Now consider the infinite line $y = c, x \in (-\infty, \infty)$. It's image is $\{ \sinh x \cos c + i \cosh x \sin c | x \in (-\infty, \infty) \}$. Here u and v satisfy the equation for a hyperbola $u^2/c^2 - v^2/\cosh^2 c = 1$. As c varies from 0 to $\pi/2$ or from $\pi/2$ to π , the semi-hyperbola cover the upper- half of ω -plane. Thus the mapping is 2-to-1. We look for branch points of $\sinh^{-1} \omega$. $\omega = \sinh z \Rightarrow z = \sinh^{-1} \omega = \frac{1}{2} \log \left(\frac{\omega + \sqrt{\omega^2 + 1}}{\omega - \sqrt{\omega^2 + 1}} \right)$ The branch points are at $\omega = \pm i$. Since $\omega + \sqrt{\omega^2 + 1}$ is non zero and finite in the finite complex plane, the logarithm does not introduce any branch in the finite plane. Thus the only branch point in the upper-half of ω -plane is at $\omega = i$. Any branch cut that connects $\omega = i$ with the boundary of $\text{Im } \omega \leq 0$ will separate the branches under the inverse mapping. We consider the line $y = \pi/4$. The image under the mapping is the upper-half of the hyperbola $2u^2 - 2v^2 = 1$. Consider the segment $x = 1$. The image under the mapping is the upper-half of the ellipse. $u^2/c^2 + v^2/\cosh^2 c = 1$

79 Example 23 : Construct a Riemann Surface for $\cos^{-1} z$. Solution : The function $\omega = \cos^{-1} z = -i \log [z + (z^2 - 1)^{1/2}]$ has two sources of multi-valuedness; one due to the square root function $(z^2 - 1)^{1/2}$ and the other due to the logarithm, if any. First we construct the branch of the square root $(z^2 - 1)^{1/2} = (z + 1)^{1/2} (z - 1)^{1/2}$ We see that there are branch points at $z = -1$ and $z = 1$.

In particular we want the $\cos^{-1} z$ to be defined for $z = x, x \in [-1, 1]$. Hence we introduce the branch cuts on the lines $(-\infty, -1]$ and $[1, \infty)$. Let $z + 1 = re^{i\theta}, z - 1 = \rho e^{i\phi}$ With the given branch cuts, the angles have the possible ranges $-\pi \leq \theta < \pi, 0 \leq \phi < 2\pi$ Now we must determine if the logarithm introduces additional branch points. The only possibilities for branch points are where the argument of the logarithm is zero. $z + (z^2 - 1)^{1/2} = 0$ or, $z^2 = z^2 - 1 \Rightarrow 0 = -1$ We see that the argument of the logarithm can not be zero and thus there are no additional branch points. Here the Riemann surface consists of two sheets S_0 and S_1 joined on the real axis $(-\infty, -1] \cup [1, \infty)$: $S_0: 0 \leq \theta < \pi; \leq \phi < 2\pi; S_1: \pi \leq \theta < 2\pi; \pi \leq \phi < 3\pi$

A particular branch of this function can be obtained by first taking $z + 1 = re^{i\theta}, -\pi \leq \theta < \pi; z - 1 = \rho e^{i\phi}, 0 \leq \phi < 2\pi$ Then adding these relations, we find $z = (re^{i\theta} + \rho e^{i\phi})/2$ and the function $z + (z^2 - 1)^{1/2}$ reduces to $z z r e^{i\theta} + \rho e^{i\phi} = z r e^{i\theta} + \rho e^{i\phi}$ Fig. 53 $-1 \leq y < 1$

$80 = z r e^{i\theta} + \rho e^{i\phi}$ Then $\cos^{-1} z = \frac{1}{2i} \log \frac{z + (z^2 - 1)^{1/2}}{z - (z^2 - 1)^{1/2}}$ If a point lying on the sheet S_0 is allowed to travel a path making a complete round about only the branch point $z = 1$, it enters to the sheet S_1 from the sheet S_0 . In this case the value of ϕ changes by 2π while the value of θ remains unchanged. The change in $(\phi - \theta)/2$ is π . So in this case, $\cos^{-1} z = \frac{1}{2i} \log \frac{z + (z^2 - 1)^{1/2}}{z - (z^2 - 1)^{1/2}}$ on S_1 . Similarly we can analyse the case when the point on S_0 encloses only the branch point $z = -1$ while travelling a complete round. Some standard branch cuts of elementary functions. Function Branch cuts z , non integral s with $\text{Re } s > 0$ $(-\infty, 0)$ z^s , non integral s with $\text{Re } s \leq 0$ $(-\infty, 0)$ e^z none $\log z$ $(-\infty, 0)$ $\sin^{-1} z, \cos^{-1} z$ $(-\infty, -1]$ and $[1, \infty)$ $\tan^{-1} z$ $y \leq -1, x = 0$ and $y \geq 1, x = 0$ $\text{cosec}^{-1} z, \sec^{-1} z$ $(-1, 1)$ $\cot^{-1} z$ $[-i, i]$ $\sinh^{-1} z$ $y > \pi/2, x = 0$ and $y < -\pi/2, x = 0$ $\cosh^{-1} z$ $(-\infty, 1)$ $\text{cosech}^{-1} z$ $-1 < y < 1, x = 0$ and $y > 1, x = 0$ $\text{coth}^{-1} z$ $[-1, 1]$

81 Exercises 1.

- Find the principal value of each of the following complex quantities : (a) $(1 - i)^{1+i}$ (b) 3^{3-i} (c) 2^{2-i}
 i 2. Give the number of branches and locations of the branch points for the functions. (a) $\cos(z^{1/2})$ (b) $(z + i)^{-z}$.
 Determine the branch points of the function $\omega = \{(z^2 - z)(z + 2)\}^{1/3}$ 4. Find the branch points of $(z^{1/2} - 1)^{1/2}$ in the finite complex plane. Introduce branch cuts to make the function single-valued. 5. Let D be the complex z -plane with a cut along the segment $[-1, 1]$, determine the regular branches of the function $f(z) = z^z$ 6. Split the function $f(z) = z^z (-z)^{-z} = 2^{2z} 4^{z^2}$ into two regular branches in the domain $D: \setminus \{-3, -2, -1, 0, 1, 2, 3\}$ 7. Evaluate (i) $\int_0^1 x^\alpha dx$ (ii) $\int_0^1 \log x dx$ (iii) $\int_0^1 \log x dx$ 8. Prove that $\int_0^1 \log \sin x dx = -\pi/2$ 9. Construct a Riemann surface for the following functions : (i) $\omega = z^{1/3}$ (ii) $\omega = (z^2 + 1)^{1/2}$ (iii) $\omega = \log(z - z^2)$ (iv) $\omega = \sin^{-1} z$.

82 $\circ -i i 10$. Let $f(z)$ have branch points at $z = 0$ and $z = \pm i$ but nowhere else in the extended complex plane. How does the value and argument of $f(z)$ change while traversing the contour given in the figures 51(a) (b). Do the branch cuts make the function single valued? \circ Fig. 54 (b) $-i$ i Fig. 54 (a) $\<$; $\<$; $\<$; $\<$; $\<$; $\<$; $\<$; $\<$; $\<$; $\<$;

83 Unit 5 Conformal Equivalence Structure 5.0 Objectives 5.1 Riemann Mapping Theorem 5.2 The Schwarz Reflection Principle 5.3 The Schwarz-Christoffel Transformation 5.4 Examples : Triangles / Rectangles 5.0 Objectives of this Chapter The concept of conformal equivalence of two regions will be introduced in this chapter. The main theorem of this chapter is Riemann mapping theorem. Also Hurwitz's theorem, Schwarz lemma, Schwarz reflection principle, Schwarz-Christoffel transformation will be studied and their applications will be shown through a few examples. 5.1 Riemann Mapping Theorem In the family of analytic functions that concern geometrical orientation, conformal mapping plays a leading role. As its consequences we shall present here a most important result named after G. F. B Riemann, known as "Riemann mapping theorem". Throughout $H(G)$ will denote the family of analytic functions defined on the region G . Definition : Conformal Equivalence : Two regions R_1 and R_2 are said to be conformally equivalent if there exists a $f \in H(R_1)$ such that f is one-to-one in R_1 and $f(R_1) = R_2$ i.e. if there exists a conformal mapping one to one of R_1 onto R_2 . Clearly, this is an equivalence relation (reflexive, symmetric and transitive).

Theorem 5.1 [Hurwitz's Theorem] Let G be a region and $\{f_n\}$ be a sequence in $H(G)$ that converges uniformly to $f \in H(G)$. Suppose $f \neq 0$, $D(a, R) \subset G$ and $f(z) \neq 0$ on $\gamma : |z-a| = R$. Then there exists an integer N such that for $n \geq N$, f_n and f have the same number of zeros in $D(a, R)$.

84 Proof. Since $f(z)$ is never zero on the circle γ , we have $\inf_{z \in \gamma} |f(z)| = \delta > 0$. Again, $f_n \rightarrow f$ uniformly on γ , so there is an integer N such that for $n \geq N$
 $\sup_{z \in \gamma} |f_n(z) - f(z)| < \delta$
 $|f_n(z)| > |f(z)| - \delta > \delta/2$ and thus on the circle γ , $|f_n(z)| > \delta/2$ and thus on the circle γ , $|f_n(z) - f(z)| < \delta/2 \leq \delta/2$ for $n \geq N$. Using Rouché's theorem we find that f_n and f have the same number of zeros inside the circle

$\gamma : |z-a| = R$ for $n \geq N$. By means of the above theorem, we can easily prove Corollary 1. Let G be a region and $\{f_n\}$ be a sequence in $H(G)$ such that each f_n never vanishes in G . Suppose $f_n \rightarrow f$ uniformly in $H(G)$. Then $f(z)$ never vanishes in G , unless $f \equiv 0$. Some useful results (i)

If $f(z)$ is analytic at z_0 and $f'(z_0) \neq 0$, then there is a neighbourhood of z_0 in which $f(z)$ is univalent. (ii) An univalent analytic function f on a domain G has a non-zero derivative at every point of G , i.e., $f'(z) \neq 0$ on G . (iii) The inverse of an univalent analytic function is analytic. (iv) Any domain in \mathbb{C} , that is conformally equivalent to a simply connected domain must itself be simply connected. (v) A domain D in \mathbb{C} is simply connected if and only if every analytic function in D has a primitive in D . Schwarz Lemma Let $f : D(0, 1) \rightarrow D(0, 1)$ be an analytic function which maps the unit disc $D(0, 1)$ to itself. If $f(0) = 0$, then (i) $|f(z)| \leq |z|$ for $0 \leq |z| < 1$ (ii) $|f'(0)| \leq 1$ (iii) if equality holds in (i) for at least one $z \in D(0, 1) - \{0\}$, or, if equality holds in (ii), then $f(z) = \lambda z$, where λ is a constant, $|\lambda| = 1$. Proof : Let us consider the function $g(z) = f(z)/z$ () =

85 which is analytic in the disc $D(0, 1) - \{0\}$ and it has removable singularity at $z = 0$, since $f(0) = 0$. It can be made analytic at $z = 0$ if we define $g(0) = \lim_{z \rightarrow 0} f(z)/z = f'(0)$ (55) For $|z| = r$, where $0 < r < 1$ $g(z) = f(z)/z$ () = $\<$; $\<$; $\<$; $\<$; $\<$; $\<$; $\<$; $\<$; $\<$; $\<$ By the Maximum Modulus Principle, $|g(z)| \leq 1/r$ for the entire disc $|z| \leq r$. We fix $z \in D(0, 1) - \{0\}$ and let $r \rightarrow 1$. Then $|g(z)| \leq 1$. This is true for all $z \in D(0,1) - \{0\}$ and we get $|f(z)| \leq |z|$ (56) i.e. $|f(z)| \leq |z|$, $0 < |z| < 1$. Since $f(0) = 0$, we have $|f(z)| \leq |z|$ for $0 \leq |z| < 1$.

So, (i) is proved and we find from (55) that $|g(0)| = |f(1(0))| \leq 1$ which proves (ii) To prove (iii), we observe that if at a point $z \neq 0$ ($|z| > 1$) $|g(z)| = 1$ i.e. $|g(z)|$ attains its maximum at an internal point and hence by the maximum modulus principle $g(z) = \lambda$, a constant and that $|\lambda| = 1$, so $f(z) = \lambda z$. Theorem 5.2 Let $a \in D(0, 1)$. Then ϕ_a defined by $\phi_a(z) = \frac{z+a}{1+\bar{a}z}$ maps $D(0, 1)$ onto $D(0, 1)$. Proof. Clearly, ϕ_a is a bilinear transformation, it is analytic in the whole complex plane except the point $1/\bar{a}$ (which is the inverse point of the point a with respect to the circle $|z| = 1$, and hence lies outside $|z| = 1$). We observe that $\phi_a(\phi_a^{-1}(z)) = z$

$$\phi_a^{-1}(\phi_a(z)) = \frac{\phi_a(z) - a}{1 - \bar{a}\phi_a(z)} = \frac{\frac{z+a}{1+\bar{a}z} - a}{1 - \bar{a}\frac{z+a}{1+\bar{a}z}} = \frac{z+a - a(1+\bar{a}z)}{1+\bar{a}z - \bar{a}(z+a)} = \frac{z+a - a - a\bar{a}z}{1+\bar{a}z - \bar{a}z - \bar{a}a} = \frac{z(1 - |a|^2)}{1 - |a|^2} = z$$

similarly.

86

Thus ϕ_a maps $D(0, 1)$ onto $D(0, 1)$ in a one to one way. Now let θ be a real number. Then $\phi_{\theta} \phi_{\theta^{-1}}$

$$\phi_{\theta} \phi_{\theta^{-1}}(z) = \frac{\phi_{\theta^{-1}}(z) + \theta}{1 + \bar{\theta} \phi_{\theta^{-1}}(z)} = \frac{\frac{z + \theta^{-1}}{1 + \bar{\theta}^{-1}z} + \theta}{1 + \bar{\theta} \frac{z + \theta^{-1}}{1 + \bar{\theta}^{-1}z}} = \frac{z + \theta^{-1} + \theta(1 + \bar{\theta}^{-1}z)}{1 + \bar{\theta}^{-1}z + \bar{\theta}(z + \theta^{-1})} = \frac{z + \theta^{-1} + \theta + \theta \bar{\theta}^{-1}z}{1 + \bar{\theta}^{-1}z + \bar{\theta}z + \bar{\theta} \theta^{-1}} = \frac{z(1 + |\theta|^2) + \theta^{-1} + \theta}{1 + \bar{\theta}^{-1}z + \bar{\theta}z + 1} = \frac{z(1 + |\theta|^2) + \theta^{-1} + \theta}{2}$$

i.e. $\phi_{\theta} \phi_{\theta^{-1}}(z) = z$ i.e., ϕ_a maps $|z| = 1$ on $|z| = 1$. Thus, ϕ_a maps $D(0, 1)$ onto $D(0, 1)$. A maximal problem Let α, β be two complex numbers with $|\alpha| > 1, |\beta| > 1$ and f be analytic on $D(0, 1)$ satisfying $f(\alpha) = \beta$. What is the maximum possible value of $|f(1(0))|$ among such mappings? Solution : Let $g = \phi_{\beta} \circ f \circ \phi_{\alpha}^{-1}$ where ϕ_{β} is defined as in theorem 5.2 (57) Then g maps $D(0, 1)$ to $D(0, 1)$ and satisfies $g(0) = \phi_{\beta}(f(\phi_{\alpha}^{-1}(0))) = \phi_{\beta}(f(\alpha)) = \phi_{\beta}(\beta) = 0$ Thus g satisfies all the conditions of Schwarz's lemma and hence $|g(1(0))| \leq 1$. To obtain an explicit form of $g(1(0))$, we use (57) and apply the chain rule $g(1(0)) = (\phi_{\beta} \circ f) \circ \phi_{\alpha}^{-1}(1(0)) = \phi_{\beta}(f(\phi_{\alpha}^{-1}(1(0)))) = \phi_{\beta}(f(\alpha)) = \phi_{\beta}(\beta) = 0$ But $|g(1(0))| \leq 1$, therefore $|f(1(0))| \leq |\alpha|$ (58)

Equality in (58) occurs only when $|g(1(0))| = 1$. In that case by virtue of Schwarz

87 lemma there is a constant $\lambda, |\lambda| = 1$ so that $g(z) = \lambda z$. Hence, $f(z) = \phi_{\beta}(\lambda \phi_{\alpha}^{-1}(z))$, $z \in D(0, 1)$ (59) We now present an important consequence of Schwarz's lemma, which may be seen as the converse form of theorem 5.2. Theorem 5.3 : Let $f : D(0, 1) \rightarrow D(0, 1)$ be any conformal map of the unit disc onto itself and $f(a) = 0, a \in D(0, 1)$. Then there is a constant $\lambda, |\lambda| = 1$ such that $f(z) = \lambda \phi_a(z)$ where ϕ_a is defined as in theorem 5.2. Proof. Since f is a conformal map from $D(0, 1)$ to $D(0, 1)$, we can have inverse of f, g defined by $g(f(z)) = z$, which is analytic too. Applying the chain rule $g(1(0)) = f^{-1}(0) = a$ (60) But according to inequality (58), f and g have to satisfy $|f(a)| = |f(1(0))| \leq |a|$ (61) (since, $f(a) = 0$ and $g(0) = a$). From (60), (61) it follows that $|f(a)| = (1 - |a|^2)^{-1/2}$. Hence applying the result (59) we find that $f(z) = \lambda \phi_a(z)$ for some λ with $|\lambda| = 1$. Lemma 5.1 : Let G be a simply connected region and $\{f_n\}$ be a sequence of injective analytic mappings (conformal mappings) of G into \mathbb{C} which converges uniformly on every compact subset of G , then the limit function f is either constant or injective. Proof. Suppose f is not constant and not injective. Then there exist two points ζ and $\eta \in G, \zeta \neq \eta$ such that $f(\zeta) = f(\eta) = \omega$, say. Let $g_n(z) = f_n(z) - \omega$. We can find a positive $\delta, \delta > 0$ so that the discs $D(\zeta, \delta)$ and $D(\eta, \delta)$ are included in G . Now $g_n(z) = f_n(z) - \omega$ never vanishes on the circles $|z - \zeta| = \delta$ and $|z - \eta| = \delta$, where $g_n(z) \rightarrow f(z) - \omega$. Applying Hurwitz's theorem, for large n , there exists ζ_n lying inside the circle $|z - \zeta| = \delta$ with $g_n(\zeta_n) = 0$ as $g_n \rightarrow g$ uniformly in G . Similarly, for all large n , there is η_n within $|z - \eta| = \delta$ with $g_n(\eta_n) = 0$. But by construction, $D(\zeta, \delta) \cap D(\eta, \delta) = \emptyset$ and hence $\zeta_n \neq \eta_n$. Thus $g_n(\zeta_n) = g_n(\eta_n) = 0, \zeta_n \neq \eta_n$ that is, $f_n(\zeta_n) = f_n(\eta_n), \zeta_n \neq \eta_n$

88 contradicting the injectivity of each f_n and the proof follows. NOTE : There is no conformal map f of the unit disc $D(0, 1)$ onto the whole complex plane \mathbb{C} because then the inverse function $f^{-1} : \mathbb{C} \rightarrow D(0, 1)$ would be a bounded entire function which is not constant, contradicting the Liouville's theorem. Open mapping theorem : Let G be a region and suppose that f is a non-constant analytic function on G . Then for any open set U in $G, f(U)$ is open. Proof : Omitted. Uniform boundedness : A sequence of functions $\{f_n\}$ defined on a set D is said to be uniformly bounded on D if there exists a constant $M > 0$ such that $|f_n(z)| \leq M$ for all n and for all $z \in D$. Normal family : Let F be a family of functions in a region G . The family F is said to be normal in G if every sequence $\{f_n\}$ of functions $f_n \in F$ contains a subsequence $\{f_{n_k}\}$ which converges uniformly on every compact subset of G . Montel's theorem : A family F in $H(G)$ is normal

if and only if F is uniformly bounded

on every compact subset of G . Proof : Omitted.

Theorem 5.4 : [Riemann Mapping Theorem] Let G be a simply connected region,

except for \mathbb{C} itself and let $a \in G$. Then there is a unique conformal map $f : G \rightarrow D(0, 1)$ of G onto the unit disc which satisfies $f(a) = 0$ and $f'(a) \neq 0$. Proof. Let us first prove that f is unique. If there was another conformal map $g : G \rightarrow D(0, 1)$ with the given properties, then $f \circ g^{-1} : D(0, 1) \rightarrow D(0, 1)$ would be a conformal map and also $(f \circ g^{-1})(0) = f(a) = 0$. Hence, applying Theorem 5.3, we find that there is a constant λ with $|\lambda| = 1$ $(f \circ g^{-1})(z) = \lambda z$. Deriving the derivative at the origin, we find $(f \circ g^{-1})'(0) = \lambda$, $-1 \leq \lambda \leq 1$, $-1 \leq \lambda \leq 1$ $f \circ g^{-1} = f \circ g^{-1} \circ f \circ g^{-1}$, $f \circ g^{-1} = f \circ g^{-1} \circ f \circ g^{-1}$ $f \circ g^{-1} = f \circ g^{-1} \circ f \circ g^{-1}$ $f \circ g^{-1} = f \circ g^{-1} \circ f \circ g^{-1}$ $f \circ g^{-1} = f \circ g^{-1} \circ f \circ g^{-1}$ from which it follows that λ is positive. But also $|\lambda| = 1$, so $\lambda = 1$. Thus $f \circ g^{-1}$ is an identity map and $f = g$. The proof of existence is divided into several stages. Lemma 5.2 Let G be a simply connected region other than \mathbb{C} . Then there exists an injective analytic map f on G with $f(G) \subset D(0, 1)$. Proof. We choose a point $b \in \mathbb{C} \setminus G$. Since G is simply connected there exists a $g : G \rightarrow \mathbb{C}$ analytic with

$$g^2(z) = z - b.$$

89 Here g is injective since

$$g(z_1) = g(z_2) \Rightarrow g^2(z_1) = g^2(z_2) \text{ i.e. } z_1 - b = z_2 - b \Rightarrow z_1 = z_2.$$

By

open mapping theorem $g(G)$ is open. Let us pick $\omega \in g(G)$ and choose $r > 0$ so that $D(\omega, r) \subset g(G)$. Then $D(-\omega, r) \subset \mathbb{C} \setminus g(G)$. For, if there exists a point $\omega \in D(-\omega, r) \cap g(G)$, then $\omega = g(z_1)$ for some $z_1 \in G$ and also $-\omega \in D(\omega, r) \subset g(G)$, so that $-\omega =$

$$g(z_2) \text{ for some } z_2 \in$$

G . Again,

$$g(z_1) = -$$

$$g(z_2) \Rightarrow g^2(z_1) = g^2(z_2) \text{ or, } z_1 - b = z_2 - b \text{ i.e. } z_1 = z_2 \text{ or, } g(z_1) = g(z_2) = -g(z_1) \Rightarrow$$

$$g(z_1) = 0 \Rightarrow 0 = g^2(z_1) =$$

$$z_1 - b$$

$$\text{i.e. } z_1 = b \in$$

$\mathbb{C} \setminus G$ contradicting $z_1 \in G$.

We take $f(z) = z^2 + 2b$

$$g(z) = z^2 + 2b$$

$\omega \in D(0, 1)$ Then f is injective analytic map on G (by construction $|g(z) + \omega| \geq r$ for $z \in G$) and also satisfies $f(z) \in D(0, 1)$ for $z \in G$. Lemma 5.3 : Let G be a simply connected region other than \mathbb{C} itself and let $a \in G$ be fixed. Then there exists a conformal map $f : G \rightarrow D(0, 1)$ of G onto the unit disc with the properties $f(z) = 0$ and $f'(a) \neq 0$. Proof : Let F denote the family of analytic functions $f : G \rightarrow \mathbb{C}$ such that either $f \equiv 0$ or f is injective, and $f(G) \subset D(0, 1)$, $f(a) = 0$ and $f'(a) \neq 0$. Let us consider the function $\psi(z) = (f(z) - f(a)) / (z - a)$ where $f(z)$ is given by (62) of lemma 5.2 and we find that $\psi(G) \subset D(0, 1)$, $\psi(a) = 0$ and $\psi'(a) \neq 0$. So F is non empty and by Montel's theorem it is normal. Applying Lemma 1 we see that all functions in the closure of F in $H(G)$ are either constant or injective. Now since all functions in F take the value zero at a , the same is true for all functions in the closure of F . Likewise the only constant function in the closure is

0 while the other functions in the closure satisfy $f(G) \subset D(0, 1)$. Since $f(G)$ is open, by open mapping theorem, $f(G) \subset D(0, 1)$. Again since the $f \rightarrow f'(a)$ is continuous, all functions in the closure of F must satisfy $f'(a) \geq 0$. The functions in the closure, that are not identically zero, are injective, so $f'(a) \neq 0$ unless $f \equiv 0$. These observations prove that the set F is closed in $H(G)$. Hence F is compact in $H(G)$. Since the map $f \rightarrow f'(a) : F \rightarrow \mathbb{R}$ is a continuous function on a compact set, it must attain its maximum value, as we are not considering constant function (here it is zero). Let $f \in F$ be a function with $f'(a)$ maximum. We now show that $f(G) = D(0, 1)$. On the contrary, suppose that $f(G) \neq D(0, 1)$ and choose $w \in D(0, 1) \setminus f(G)$. Using the property that every non-vanishing analytic function in a simply connected region has an analytic square root, we take a function $h \in H(G)$ with $h^2(z) = z - w$ (63) Now as the bilinear transformation $\phi(z) = \frac{z - w}{z + w}$ maps $D(0, 1)$ onto $D(0, 1)$ and as $f \in F$,

$h \circ f : G \rightarrow \mathbb{C}$ defined by

$$g(z) = h \circ f(z) = \frac{f(z) - w}{f(z) + w}$$

Then clearly, $g(G) \subset D(0, 1)$, $g(a) = 0$ and g is analytic injective and $g'(a) \neq 0$, since $g'(a) = \frac{f'(a)}{2f(a)}$ (64) So, $g \in F$

So, $g \in F$ and $g(G) \subset D(0, 1)$. Let $g : G \rightarrow \mathbb{C}$ defined by

$$g(z) = \frac{f(z) - w}{f(z) + w}$$

Then clearly, $g(G) \subset D(0, 1)$, $g(a) = 0$ and g is analytic injective and $g'(a) \neq 0$, since $g'(a) = \frac{f'(a)}{2f(a)}$ (64) So, $g \in F$

So, $g \in F$ and $g(G) \subset D(0, 1)$. Let $g : G \rightarrow \mathbb{C}$ defined by

$$g(z) = \frac{f(z) - w}{f(z) + w}$$

Then clearly, $g(G) \subset D(0, 1)$, $g(a) = 0$ and g is analytic injective and $g'(a) \neq 0$, since $g'(a) = \frac{f'(a)}{2f(a)}$ (64) So, $g \in F$

93 where $A, B \in \mathbb{C}$, that maps the upper plane $\text{Im } z < 0$ onto the interior of P , with $f(x_1) = w_1, \dots, f(x_{k-1}) = w_{k-1}, f(\infty) = w_k$. (68) Remarks : (i) We do not need to have specific information on w_k and α_k . While travelling the polygon anticlockwise direction we made a left turn of an angle $\pi - \alpha_j$ at the vertex w_j . (ii) Sometimes certain infinite regions can be thought of as infinite polygons. In this case it is convenient to take w_k as the point at infinity, as we need no information on α_k . (iii) It can be shown that Schwarz-Christoffel transformation can be uniquely determined by three points as in the case of bilinear transformation. One of these is used by taking $f(\infty) = w_k$. We can therefore have the freedom to choose two points say, x_1 and x_2 satisfying $x_1 < x_2$. (iv) Note that the integral involved may be impossible to calculate theoretically. In practical problems numerical techniques are often used to evaluate the integral. In first part of the proof we take $f(x_k) = w_k, x_k = \text{finite}$. Proof. By Riemann mapping theorem such a mapping exists. We shall prove that its form is given by (67). So $f(z)$ is analytic for $\text{Im } z < 0$ and $f'(z) \neq 0$ in the upper half plane. From these it is clear that $d dz f(z) / f(z) \log(\dots)$ is analytic in the upper half plane. To construct the function $f(z)$ our aim is to establish that $f'(z)/f(z)$ is analytic for $\text{Im } z \geq 0$ save for the pre-image points of the vertices of the polygon lying on the real axis. Let l be a side of the polygon P , which makes an angle θ (positive sense) with the real-axis and ζ be any point on l but not a vertex of the polygon P . Then for any ω on $l, (\omega - \zeta)e^{-i\theta}$ is real and there is a point z on the real axis of the z -plane so that $f(z) = \omega$ and a corresponding point $z = a$ for ζ on the same line. Hence $\{f(z) - \zeta\}e^{-i\theta}$ is real and continuous on the segment γ of the real axis of the z -plane corresponding to the straight line l of the ω -plane. Moreover, this function is also analytic for $\text{Im } z < 0$, thus following the Schwarz reflection principle we can continue this function analytically across γ to the lower half plane $\text{Im } z > 0$. In particular, this function is analytic in a neighbourhood of the point $z = a$ and can be expanded in the form of the Taylor series. Fig. 55 a θ w -plane l

94 $\{f(z) - \zeta\}e^{-i\theta} = \sum_{k=1}^{\infty} c_k (z - a)^{k-1}$ where $c_1 = f'(a) \neq 0$, maintaining the status quo that $f(a) = \zeta$ and the function f maps the segment γ onto the straight line l . Now $f'(z) = e^{i\theta} \{c_1 + 2c_2(z - a) + \dots\}$ and $\log f'(z) = i\theta + \log\{c_1 + 2c_2(z - a) + \dots\}$ So, $d dz f(z) \log(\dots)$ is analytic in a neighbourhood of $z = a$ and real on a real line segment intercepted by the neighbourhood. Let us consider the case when the point ζ is the corresponding point at infinity on γ (in this case γ is divided into two parts, each of infinite length). Here the Taylor series expansion in the neighbourhood of point at infinity $\{f(z) - \zeta\}e^{-i\theta} = \sum_{k=1}^{\infty} c_k z^{-k}$ where each $c_k \in \mathbb{R}$ is real and $c_1 \neq 0$ (with the same reason mentioned in the finite case). So $f'(z) =$

$f(z) = e^{i\theta} \{c_1 z^{-1} + c_2 z^{-2} + c_3 z^{-3} + c_4 z^{-4} + c_5 z^{-5} + \dots\}$ and we find that $f''(z) = -c_1 z^{-2} - 2c_2 z^{-3} - 3c_3 z^{-4} - 4c_4 z^{-5} - 5c_5 z^{-6} - \dots$

$z f'(z) = e^{i\theta} \{-c_1 + c_2 z^{-1} + c_3 z^{-2} + c_4 z^{-3} + c_5 z^{-4} + \dots\}$

$z f''(z) = e^{i\theta} \{2c_1 z^{-2} - 3c_2 z^{-3} + 4c_3 z^{-4} - 5c_4 z^{-5} + 6c_5 z^{-6} - \dots\}$

$z^2 f''(z) = e^{i\theta} \{2c_1 z^{-1} - 3c_2 + 4c_3 z^{-1} - 5c_4 z^{-2} + 6c_5 z^{-3} - \dots\}$

$z^3 f''(z) = e^{i\theta} \{2c_1 z - 3c_2 z + 4c_3 z^2 - 5c_4 z^3 + 6c_5 z^4 - \dots\}$

$z^k f''(z) = e^{i\theta} \{2c_1 z^{k-1} - 3c_2 z^{k-2} + 4c_3 z^{k-3} - 5c_4 z^{k-4} + 6c_5 z^{k-5} - \dots\}$

(69) $d dz f(z) \log(\dots)$ is analytic in a neighbourhood of the point at infinity and is real when z is real. In the polygon P , let l_1 be an adjacent side to l making an angle α_1 at their point of intersection w_1 . The corresponding point of w_1 on the real axis is x_1 . Here

95 the function $f(z)$ is not analytic in a neighbourhood of x_1 , we choose the branch of the argument so that $\pi < \arg z(x_1 - \omega_1) \leq 2\pi$ introducing a branch cut along the axis $\{x_1 + iy : y \leq 0\}$ [$f(z)$ is not continuous on this branch cut]. Here $\arg\{(\omega_1 - \omega)e^{-i\theta}\}$ is equal to zero or $\alpha_1\pi$ according as ω lies on or to the left of x_1 . So the function $\{f(z) - f(x_1)\} / (\omega_1 - z)^{\alpha_1}$ is real and continuous on the segment of the real axis corresponding to the consecutive sides w_1w_2 and w_2w_3 . Again this function is analytic for $\text{Im } z > 0$ since $f(z) - f(x_1)$ is analytic and non zero there. Expanding $\{f(z) - f(x_1)\} / (\omega_1 - z)^{\alpha_1}$ in Taylor's series in a neighbourhood of x_1 we find $\{f(z) - f(x_1)\} / (\omega_1 - z)^{\alpha_1} = \sum_{k=0}^{\infty} c_k (z - x_1)^k$ where each c_k is real and $c_0 \neq 0$. On simplifying, we find $f(z) - f(x_1) = (\omega_1 - z)^{\alpha_1} \sum_{k=0}^{\infty} c_k (z - x_1)^k = \sum_{k=0}^{\infty} c_k (\omega_1 - z)^{\alpha_1} (z - x_1)^k$ where c_0 is a constant multiple of c_1 , hence not equal to zero. Now we have $f'(z) = \sum_{k=0}^{\infty} c_k (\alpha_1 - k) (\omega_1 - z)^{\alpha_1 - k} (z - x_1)^k + \sum_{k=1}^{\infty} c_k k (\omega_1 - z)^{\alpha_1} (z - x_1)^{k-1}$ where c_0 is a constant multiple of c_1 , hence not equal to zero. Now we have $f'(z) = \sum_{k=0}^{\infty} c_k (\alpha_1 - k) (\omega_1 - z)^{\alpha_1 - k} (z - x_1)^k + \sum_{k=1}^{\infty} c_k k (\omega_1 - z)^{\alpha_1} (z - x_1)^{k-1}$ Fig. 56 Fig. 57 new position after rotation through an angle θ clockwise $\alpha_1\pi$ $\alpha_1\pi$ θ $\alpha_1\pi$ w_1w_2

96 where $F(z)$ is analytic and not zero in a neighbourhood of $z = x_1$ and we obtain $d dz f z z x F z F z \log (\omega_1 - z)^{-\alpha_1} = -\alpha_1 (\omega_1 - z)^{-\alpha_1 - 1} f(z) dz$ This shows that if the polygon P has an angle $\alpha_1\pi$ at a point ω_1 then $d dz f z \log (\omega_1 - z)^{-\alpha_1}$ will have a simple pole of residue $\alpha_1 - 1$ at its corresponding point x_1 . Now if the point at infinity be the corresponding point to ω_1 at which the polygon P has an angle $\alpha_1\pi$, then we can express $\omega_1 = \theta_1 + i\theta_2$ $\dots - / f z e c z c z i = + +$ or, $f z e c z c z c i (\omega_1 - z)^{-\alpha_1} = + + \omega_1 \alpha_1 \theta_1 \theta_2 \dots = + + + f z e c z c z c e c z c z c i i (\omega_1 - z)^{-\alpha_1} = + + + \theta_1 \alpha_1 \theta_2 \alpha_2 \alpha_3 \dots$

$+ e c z c z c i \theta_1 \alpha_1 \alpha_2 \alpha_3 \dots = + + + + + + f z e c z c z c e c z c z c i i (\omega_1 - z)^{-\alpha_1} = + + + + \alpha_1 \alpha_2 \alpha_3 \dots = + + + +$
 $+ - (\omega_1 - z)^{-\alpha_1} \dots e c z c z c i \theta_1 \alpha_1 \alpha_2 \alpha_3 \dots d dz f z f z f z z c z c z c \log (\omega_1 - z)^{-\alpha_1} = + + + + + \alpha_1 \alpha_2 \alpha_3 \dots = + + + + - (-) \dots \alpha_1 \alpha_2 \alpha_3 \dots z c z c = + + = \bullet \sum - \sim \alpha_1 2 1 z c z k k k (71)$ Now since x_2, x_3, \dots, x_k are the corresponding points lying on the real-axis of the z -plane, to the vertices w_2, w_3, \dots, w_k respectively of the polygon P with angles $\alpha_2\pi, \dots, \alpha_k\pi$ there, the function $d dz f z \log (\omega_j - z)^{-\alpha_j}$ will have simple poles with residue $\alpha_j - 1$ at $x_j, j = 2, \dots, k$. Thus we see that this function is analytic for $\text{Im } z > 0$ and continuous on $\text{Im } z = 0$ except the points x_1, x_2, \dots, x_k and using the Schwarz reflection principle it can be continued analytically across the real axis. Hence $d dz f z \log (\omega_j - z)^{-\alpha_j}$ possesses only simple poles at x_1, x_2, \dots, x_k as its only singularities and can be expressed as $d dz f z z x z x z x G z k k \log (\omega_j - z)^{-\alpha_j} = + + + + \alpha_j \alpha_j (72)$ where $G(z)$ is a polynomial. When $|z|$ is large enough $\alpha_j i i i i z x z x z x z i k - - \dots, \dots, 1 1 1 1 2 2 = + + + =$ So, $d dz f z z x z x z x G z i i i k i k i k \log (\omega_j - z)^{-\alpha_j} / (-) / (-) / \dots (\omega_j - z)^{-\alpha_j} = + + + + = \sum \sum \sum \alpha_j \alpha_j = + + = \bullet \sum - (\alpha_j - 1) 2 z d z G z i i i (73)$ Using the property of the sum of the exterior angles of a polygon, $(1 - \alpha_1)\pi + (1 - \alpha_2)\pi + \dots + (1 - \alpha_k)\pi = 2\pi$. Comparing (73) with (69) we get $G(z)$ identically zero. Finally integrating equation (72), we find the desired mapping $f(z)$ as $f z A s x s x s x ds B k z z k (\omega_j - z)^{-\alpha_j} \dots (-) - - - -1 = + 1 1 2 1 1 0 1 2 \alpha_j \alpha_j (74)$ Role of constants A and B (i) $|A|$ controls the size of the polygon (ii) $\arg A$ and B help to select the position, if any, in determining orientation and translation respectively. An useful observation In some occasions we urge to make the evaluation process of the integral in (74) simple. For this sake, we consider the point at infinity corresponds to the vertex w_k where the polygon P has an angle $\alpha_k\pi$. Then we can express [see eq. (71)] $d dz f z z c z k i i \log (\omega_k - z)^{-\alpha_k} = + \bullet \sum \alpha_k (75)$

97 $\alpha_3\pi, \dots, \alpha_k\pi$ there, the function $d dz f z \log (\omega_j - z)^{-\alpha_j}$ will have simple poles with residue $\alpha_j - 1$ at $x_j, j = 2, \dots, k$. Thus we see that this function is analytic for $\text{Im } z > 0$ and continuous on $\text{Im } z = 0$ except the points x_1, x_2, \dots, x_k and using the Schwarz reflection principle it can be continued analytically across the real axis. Hence $d dz f z \log (\omega_j - z)^{-\alpha_j}$ possesses only simple poles at x_1, x_2, \dots, x_k as its only singularities and can be expressed as $d dz f z z x z x z x G z k k \log (\omega_j - z)^{-\alpha_j} = + + + + \alpha_j \alpha_j (72)$ where $G(z)$ is a polynomial. When $|z|$ is large enough $\alpha_j i i i i z x z x z x z i k - - \dots, \dots, 1 1 1 1 2 2 = + + + =$ So, $d dz f z z x z x z x G z i i i k i k i k \log (\omega_j - z)^{-\alpha_j} / (-) / (-) / \dots (\omega_j - z)^{-\alpha_j} = + + + + = \sum \sum \sum \alpha_j \alpha_j = + + = \bullet \sum - (\alpha_j - 1) 2 z d z G z i i i (73)$ Using the property of the sum of the exterior angles of a polygon, $(1 - \alpha_1)\pi + (1 - \alpha_2)\pi + \dots + (1 - \alpha_k)\pi = 2\pi$. Comparing (73) with (69) we get $G(z)$ identically zero. Finally integrating equation (72), we find the desired mapping $f(z)$ as $f z A s x s x s x ds B k z z k (\omega_j - z)^{-\alpha_j} \dots (-) - - - -1 = + 1 1 2 1 1 0 1 2 \alpha_j \alpha_j (74)$ Role of constants A and B (i) $|A|$ controls the size of the polygon (ii) $\arg A$ and B help to select the position, if any, in determining orientation and translation respectively. An useful observation In some occasions we urge to make the evaluation process of the integral in (74) simple. For this sake, we consider the point at infinity corresponds to the vertex w_k where the polygon P has an angle $\alpha_k\pi$. Then we can express [see eq. (71)] $d dz f z z c z k i i \log (\omega_k - z)^{-\alpha_k} = + \bullet \sum \alpha_k (75)$

98 in the neighbourhood of the point at infinity. Again considering the expression of $d dz f z \log (\omega_j - z)^{-\alpha_j}$ in the neighbourhood of the points corresponding to the vertices w_1, w_2, \dots, w_{k-1} [see eq. (70)]. $d dz f z z x z x z x G z k k \log (\omega_j - z)^{-\alpha_j} = + + + + \alpha_j \alpha_j (75.1)$ where $G(z)$ is a polynomial. If $|z|$ is large enough, proceeding as earlier $d dz f z z x z x z x G z i i i k i k i k \log (\omega_j - z)^{-\alpha_j} / (-) / (-) / (-) / (-) = + + + + 1 1 1 1 2 2 1 1 3 1 1 1 = + + + \sum \sum \sum \alpha_j \alpha_j = + + + \bullet \sum - \sim (\alpha_j - 1) \alpha_j k i z d z G z 1 1 2 (76)$ Comparing (76) with (75), $G(z)$ turns out to be identically zero and hence integrating (75.1) we obtain $f z A s x s x s x ds B k z z k (\omega_j - z)^{-\alpha_j} \dots (-) - - - -1 -1 = + 1 1 2 1 1 1 2 0 \alpha_j \alpha_j$ where the role of the constants A and B remain as before. 5.4 Examples : Triangles / Rectangles The Schwarz-Christoffel transformation is expressed in terms of the points x_j , not in terms of their images i.e., the vertices of the polygon. Not more than three points (x_j) can be chosen arbitrarily. If the point at infinity be one of the x_j 's then only two finite points on the real-axis are free to be chosen, whether the polygon is a triangle or a rectangle etc. Triangle Let the polygon be a triangle with vertices w_1, w_2 and w_3 . The S-C transformation is written as $w A s x s x s x ds B z z = + (-) (-) (-) - - - 1 1 2 1 3 1 1 2 3 0 \alpha_1 \alpha_2 \alpha_3 (77)$ where α_1, α_2 and α_3 are the internal angles at the respective vertices. Fig. 58 Fig. 59 w -plane z -plane $x_1 x_2 x_3 w_1 w_2 w_3 \alpha_3 \pi \alpha_1 \pi \alpha_2 \pi$

99 Here we have chosen all the three finite points x_1, x_2, x_3 on the real-axis. The constants A, B control the size and position of the triangle respectively. If we take the vertex w_3 as the image of the point at infinity, the S-C transformation becomes $w = A s^2 + B z + \dots$ (78) Here x_1 and x_2 can be chosen arbitrarily. Example 1 : Find a

Schwarz-Christoffel transformation that maps the upper half-plane to the inside of the triangle with vertices $-1, 1$

and $\sqrt{3}i$. Solution : Following our notation, we write $w_1 = -1, w_2 = 1$ and $w_3 = \sqrt{3}i$ so that $\alpha_1 = \alpha_2 = \alpha_3 = 1/3$. We choose the form (78) of S-C transformation and consider the mapping. $f(z) = A s^2 + B z + \dots$ [here $f(\infty) = \sqrt{3}i$] We may choose $x_1 = -1$ and $x_2 = 1$, so that $f(-1) = -1$ and $f(1) = 1$. Therefore $f(z) = A s^2 + B z + \dots$ It then follows that $A = \dots$ and $B = \dots$ We obtain $A = \dots$ and $B = \dots$ Hence Fig. 60 Fig. 61 $-1, 1, \sqrt{3}i$

100 $f(z) = A s^2 + B z + \dots$ Example 2 : Using Schwarz-Christoffel transformation map the upper half-plane onto an equilateral triangle of side 5 units. Solution : It is convenient to choose three arbitrary points $x_1 = -1, x_2 = 1$ and $x_3 = \infty$ which are mapped into the vertices of the equilateral triangle, so we take S-C transformation (78). $f(z) = A s^2 + B z + \dots$ Here, $f(-1) = w_1 = 0$ and $f(1) = w_2 = 5$. So that $A = \dots$ Hence the desired transformation is $f(z) = A s^2 + B z + \dots$ Alternative : We take $z_0 = -1, A = 1, B = 0$ and find S-C transformation as, (choosing one of x_i 's as point at infinity) $w = A s^2 + B z + \dots$ (79) taking $x_1 = -1$ and $x_2 = 1$. Then $\sim(z) \sim f(w) = \dots$, and the image of the point $z = -1$ is the point $w = 0$. When $z = 1$ in the integral we can write $s = x$, where $-1 < x < 1$. Then $x + 1 > 0$ and $\text{Arg}(x+1) = 0$, while $|x-1| = 1-x$ and $\text{Arg}(x-1) = \pi$. Hence $\sim(z) = \dots$ Fig. 62 Fig. 63 $(-1, 1, w_3 = \pi/3 - \pi/3)$ $101 = \dots$ substituting $x = \sqrt{t} = \dots$ We choose w_2 as, $w_2 = 5 = \dots$ where $k \in \mathbb{R} = -5, \dots$ To find w_3 let us first calculate for $\sim(w_3) = \dots$

x
 $\int dx \dots = \dots$
 $\int dx \dots = \dots$
 $\int dx \dots = \dots$

Now, the value of $\sim(w_3)$ can also be represented by the integral $\int \dots$ when z tends to infinity along the negative real axis. Thus from the above relation, we have $\sim(w_3) = \dots$ i.e., $\sim(w_3) = \dots$ So, $w_3 = 5e^{i\pi/3} = \dots$

102 Therefore, the three vertices of the equilateral triangle are $w_1 = 0, w_2 = 5$ and $w_3 = 5e^{i\pi/3}$. Clearly each of its side is of length 5 unit. The desired transformation is then $f(z) = A s^2 + B z + \dots$ which is same as obtained in the first process. Remark : Following the above technique we can determine a S-C transformation from $\text{Im } z \geq 0$ onto a triangle, in particular, whose one side opposite to an angle is given. Rectangle :

Example 3 : Find a S-C transformation that maps the upper half of the z -plane to the inside of the rectangle in the w -plane

with vertices $-a, a, a + ib$ and $-a + ib$ which are the preimages of $-1, 1, \alpha$ and $-\alpha$ respectively. Solution : Let us first make the identification of the vertices of the rectangle $w_1 = -a + ib, w_2 = -a, w_3 = a, w_4 = a + ib$ $\alpha_1 = \alpha_2 = \alpha_3 = \alpha_4 = 1/2$ We choose $x_1 = -\alpha, x_2 = -1, x_3 = 1, x_4 = \alpha$ where $\alpha < 1$ will be determined later. We are attempting to benefit from the symmetry here, which requires the image $z = 0$ to be $w = 0$. So taking $z_0 = 0$ we get $B = 0$ in the formula (74) for S-C transformation, which reduces to $f(z) = A s^2 + \dots$ Fig. 64 Fig. 65 $\wedge \wedge \wedge \wedge \wedge \wedge \wedge \wedge$
 $\wedge \wedge -a + ib, a + ib, -a, a$

103 $\equiv A ds s s z z [(-)(-)] ((,)) 1 2 2 2 0 \alpha \varphi \alpha$ (80) The constant A may be found by using the fact that $f(1) = a$ i.e., $a A ds s s a ds s s = [(-)(-)] / [(-)(-)] 1 1 2 2 2 0 1 2 2 2 0 1 \alpha \alpha$ or $A = a/\varphi(\alpha)$, say (81) To find α , we apply $f(\alpha) = a + ib$, $a ib a ds s s a + = \varphi \alpha \alpha () [(-)(-)] 1 2 2 2 0 = + a ds s s i ds s s \varphi \alpha \alpha \alpha () [(-)(-)] [(-)(-)] 1 1 2 2 2 2 2 1 0 1$ from which, equating imaginary parts, we arrive at $b ds s s \varphi \alpha \alpha \alpha () [(-)(-)] = 2 2 2 1 1$ Since a and b are known, this equation determines α , which gives rise to the evaluation of $\varphi(\alpha)$ i.e. A is completely known. Note : The function $\varphi(z, \alpha)$, given in (80), which involves z as the upper limit of an integral, is called an elliptic integral of the first kind and it is not an elementary function. The real definite integral $\varphi(\alpha)$ in (81) is called a complete elliptic integral of the first kind.

Example 4 :

Find a Schwarz-Christoffel transformation that maps the upper half of the z-plane to the

vertical semi-infinite strip $-\pi/2 < u < \pi/2$,

$v > 0$ of the w-plane. Solution : Fig. 66 Fig. 67 w-plane z-plane $-1 1 0 - - \pi 2 - \pi 2$

104 Here we take $x_1 = -1$, $x_2 = 1$ and $x_3 = \infty$ and the image points are $w_1 = -\pi/2$ and $w_2 = \pi/2$ respectively, so that a S-C transformation can be written as $f z A s s ds B z z () () (-) -1/ -1/ = + + 1 1 2 2 0 = + A s B z z 1 1 2 1 2 0 (-) / ds = + \sim \log - \sim A iz z B 1 2$ Using $f(-1) = -\pi/2$ and $f(1) = \pi/2$ we find $f z i iz z () - \log - , = + 1 2$ Choosing a suitable branch of the logarithm.

105 Unit 6 Entire and Meromorphic Functions Structure 6.0 Objectives 6.1 Entire function 6.2 Infinite Products 6.3 Infinite product of functions 6.4 Weierstrass Factorization 6.5 Counting zeros of analytic functions 6.6 Convex functions 6.7 Order of an entire function 6.8 The function $n(r)$ 6.9 Convergence exponent 6.10 Canonical Product 6.11 Hadamard's Factorization Theorem 6.12 Consequences of Hadamard's Theorem 6.13 Meromorphic functions 6.14 Partial Fraction Expansions of Meromorphic Functions 6.15 Partial Fraction Expansion of Meromorphic functions Using Residue theorem 6.16 The Gamma Function 6.17 A few properties of $\Gamma\Gamma\Gamma\Gamma(z)$ 6.0 The Objectives of the Chapter In this chapter we shall study entire functions, their growth properties and meromorphic functions. Infinite products and their convergence will be discussed. Properties of zeros of

106 an entire function, convex functions, gamma function and its important properties will also be discussed. 6.1 Entire function A function $f(z)$

analytic in the finite complex plane is said to be entire (or

sometimes integral) function. Clearly, the sum, difference and product of two or more entire functions are entire functions. Examples : The polynomial function $P(z) = a_0 + a_1 z + \dots + a_n z^n$, exponential function e^z , $\sin z$, $\cos z$ etc. are entire functions. Let us consider the first example, the polynomial function. It is evident that $P(z)$ can be uniquely expressed as a product of linear factors in

the form

$$A z z z z z z n 0 1 2 0 1 1 1 0 - - - \neq , \text{ if } a \text{ or } A z z z z z$$

a
a
a

$p p n p p p 1 1 1 0 0 1 2 0 1 1 - - - = = \neq - - \zeta \zeta \zeta , , , \text{ if } a$ (82) where A_0 (or, A_p) is constant and $z = z_1, z_2, \dots, z_n$ (or, $z = 0, \zeta_1, \zeta_2, \dots, \zeta_{n-p}$) are the zeros of $P(z)$, multiple zeros are counted according to their multiplicities. There arises a natural question : whether any entire function can be expressed in a similar manner in terms of its zeros. The observations are as follows : (i) There may exist entire function which never vanishes, (ii) If an entire function possesses finite number of zeros, then it is always possible to express it in the form (82) stated above. But when the number of zeros are infinite the form (82) reduces to a product of infinite number of linear factors which need not always be convergent. We first consider infinite products of complex numbers and functions. 6.2 Infinite Products An infinite product is an expression of the form $p n n = \bullet \prod 1$ (83)

107 where $p_1, p_2, \dots, p_n, \dots$ are non-zero complex factors. If we allow any of the factors be zero, it is evident that the infinite product would be zero regardless of the behaviour of the other terms.

Let $P_n = p_1 p_2 \dots p_n$. If P_n tends to a finite limit (non-zero) p as n tends to infinity, we say that the infinite product (83) is convergent and write as $p = \prod_{n=1}^{\infty} p_n$ (84). An infinite product which does not tend to a non-zero finite limit as n tends to infinity is said to be divergent. To find the necessary condition for convergence for the infinite product $p = \prod_{n=1}^{\infty} p_n$, say (84) holds, then writing $p_n = 1 + a_n$ we conclude in view of (84) that

$$\lim_{n \rightarrow \infty} a_n = 0 \quad (85)$$

Thus, $\lim_{n \rightarrow \infty} a_n = 0$ (85) is a necessary condition for convergence of the infinite product (83). It is then better to write the product as $\prod_{n=1}^{\infty} (1 + a_n)$ (86) so that $a_n \rightarrow 0$ as $n \rightarrow \infty$ is a necessary condition for convergence. Theorem 6.1: The infinite product (86) converges if and only if $\sum_{n=1}^{\infty} a_n$ converges.

We use the principal branch of the log function and omit, as usual, the terms with $a_n = -1$.

Proof. Let $P_n = \prod_{k=1}^n (1 + a_k)$ and $S_n = \sum_{k=1}^n a_k$. Then $\log P_n = S_n$ and $P_n = e^{S_n}$. Now if the given series is convergent i.e. $S_n \rightarrow S$ as $n \rightarrow \infty$, P_n tends to the limit $P = e^S (\neq 0)$. This proves the sufficiency of the condition. Conversely, assume that the product converges i.e. $P_n \rightarrow P (\neq 0)$ as $n \rightarrow \infty$. We shall show, by virtue of $P_n = e^{S_n}$, that the series (87) converges to some value of $\log P$, not necessarily the principal value of $\log P$. For $n > N$, $P_n \rightarrow P$ and $P_{n+1} \rightarrow P$, so $\frac{P_{n+1}}{P_n} \rightarrow \frac{P}{P} = 1$. Now there exists an integer K such that $\log \frac{P_{n+1}}{P_n} = \log \frac{1 + a_{n+1}}{1 + a_n} = -2\pi i k_n$ (88). To establish the convergence of the sequence $\{k_n\}$, we form the difference $(k_{n+1} - k_n) \log P$.

$$\log \frac{P_{n+1}}{P_n} - \log \frac{P_n}{P_{n-1}} = \log \frac{1 + a_{n+1}}{1 + a_n} - \log \frac{1 + a_n}{1 + a_{n-1}}$$

and that $k_{n+1} - k_n = \frac{a_{n+1} - a_n}{1 + a_n}$ tends to zero as $n \rightarrow \infty$, and let the limit of the sequence $\{k_n\}$ be k . Taking limit in (88), we find that $\log P = 2\pi i k$ and so the condition assumed is necessary. Definition: An infinite product $\prod_{n=1}^{\infty} (1 + a_n)$ is absolutely convergent if and only if $\sum_{n=1}^{\infty} |a_n|$ is convergent. Theorem 6.2: The infinite product (86) converges absolutely if and only if the series $\sum a_n$ converges absolutely. Proof: If $\sum a_n$ converges absolutely, then in particular $a_n \rightarrow 0$ as $n \rightarrow \infty$. Also, if $\sum |a_n|$ converges absolutely then $\sum a_n$ converges absolutely and thus in

either of the cases $a_n \rightarrow 0$ and we can take $|a_n| \leq \frac{1}{2}$ for sufficiently large n . Then by elementary calculation,

$$|a_n| \leq \frac{1}{2} \implies |1 + a_n| \leq \frac{3}{2} \implies \log |1 + a_n| \leq \frac{3}{2} |a_n|$$

, $n =$ large enough. It follows that $\sum |a_n|$ converges and $\sum a_n$ converges absolutely.

6.3 Infinite product of functions So far we have considered infinite product of complex numbers. Now we shall study infinite products whose factors are functions of a complex variable. Some of the factors (finite in number) may vanish on a region considered. In that case we consider the infinite product omitting those factors. The theorems proved earlier hold good in this case too with some modifications. Definition: (Uniform convergence of infinite products) An infinite product $\prod_{n=1}^{\infty} (1 + a_n(z))$ (89) where the functions $a_n(z)$ are defined on a region D ,

is said to be uniformly convergent on D if the sequence of partial products

$P_n(z) = \prod_{k=1}^n (1 + a_k(z))$ converges uniformly to a non-zero limit on D . Theorem 6.3: An infinite product (89) is uniformly convergent on a domain D if the series $\sum_{n=1}^{\infty} |a_n(z)|$ converges uniformly and has a bounded sum there.

Proof: Let M be the upper bound of the sum $\sum_{n=1}^{\infty} |a_n(z)|$ on D . Then $|1 + a_n(z)| \leq 1 + |a_n(z)| \leq 1 + \frac{M}{n}$

$$\prod_{n=1}^{\infty} (1 + a_n(z)) \leq \prod_{n=1}^{\infty} \left(1 + \frac{M}{n}\right)$$

Let us consider the sequence $\{Q_n\}$ with $Q_n = \prod_{k=1}^n (1 + \frac{M}{k})$

We observe $Q_n \leq Q_{n+1}$

z

$a_n z^n$

$n \rightarrow \infty$... $\sum_{n=1}^{\infty} a_n z^n$ is uniformly convergent.

e^z

$M_n(z)$

Now since

the series $\sum a_n z^n$

is uniformly convergent, the series $\sum_{n=1}^{\infty} \frac{Q_n(z)}{Q_n(z)}$ is uniformly convergent.

Thus the sequence $\{Q_n\}$ tends to a limit. Again $P(z) = \prod_{n=1}^{\infty} (1 - \frac{z}{z_n}) e^{\frac{z}{z_n} + \frac{z^2}{2z_n^2} + \dots}$ so the result follows.

Theorem 6.4 : An infinite product $\prod_{n=1}^{\infty} (1 + a_n)$ converges uniformly and absolutely in a closed bounded domain D if each function $a_n(z)$ satisfies $|a_n(z)| \leq M_n$ for all $z \in D$ and M_n is independent of z and moreover $\sum M_n$ is convergent. Proof : Given $\sum M_n$ is convergent, so the infinite product $\prod_{n=1}^{\infty} (1 + M_n)$ converges by theorem 6.2

Now, for

$n < m$

$z \in D$

$z \in D$

$a_n z^n$

$\prod_{k=1}^m (1 + a_k(z)) = \prod_{k=1}^m (1 + a_k(z)) \dots$ (90) Again, $\prod_{k=1}^m (1 + a_k(z)) \dots$

a_n

z

$a_n z^n$

$a_n z^n$

a_n

z

k

m

n

$\prod_{k=1}^m (1 + a_k(z)) \dots$

$a_n z^n$

m

$\prod_{k=1}^m (1 + a_k(z)) \dots$ Taking moduli $\prod_{k=1}^m (1 + |a_k(z)|) \dots$

$\prod_{k=1}^m (1 + |a_k(z)|) \dots$

$\prod_{k=1}^m (1 + |a_k(z)|) \dots$ Utilising this in (90) we obtain

$\prod_{k=1}^m (1 + |a_k(z)|) \dots$

$\prod_{k=1}^m (1 + |a_k(z)|) \dots$

$\prod_{k=1}^m (1 + |a_k(z)|) \dots$ Now as the infinite product $\prod_{k=1}^{\infty} (1 + |a_k(z)|)$ is

convergent, we choose m large enough so that r.h.s in (91) is less than ϵ and hence $|Q_n(z) - Q_m(z)| < \epsilon$, when $n > m$

Thus the sequence $\{Q_n(z)\}$ converge uniformly, since m depends only on ϵ . Finally, absolute convergence of the

infinite product follows on utilising Th. 6.2 Example 1 : Test for convergence of the infinite product $\prod_{n=1}^{\infty} (1 - \frac{z}{2^n})$

Solution : The terms of the product vanish when $z = \pm 2^n, \dots$ etc. Here $a_n z^n = -\frac{z}{2^n}$ and Now

since the series $\sum \frac{1}{2^n}$ is convergent, the given infinite product is uniformly and absolutely convergent in the entire plane

excluding the points $z = \pm 2^n, \dots$ etc. Example 2 : Discuss the convergence of the infinite product $\prod_{n=1}^{\infty} (1 + \frac{z}{n^2})$

Solution : Let $P_n(z) = \prod_{k=1}^n (1 + \frac{z}{k^2})$ and we consider a bounded closed domain D which does not contain

the points $z = \pm 2^n, \dots$. The sequence $\{P_n(z)\}$ converges uniformly in D (see example 1). Again

let

F

z

z

$z z z$

$\prod_{k=1}^n (1 + \frac{z}{k^2}) = \prod_{k=1}^n (1 + \frac{z}{k^2}) \dots$, then $F z P z z z$

n

$P z$

n
 n

$$\sum_{n=1}^{\infty} \frac{1}{n^2} = \frac{\pi^2}{6}$$

and

F

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and obviously the sequences F_2, F_4, F_6, \dots and F_1, F_3, F_5, \dots

converge uniformly in D . Hence the given infinite product converges uniformly in D . To test for the absolute

convergence of the given product we notice that $\sum_{n=1}^{\infty} \frac{1}{n} = \infty$

and it is divergent since the series on the right is divergent and $|z|$ is finite. Therefore the given product does not converge absolutely. Considering the theorem 4.4 on uniformly convergent sequence of analytic functions [(14) Page-72] we get the following theorem : Theorem 6.5 : If an infinite product $\prod (1 + f_n(z))$ converges uniformly to $f(z)$ in a bounded closed domain D and if each function $f_n(z)$ is analytic in D , then $f(z)$ is also analytic in D . 6.4

Weierstrass' Factorization Theorem 6.6 : If $f(z)$ is an entire function and never vanishes on C , then $f(z)$ is of the form $f(z) = e^{g(z)}$, or, more generally, $f(z) = ce^{g(z)}$, $c \neq 0$, constant. where $g(z)$ is also an entire function. Proof : Since f is entire and never vanishes on C , f'/f is also entire and is thus the derivative of an entire function $g(z)$. [follows from Result 1, PG(MT) 02-complex analysis [14, page-54]. Then $f' = f'fg$ i.e. $f' = f'fg$ Now, $(f'/f)' = f''/f - (f')^2/f^2 = 0$ Hence, $f(z) = ce^{g(z)}$

proving the result. Assume now that f possesses finitely many zeros, a zero of order $m < 0$ at the origin, and the non-zero

ones, possibly repeated are a_1, \dots

a_n . Then $f(z)$

$z^m e^{g(z)}$

$f(z) = z^m e^{g(z)} \prod_{k=1}^n (1 - \frac{z}{a_k})$ where g is

entire. This is clear, since if we divide f by the

factors which produce zero at the points $z = 0, a_1, \dots, a_n$ we get an entire function with no zeros. However we cannot expect, in general, such a simple formula to hold

in the case of

infinitely many zeros. Here we have to take care of convergence problems

for an infinite product. In fact the obvious generalization.

Theorem 6.7 (Weierstrass' Factorization Theorem) :— Let $\{a_n\}$ be a sequence of complex numbers with the property $a_n \rightarrow \infty$ as $n \rightarrow \infty$. Then it is possible to construct an entire function $f(z)$ with zeros precisely at these points. Proof : We need Weierstrass' primary factors to construct the desired function. The expressions

$E(z, 0) = 1 - z$, $E(z, p) = (1 - z) e^{z + \frac{z^2}{2p} + \dots + \frac{z^p}{p}}$, $p = 1, 2, \dots$, are called Weierstrass' primary factors. Each primary factor is an entire function having only one

simple zero at $z = 1$. Now, when $|z| > 1$

$|z| > 1$
we have, log

$E(z, p) = \log(1 - z) +$

$$\frac{z}{p} + \frac{z^2}{2p^2} + \dots + \frac{z^p}{p^p}$$

$$= \log(1 - z) + \frac{z}{p} + \frac{z^2}{2p^2} + \dots + \frac{z^p}{p^p}$$

p

Here we have taken the principal branch of $\log(1 - z)$. Hence if $z \in D$

$z \in D$

$$\log(1 - z) + \frac{z}{p} + \frac{z^2}{2p^2} + \dots + \frac{z^p}{p^p} \leq \frac{1}{p} + \frac{1}{2p^2} + \dots + \frac{1}{p^p} \leq \frac{1}{p-1} \quad (92)$$

We may suppose that the origin is not a zero of the entire function $f(z)$ to be constructed so that $a_n \neq 0$ for all n . For, if origin is a zero of $f(z)$ of order m we

need only multiply the constructed function by z^m . We also arrange the zeros in order of non-decreasing modulus (if several distinct points a_n have the same modulus, we take them in any order) so that $|a_1| \leq |a_2| \leq \dots$. Let $|a_n| = r_n$. Since $r_n \rightarrow \infty$ we can always find a sequence of positive integers $m_1, m_2, \dots, m_n, \dots$ such that the series $\sum_{n=1}^{\infty} \frac{r_n^{m_n}}{z^{m_n}}$ converges for all positive values of r .

114 In fact, we may take $m_n = n$ since for any given value of r , we have $\frac{r^n}{z^n} < \frac{1}{2^n}$ for all sufficiently large n and the series

is therefore convergent. Next we take an arbitrary positive number R and choose the integer N such that $\frac{r^N}{R^N} \leq \frac{1}{2}$. Hence, when $z \leq R$ and $n > N$ we have, $\frac{r^n}{z^n} \leq \frac{1}{2^n}$ and so by (92), $\log \left| \prod_{n=N+1}^{\infty} \left(1 - \frac{a_n}{z}\right)^{m_n} \right| < \sum_{n=N+1}^{\infty} \frac{r^n}{z^n} < \frac{1}{2}$ and so by (92), $\log \left| \prod_{n=N+1}^{\infty} \left(1 - \frac{a_n}{z}\right)^{m_n} \right| < \frac{1}{2}$

By Weierstrass' M-test the series

$\log \left| \prod_{n=1}^{\infty} \left(1 - \frac{a_n}{z}\right)^{m_n} \right|$ converges absolutely and uniformly when $z \leq R$ and so the infinite product $E(z) = \prod_{n=1}^{\infty} \left(1 - \frac{a_n}{z}\right)^{m_n}$, converges absolutely and uniformly in the disc $|z| \leq R$, however large R may be. Hence the above product represents an entire function, say $G(z)$. Thus, $G(z) = \prod_{n=1}^{\infty} \left(1 - \frac{a_n}{z}\right)^{m_n}$ (93)

With the same value of R , we choose another integer k such that $\frac{r^k}{R^k} \leq \frac{1}{2}$. Then each of the functions of the sequence $E_k(z) = \prod_{n=1}^k \left(1 - \frac{a_n}{z}\right)^{m_n} = \left(1 - \frac{a_1}{z}\right)^{m_1} \dots \left(1 - \frac{a_k}{z}\right)^{m_k}$, vanish at the points a_1, \dots, a_k and nowhere else in $|z| \leq R$. Hence by Hurwitz's theorem the only zeros of G in $|z| \leq R$ are a_1, \dots, a_k . Since R is arbitrary, this implies that the only zeros of G are the points of the sequence $\{a_n\}$. Now, if origin is a zero of order m of the required

entire function $f(z)$, then $f(z)$ is

of the form $f(z) = z^m G(z)$. Again, for any entire function

$g(z)$, $e^{g(z)}$ is also an entire function without any zero. Hence the general form of the required entire function $f(z)$ is

$$f(z) = z^m e^{g(z)} \prod_{n=1}^{\infty} \left(1 - \frac{a_n}{z}\right)^{m_n} e^{\frac{a_n}{z}}$$

where

$$\sum_{n=1}^{\infty} \frac{a_n^{m_n}}{z^{m_n}} < \infty \quad (94)$$

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Remark : As there are many possible sequences $\{m_n\}$ in the construction of the function $G(z)$ and ultimately of $f(z)$, the form of the function $f(z)$ achieved is not unique. 6.5

Counting zeros of analytic functions The rate of growth of an entire function is closely related to the density of zeros. We have a quite effective

formula in this regard due to J.L.W.V. Jensen, a Danish mathematician who discovered it in the year 1899. Theorem 6.8

[Jensen's Formula] :- Let $f(z)$ be analytic on $|z| \leq R$, $f(0) \neq 0$ and $f(z) \neq 0$ on $|z| = R$. If a_1, \dots, a_n be the zeros of $f(z)$ within the circle $|z| = R$, multiple zeros being repeated according to their multiplicities, then $\log \left| \frac{f(z)}{z^k} \right| = \log |f(z)| - k \log |z|$

Proof : Let $\varphi(z) = \frac{f(z)}{z^k}$. The zeros of the denominator of $\varphi(z)$ are also the zeros of $f(z)$ of the same order. Hence the zeros of $f(z)$ cancels the poles a_n in the product and so $\varphi(z)$ is analytic on $|z| \leq R$. Also, $\varphi(z) \neq 0$ on $|z| \leq R$. For, if $\frac{f(z)}{z^k} = 0$ then $\frac{f(z)}{z^k} = 0$ is the inverse point of a_k with respect to the circle $|z| = R$ and so lies outside the circle. Again, $\varphi(z) = \prod_{n=1}^{\infty} \left(1 - \frac{a_n}{z}\right)^{m_n} e^{\frac{a_n}{z}}$

Now, when $|z| = R$ we have, $\log \left| \frac{f(z)}{z^k} \right| = \log |f(z)| - k \log |z| = \log |f(z)| - k \log R$

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124 Hence, $|\varphi(z)| = |f(z)|$ on $|z| = R$.

Since $\varphi(z)$ is analytic and non-zero on $|z| \leq R$, $\log \varphi(z)$ is also analytic on $|z| \leq R$ and consequently $\text{Re} \log \varphi(z) = \log |\varphi(z)|$ is harmonic on $|z| \leq R$. Hence by Gauss' mean value theorem, $\log \left| \frac{f(z)}{z^k} \right| = \log |f(z)| - k \log |z|$ (98)

116 From (97) we have, $\varphi(z) = \frac{1}{2\pi} \int_0^{2\pi} f(re^{i\theta}) d\theta$. Hence from (98) we get, $\log \varphi(z) = \frac{1}{2\pi} \int_0^{2\pi} \log |f(re^{i\theta})| d\theta$. Note: We observe that Jensen's formula can also be expressed as $\log \varphi(z) = \frac{1}{2\pi} \int_0^{2\pi} \log |f(re^{i\theta})| d\theta - \sum_{i=1}^n \log \frac{r}{|a_i|}$ (99) or as, $\log \varphi(z) = \frac{1}{2\pi} \int_0^{2\pi} \log |f(re^{i\theta})| d\theta - \sum_{i=1}^n \log \frac{r}{|a_i|}$ (100) where $|a_i| = r_i, i = 1, \dots, n$. Theorem 6.9 (Jensen's inequality):— Let $f(z)$ be analytic on $|z| \leq R, f(0) \neq 0$ and $f(z) \neq 0$ on $|z| = R$. If a_1, \dots, a_n be the zeros of $f(z)$ within $|z| = R$, multiple zeros being repeated according to their multiplicities, and $|a_i| = r_i, i = 1, \dots, n$, then $\log |f(0)| \leq \frac{1}{2\pi} \int_0^{2\pi} \log |f(re^{i\theta})| d\theta - \sum_{i=1}^n \log \frac{r}{|a_i|}$ (101) where $M(r) = \max_{|z|=r} |f(z)|$. Proof: As in Jensen's formula (theorem 6.8) we have, $\log |f(z)| = \frac{1}{2\pi} \int_0^{2\pi} \log |f(re^{i\theta})| d\theta - \sum_{i=1}^n \log \frac{r}{|a_i|}$ on $|z| = R$ and so by the maximum modulus theorem, $|f(z)| \leq M(R)$ for $|z| \leq R$. In particular, $|f(0)| \leq M(R)$ i.e. $\log |f(0)| \leq \log M(R)$.

Theorem 6.10 (Poisson-Jensen formula):— Let $f(z)$ be analytic on $|z| \leq R, f(0) \neq 0$ and $f(z) \neq 0$ on $|z| = R$. If a_1, \dots, a_n be the zeros of $f(z)$ within the circle $|z| = R$, multiple zeros being repeated according to their multiplicities, then for any $z = re^{i\theta}, r < R$, $\log |f(z)| = \frac{1}{2\pi} \int_0^{2\pi} \log |f(re^{i\theta})| d\theta - \sum_{i=1}^n \log \frac{r}{|a_i|} + \sum_{i=1}^n \log \frac{r}{|a_i|} \frac{R^2 - r^2}{|R^2 - \bar{a}_i z|}$.

117 Proof: Let $\varphi(z) = \frac{f(z)}{f(0)}$. Then, as in Jensen's formula we have, $|\varphi(z)| = |f(z)|/|f(0)|$ on $|z| = R$. Since $\varphi(z)$ is analytic and non-zero on $|z| \leq R$, $\log \varphi(z)$ is also analytic on $|z| \leq R$ and consequently $\log |\varphi(z)|$ is harmonic on $|z| \leq R$. So, by Poisson's integral formula, $\log |\varphi(z)| = \frac{1}{2\pi} \int_0^{2\pi} \log |f(re^{i\theta})| d\theta - \sum_{i=1}^n \log \frac{r}{|a_i|} + \sum_{i=1}^n \log \frac{r}{|a_i|} \frac{R^2 - r^2}{|R^2 - \bar{a}_i z|}$ (102) Now, $\log |\varphi(z)| = \log |f(z)| - \log |f(0)|$ on $|z| = R$ we get from (102) $\log |f(z)| = \frac{1}{2\pi} \int_0^{2\pi} \log |f(re^{i\theta})| d\theta - \sum_{i=1}^n \log \frac{r}{|a_i|} + \sum_{i=1}^n \log \frac{r}{|a_i|} \frac{R^2 - r^2}{|R^2 - \bar{a}_i z|} + \log |f(0)|$ (103)

6.6 Convex functions The property of convexity plays an important role in function theory because in several cases some lead factors associated with entire, meromorphic and subharmonic functions appear to be convex functions. A real-valued function φ defined on the interval $I = [a, b]$ is said to be convex if for any two points s, u in $[a, b]$ $\varphi(\lambda s + (1-\lambda)u) \leq \lambda \varphi(s) + (1-\lambda)\varphi(u)$ for $0 \leq \lambda \leq 1$ (104) Geometrically, the condition (104) is equivalent to the condition that if $s < x < u$, then the point $(x, \varphi(x))$ should lie below or on the chord joining the points $(s, \varphi(s))$ and $(u, \varphi(u))$ in the plane. Analytical condition for $\varphi(x)$ to be convex in $[a, b]$:— Let the coordinates of the points A, B, C on the curve $y = \varphi(x)$ as shown in the adjoining figure be $(s, \varphi(s)), (u, \varphi(u))$ and $(x, \varphi(x))$ respectively where $s < x < u$.

118 Equation of the chord AB is $y - \varphi(x) = \frac{\varphi(u) - \varphi(s)}{u - s} (x - s)$ or, $y - \varphi(x) = \frac{\varphi(u) - \varphi(s)}{u - s} (x - s) + \varphi(x)$ (105) Let the coordinates of any point D on the chord AB be (x, y) . According to definition $\varphi(x)$ will be convex if and only if $CN \leq DN$ i.e., if and only if $\varphi(x) \leq y$; i.e. if and only if $\varphi(x) \leq \frac{\varphi(u) - \varphi(s)}{u - s} (x - s) + \varphi(x)$ (106) for $s < x < u$. We now state two results on convex functions without proof. Result 1. A differentiable function $f(x)$ on $[a, b]$ is convex if and only if $f'(x)$ is increasing in $[a, b]$. Result 2. A sufficient condition for $f(x)$ to be convex is that $f''(x) \geq 0$.

The maximum modulus function: Let $f(z)$ be a non-constant analytic function in $|z| > R$. Then for $0 < r < R$ we define the maximum modulus function $M(r, f)$ or, simply $M(r)$ by $M(r) = \max_{|z|=r} |f(z)|$. By maximum modulus theorem we can also write $M(r) = \max_{|z|=r} |f(z)|$. Result: Let $f(z)$ be a non-constant analytic function in $|z| > R$. Then $M(r)$ is a strictly increasing function of r in $0 < r < R$. Proof: Let $0 < r_1 < r_2 < R$. Since $f(z)$ is analytic in $|z| \leq r_2$, the maximum value of $|f(z)|$ for $|z| \leq r_2$

is attained on $|z| = r_2$. Let z_2 be a point on $|z| = r_2$ such that $|f(z_2)| = M(r_2)$. Similarly, the maximum value of $|f(z)|$ for $|z| \leq r_1$ is attained on $|z| = r_1$. Let z_1 be a point on $|z| = r_1$ such that $|f(z_1)| = M(r_1)$. Since $r_1 < r_2, z_1$ is an interior point of the closed region $|z| \leq r_2$. Hence by maximum modulus theorem, $|f(z_1)| < M(r_2)$; i.e. $M(r_1) < M(r_2)$.

This proves the result. $(x, \varphi(x)) (s, \varphi(s)) D(x, y) B(u, \varphi(u)) A$

119 Corollary: Let $f(z)$ be a non-constant entire function. Then its maximum modulus function $M(r) \rightarrow \infty$ as $|z| = r \rightarrow \infty$. For, if $M(r)$ is bounded, then by Liouville's theorem $f(z)$ would be a constant function. Theorem 6.11 [Hadamard's three-circles theorem]. Let $0 < r_1 < r < r_3$ and suppose that $f(z)$ is analytic on the closed annulus $r_1 \leq |z| \leq r_3$. If $M(r) = \max_{|z|=r} |f(z)|$, then $M(r) \leq M(r_1)^{\frac{r_3 - r}{r_3 - r_1}} M(r_3)^{\frac{r - r_1}{r_3 - r_1}}$ (107) Proof: Let us consider the function $\varphi(z) = z^\alpha f(z)$, where α is a real constant to be chosen later.

If $\alpha \neq$ an integer, $\phi(z)$ is multi-valued in $r_1 \leq |z| \leq r_2$ and so we cut the annulus along the negative part of the real axis. Thus we obtain a simply connected region G in which the principal branch of $\phi(z)$ is analytic. Hence the maximum modulus of this branch of $\phi(z)$ in G is attained on the boundary of G . Since α is real, all the branches of $\phi(z)$ have the same modulus. If we consider another branch of $\phi(z)$ which is analytic in another cut annulus it is clear that the principal branch of $\phi(z)$ can not attain its maximum value on the cut. Hence maximum of $|\phi(z)|$ is attained on at least one of the bounding circles $|z| = r_1$ or $|z| = r_2$. Thus, $\max_{r_1 \leq |z| \leq r_2} |\phi(z)| \leq \max\{M(r_1), M(r_2)\}$. Hence on $|z| = r$, $M(r) \leq \max\{M(r_1), M(r_2)\}$. (108) We now choose α such that $M(r_1) = M(r_2)$. Then $\alpha = -\log(M(r_1)/M(r_2)) / \log(r_1/r_2)$. Substituting this value of α in (108) we get, $M(r) \leq M(r_1) (r/r_1)^\alpha$.

$$M(r) \leq M(r_1) \left(\frac{r}{r_1}\right)^\alpha \leq M(r_1) \left(\frac{r}{r_1}\right)^{\frac{\log(M(r_1)/M(r_2))}{\log(r_1/r_2)}} \quad (109)$$

That is, $M(r) \leq M(r_1) \left(\frac{r}{r_1}\right)^{\frac{\log(M(r_1)/M(r_2))}{\log(r_1/r_2)}}$ [since $a \log b = b \log a$]

Note : Equality in (107) occurs when $\phi(z)$ is a constant, i.e. when $f(z)$ is of the form cz^α for some real α and c is a constant. Corollary : $\log M(r)$ is a convex function of $\log r$. Proof : Let $f(z)$ be analytic in the closed annulus $0 < r_1 \leq |z| \leq r_2$. If $r_1 < r < r_2$ we have, by Hadamard's three-circles theorem,

$$\log M(r) \leq \frac{\log M(r_1) \log(r_2/r) + \log M(r_2) \log(r/r_1)}{\log(r_2/r_1)} \quad (109)$$

Taking logarithms we get $(\log \log M(r)) \log(r) \leq \frac{(\log \log M(r_1)) \log(r_2/r) + (\log \log M(r_2)) \log(r/r_1)}{\log(r_2/r_1)}$. That is, $\log(\log M(r)) \log(r)$ is a convex function of $\log r$.

The inequality (109) shows that $\log M(r)$ is a convex function of $\log r$. 6.7 Order of an entire function An entire function $f(z)$ is said to be of finite order if there is a positive number A such that as $|z| = r \rightarrow \infty$, the inequality $M(r) \leq e^{r^A}$ holds. The lower bound ρ of such numbers A is called the order of the function. f is said to be of infinite order if it is not of finite order. From the definition it is clear that order of an entire function is non-negative. Result : Let f be an entire function of order ρ and $M(r) = \max\{|f(z)| : |z| = r\}$. Then

$$\rho = \limsup_{r \rightarrow \infty} \frac{\log \log M(r)}{\log r} \quad (110)$$

Proof : By hypothesis, given $\epsilon > 0$ there exists $r_0(\epsilon) > 0$ such that $M(r) \leq e^{r^{\rho + \epsilon}}$ for $r > r_0$ while $M(r) \geq e^{r^{\rho - \epsilon}}$ for an increasing sequence $\{r_n\}$ of values of r , tending to infinity. In other words, $\log \log M(r) \leq (\rho + \epsilon) \log r$ and $\log \log M(r) \geq (\rho - \epsilon) \log r$ for a sequence of values of $r \rightarrow \infty$ (111) and (112) precisely means $\rho = \limsup_{r \rightarrow \infty} \frac{\log \log M(r)}{\log r}$. Example 3 : Determine the order of the functions. (i)

- (i) $p(z) = a_0 + a_1 z + \dots + a_n z^n, a_n \neq 0$. (ii) $e^{kz}, k \neq 0$. (iii) $\sin z$ (iv) $\cos z$ Solution : (i)
- $$M(r) \leq a_0 + a_1 r + \dots + a_n r^n$$
- Hence, $M(r) \leq a_n r^n$ for $r \geq 1$. Since ultimately $r \rightarrow \infty$, the choice is justified). $M(r) \leq B r^n$, where $B = a_0 + a_1 r + \dots + a_n r^n$. Hence $\log M(r) \leq \log B + n \log r \leq \log r + n \log r$ (Taking r sufficiently large). $\rho \leq n + 1$. Now, $\rho = \limsup_{r \rightarrow \infty} \frac{\log \log M(r)}{\log r} \leq n + 1$. i.e. $\rho \leq n + 1$. But by definition $\rho \geq 0$. Hence $\rho = 0$ (ii) Here $M(r) = e^{|k|r}$ and hence $\rho = \limsup_{r \rightarrow \infty} \frac{\log \log M(r)}{\log r} = \limsup_{r \rightarrow \infty} \frac{\log |k| + r}{\log r} = \infty$.

125 $n \times x \times dx \text{ Ar } r () \delta gt; + \rho \epsilon 0 2$ for all large r , A being a constant independent of r . Since $n(x)$ is non-negative and non-decreasing function of x , $n \times x \times dx \text{ n } \times x \times dx \text{ r } r () () \leq \delta gt; 0 2 2 \text{ Ar } \rho + \epsilon$ and also $n \times x \times dx \text{ n } \times x \times dx \text{ r } r r r () () \log \geq = 2 2 2$ Hence, $n \text{ r } n \times x \times dx \text{ Ar } r r () \log () , 2 2 \leq \delta gt; + \rho \epsilon$ i.e., $n \text{ r } A \text{ r } () \log \delta gt; + 2 \rho \epsilon$ for all large r . Hence, $n(r) = O(r \rho + \epsilon)$.

6.9 Convergence exponent (or, exponent of Convergence) Let $f(z)$ be an entire function with zeros at the points a_1, a_2, \dots , arranged in order of non-decreasing modulus, multiple zeros being repeated according to their multiplicities and $|a_i| = r_i, i = 1, 2, \dots$. We define convergence exponent ρ_1 of the zeros of $f(z)$ by the equation $\rho_1 = \limsup \log \log n \text{ n } r$ (117) or, equivalently by $\rho_1 = \limsup \log () \log n \text{ n } r$ (118) The convergence exponent has the following property.

Theorem 6.14 : Let $f(z)$ be an entire function with zeros at $a_1 a_2, \dots$, arranged in order of non-decreasing modulus, multiple zeros being repeated according to their multiplicities and $|a_i| = r_i$. If the convergence exponent ρ_1 of the zeros of $f(z)$ is finite, then the series $\sum 1/r_n^\alpha$ converges when $\alpha > \rho_1$ and diverges when $\alpha < \rho_1$. If ρ_1 is infinite, the above series diverges for all positive values of α . Proof : Let ρ_1 be finite and $\alpha < \rho_1$. Then, $\rho \alpha 1 1 1 2 \delta gt; + ()$. Hence, $\log \log () n \text{ r } n \delta gt; + 1 2 1 \rho \alpha$ for all large n .

126 or, $\log \log () n \text{ r } n \delta gt; + 1 2 1 \rho \alpha$ i.e. $n \text{ r } n \delta gt; + 1 2 1 1 2 () ; \dots, \rho \alpha \rho \alpha r n n n$ where $p n p \alpha \rho \alpha \rho \alpha \rho \alpha \rho \alpha \rho \alpha \delta gt; = = = - + \delta gt; + + - + + 2 0 1 1 1 1 1 1 1 , .$ Hence, $1 1 1 r n n p \alpha \delta gt; +$ for all large n . Hence, $1 1 r n n \alpha = \bullet \sum$ converges. Next, let $\alpha > \rho_1$. Then, $\log \log n \text{ r } n \delta gt; \alpha$ for a sequence of values of n , tending to infinity. That is, $\log \log n \text{ r } n \delta gt; \alpha$ or, $1 1 r n n \alpha \delta gt; (119)$ for a sequence of values of n tending to infinity. Let N be such a value of n for which (119) holds and m be the least integer $< N 2$. Then, as r_n is non-decreasing, $1 1 1 1 1 1 1 r r r r r n N m N m N N N n N m N \alpha \alpha \alpha \alpha \alpha = + + + \geq + + - - + = - \sum = + \delta gt; \delta gt; \delta gt; m r m r m N N N 1 1 2 \alpha \alpha$. Since N may be as large as we please, by Cauchy's principle of convergence, the series $\sum 1/r_n^\alpha$ diverges. If ρ_1 is infinite, then for any positive value of α , $\log \log n \text{ r } n \delta gt; \alpha$ for a sequence of values

127 of n tending to infinity; i.e., $n \text{ r } n \delta gt; \alpha$ for a sequence of values of n tending to infinity. Hence as before, the series $\sum 1/r_n^\alpha$ diverges for any positive α . Note 1. Observe that ρ_1 may also be defined as the lower bound of the positive numbers α for which the series $\sum 1/r_n^\alpha$ is convergent. If $f(z)$ has no zeros we define $\rho_1 = 0$ and if $\sum 1/r_n^\alpha$ diverges for all positive α , then $\rho_1 = \bullet$. Note 2. If ρ_1 is finite, the series $\sum 1/r_n^{\rho_1}$ may be convergent or divergent. For example, if $r_n = n$, then $\rho_1 = \limsup \log \log n \text{ n } r$ and $\sum 1/n^{\rho_1} = \bullet \sum \sum$ diverges. Again, if $r_n = n(\log n)^2$, then, $\rho_1 2 1 = + = \limsup \log \log \log \log , n \text{ n } n n$ and $\sum 1/n^{\rho_1} = \bullet \sum \sum (\log)$ converges.

Theorem 6.15 : If $f(z)$ is an entire function with finite order ρ and r_1, r_2, \dots , are the moduli of the zeros of $f(z)$, then $\sum 1/r_n^\alpha$ converges if $\alpha > \rho$. Proof : We choose β such that $\rho < \beta < \alpha$. Since for any $\epsilon < 0, n(r) = O(r^{\rho + \epsilon}), n(r) > Kr^\beta$ (120) for all large r, K being a constant. Putting $r = r_n, n$ large, (120) gives $n > Kr_n^\beta > \beta$, i.e., $r_n < n^{1/\beta}$ or, $1/r_n < B/n^{\alpha/\beta} > /$ for all large n, B being a constant. Since $\alpha/\beta > \rho$, $\sum 1/r_n^\alpha$ converges.

128 Corollary : Since convergence exponent ρ_1 is the lower bound of positive numbers α for which $\sum 1/r_n^\alpha$ is convergent, it follows that $\rho_1 \leq \rho$. Note : ρ_1 may be 0 or \bullet . For example if $r_n = e^n, \rho_1 = 0$ and if $r_n = \log n$, then $\rho_1 = \bullet$. For the function $f(z) = e^z, \rho = 1$ and $\rho_1 = 0$ so that $\rho_1 > \rho$. But for $\sin z$ or $\cos z, \rho = \rho_1 = 1$. Result : If the convergence exponent ρ_1 of the zeros of an entire function $f(z)$ is greater than 0, then $f(z)$ has infinite number of zeros. Proof : If possible, suppose $f(z)$ has finite number of zeros with moduli r_1, \dots, r_N . The series $\sum 1/r_n^\alpha$, being of finite number of terms, converges for every $\alpha > 0$. Hence $\rho_1 = 0$, a contradiction. Hence $f(z)$ has infinite number of zeros. Note : For an entire function with finite number of zeros, $\rho_1 = 0$. Example : Find the convergence exponent of the zeros of $\cos z$. Solution : First method : The zeros of $\cos z$ are $\pi/2, 3\pi/2, 5\pi/2, \dots$. Now, $1/2^2, 1/3^2, 1/5^2, \dots$. The series converges when $\alpha > 1$ and diverges when $\alpha < 1$. Hence the lower bound of the positive numbers α for which $\sum 1/r_n^\alpha$ converges is 1 i.e., $\rho_1 = 1$. Second method : The zeros of $\cos z$ are $(2n+1)\pi/2, n \in \mathbb{Z}$. Let $a_n = (2n+1)\pi/2$. Then, $1/a_n^\alpha = 1/((2n+1)\pi/2)^\alpha = 2^\alpha / ((2n+1)^\alpha \pi^\alpha) = 2^\alpha / (2^\alpha n^\alpha \pi^\alpha) = 1/n^\alpha \pi^\alpha$. Hence, $\sum 1/a_n^\alpha = \sum 1/n^\alpha \pi^\alpha$. Hence, $\rho_1 = \limsup \log \log n \text{ n } r = - + = -$

$f(z) = p(z)/q(z)$ where p and q are polynomials with no common zeros. If the degree of p is less than or equal to the degree of q , then f has only a finite number of poles and the point at infinity is not a pole. On the other hand, if the degree of p is greater than the degree of q , then (taking degree of $p(z) = m$ and degree of $q(z) = n$).

$$f(z) = \frac{z^m + a_{m-1}z^{m-1} + \dots + a_1z + a_0}{z^n + b_{n-1}z^{n-1} + \dots + b_1z + b_0}$$

where degree of $r(z) \leq n - 1$. This shows that the point at infinity is a pole of order $(m - n)$ and there lie a finite number of poles in the unextended plane. These establish that $f(z)$ is meromorphic. Theorem 6.20 : [Partial fraction decomposition]. Let $p(z)$, $q(z)$ be two polynomials with no common zeros and that $0 \leq \deg(p) < \deg(q)$. Let a_1, \dots, a_k be the zeros of $q(z)$ with multiplicities $\alpha_1, \dots, \alpha_k$. Then $p(z)/q(z)$ can be expressed uniquely as
$$\frac{p(z)}{q(z)} = \sum_{i=1}^k \frac{c_i}{(z - a_i)^{\alpha_i}} + g(z) \quad (131)$$
 Proof. The decomposition is unique. We assume that the relation (131) exists. Let $r > 0$ be small enough. Then for $z \in N(a_i, r)$, (131) can be rewritten as
$$\frac{p(z)}{q(z)} - \frac{c_i}{(z - a_i)^{\alpha_i}} = \sum_{j=1}^{i-1} \frac{c_j}{(z - a_j)^{\alpha_j}} + g(z) \quad (132)$$
 since $N(a_i, r)$ does not contain any zero of $q(z)$ other than a_i , $g(z)$ is analytic at $z = a_i$. Multiplying both sides of (132) by $(z - a_i)^{\alpha_i}$, we obtain
$$\frac{p(z)}{q(z)} (z - a_i)^{\alpha_i} - c_i = \sum_{j=1}^{i-1} c_j (z - a_i)^{\alpha_i - \alpha_j} + g(z) (z - a_i)^{\alpha_i} \quad (133)$$
 Now the function $\frac{p(z)}{q(z)} (z - a_i)^{\alpha_i} - c_i$ is analytic for all z belonging to $N(a_i, r)$ and hence can be expanded in a Taylor series in a neighbourhood of a_i in $N(a_i, r)$
$$\frac{p(z)}{q(z)} (z - a_i)^{\alpha_i} - c_i = \sum_{n=0}^{\infty} \alpha_n (z - a_i)^n \quad (134)$$
 Combining (133) and (134), we write
$$\frac{c_j}{(z - a_j)^{\alpha_j}} - \frac{c_i}{(z - a_i)^{\alpha_i}} = \sum_{n=0}^{\infty} \alpha_n (z - a_i)^{n - \alpha_j} \quad (135)$$
 Comparing the coefficients we find $c_j - c_i = \sum_{n=0}^{\infty} \alpha_n (a_i - a_j)^{n - \alpha_j}$, ..., uniquely Existence of the decomposition. The principal part associated to each pole a_i is $\frac{c_i}{(z - a_i)^{\alpha_i}} = \sum_{j=1}^{\alpha_i} \frac{c_{ij}}{(z - a_i)^j}$ Now if we subtract all the principal parts we find the function $f(z) - \sum_{i=1}^k \frac{c_i}{(z - a_i)^{\alpha_i}} = g(z)$ is analytic in the extended plane. Now each of the terms $\frac{c_i}{(z - a_i)^{\alpha_i}}$ converges to zero for $z \rightarrow \infty$, and also $p(z)/q(z)$ converges to zero for $z \rightarrow \infty$ since $\deg(q) > \deg(p)$. This shows that $f(z) \rightarrow 0$ for $z \rightarrow \infty$. But then f is necessarily bounded and hence constant by Liouville's theorem. A constant function tending to zero as $z \rightarrow \infty$ must be identically zero. Example 4 : Consider the rational function

$$f(z) = \frac{z^5 + 2z^3 + 3z^2 + 4z + 1}{z^6 + z^5 + z^4 + z^3 + z^2 + z + 1}$$

Partial Fraction Expansion of Meromorphic Functions Let $f(z)$ be a meromorphic function and z_0 be a pole of order m with the principal part $\frac{c_{-m}}{(z - z_0)^m} + \frac{c_{-m+1}}{(z - z_0)^{m-1}} + \dots + \frac{c_{-1}}{(z - z_0)}$ Then $f(z)$ can be written as [see § 6.2, (14)] $f(z) = p(z) + g(z)$ where $g(z)$ is an entire function. Now

if, in general, z_1, z_2, \dots, z_n are the poles of a meromorphic function f with the corresponding principal parts P_1, P_2, \dots, P_n then f can be expressed as $f(z) = \sum_{j=1}^n P_j(z) + \psi(z)$ (136) where $\psi(z)$ is an entire function. But the question arises whether it is possible to construct a meromorphic function possessing poles at the sequence of points $\{z_n\}$ with corresponding principal parts P_1, P_2, \dots . Because in this case the series $\sum P_j(z)$ in (136) turns out to be an infinite series $\sum_{j=1}^{\infty} P_j(z)$, which needs to be convergent.

137 Gösta Mittag Leffler (1846-1927), German in origin but his several generations lived in Sweden, overcame this difficulty by introducing a polynomial $p_n(z)$ dependent on z_n and $P_n(z)$ so that the series $\sum_{n=1}^{\infty} \frac{P_n(z)}{p_n(z)}$ is uniformly convergent in any compact set K not containing any points of the sequence $\{z_n\}$. Theorem 6.21 [The Mittag Leffler Theorem]: Given a sequence of distinct complex numbers $\{z_n\}$, $z_n \neq 0, 1, 2, \dots$, $\lim_{n \rightarrow \infty} |z_n| = \infty$ and a sequence of rational functions $\{P_n(z)\}$, $P_n(z) = \frac{c_n}{z - z_n + k_n}$, $c_n \neq 0, k_n \in \mathbb{Z}, \dots$ (137) there exists a meromorphic function $f(z)$ having poles at the points z_n and only there with $P_n(z)$ as its principal part at z_n and can be represented in the form of an expansion

$f(z) = \sum_{n=1}^{\infty} \frac{P_n(z)}{p_n(z)} + h(z)$ where $h(z)$ is an arbitrary entire function

and $p_n(z)$ is suitable partial sum of Taylor's expansion of the singular part which is analytic in the open disc $|z| < |z_n|$. Proof. Without loss of generality we assume that $z = 0$ is not a pole of $f(z)$. Now $P_k(z)$ is analytic for $|z| < |z_k|$ and can be expanded in this neighbourhood of z : $P_k(z) = \sum_{j=0}^{\infty} a_{kj} z^j$ and hence this series converges uniformly in the disk $|z| \leq 2$. Let $p_k(z) = \sum_{j=0}^k a_{kj} z^j$ be a partial sum of this expansion such that $|P_k(z) - p_k(z)| < \frac{1}{2}$ for $|z| \leq 2$. Let R be an arbitrary large positive number and since $|z_n| \rightarrow \infty$ as $n \rightarrow \infty$ we can find an $N(R)$ so large that $|z_n| > 2R$ when $n \geq N(R)$. Therefore in the circle

$|z| \leq R$

z

$|z| \leq R < |z_n| < 2R$

z

$|z| \leq R < |z_n| < 2R$

z

p

z

n

$n < |z_n| < 2R$

n

N

$R < |z_n| < 2R$

138

the

first sum in the r.h.s is finite and the second sum $\sum_{n=N}^{\infty} \frac{P_n(z)}{p_n(z)}$ is absolutely and uniformly convergent by comparison with the convergent series $\sum_{n=N}^{\infty} \frac{1}{|z_n|^{2k}}$. Therefore $\sum_{n=1}^{\infty} \frac{P_n(z)}{p_n(z)}$ is analytic in $|z| < R$ except at the poles belonging to the sequence $\{z_n\}$. It is thus a meromorphic function with the poles at z_1, z_2, \dots and with the principal parts $P_1(z), P_2(z), \dots$ at each point z_n respectively. Now if $f(z)$ possesses the same poles only with the same principal parts then $f(z) - \sum_{n=1}^{\infty} \frac{P_n(z)}{p_n(z)}$ is an entire function $h(z)$, say. This completes the proof. Example 5: Prove that $\pi \cot \pi z = \frac{1}{z} + \sum_{n=1}^{\infty} \left(\frac{1}{z-n} + \frac{1}{z+n} \right)$

'Solution: The given function $\pi \cot \pi z$ has simple poles at $z = 0, \pm 1, \pm 2, \dots$ with residue 1.

Here, $\frac{1}{z-n} + \frac{1}{z+n} = \frac{2z}{z^2 - n^2}$

z

z

$\frac{1}{z-n} + \frac{1}{z+n} = \frac{2z}{z^2 - n^2}$, (138) Let $|z| < R$

$R < |z_n| < 2R$

be so large that $R < |z_n| < 2R$ when $n \geq N(R)$. Then from (138), we find $\frac{1}{z-n} + \frac{1}{z+n} = \frac{2z}{z^2 - n^2}$

$N < |z_n| < 2N$

Now, since $\sum \frac{1}{N^2}$ is convergent, the series $\sum_{n=N}^{\infty} \frac{1}{z^2 - n^2}$

$= - \sum_{n=N}^{\infty} \frac{1}{n^2}$

converges uniformly on any compact set (lying in $|z| > R$) not containing any of the points $z = \pm 1, \pm 2, \dots$. Therefore applying the

Mittag-Leffler theorem we can express $\pi \cot(\pi z)$

$$\pi \cot(\pi z) = \frac{1}{z} + \sum_{n=1}^{\infty} \left(\frac{1}{z-n} + \frac{1}{z+n} \right) \quad (139)$$

$$+ \dots + \frac{1}{z-n} + \frac{1}{z+n} \quad (139)$$

where $h(z)$ is an entire function.

Differentiating term-wise, we obtain $\pi^2 \operatorname{cosec}^2(\pi z)$

$$= \frac{1}{z^2} + \sum_{n=1}^{\infty} \left(\frac{1}{(z-n)^2} + \frac{1}{(z+n)^2} \right) + h'(z)$$

$$= \frac{1}{z^2} + \sum_{n=1}^{\infty} \left(\frac{1}{(z-n)^2} + \frac{1}{(z+n)^2} \right) + h'(z)$$

$$= \frac{1}{z^2} + \sum_{n=1}^{\infty} \left(\frac{1}{(z-n)^2} + \frac{1}{(z+n)^2} \right) + h'(z)$$

$$= \frac{1}{z^2} + \sum_{n=1}^{\infty} \left(\frac{1}{(z-n)^2} + \frac{1}{(z+n)^2} \right) + h'(z)$$

We notice that the functions $f(z)$ and $\psi(z)$ are both periodic with period 1 and consequently $h'(z)$ is also periodic with the same period. Let $z = x + iy$. Consider the strip $0 \leq x \leq 1$. In fact, the convergence of the series in (140) is uniform for $y \geq 1$, say and the limit tends to 0 as $y \rightarrow \infty$ (this can be seen on taking the limit in each term of the series). Again,

$$\sin(x + iy) = \sin x \cosh y + i \cos x \sinh y$$

$$\cos(x + iy) = \cos x \cosh y - i \sin x \sinh y$$

$$\cot(x + iy) = \frac{\sin(x + iy)}{\cos(x + iy)} = \frac{\sin x \cosh y + i \cos x \sinh y}{\cos x \cosh y - i \sin x \sinh y} = \frac{\sin x \cosh y + i \cos x \sinh y}{\cos x \cosh y - i \sin x \sinh y}$$

and so $\cot(x + iy) \rightarrow 0$ as $y \rightarrow \infty$. From these we conclude that $h'(z)$ is bounded in the period strip $0 \leq x \leq 1$ and due to its periodicity it is bounded in the entire plane. By Liouville's theorem it then reduces to a constant.

Now since $\lim_{y \rightarrow \infty} \cot(x + iy) = 0$, $h'(z)$ is indeed zero and $h(z) = c$, a constant.

Then from (139), $\pi \cot(\pi z) = \frac{1}{z} + \sum_{n=1}^{\infty} \left(\frac{1}{z-n} + \frac{1}{z+n} \right) + c$

$$\pi \cot(\pi z) = \frac{1}{z} + \sum_{n=1}^{\infty} \left(\frac{1}{z-n} + \frac{1}{z+n} \right) + c$$

$$= \frac{1}{z} + \sum_{n=1}^{\infty} \left(\frac{1}{z-n} + \frac{1}{z+n} \right) + c$$

For $z = 1/2$, $\pi \cot(\pi/2) = 0 = \frac{1}{1/2} + \sum_{n=1}^{\infty} \left(\frac{1}{1/2-n} + \frac{1}{1/2+n} \right) + c$

$$= \frac{1}{1/2} + \sum_{n=1}^{\infty} \left(\frac{1}{1/2-n} + \frac{1}{1/2+n} \right) + c$$

$$= \frac{1}{1/2} + \sum_{n=1}^{\infty} \left(\frac{1}{1/2-n} + \frac{1}{1/2+n} \right) + c$$

$$= \frac{1}{1/2} + \sum_{n=1}^{\infty} \left(\frac{1}{1/2-n} + \frac{1}{1/2+n} \right) + c$$

Remark : Here it is proved incidentally that $\pi \cot(\pi z) = \frac{1}{z} + \sum_{n=1}^{\infty} \left(\frac{1}{z-n} + \frac{1}{z+n} \right) + c$ [see equation (140)] We can now utilize the identity (141) to calculate easily some familiar sums. Here the l.h.s of (141) has the Laurent series expansion in the neighbourhood of

$z = 0$. $\pi \cot(\pi z) = \frac{1}{z} - \frac{\pi^2 z}{3} + \frac{7\pi^4 z^3}{15} - \dots$ Note that the series on the r.h.s of (141) converges uniformly near $z = 0$. By Th. 4.14 [14] it converges uniformly together with all derivatives. Again

$$\pi \cot(\pi z) = \frac{1}{z} - \frac{\pi^2 z}{3} + \frac{7\pi^4 z^3}{15} - \dots$$

$$= \frac{1}{z} - \frac{\pi^2 z}{3} + \frac{7\pi^4 z^3}{15} - \dots$$

$$= \frac{1}{z} - \frac{\pi^2 z}{3} + \frac{7\pi^4 z^3}{15} - \dots$$

$$= \frac{1}{z} - \frac{\pi^2 z}{3} + \frac{7\pi^4 z^3}{15} - \dots$$

$$= \frac{1}{z} - \frac{\pi^2 z}{3} + \frac{7\pi^4 z^3}{15} - \dots$$

$$= \frac{1}{z} - \frac{\pi^2 z}{3} + \frac{7\pi^4 z^3}{15} - \dots$$

$$= \frac{1}{z} - \frac{\pi^2 z}{3} + \frac{7\pi^4 z^3}{15} - \dots$$

Solution : Here the given function $\pi \tan(\pi z)$ possesses simple poles at $z = \pm \frac{1}{2}, \pm \frac{3}{2}, \dots$, with residue -1 . Then, $\pi \tan(\pi z) = \sum_{n=0}^{\infty} \left(\frac{1}{z - (n + \frac{1}{2})} - \frac{1}{z + (n + \frac{1}{2})} \right) + h(z)$

$$\pi \tan(\pi z) = \sum_{n=0}^{\infty} \left(\frac{1}{z - (n + \frac{1}{2})} - \frac{1}{z + (n + \frac{1}{2})} \right) + h(z)$$

$$= \sum_{n=0}^{\infty} \left(\frac{1}{z - (n + \frac{1}{2})} - \frac{1}{z + (n + \frac{1}{2})} \right) + h(z)$$

$$= \sum_{n=0}^{\infty} \left(\frac{1}{z - (n + \frac{1}{2})} - \frac{1}{z + (n + \frac{1}{2})} \right) + h(z)$$

$$= \sum_{n=0}^{\infty} \left(\frac{1}{z - (n + \frac{1}{2})} - \frac{1}{z + (n + \frac{1}{2})} \right) + h(z)$$

converges uniformly on any compact set not containing any of the poles of the given function. By Mittag-Leffler theorem, $\pi \tan(\pi z) = \sum_{n=0}^{\infty} \left(\frac{1}{z - (n + \frac{1}{2})} - \frac{1}{z + (n + \frac{1}{2})} \right) + h(z)$ where $h(z)$ is an arbitrary entire function. Now proceeding as in example 5, we can have the desired result. Example 7 : Establish that $\sum_{n=0}^{\infty} \frac{1}{(2n+1)^2} = \frac{\pi^2}{8}$

$$\sum_{n=0}^{\infty} \frac{1}{(2n+1)^2} = \frac{\pi^2}{8}$$

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$$\sum_{n=0}^{\infty} \frac{1}{(2n+1)^2} = \frac{\pi^2}{8}$$

$\pm +$

$\geq + + 1 2 1 2 1 2 1 2 \pi \pi \pi \pi = - + \geq - + - + - + 1 1 1 1 2 1 2 1 e e e e n () () \pi \pi \pi \pi$ (148) The given function can be rewritten as $z z z z z z z \sin \sin \cos \cot - = - 1 1 B n A n P n Q n x y o b n - b n - b n$

145 I. Bound on the sides $A n Q n \& B n P n$ of the square $C n$: Using (147), we obtain $1 1 1 1 1 1 2 2 z z z z e e e e b y y y y y n - \leq - = - + - + \rightarrow \rightarrow \bullet - - \cot \cot .$ as n II. Bound on the sides $A n B n \& Q n P n$ of $C n$: Here we apply (148) to achieve $1 1 1 1 1 1 1 1 1 2 2 z z z z e e b y e e n - \leq - \leq - + - + \rightarrow + - \rightarrow \bullet \cot \cot . \pi \pi \pi \pi$ as n Thus, $z z z z z e e z C n$

$n \sin \sin \cos , , , \dots - \leq + - \epsilon = \pi \pi 1 1 1 2$ This shows that the function $F(z)$ is bounded on the sequence of contours $\{C n\}$ and we can apply (144) to prove

$z z z z z z z n n n n$

n

$\sin \sin \cos - = + - + + + - = \bullet \sum 3 2 1 1 1 1 1 \alpha \alpha \alpha \alpha = + - = \bullet \sum 3 2 2 2 2 1 z z n n$

α Exercises 1. Obtain partial fraction expansion of $\operatorname{cosec} z$. 2. Prove that $\sec () () z n z n n n = - - - - = \bullet \sum 1 2 1 1 2 2 2 2 1 \pi \pi 3$. Show that $\tan z z z n n = - - - = \bullet \sum 2 1 2 2 2 2 1 \pi$

146 and hence deduce $1 1 3 1 5 8 2 2 2 + + + = \pi 6.16$ The Gamma Function The gamma function $\Gamma(z)$ was introduced by Swedish Mathematician L. Euler (1707- 1783), in 1729 while he was seeking for a function of a real variable x which is continuous for positive x and reduces to $x!$ when x is a positive integer. Gamma function is widely used in the fields of probability and statistics, as well as combinatorics. Gamma function $\Gamma(z)$ can be introduced in either of the ways : (i) in terms of infinite product (ii) in the form of infinite integral (iii) in limit formula We establish the form (i) first considering the fact that it possesses simple poles at $z = 0, -1, -2, \dots$ and nowhere vanishes in the entire plane and satisfies $z\Gamma(z) = \Gamma(z + 1), \Gamma(1) = 1$ (149) To construct $\Gamma(z)$ we claim that $f(z) = 1/\Gamma(z)$ is entire with simple zeros at $z = -n (n = 0, 1, 2, \dots)$. Again we know that $k = 1$ is the largest non-negative integer for which $1 1 n k n = \bullet \sum$ diverges. Then utilizing the Weierstrass Factorization theorem $f(z)$ can be represented as $f z z e$

$z n e g z n z n () () = + = \bullet - \prod 1 1$ where $g(z)$ is an entire function,

so that gamma function will be of the form $\Gamma () () / z e z z n e g z z n = + - - \bullet \prod 1 1 1$ (150) Now we find $g(z)$ so that (149) hold. We write (150) in the form

147 $\Gamma () \lim ()$

$z e z z m e n g z z m n = + \rightarrow \bullet - - \prod 1 1 = - + + + = \rightarrow \bullet \rightarrow \bullet \sum \lim ! \exp () () () \lim () , n$

$n n n n$

g

z

z m

z

z

$z n z 1 1 \Gamma$ say (151) $z z z n z g z z m z z z n z z z$

n

$n g z z$

m

n

$n n n$

$\Gamma \Gamma () () ! \exp () () () () ! \exp () + = - + + + + + - + + + \sum \sum 1 1 1 2 1 1 1 1 1 = + + + - - \sum () \exp () () z$

n

$g z g z m n 1 1 1 1 = + + + - - \sum 1 1 1 1 1 z$

$n n g z$

g z

$m n \exp () () = + + + - - + \sum 1 1 1 1 1 z n g z g z$

$m n n$

$\exp () () \log$ Now from the relation

$z z z z z z n n$

$n \Gamma \Gamma \Gamma \Gamma () () \lim () () , + = + \rightarrow \bullet 1 1$ we find that $z z z z n g z g z m n n$

n

$$\Gamma(z) = \lim_{n \rightarrow \infty} \frac{n! n^z}{z(z+1)\dots(z+n)}$$

by Gauss's formula. In the expression (155) for $\Gamma(z)$ the infinite product is uniformly convergent on every compact subset of $\mathbb{C} - \{0, -1, -2, \dots\}$. So calculating $\Gamma'(z)/\Gamma(z)$ we find that $\Gamma'(z)/\Gamma(z) = -\gamma - \sum_{n=1}^{\infty} \left(\frac{1}{z+n} - \frac{1}{n}\right)$. This function $\Gamma'(z)/\Gamma(z)$ is denoted by $\psi(z)$ and named as Gaussian psi function and it is seen from its expression that ψ is meromorphic in \mathbb{C} with simple poles at $z = 0, -1, -2, \dots$ and $\text{Res}(\psi; -n) = -1$ for $n = 0, 1, 2, \dots$ Example 10 : Show that (i) $\psi(1) = -\gamma$ (ii) $\psi'(z) = -\sum_{n=0}^{\infty} \frac{1}{(z+n)^2}$ (iii) $\psi(z) = -\gamma - \sum_{n=1}^{\infty} \left(\frac{1}{z+n} - \frac{1}{n}\right)$

150 so, $\psi(1) = -\gamma$ (ii) $\psi'(z) = -\sum_{n=0}^{\infty} \frac{1}{(z+n)^2}$ (iii) $\psi(z) = -\gamma - \sum_{n=1}^{\infty} \left(\frac{1}{z+n} - \frac{1}{n}\right)$

$$\psi(z) = -\gamma - \sum_{n=1}^{\infty} \left(\frac{1}{z+n} - \frac{1}{n}\right)$$

$$\psi'(z) = -\sum_{n=0}^{\infty} \frac{1}{(z+n)^2}$$

$$\psi'(1) = -\sum_{n=0}^{\infty} \frac{1}{(1+n)^2} = -\sum_{n=1}^{\infty} \frac{1}{n^2} = -\frac{\pi^2}{6}$$

$$\psi(z) = 1/z + \sum_{n=1}^{\infty} \left(\frac{1}{z+n} - \frac{1}{n}\right) = -\gamma - \sum_{n=1}^{\infty} \left(\frac{1}{z+n} - \frac{1}{n}\right)$$

$$\psi'(z) = -\sum_{n=0}^{\infty} \frac{1}{(z+n)^2}$$

$$\psi'(1) = -\sum_{n=0}^{\infty} \frac{1}{(1+n)^2} = -\sum_{n=1}^{\infty} \frac{1}{n^2} = -\frac{\pi^2}{6}$$

by $\pi \cot \pi z = \sum_{n=-\infty}^{\infty} \frac{1}{z+n}$

6.17 A Few Properties of $\Gamma(z)$

We have $\Gamma(z)\Gamma(1-z) = \frac{\pi}{\sin \pi z}$

$$\Gamma(z)\Gamma(1-z) = \frac{\pi}{\sin \pi z}$$

z z

n

$$\pi \pi 1 2 2 1 = - z z$$

π

$$\pi \sin or, 1 \Gamma \Gamma () [()] \sin z z z z - - =$$

$$\pi \pi \text{ i.e. } 1 \Gamma \Gamma () () \sin , z z$$

$$z - =$$

π

$$\pi \text{ [using } z\Gamma(z) = \Gamma(z + 1) \text{ i.e., } -z\Gamma(-z) = \Gamma(1 - z)] \text{ (157)}$$

In particular, $\Gamma 1 2 2 = \pi$ and $\Gamma 1 2 = \pi$ (minus sign is excluded since $\Gamma 1 2$ is positive by (155)). Likewise using $\Gamma(z + 1) = z\Gamma(z)$ we find $\Gamma 3 2 1 2 1 2 1 2 = = \pi \Gamma 5 2 3 2 3 2 3 2 1 2 = = \cdot \pi \Gamma 7 2 5 2 5 2 5 2 3 2 1 2 = = \cdot \cdot \pi$ and in general $\Gamma n n n n + = - = 1 2 1 3 2 1 2 1 2 \cdot () , (, , \pi i$.

e. $\Gamma n n n n + = 1 2 2 2 2 / () ! () \pi$ (158) If n is a positive integer repeated use of (149) produce $\Gamma () ! n n + = 1$ The Γ -function can therefore be considered as an extension of the factorial function to the complex plane.

152 Legendre's Duplication Formula Let us consider the Gauss's formula $\Gamma \Gamma () \lim ! () () \lim (,)$,

z

n n

z z

z n

$$z n n z n = + + = \rightarrow \bullet \rightarrow \bullet 1 \text{ say Then, } \Gamma (,) () ! () () () 2 2 2 2 2 2 1 2 2 2 2 z$$

n n n

z

$$z z n n z z = + + + = + + + - 2 1 2 2 2 2 1 2 2 2 2 2 1 2 n z n n n$$

z z

z z

$$n ! () () () \Gamma \pi \text{ [Replacing } (2n)! \text{ by (158)] } = + + + + + - - 2 1 2 1 2 1 2 3 2 1 2 2 1 2$$

z

z

n

n n

z z

z z

n z

z z

$$n ! () () () \Gamma \pi = + + + + - - 2 1 2 1 1 2 3 2 1 2 2 1 z z$$

n n z z

z

$$n \pi \Gamma \Gamma (,) = + + + + - 2 1 2 1 2 1 2 2 1 1 2 z z$$

n n

z

n n n z

n

n

π

$$\Gamma \Gamma \Gamma \Gamma (,) , () / \text{ and } \Gamma \Gamma \Gamma \Gamma$$

$$\Gamma () \lim (,) ()$$

$$\lim () / 2 2 2 2 1 2 1 2 1 2 2 1 1 2$$

z

z

$$n z z n n n z n n n z$$

$$n = + + + + \rightarrow \bullet - \rightarrow \bullet$$

Integrating by parts ($n - 1$ times) = + + = $n n$

z

z

z

n

$F z$

$n! \dots () () 2 1$

Now to prove (ii) we show that $\lim_{n \rightarrow \infty} \int_C z^n e^{-z} dz = 0$

t

$\int_C z^n e^{-z} dz \rightarrow 0$ because the integral on the right converges. This completes the proof of (ii). Finally

combining the results (i) and (ii) with the Gauss's formula (156) we get $\Gamma(z) = \int_0^\infty t^{z-1} e^{-t} dt$

for $\text{Re}(z) > 0$. Consequently, $\Gamma(z) = \int_0^\infty t^{z-1} e^{-t} dt$

for $\text{Re}(z) > 0$. Therefore, $\Gamma(z) = \int_0^\infty t^{z-1} e^{-t} dt$

for $\text{Re}(z) > 0$.

$n e^{-t} dt$

n

$z^n e^{-z} dz \rightarrow 0$ because the integral on the right converges.

This completes the proof of (ii). Finally combining the results (i) and (ii) with the Gauss's formula (156) we get $\Gamma(z) = \int_0^\infty t^{z-1} e^{-t} dt$ for $\text{Re}(z) > 0$.

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Instead of being guided by any presumption about ability level, it would perhaps stand to reason if receptivity of a learner is judged in the course of the learning process. That would be entirely in keeping with the objectives of open education which does not believe in artificial differentiation.

Keeping this in view, study materials of the Post-Graduate level in different subjects are being prepared on the basis of a well laid-out syllabus. The course structure combines the best elements in the approved syllabi of Central and State Universities in respective subjects. It has been so designed as to be upgradable with the addition of new information as well as results of fresh thinking and analysis.

The accepted methodology of distance education has been followed in the preparation of these study materials. Co-operation in every form of experienced scholars is indispensable for a work of this kind.

We, therefore, owe an enormous debt of gratitude to everyone whose tireless efforts went into the writing, editing and devising of proper lay-out of the materials. Practically speaking, their role amounts to an involvement in 'invisible teaching'.

For, whoever makes use of these study materials would virtually derive the benefit of learning under their collective care without each being seen by the other.

The more a learner would seriously pursue these study materials, the easier it will be for him or her to reach out to larger horizons of a subject. Care has also been taken to make the language lucid and presentation attractive so that they may be rated as quality self-learning materials. If anything remains still obscure or difficult to follow, arrangements are there to come to terms with them through the counselling sessions regularly available at the network of study centres set up by the University. Needless to add, a great deal of these efforts is still experimental—in fact, pioneering in certain areas. Naturally, there is every possibility of some lapse or deficiency here and there. However, these do admit of rectification and further improvement in due course. On the whole, therefore, these study materials are expected to evoke wider appreciation the more they receive serious attention of all concerned. Chandan Basu

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NETAJI S U B H A S O P E N U N I V E R S I T Y PG (MT)–XA (1) Advanced Differential Geometry NETAJI SUBHAS OPEN UNIVERSITY Unit 1 : DIFFERENTIABLE MANIFOLD 7-64 Unit 2 : LIE GROUP 65-91 Unit 3 : LINEAR CONNECTION 92-102 Unit 4 : RIEMANNIAN MANIFOLD 103-125

7 UNIT - 1 1.1 Calculus on \mathbb{R}^n : Let \mathbb{R} denote the set of real numbers. For an integer $n \in \mathbb{N}$, let \mathbb{R}^n be the cartesian product $\mathbb{R} \times \mathbb{R} \times \dots \times \mathbb{R}$ n times of the set of all ordered n -tuples (x_1, x_2, \dots, x_n) of real numbers. Individual n -tuple will be denoted at times by a single letter, e.g. x, y, z, \dots and so on. Co-ordinate functions : Let $f : \mathbb{R}^n \rightarrow \mathbb{R}$. Then, the functions $f_i : \mathbb{R}^n \rightarrow \mathbb{R}$ defined by $f_i(x_1, \dots, x_n) = x_i$. We are now going to define a function to be differentiable of class C^k . A real-valued function $f : U \subset \mathbb{R}^n \rightarrow \mathbb{R}$, U being an open set of \mathbb{R}^n , is said to be of class C^k if i) all its partial derivatives of order less than or equal to k exist and ii) are continuous functions at every point of U . By class C^0 , we mean that f is merely continuous from U to \mathbb{R} . By class C^k , we mean that that partial derivatives of all orders of f exist and are continuous at every point of U . In this case, f is said to be a smooth function. Note : By class C^k on U , we mean that f is real analytic on U i.e. expandable in a power series about each point on U . A C^k function is a C^0 function but the converse is not true. Exercise : 1. Let $f : \mathbb{R} \rightarrow \mathbb{R}$ be defined by $f(x) = x^2 \sin(1/x)$ for $x \neq 0$ and $f(0) = 0$. Show that f is a differentiable function of class C^1 . Solution : Note that $f'(x) = 2x \sin(1/x) - \cos(1/x)$ for $x \neq 0$ and $f'(0) = 0$. Apply L'Hospital's Rule, on taking, $h \rightarrow 0$ we see that $f'(0) = 0$. 2. Let $f : \mathbb{R}^2 \rightarrow \mathbb{R}$ be defined by $f(x, y) = x^2 y^2$. Show that f is of class C^2 at $(0, 0)$. Applying L' Hospital rule successively, we find $f''(0, 0) = 0$.

9 Proceeding in this manner, we can show that $f^{(n)}(x) = 0$ for $n \geq 2$. Hence f is a function of class C^1 . A mapping $f: U \rightarrow V$ of an open set $U \subset \mathbb{R}^n$ to an open set $V \subset \mathbb{R}^n$ is called a homeomorphism if i) f is bijective i.e. one to one and onto, as well as ii) f, f^{-1} are continuous. Exercise : 2. Let $f: \mathbb{R}^2 \rightarrow \mathbb{R}^2$ be such that $f(x, y) = (x^2 + y^2, 3x - y)$. Show that f is a homeomorphism on \mathbb{R}^2 . 3. Let $f: \mathbb{R}^2 \rightarrow \mathbb{R}^2$ be defined by $f(x, y) = (x^2 + y^2, 3x - y)$. Test i) whether f is a differentiable function of class C^1 or not ii) whether f is a homeomorphism or not. [Ans. : i) f is of class C^1 . ii) f is homeomorphism] Solution : 2. Note that $f(x, y) = (x^2 + y^2, 3x - y)$ if and only if $x = y^2$. Hence f is one one. Let $y = x^2 + 3x - y$ and hence $f^{-1}: \mathbb{R}^2 \rightarrow \mathbb{R}^2$ is defined as $f^{-1}(x, y) = (x^2 + y^2, 3x - y)$. Again, $f^{-1}(x, y) = (x^2 + y^2, 3x - y)$ and $f^{-1}(x, y) = (x^2 + y^2, 3x - y)$. Thus f is onto. Consequently f is bijective. $f: U \rightarrow V \subset \mathbb{R}^n$ is called a homeomorphism if i) f is a homeomorphism of U onto V and ii) f, f^{-1} are of class C^k . when f is a C^k - diffeomorphism, we simply say diffeomorphism. Exercise : 4. Consider the mapping $f: \mathbb{R}^2 \rightarrow \mathbb{R}^2$ given by $f(x, y) = (x^2 + y^2 \cos y, x^2 + y^2 \sin y)$. Show that f is one-to-one on a sufficiently small neighbourhood of each point $(x, y) \in \mathbb{R}^2$ with $x \neq 0$. Solution : The given mapping $f: \mathbb{R}^2 \rightarrow \mathbb{R}^2$ is given by $f(x, y) = (x^2 + y^2 \cos y, x^2 + y^2 \sin y)$. Then, we have $f(x, y) = (x^2 + y^2 \cos y, x^2 + y^2 \sin y)$. Hence each $f^{-1}(x, y) = (x, y)$ is continuous for all values of x and y in \mathbb{R}^2 . Thus f is continuously differentiable on \mathbb{R}^2 . Again the Jacobian is given by $J = \begin{pmatrix} 2x + 2y \cos y & -y^2 \sin y \\ 2x & 2y \cos y + 2x \sin y \end{pmatrix}$ if and only if $x \neq 0$ in \mathbb{R}^2 . Consequently, f is one-to-one on a sufficiently small neighbourhood of each point $(x, y) \in \mathbb{R}^2$ with $x \neq 0$. A mapping $f: U \rightarrow V$ of an open set $U \subset \mathbb{R}^n$ onto an open set $V \subset \mathbb{R}^n$ is called a C^k - diffeomorphism, $k \geq 1$ if i) f is a homeomorphism of U onto V and ii) f, f^{-1} are of class C^k . when f is a C^k - diffeomorphism, we simply say diffeomorphism. Exercise : 5. Let $f: \mathbb{R}^2 \rightarrow \mathbb{R}^2$ be defined by $f(x, y) = (x^2 + y^2, x^2 - y^2)$. Determine whether f is a diffeomorphism or not. 6. Let $f: \mathbb{R}^2 \rightarrow \mathbb{R}^2$ be defined by $f(x, y) = (x^2 + y^2, x^2 - y^2)$. Show that f is a diffeomorphism. [Ans. : 5. f is a diffeomorphism] For $i = 1, 2$; let $u_i: \mathbb{R}^2 \rightarrow \mathbb{R}$ be the coordinate functions on \mathbb{R}^2 i.e. for every $p \in \mathbb{R}^2$, $u_i(p) = p_i$ where $p = (x, y) \in \mathbb{R}^2$. Such u_i 's are continuous functions from \mathbb{R}^2 to \mathbb{R} . We call this n -tuple of functions (u_1, \dots, u_n) the standard co-ordinate system of \mathbb{R}^n . If $f: U \subset \mathbb{R}^n \rightarrow \mathbb{R}^n$ is a mapping defined on $U \subset \mathbb{R}^n$, then, f is determined by its co-ordinate functions $f_i = u_i \circ f$ where $f_i = u_i \circ f$. 1.2) $f = (f_1, \dots, f_n)$ and each $f_i: U \rightarrow \mathbb{R}$ are real valued functions, defined on an open subset U of \mathbb{R}^n . Thus for every $p \in U$, $f(p) = (f_1(p), \dots, f_n(p))$ where $f_i(p) = u_i(f(p))$. 1.3) consequently f is of class C^k if each of its co-ordinate functions $f_i = u_i \circ f$ is of class C^k . \mathbb{R}^n is a Hausdorff, second countable space. If every point of M has a neighbourhood homeomorphic to an open set in \mathbb{R}^n , then M is said to be a manifold.

17.3.1) i) $f: U \rightarrow V$ and ii) the mapping $\gamma: I \rightarrow M$ is of class C^k . By a differentiable mapping, we shall mean, unless otherwise stated, a mapping of class C^1 . If (α, β) and (γ, δ) are respectively the local coordinate systems defined in a neighbourhood U of p of M and V of $f(p)$ of N , then it can be shown, as done earlier (3.2) $\gamma \circ f \circ \alpha^{-1} \circ \beta^{-1}$ where g is a differentiable function defined on $V \rightarrow N$ and (3.3) $g \circ f \circ \alpha^{-1} \circ \beta^{-1}$. Let M and N be two n -dimensional differentiable manifolds. A mapping $f: M \rightarrow N$ is called a diffeomorphism if i) f and f^{-1} are differentiable mappings of class C^1 ii) f is a bijection. In such cases, M and N are said to be diffeomorphic to each other.

Exercise: 1. Let M and N be two differentiable manifolds with $M=N=R$. Let (U, α) and (V, β) be two charts on M and N respectively, where $U = R^1$ and $V = R^1$ be the identity mapping and $V = R^1 \rightarrow R^1$ be the mapping defined by $\beta \circ f \circ \alpha^{-1}(x) = x^3$. Show that the two structures defined on R are not C^1 -related even though M and N are diffeomorphic where $f: M \rightarrow N$ is defined by $f(t) = t^3$. Hint: Note that, $(\alpha^{-1})^{-1} \circ \beta \circ f \circ \alpha^{-1}(x) = x^3$ and $(\beta^{-1})^{-1} \circ \alpha \circ f \circ \beta^{-1}(x) = x^3$. Thus f is of class C^1 but f^{-1} is not of class C^1 . Again $(\alpha^{-1})^{-1} \circ \beta \circ f \circ \alpha^{-1}(x) = x^3$ Also $f \circ \alpha^{-1} \circ \beta^{-1}(x) = x^3$ if and only if $y = x^3$. Thus f is one-one. Finally $f \circ \alpha^{-1} \circ \beta^{-1}(x) = x^3$, so that $f \circ \alpha^{-1} \circ \beta^{-1}(x) = x^3$ and $f \circ \alpha^{-1} \circ \beta^{-1}(x) = x^3$. Thus f is a bijection. Note: A diffeomorphism f of M onto itself is called a transformation of M . A real-valued function on M ; i.e. $f: M \rightarrow R$ is said to be a differentiable function of class C^1 , if for every chart (U, α) containing p of M , the function $f \circ \alpha^{-1} \circ \beta^{-1} \circ \gamma^{-1}: I \rightarrow R$ is of class C^1 . We shall often denote by $F(M)$, the set of all differentiable functions on M and will sometimes denote by $F(p)$, the set of functions on M which are differentiable at p of M . R^n is a $F(M)$ if $U = M$ and $\alpha = \text{id}$.

19. It is to be noted that such $F(M)$ is i) an algebra over R ii) a ring over R iii) an associative algebra over R and iv) a module over R . Where the defining relations are a) $(f+g)(p) = f(p) + g(p)$ b) $(fg)(p) = f(p)g(p)$ c) $(f \circ g)(p) = f(g(p))$.

1.4. Differentiable Curve: We are now in a position to define a curve on a manifold. A differentiable curve through p in M of class C^1 is a differentiable mapping $\gamma:]c, d[\rightarrow M$, namely the restriction of a differentiable mapping of class C^1 of an open interval $]c, d[$ containing $[a, b]$, such that $\gamma(a) = p$, $\gamma(b) = q$. Also $\gamma'(t) \neq 0$. We write it as $\gamma:]c, d[\rightarrow M$. The tangent vector to the curve γ at p is a function $T_p: F(p) \rightarrow R^n$ defined by $T_p(f) = \lim_{t \rightarrow a} \frac{f(\gamma(t)) - f(\gamma(a))}{t - a}$ where $\gamma(a) = p$. It can be shown that it satisfies 4.5) $X(fg) = Xf \cdot g + f \cdot Xg$ and 4.6) $X(fg) = Xf \cdot g + f \cdot Xg$. Note: Each function $X_p: F(p) \rightarrow R^n$ cannot be a tangent vector to some curve at $p \in M$, unless it is a linear function and satisfies Leibnitz Product Rule. Exercises: 1. Let a curve γ on R^n be given by $\gamma(t) = (t^2, t^3, t^4)$, $t \in]1, 2[$. Find the tangent vector to the curve γ at the point $(1, 1, 1)$. 2. If C is a constant function on M and X is a tangent vector to some curve γ at $p \in M$, then $X \cdot C = 0$. [Ans. i) $(1, 1, 1)$ ii) use 4.5), 4.6) and the definition of constant function. Let us define $X = \frac{d}{dt} \gamma(t)$. 4.8) $X \cdot (fg) = Xf \cdot g + f \cdot Xg$. If we denote the set of tangent vectors to M at p by $T_p(M)$, then from 4.7) and 4.8) it is easy to verify that $T_p(M)$ is a vector space over R . We are now going to determine the basis of such vector space. For each $i = 1, \dots, n$, we define a mapping $\gamma_i:]c, d[\rightarrow M$ by $\gamma_i(t) = (t, \dots, t, \dots, t)$ where t is in the i -th position. Note that $\gamma_i'(t) = (1, \dots, 1, \dots, 1)$ where the 1 is in the i -th position. Let us define a differentiable curve $\gamma:]c, d[\rightarrow M$ by $\gamma(t) = (t^2, t^3, \dots, t^i, \dots, t^i, \dots, t^2)$, for fixed i . Then $\frac{d}{dt} \gamma(t) = (2t, 3t^2, \dots, t^{i-1}, \dots, t^{i-1}, \dots, 2t)$ by chain rule. For fixed i , by (4.3) $\{ \gamma_i'(t) \}$ is a basis for $T_p(M)$. Thus we can claim that each $X = \sum_{i=1}^n \alpha_i \gamma_i'(t)$, $\alpha_i \in R$, is a tangent vector to the curve γ defined above, at $p = \gamma(t)$. Again from the definition of the tangent vector, $X(f) = \lim_{t \rightarrow a} \frac{f(\gamma(t)) - f(\gamma(a))}{t - a}$ where $\gamma(a) = p$. By chain rule $X(f) = \sum_{i=1}^n \alpha_i \frac{d}{dt} \gamma_i(t)$ where $\gamma_i(t) = (t, \dots, t, \dots, t)$ where t is in the i -th position. We write it as $X = \sum_{i=1}^n \alpha_i \frac{d}{dt} \gamma_i(t)$ where $\gamma_i(t) = (t, \dots, t, \dots, t)$ where t is in the i -th position. Thus each $\alpha_i: M \rightarrow R$, $i = 1, \dots, n$, is a differentiable function and every tangent vector, say X_p , to some curve, say $\gamma(t)$ at $p = \gamma(t)$ can be expressed as a linear combination of the tangent vector $\gamma_i'(t)$, $i = 1, \dots, n$, to the curve γ defined in (4.10).

23 If possible, for a given linear combination of the form $\sum_{i=1}^n \lambda_i \mathbf{p}_i(x)$, where λ_i 's are functions on M , let us define a curve γ by $\gamma(t) = \sum_{i=1}^n \lambda_i(t) \mathbf{p}_i(x)$. If we assume that $\sum_{i=1}^n \lambda_i \mathbf{p}_i(x) = 0$ then, $\sum_{i=1}^n \lambda_i \mathbf{p}_i(x) = 0$ where $x \in M$. Thus the set $\{\mathbf{p}_1, \dots, \mathbf{p}_n\}$ is linearly independent. Hence we state Theorem 1: If (x, y) is a local coordinate system in a neighbourhood U of $p \in M$,

then, the basis of the tangent space $T_p(M)$ is given by $\mathbf{p}_1, \dots, \mathbf{p}_n$. Let us define $T(M) = \cup_{p \in M} T_p(M)$. This $T(M)$ is called the tangent space of M .

24 1.5. Vector Field : In classical notation, if to each point p of R^3 or in a domain U of R^3 , a vector $(\cdot)_p$ is specified, then, we say that a vector field is given on R^3 or in a domain U of R^3 . A vector field X on M is a correspondence that associates to each point $p \in M$, a vector $X_p \in T_p(M)$. In fact, if $f \in F(M)$, then Xf is defined to be a real-valued function on M , defined as follows 5.1) $(Xf)(p) = X_p f$. A vector field X is called differentiable if Xf is so for every $f \in F(M)$. Using (4.11) of 1.4, a vector field X may be expressed as 5.2) $X = \sum_{i=1}^n X^i \frac{\partial}{\partial x^i}$ where X^i 's are differentiable functions on M . Let (M) denote the set of all differentiable vector fields on M . We define 5.3) $[X, Y] = X(Y) - Y(X)$.

It is easy to verify that (M) is a vector space over R . Also, for every $f \in F(M)$, fX is defined to be a vector field on M , defined as 5.4) $(fX)(p) = f(p)X_p$. Let us define a mapping as $[\cdot, \cdot] : F(M) \times F(M) \rightarrow F(M)$ as 5.5) $[X, Y] = X(Y) - Y(X)$. Such a bracket is known as Lie bracket of X, Y . Exercises : 1. Show that for every X, Y, Z in (M) , for every f, g in $F(M)$, i) $[X, Y] = -[Y, X]$ ii) $[bX, Y] = b[X, Y]$ iii) $[X, Y + Z] = [X, Y] + [X, Z]$ iv) $[X, Y + Z] = [X, Y] + [X, Z]$ v) $[X, X] = 0$ vi) $[X, Y] = -[Y, X]$ vii) $X[Y, Z] - Y[X, Z] + Z[X, Y] = 0$: Jacobi Identity viii) $[fX, gY] = (fg)[X, Y] + \{f(Xg) - g(Yf)\}X - \{g(Yf)X - f(Xg)Y\}$ a) $[X, fY] = f[X, Y] + (Xf)Y$ b) $fX, Y = f[X, Y] - (Yf)X$.

In terms of a local co-ordinate system (x^1, \dots, x^n) , $L = \sum_{i,j=1}^n g_{ij} dx^i dx^j$, where $X = \sum_{i=1}^n X^i \frac{\partial}{\partial x^i}$, $Y = \sum_{j=1}^n Y^j \frac{\partial}{\partial x^j}$. Complete $[X, Y]$ where i) $X = x^1 \frac{\partial}{\partial x^1} + x^2 \frac{\partial}{\partial x^2}$, $Y = x^1 \frac{\partial}{\partial x^1} + x^2 \frac{\partial}{\partial x^2}$. Prove that i) (M) is a $F(M)$ module Hints : 1. viii) Note that $\{f(Yh)\}(p) = f(p)(Yh)_p$ by (5.4) of 1.5) = $f(p)Y_p h$ by (5.1) of 1.5) Again, $\{(fY)\}(p) = (fY)(p)h$ by (5.1) = $f(p)Y_p h$ by (5.4) Thus $\{f(Yh)\}(p) = \{(fY)h\}(p)$, $p f(Yh) = (fY)h$ Use the above result, 5.5) of 1.5 & (4.6) of 1.4, the result follows after a few steps.

26 . 1.6. Integral Curve : In this article, we are going to give the geometrical interpretation of a vector field. Let Y be a vector field on M . The assignment of the vector Y_p at each point $p \in U \subset M$, is given by $Y : p \rightarrow Y_p \in T_p(M)$. A curve is an integral curve of Y if the range of γ is contained in U and for every $a < t < b$ in the domain $[a, b]$ of γ , the tangent vector to γ at $(t, \gamma(t)) = p$ coincides with Y_p . i.e. $\dot{\gamma}(t) = Y_{\gamma(t)}$. $\gamma : (a, b) \rightarrow M$, $\dot{\gamma}(t) = Y_{\gamma(t)}$ by (4.4) of 1.4 Using 4.11) 1.4 one can write $\dot{\gamma}(t) = \sum_{i=1}^n \dot{\gamma}^i(t) \mathbf{p}_i(\gamma(t))$ where $\dot{\gamma}^i$'s are functions on M . $\int \sum_{i=1}^n \dot{\gamma}^i(t) \mathbf{p}_i(\gamma(t)) dt = \int \sum_{i=1}^n \dot{\gamma}^i(t) dx^i$ or $\int \sum_{i=1}^n \dot{\gamma}^i(t) dx^i = \int \sum_{i=1}^n \dot{\gamma}^i(t) dx^i$ Using (4.3) of 1.4 we get

27 $\int \sum_{i=1}^n \dot{\gamma}^i(t) dx^i = \int \sum_{i=1}^n \dot{\gamma}^i(t) dx^i$ Hence they are related by 6.1) $dx^i dt = \dot{\gamma}^i(t) dt$. Exercises : 1. Find the integral curve of a zero vector. 2. Find the integral curve of the following vector field i) $X = x^1 \frac{\partial}{\partial x^1} + x^2 \frac{\partial}{\partial x^2}$ on R^2 ii) $X = x^1 \frac{\partial}{\partial x^1} + x^2 \frac{\partial}{\partial x^2}$ on R^2 iii) $X = x^1 \frac{\partial}{\partial x^1} + x^2 \frac{\partial}{\partial x^2}$ on R^2

Solution : 2.i) From (6.1) of 1.6, we see that $dx dt x 1 1 ? , dx dt x 2 2 ?$ or $dx x dt 1 1 ? , dx x dt 2 2 ?$ Integrating $\log x t 1 ? ? C$, $\log x t D 2 ? ?$ say, where C, D are integrating constant. When $t = 0$, if $x p 1 1 ? , x p 2 2 ?$, then from $x C e t 1 ?$ and $x D e t 2 ?$ we find that $p 1 = C, p 2 = D$ Thus $?, p e p e t t 1 2 b g$ is the integral curve of X passing through the point $1 2 (,) p p$ 28 . 1.7 Differential of a mapping : Let $f : M \rightarrow N$ be a differentiable mapping of an n -dimensional manifold M to an m -dimensional manifold N . Let $F(p)$ denote the set of all differentiable functions at $p \in M$ and $F f p () b g$ denote the set of all differentiable functions at $f p \in N () . ?$ Such a map f , induces a map $f F f p F p * : () () b g ?$, usually called pull back map. and is defined by 7.1) $f g g f * () , ? ? ? ? () g F f p ?$ called the pull back of g by f , which satisfies 7.2) $f a g b h a f g b f h * * * () () ? ? ? f g h f g f h * * * () () ?$ where $?, () g h F f p ?$ and $, ? a b R$ The map f , also induces a linear mapping $f T M p * : () ? () () f p T N$ such that 7.3) $?, ? ? ? * * () () p p p f X g X g f X f g ? ? ?$ called the push forward of X by f . Such $f *$ is also called derived linear map or Jacobian map or differential map of f on $T p (M)$ $f *$? push forward objects defined on objects defined on $f *$? pull back $N M f f *$

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f M N p f p ? () T p (M) T f(p) (N) { 29		

Let us write 7.4) $f X f X p f p * * () () ?$ We can also define push forward of X by f , geometrically, in the following manner : Given a differential mapping $f M N : , ?$ the differential of f at $p \in M$ is the linear mapping f

$T M p * : () ? () () f p T N$ defined as follows : For each $X p \in T p (M)$, we choose a curve $?() t$ in M such that $X p$ is the tangent vector to the curve $?()$

t at p
 $t ? ? () . 0$ Then $f X p * ()$ is defined to be the tangent vector to the curve $f t ? () b g$ at $f p f t () () ? ? 0 b g$ Exercises : 1. If f is a differentiable map from a manifold M into another manifold N and g is a differentiable map from N into another manifold L , then, show that i) $() * * * g f g f ? ? ?$ ii) $() * * * g f f g ? ? ?$ 2. If f is a transformation of M and g is a differentiable function on M , prove that i) $f X Y f X Y * * [,] ?$ ii) $f f X g X g f g * * * () b g ?$ iii) $f g X g f f X * * () () ? ? ? 1$ for all vector fields X, Y on M . Solution : 1. By definition,

$f X p * ()$ is the tangent vector to the curve $f t ? () b g$ at $f p f$
 $t () () ? ? 0 b g$ where
 $X p$ is the tangent vector to the curve $?() t$ at p
 $t ? ? () . 0$ Hence by (4.4) of 1.4

30 $f g p * () X d i ? d d t g f t t t () ? b g L N M O Q P ? 0 g F f p ? () b g ? L N M O Q P ? d d t g f t t t () ? ? b g 0 = X p () g$
 $f ?$ by 4.4) of 1.4 Hints 3. Given that $f : M \rightarrow M$ is a transformation and hence for every $p \in M, f p q () , ?$ say.. Thus, $p f q ? ? 1 ()$ consequently, from 7.3) of 1.7, we find that $f X g f p p * () d i { } ? X g f p p () () , ? n s ? ? p M$ or $f X g q X g f f q p p * () () () d i { } n s ? ? ? 1$ or $f X g X g f f * () () b g b g ? ? ? 1$ Using this relation, one can deduce the three results. We are now going to give a matrix representation of the linear mapping $f *$. Theorem 1 : If f is a mapping from an n -dimensional manifold M to an m -dimensional manifold N , where $(,) x x n 1 ?$ is the local co-ordinate system in a neighbourhood of a point p of M and $(,) y y m 1 ?$ is the local co-ordinate system in a neighbourhood of $f p ()$ of N , then $f x f x y i p j i p j m j f p * () ? ? ? ? ? F H I K ? F H G I K J ? ? 1$ where $f y f j j ? ?$ Proof : We write $f x a y i p i j j f p j m * () , ? ? ? ? F H I K ? F H G I K J ? ? 1$ in $?, \dots,$

31 where $a s i j$, are unknown to be determined or $f x y a y y i k i j j f p k j m * () ? ? ? ? F H I K R S T U V W ? F H G I K J ? ? 1$ where each $y F f p k ? ? () k m ? 1, \dots$, using 7.3) of 1.7, we find $?, ? x y f a i p k i j j k j m F H I K ? ? ? () ? 1$ or $?, ? x f a i p k i k F H I K ?$ or $?, ? x a k i p i k F H G I K J ?$ by (4.9) of 1.4 Thus $f x f x y i p j i p j f p j m * () ? ? ? ? ? F H I K ? F H G I K J F H G I K J ? ? 1$ Note : 1. The matrix of $f *$, denoted by $(f *)$ is

given by $() *$
 f
 f
 $x f x f x f x f x f x f x f$
 x
 n

$n \times m$ matrix M is invertible if $\det M \neq 0$. The kernel of f^* is the set of $X \in T_p(M)$ for which $f^*(X) = 0$. The image of f^* is the set of $Y \in T_p(N)$ for which, there exists $X \in T_p(M)$ such that $f^*(X) = Y$. Now from a known theorem
 $\dim(\ker f^*) + \dim(\text{Range } f^*) = \dim T_p(M)$. We write it as 7.5) $\dim(\ker f^*) + \dim(\text{Range } f^*) = \dim T_p(M)$ for each $p \in M$. The $\dim(\text{Range } f^*)$ is called the rank f^* . If $\text{rank } f^* = \dim T_p(M)$ we say i) f is an immersion if $\dim M < \dim N$ and $f(M)$ is an immersed submanifold of N ii) f is an imbedding if f is one to one and an immersion and then $f(M)$ is an imbedded submanifold of N iii) f is a submersion if $\dim M > \dim N$. Exercises : 1. Show that $f: \mathbb{R}^2 \rightarrow \mathbb{R}^2$ given by $f(t) = (a \cos t, a \sin t)$ is an immersion. 2. Find (f^*) in the following cases i) $f: \mathbb{R}^2 \rightarrow \mathbb{R}^2$ given by $f(x, y) = (x^2, y^2)$,
 $x \times x \times 1 \ 2 \ 2 \ 2 \ 1 \ 2 \ 2 \ 3 \ ? \ b \ g$ ii) $f: \mathbb{R}^2 \rightarrow \mathbb{R}^2$ given by $f = x \ e \ x \ x \ e \ x \ x \times 1 \ 2 \ 1 \ 2 \ 2 \ 2 \ ? \ ?$,
 $c \ h$ at $(0, 0)$ where (x, y) are the local co-ordinates on \mathbb{R}^2 . 1.8 f-related vector Field : Let X and Y be fields on M and N respectively. Then, for $p \in M$, let $p \in X \in T_p(M)$ and $f(p) \in Y \in T_p(N)$ and such that 8.1) $f^*(X) = Y$ where $f: M \rightarrow N$ is a differentiable mapping and f^* is already defined in the previous article. In such a case, we say that the two vector fields X, Y are f-related.
3.3 For $g \in F_p$ we see that $f^*(X) = Y$ using 7.3) of 1.7 and (5.1) of 1.5 we find that $X \cdot g = f^*(Y \cdot g)$. Then 8.2) $(f^*)^2(X) = X \cdot g$. If f is a transformation on M and $f^*(X) = X$ we say that, X is f-related to itself or X is invariant under f . We also write it as 8.3) $f^*(X) = X$. Exercises : 1. Let X, Y be two f-related vector fields on M and N respectively.. Show that the vector fields $[X_1, X_2]$ and $[Y_1, Y_2]$ are also f-related. 2. Prove that two vector fields X, Y respectively on M and N are f-related if and only if $f^*(X) = Y$ where $f: M \rightarrow N$ is a C map. 3. If f is a transformation on M , show that, for every $X \in T_p(M)$, there exists a unique f-related vector field to X . Solution : 1. From the definition of the Lie bracket, we see that $[f^*(X), f^*(Y)] = f^*([X, Y])$.
 $g \cdot f^*(X) = f^*(X) \cdot g$
 $b \cdot g = g \cdot b$
 $Y \cdot g = X \cdot g$ by (8.2) above
 $Y \cdot g = X \cdot g$ by (8.2) above
 $l \cdot q$
34 $[f^*(X), f^*(Y)] = f^*([X, Y])$
 $l \cdot q$ from the definition of the Lie Bracket. Hence from 8.2), one claims that $[X_1, X_2]$ and $[Y_1, Y_2]$ are f-related. 1.9 One parameter group of transformations on a manifold : Definition Let a mapping $\tau: \mathbb{R} \rightarrow M$ is defined by $\tau(t) = (x(t), y(t))$ which satisfy i) for each $t \in \mathbb{R}$, $\tau(t) \in M$ is a transformation on M and $0 \in \mathbb{R}$ ii) for all $t, s, t + s \in \mathbb{R}$
 $\tau(t) \cdot \tau(s) = \tau(t + s)$. Then the family $\tau(t)$ of mappings is called a one-parameter group of transformations on M . Exercise : 1. Let $\tau(t)$ be a one-parameter group of mappings on M . Show that i) $\tau(t) \cdot \tau(s) = \tau(t + s)$ ii) $\tau(t)$ form an abelian group. Let us set 9.1) $\tau(t) = (x(t), y(t))$. Then $\tau(t)$ is a differentiable curve on M such that $\tau(0) = p$ by Def. (i) above Such a curve is called the orbit through p of M . The tangent vector, say X_p to the curve $\tau(t)$ at p is therefore 9.2) $X_p = \frac{d}{dt} \tau(t) \Big|_{t=0}$. $l \cdot n \cdot m \cdot o \cdot q \cdot p$ $\lim_{t \rightarrow 0} \tau(t) = p$. $\tau(t) = (x(t), y(t))$
35 In this case, we say that $\tau(t)$ induces the vector field X and X is called the generator of τ . The curve $\tau(t)$ defined by 9.1) is called the integral curve of X . Exercises : 2. Show that the mapping $\tau: \mathbb{R} \rightarrow \mathbb{R}^3$ defined by $\tau(t) = (e^t, e^{-t}, t)$ is a one-parameter group of transformations on M and the generator is given by $X = x \frac{\partial}{\partial x} - y \frac{\partial}{\partial y} + \frac{\partial}{\partial z}$. Let $M = \mathbb{R}^2$ and let $\tau: \mathbb{R} \rightarrow M$ be defined by $\tau(t) = (x(t), y(t))$. Show that τ defines a one-parameter group of transformation on \mathbb{R}^2 and find its generator.. Note : Since every 1-parameter group of transformations induces a vector field on M , the question now arises whether every vector field on M generates one parameter group of transformations. This question has been answered in the negative. Example : Let $X = x \frac{\partial}{\partial x} + 1 \frac{\partial}{\partial y}$ on $M = \mathbb{R}^2$. Then, $\frac{dx}{dt} = x$, $\frac{dy}{dt} = 1$. Thus $e^t = A x$, $x = B e^t$, where A, B are integrating constant. Let $x = 1$, $x = 2$ for $t = 0$. Then, $A = e^{-1}$, $B = 2e^{-1}$.
36 Consequently the integral curve of X is $(t) = (\log t, t)$ which is not defined for all values of t in \mathbb{R} . Thus, if $\tau(t) = (t, e^t)$, $\tau(t)$ does not generate one parameter group of transformations. Problem 7 leads us to the following definition : Let I be an open interval (a, b) and U be a nbd of p of M . Let a mapping $\tau: I \rightarrow U$ denoted by $\tau(t) = (x(t), y(t))$ be such that i) U is an open set of M ii) for each $t \in I$, $\tau(t)$ is a transformation of U onto an open set $V(t)$ of M and $0 \in \mathbb{R}$ iii) if $t, s, t + s$ are in I and if $s \in V(t)$ then $\tau(t) \cdot \tau(s) = \tau(t + s)$. Such a family $\tau(t)$ of mappings is called a local one parameter group of transformations, defined on I . We are now going to establish the following theorem Theorem 1 : Let X be a vector field on a manifold M . Then, X generates a local one-parameter group of

transformations in a neighbourhood of a point of M . Proof : Let (\dots)
 $x \times x \times n \times 1 \times 2$ be a local coordinate system in a neighbourhood U of
 p of M such that $\varphi(U, \dots) \cong \mathbb{R}^n$, where (U, φ) is the chart at p of M . Then $x \times p \times U \times \varphi^{-1}(\dots)$, $\dots \times \mathbb{R}^n \times 1, \dots$, Let X
 $x \times x \times i \times \dots \times (\dots) \times 1$
 X be a given vector field on U , the neighbourhood of $p \in M$, where each X_i , the components of X , are differentiable
functions on U of M . Then, for every X on M , we have a φ -related vector field on $\mathbb{R}^n(U) \cong \mathbb{R}^n$ with $\varphi_*(p) = (\dots) \times 0 \times n$
 $\times 1 \times U \times \mathbb{R}^n$. Let \tilde{X}_i be the components of the φ -related vector field on $U \cong \mathbb{R}^n$. Then by the existence theorem of
ordinary differential equations, for each $\varphi(p) \in U \cong \mathbb{R}^n$, there exists a $\varphi^{-1}(0)$ and a neighbourhood $V \subset U$ of $\varphi(p)$, $V \subset U \subset \mathbb{R}^n$
such that, for each $q \in V \cong \mathbb{R}^n$, $q \in \varphi(U)$, say, there exists n -tuple of C^1 functions f_1, \dots, f_n defined on $I \times \mathbb{R}^n$
 $I \times \mathbb{R}^n \times 1$ and mapping $f: I \times V \subset U \times \mathbb{R}^n \times 1, \dots$, which satisfies the system of first order differential equations $\frac{d}{dt} \varphi^{-1} \circ \tilde{X}_i$
 $(t, \dots) = f_i(t, \dots)$, with the initial condition $\varphi^{-1}(0) = (0, \dots)$. Let us write $\tilde{X} = \sum_{i=1}^n f_i \frac{\partial}{\partial x_i}$. We are
going to show $\tilde{X} = \varphi_* X$. Note that if $t, s, t+s$ are in I and if $q \in V \subset U \subset \mathbb{R}^n$ then each $f_i(t, \dots), f_i(t+s, \dots)$
 $f_i(s, \dots)$ are defined on $I \times U \subset \mathbb{R}^n$. Now let us set $g = \varphi_* X - \tilde{X}$. For simplicity,
we write $g = \sum_{i=1}^n g_i \frac{\partial}{\partial x_i}$. Then each $g_i(t, \dots)$ is defined on $I \times U \subset \mathbb{R}^n$ and hence satisfies $\frac{d}{dt} g_i(t, \dots) = 0$
condition $g_i(0, \dots) = 0$.
38 Also, let us set $h = \varphi_* X - \tilde{X}$. For simplicity, we write $h = \sum_{i=1}^n h_i \frac{\partial}{\partial x_i}$. Then each $h_i(t, \dots)$ is defined on $I \times U \subset \mathbb{R}^n$
and hence satisfies $\frac{d}{dt} h_i(t, \dots) = 0$. Hence from the uniqueness we must have $g = h = 0$. Using 3) we must have $\tilde{X} = \varphi_* X$. Thus, we claim that, for every vector field defined in a
neighbourhood U of $\varphi(p)$ of \mathbb{R}^n , there exists $\varphi^{-1} \circ \tilde{X}$ as its local 1-parameter group of transformations defined
on $I \times U \subset \mathbb{R}^n$. Let us set $V = \varphi^{-1}(V) \subset U$ and define $\varphi_t: V \rightarrow M$ as follows $\varphi_t = \varphi^{-1} \circ \tilde{X}_t \circ \varphi$. Then φ_t is a transformation of V onto $\varphi^{-1}(V)$ of M .
39 ii) if $t, s, t+s$ are in I and if $q \in V$, then $\varphi_{t+s} \circ \varphi_t^{-1} \circ \varphi_s = \varphi_{t+s}$, after a few
steps $\varphi_{t+s} \circ \varphi_t^{-1} \circ \varphi_s = \varphi_{t+s}$. Thus for the given vector field X , defined in a neighbourhood U of p of M , there exists φ_t as its
local 1-parameter group of transformations, defined on $I \times V \subset U$ of M . Note that if we define $\varphi_t = \varphi^{-1} \circ \tilde{X}_t \circ \varphi$,
 $\varphi_t \circ \varphi_s^{-1} \circ \varphi_{t+s} = \varphi_{t+s}$, then φ_t is the integral curve of X . This completes the proof. Theorem 2 : Let φ_t
be a transformation of M . If a vector field X generates φ_t as its local 1-parameter group of transformations, then, the
vector field $\varphi_t^* X$ will generate φ_t as its local 1-parameter group of transformations. Proof : Left to the reader.
Exercise : 4. Show that a vector field X on M is invariant under a transformation φ_t on M if and only if $\varphi_t^* X = X$
where φ_t is the local 1-parameter group of transformations induced by X . We
now give a geometrical interpretation of $[X, Y]$, for every vector field X, Y on M . Theorem 3 : If X generates φ_t as its local
1-parameter group of transformations, then, for every vector field Y on M , $\lim_{t \rightarrow 0} \frac{1}{t} (\varphi_t^* Y - Y) = [X, Y]$
where $q \in U$ and $(\dots) \times \dots \times t \times Y \times p \times b \times g$
40 To prove the theorem, we require some lemmas which are stated below : Lemma 1 : If $\varphi(t, p)$ is a function on $I \times M$,
where I is an open interval (\dots) such that $\varphi(0, p) = 0$, $\varphi \in M$ then, there exists a function $h(t, p)$ on $I \times M$ such that $\frac{d}{dt} \varphi(t, p) = h(t, p)$. Moreover $h(0, p) = \varphi_*(0, p)$, where $\varphi_*(0, p) = \frac{d}{dt} \varphi(t, p)|_{t=0}$. Proof : It is sufficient to define $h(t, p) = \frac{d}{dt} \varphi(t, p)$.
Hence by the fundamental theorem of calculus $\varphi(t, p) = \int_0^t h(s, p) ds + \varphi(0, p)$. Also
from above $h(0, p) = \varphi_*(0, p)$. Lemma 2 : If f is a function on M and X is a
vector field on M which induces a local 1-parameter group of transformations φ_t then there exists a function g defined
on $I \times V$, V being the neighbourhood of p of M , where
 $\frac{d}{dt} \varphi_t^* f = X \varphi_t^* f$ such that $f \varphi_t = \int_0^t X \varphi_s^* f ds + f \varphi_0$
41 Moreover, $X f = 0$ on M . Symbolically, $X f = 0$ on M . Proof : Let us set $\tilde{f} = f \varphi_t$. Then \tilde{f} is a function on $I \times M$ such that $\tilde{f}(0, p) = f(p)$. Hence by
Lemma 1, there exists a function, say, $g(t, p)$ on $I \times V$, $V \subset M$ being the neighbourhood of p of M , such
that $\frac{d}{dt} \varphi_t^* f = g(t, p)$. If $X f = 0$, then $g(t, p) = 0$. or, $\frac{d}{dt} \varphi_t^* f = 0$. Hence $\varphi_t^* f = f$.
 $\lim_{t \rightarrow 0} \frac{1}{t} (\varphi_t^* f - f) = 0$. As, $\frac{d}{dt} \varphi_t^* f = X \varphi_t^* f$, we find that $X f = 0$. Proof of the main
theorem : Let us write $\varphi_t = \varphi_t \circ \varphi_0^{-1}$.

42 Now, $(\phi^{-1})^* \circ \phi^* = \text{id}$ by Lemma 2 or $(\phi^{-1})^* \circ \phi^* = \text{id}$ by Lemma 2. From the definition we find that, $X_{\phi^{-1}(p)} = \phi_* X_p$. Taking $f = Y$, we find from above after a few steps $X_{\phi^{-1}(p)}(Y_{\phi^{-1}(p)}) = X_p(Y_p)$. Thus we write, $(\phi^{-1})^* X_p = X_{\phi^{-1}(p)}$. Note: We abbreviate the above result as $(\phi^{-1})^* X_p = X_{\phi^{-1}(p)}$. Corollary: 1. Show that $(\phi^{-1})^* \circ \phi^* = \text{id}$.

43 Proof: From the last theorem $(\phi^{-1})^* \circ \phi^* = \text{id}$, as ϕ^* is a linear mapping $(\phi^{-1})^* \circ \phi^* = \text{id}$, from a known result Using the definition of local 1-parameter group of transformations, the result follows immediately. Corollary 2: Show that $(\phi^{-1})^* \circ \phi^* = \text{id}$. Proof: Left to the reader Corollary 3: Let X, Y generate ϕ^* and ψ^* respectively, as its local 1-parameter group of transformations. Then $\phi^* \circ \psi^* = \psi^* \circ \phi^*$ if and only if $[X, Y] = 0$. Proof: Let ψ^* be a local 1-parameter group of transformations. Hence by 1.8 $(\psi^*)^* X = X$. Consequently from Theorem 3, $[X, Y] = 0$. Converse result follows from corollary 2.

A vector field X on a manifold M is said to be complete if

it induces a one parameter group of transformations on M . Theorem 4: Every vector field on a compact manifold M is complete. Proof: Let X be a given vector field on M . Then by Theorem 1, X induces ϕ^* as its

44 local 1-parameter group of transformations in a neighbourhood U of p of M and $t \in \mathbb{R}$. If p runs over M , then for each p , we get a neighbourhood $U(p)$ and $I(p)$, where all such $U(p)$ form an open covering of M . Since M is compact, every open covering $\{U(p)\}$ of M has a finite subcovering $\{U(p_1), \dots, U(p_n)\}$. If we let $\delta = \min\{\delta_1, \dots, \delta_n\}$ then, there is a δ such that for $|t| < \delta$, ϕ^* is a local 1-parameter group of transformations on M . We are left to prove that ϕ^* is defined on $\mathbb{R}M$. Case a) : $t > 0$. We write $t = k\delta + r$, $0 \leq r < \delta$, k being integer. Then $t = k\delta + r = (k-1)\delta + (\delta + r)$. Similarly for $t < 0$, we can show that ϕ^* is defined on $\mathbb{R}M$. Combining all the cases, we claim that ϕ^* is defined on $\mathbb{R}M$. Hence X induces ϕ^* as its 1-parameter group of transformations on a compact manifold M . Thus X is a complete vector field.

45 1.10 Cotangent Space: Note that $T_p(M)$ is a vector space over the field of real numbers. A mapping $\phi^* : T_p(M) \rightarrow T_p(M)$ that satisfies $\phi^*(X+Y) = \phi^*(X) + \phi^*(Y)$, $\phi^*(bX) = b\phi^*(X)$, $b \in \mathbb{R}$ and for all $X, Y \in T_p(M)$, is a linear mapping over \mathbb{R} . The linear mapping $\phi^* : T_p(M) \rightarrow T_p(M)$ denoted by $\phi^* : X \rightarrow \phi^*(X)$ is called a 1-form on M . Let $D_p(M)$ be the set of all 1-forms on M . Let us define 10.1) $D_p(M) = \{ \alpha \mid \alpha(X) = \alpha(\phi^*(X)) \}$. It can be shown that $D_p(M)$ is a vector space over \mathbb{R} , called the dual of $T_p(M)$. For every $p \in M$, $\phi^* : T_p(M) \rightarrow T_p(M)$ is a mapping $\phi^* : T_p(M) \rightarrow T_p(M)$ defined by 10.2) $\phi^*(X) = X \circ \phi$.

46 so that $\phi^* : T_p(M) \rightarrow T_p(M)$. Thus ϕ^* dual of $T_p(M)$. We write the dual of $T_p(M)$ by $T_p^*(M)$ and is the cotangent space of $T_p(M)$. Elements of $T_p^*(M)$ are called the covectors at p of M or linear functionals on $T_p(M)$. For every $f \in F(M)$, we denote the total differential of f by df and is defined as 10.3) $df(X) = Xf$. We also write it as 10.4) $(df)(X) = Xf$. Exercises: 1. Show that for every $f \in F(M)$, df is a 1-form on M . 2. If x_1, \dots, x_n are coordinate functions defined in a neighbourhood U of $p \in M$, show that each dx_i is a 1-form on U . Solution: 2 Note that $dx_i(X) = Xx_i$, (10.4) above $dx_i(X) = Xx_i = \delta_{ij}$, by (10.4) Similarly it can be shown that dx_i is a 1-form on U . Thus each dx_i is a 1-form on U . From Exercise 2 above, it is evident that each $\phi^*(dx_i)$, for $i=1, \dots, n$. We now define

47 10.5) $dx_i \circ \phi^* = dx_i$. Let $\phi^* : T_p(M) \rightarrow T_p(M)$ be such that 10.6) $\phi^*(X) = \sum_{j=1}^n \alpha_j dx_j$ where each $\alpha_j \in \mathbb{R}$. If possible, let $\phi^* : T_p(M) \rightarrow T_p(M)$ be such that $\phi^*(X) = \sum_{j=1}^n \alpha_j dx_j$ then $\phi^*(X) = \sum_{j=1}^n \alpha_j dx_j$. Proceeding in this manner we will find that $\phi^*(X) = \sum_{j=1}^n \alpha_j dx_j$. As dx_1, \dots, dx_n are linearly independent, we must have $\phi^*(X) = \sum_{j=1}^n \alpha_j dx_j$. Thus any $\phi^* : T_p(M) \rightarrow T_p(M)$ can be expressed uniquely as 10.7) $\phi^*(X) = \sum_{j=1}^n \alpha_j dx_j$. Finally if $\phi^*(X) = \sum_{j=1}^n \alpha_j dx_j$ then,

where each b_i is a 1-form, $\{b_1, \dots, b_n\}$ is a basis for $D^1(M)$ and $\{b_1, \dots, b_n\}$ is a basis for $D^1(M)$. Consequently any 1-form ω can be expressed as

53) Using (11.6) one gets after a few steps $\omega = \sum_{i=1}^n A_i b_i + \sum_{j=1}^n B_j b_j$. As $\{b_1, \dots, b_n\}$ is a basis, ω is linearly independent, so we must have $A_i = 0$ and $B_j = 0$. Consequently ω reduces to $\sum_{j=1}^n B_j b_j$.

1.12. Exterior Differentiation: The exterior derivative, denoted by d on D is defined as follows: i) $d(D^r) = D^{r+1}$ ii) for $f \in D^0$, df is the total differential iii) if $\omega \in D^r$, $\eta \in D^s$ then $d(\omega \wedge \eta) = d\omega \wedge \eta + (-1)^r \omega \wedge d\eta$ iv) $d^2 = 0$ From 11.7) of 1.11 we find that

54 Exercises: 1. Find the exterior differential of i) $xy \, dy - xy \, dx$ ii) $\cos(x) \, xy \, dx \, dz + x \, dy \, dz + y \, dz \, dx + z \, dx \, dy$ 2. Find the exterior differential of $d\omega$. A form ω is said to be closed if $d\omega = 0$. If ω is a r -form and $d\omega = 0$ for some $(r-1)$ form η then ω is said to be an exact form. Exercise: 3. Test whether ω is closed or not where

i) $\omega = F \, dx + G \, dy + H \, dz$ ii) $\omega = e^y \, dx + e^y \, dy + x \, \cos \, \sin$ Theorem 1: If ω is a 1-form, then $d\omega = 0$ iff ω is an exact form. Proof: Without any loss of generality, one may take an 1-form as $\omega = f \, dg + g \, df$, $f, g \in D^0(M)$. Using 11.5) of 1.11, we find $d\omega = (df \, g - f \, dg) + (df \, g + g \, df) = 2 \, df \, g$. Thus we get $d\omega = 0$ iff $df \, g = 0$.

55 Using (10.4) of 1.10, we get $d\omega = (df \, g - f \, dg) + (df \, g + g \, df) = 2 \, df \, g$. Thus we get $d\omega = 0$ iff $df \, g = 0$. Now $df \, g = 0$ iff $\omega = f \, dg + g \, df$ for some $f, g \in D^0(M)$. Thus we get from above $d\omega = 0$ iff $\omega = f \, dg + g \, df$.

56) Using 12.1 we get $d(f \, dx + g \, dy + h \, dz) = (df \, dx + dg \, dy + dh \, dz) + (-1)^0 (f \, dx + g \, dy + h \, dz) \wedge (dx + dy + dz)$. Using 12.1 we get $d(f \, dx + g \, dy + h \, dz) = (df \, dx + dg \, dy + dh \, dz) + (-1)^0 (f \, dx + g \, dy + h \, dz) \wedge (dx + dy + dz)$. Thus we get $d\omega = (df \, dx + dg \, dy + dh \, dz) + (-1)^0 (f \, dx + g \, dy + h \, dz) \wedge (dx + dy + dz)$.

Existence and Uniqueness of Exterior Differentiation: Without any loss of generality we may take an r -form as $\omega = f \, dx_1 \wedge \dots \wedge dx_r + \dots + g \, dx_1 \wedge \dots \wedge dx_r$. Let us define an R -linear map $d: D^r \rightarrow D^{r+1}$ as $d(f \, dx_1 \wedge \dots \wedge dx_r) = df \wedge dx_1 \wedge \dots \wedge dx_r + (-1)^i f \, dx_1 \wedge \dots \wedge dx_{i-1} \wedge dx_{i+1} \wedge \dots \wedge dx_r + \dots + (-1)^r f \, dx_1 \wedge \dots \wedge dx_{r-1} \wedge dx_{r+1}$. Clearly i) $d(D^r) = D^{r+1}$ and ii) if ω is a 0-form, then $d\omega$ is the total differential of ω . iii) Let $\omega \in D^r$ and it is enough to consider

$\omega = f \, dx_1 \wedge \dots \wedge dx_r + g \, dx_1 \wedge \dots \wedge dx_r$. Then $d\omega = df \wedge dx_1 \wedge \dots \wedge dx_r + (-1)^i f \, dx_1 \wedge \dots \wedge dx_{i-1} \wedge dx_{i+1} \wedge \dots \wedge dx_r + \dots + (-1)^r f \, dx_1 \wedge \dots \wedge dx_{r-1} \wedge dx_{r+1}$. Using 12.1 we get $d\omega = (df \, dx_1 + \dots + dg \, dx_r) \wedge dx_1 \wedge \dots \wedge dx_r + (-1)^i f \, dx_1 \wedge \dots \wedge dx_{i-1} \wedge dx_{i+1} \wedge \dots \wedge dx_r + \dots + (-1)^r f \, dx_1 \wedge \dots \wedge dx_{r-1} \wedge dx_{r+1}$.

Thus we get $d\omega = (df \, dx_1 + \dots + dg \, dx_r) \wedge dx_1 \wedge \dots \wedge dx_r + (-1)^i f \, dx_1 \wedge \dots \wedge dx_{i-1} \wedge dx_{i+1} \wedge \dots \wedge dx_r + \dots + (-1)^r f \, dx_1 \wedge \dots \wedge dx_{r-1} \wedge dx_{r+1}$. This completes the proof. Existence and Uniqueness of Exterior Differentiation: Without any loss of generality we may take an r -form as $\omega = f \, dx_1 \wedge \dots \wedge dx_r + \dots + g \, dx_1 \wedge \dots \wedge dx_r$. Let us define an R -linear map $d: D^r \rightarrow D^{r+1}$ as $d(f \, dx_1 \wedge \dots \wedge dx_r) = df \wedge dx_1 \wedge \dots \wedge dx_r + (-1)^i f \, dx_1 \wedge \dots \wedge dx_{i-1} \wedge dx_{i+1} \wedge \dots \wedge dx_r + \dots + (-1)^r f \, dx_1 \wedge \dots \wedge dx_{r-1} \wedge dx_{r+1}$. Clearly i) $d(D^r) = D^{r+1}$ and ii) if ω is a 0-form, then $d\omega$ is the total differential of ω . iii) Let $\omega \in D^r$ and it is enough to consider

$\omega = f \, dx_1 \wedge \dots \wedge dx_r + g \, dx_1 \wedge \dots \wedge dx_r$. Then $d\omega = df \wedge dx_1 \wedge \dots \wedge dx_r + (-1)^i f \, dx_1 \wedge \dots \wedge dx_{i-1} \wedge dx_{i+1} \wedge \dots \wedge dx_r + \dots + (-1)^r f \, dx_1 \wedge \dots \wedge dx_{r-1} \wedge dx_{r+1}$. Using 12.1 we get $d\omega = (df \, dx_1 + \dots + dg \, dx_r) \wedge dx_1 \wedge \dots \wedge dx_r + (-1)^i f \, dx_1 \wedge \dots \wedge dx_{i-1} \wedge dx_{i+1} \wedge \dots \wedge dx_r + \dots + (-1)^r f \, dx_1 \wedge \dots \wedge dx_{r-1} \wedge dx_{r+1}$.

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Thus we get $d\omega = (df \, dx_1 + \dots + dg \, dx_r) \wedge dx_1 \wedge \dots \wedge dx_r + (-1)^i f \, dx_1 \wedge \dots \wedge dx_{i-1} \wedge dx_{i+1} \wedge \dots \wedge dx_r + \dots + (-1)^r f \, dx_1 \wedge \dots \wedge dx_{r-1} \wedge dx_{r+1}$. This completes the proof. Existence and Uniqueness of Exterior Differentiation: Without any loss of generality we may take an r -form as $\omega = f \, dx_1 \wedge \dots \wedge dx_r + \dots + g \, dx_1 \wedge \dots \wedge dx_r$. Let us define an R -linear map $d: D^r \rightarrow D^{r+1}$ as $d(f \, dx_1 \wedge \dots \wedge dx_r) = df \wedge dx_1 \wedge \dots \wedge dx_r + (-1)^i f \, dx_1 \wedge \dots \wedge dx_{i-1} \wedge dx_{i+1} \wedge \dots \wedge dx_r + \dots + (-1)^r f \, dx_1 \wedge \dots \wedge dx_{r-1} \wedge dx_{r+1}$. Clearly i) $d(D^r) = D^{r+1}$ and ii) if ω is a 0-form, then $d\omega$ is the total differential of ω . iii) Let $\omega \in D^r$ and it is enough to consider

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Thus we get $d\omega = (df \, dx_1 + \dots + dg \, dx_r) \wedge dx_1 \wedge \dots \wedge dx_r + (-1)^i f \, dx_1 \wedge \dots \wedge dx_{i-1} \wedge dx_{i+1} \wedge \dots \wedge dx_r + \dots + (-1)^r f \, dx_1 \wedge \dots \wedge dx_{r-1} \wedge dx_{r+1}$. This completes the proof. Existence and Uniqueness of Exterior Differentiation: Without any loss of generality we may take an r -form as $\omega = f \, dx_1 \wedge \dots \wedge dx_r + \dots + g \, dx_1 \wedge \dots \wedge dx_r$. Let us define an R -linear map $d: D^r \rightarrow D^{r+1}$ as $d(f \, dx_1 \wedge \dots \wedge dx_r) = df \wedge dx_1 \wedge \dots \wedge dx_r + (-1)^i f \, dx_1 \wedge \dots \wedge dx_{i-1} \wedge dx_{i+1} \wedge \dots \wedge dx_r + \dots + (-1)^r f \, dx_1 \wedge \dots \wedge dx_{r-1} \wedge dx_{r+1}$. Clearly i) $d(D^r) = D^{r+1}$ and ii) if ω is a 0-form, then $d\omega$ is the total differential of ω . iii) Let $\omega \in D^r$ and it is enough to consider

$\omega = f \, dx_1 \wedge \dots \wedge dx_r + g \, dx_1 \wedge \dots \wedge dx_r$. Then $d\omega = df \wedge dx_1 \wedge \dots \wedge dx_r + (-1)^i f \, dx_1 \wedge \dots \wedge dx_{i-1} \wedge dx_{i+1} \wedge \dots \wedge dx_r + \dots + (-1)^r f \, dx_1 \wedge \dots \wedge dx_{r-1} \wedge dx_{r+1}$. Using 12.1 we get $d\omega = (df \, dx_1 + \dots + dg \, dx_r) \wedge dx_1 \wedge \dots \wedge dx_r + (-1)^i f \, dx_1 \wedge \dots \wedge dx_{i-1} \wedge dx_{i+1} \wedge \dots \wedge dx_r + \dots + (-1)^r f \, dx_1 \wedge \dots \wedge dx_{r-1} \wedge dx_{r+1}$.

57 be a differentiable mapping. Let $T_p(M)$ be the tangent space at p of M where $(\cdot) \cdot f_p T_N$ is its dual. Let $(\cdot) \cdot f_p T_N$ be the tangent space at $f(p)$ of N where $(\cdot) \cdot f_p T_N$ is its dual. If (x_1, \dots, x_n) and (y_1, \dots, y_m) are the local coordinate system at p of M and at $f(p)$ of N respectively, then, it is known that $\{dx_i : i=1, \dots, n\}$ and $\{dy_j : j=1, \dots, m\}$ are the basis of $T_p(M)$ and $(\cdot) \cdot f_p T_N$ respectively. Consequently $\{dx_i : i=1, \dots, n\}$ and $\{dy_j : j=1, \dots, m\}$ are the basis of $T_p(M)$ and $(\cdot) \cdot f_p T_N$ respectively. Let ω be a 1-form on N . We define an 1-form on M , called the pull-back 1 form of ω on M , denoted by $f^*\omega$, as follows 13.1) $f^*\omega = \omega \circ df_p$. So, we write 13.2) $f^*(\sum_{j=1}^m a_j dy_j) = \sum_{j=1}^m a_j (df_p)^{-1} dy_j$. Therefore we may write, for a 1 form ω on N and a vector field X on M by 13.4) $(f^*\omega)(X) = \omega(df_p(X))$. Theorem 1: If f is a mapping from an n -dimensional manifold M to an m -dimensional manifold N , where (x_1, \dots, x_n) is the local coordinate system in a neighbourhood of a point p of M and (y_1, \dots, y_m) is the local coordinate system in a neighbourhood of $f(p)$ of N , then $f^*dy_j = \sum_{i=1}^n \frac{\partial y_j}{\partial x_i} dx_i$ where $f_y f_j = \frac{\partial y_j}{\partial x_i}$. Proof: Since $(f^*\omega)_p$ is a co-vector at P on M , it can be expressed as the linear combination of the basis co-vectors dx_i at P and we take $f^*dy_j = \sum_{i=1}^n a_{ij} dx_i$ where a_{ij} are unknowns to be determined or $\{a_{ij}\}$. Using 10.5 of 1.10 we find that $f^*(\sum_{j=1}^m a_j dy_j) = \sum_{j=1}^m a_j (\sum_{i=1}^n a_{ij} dx_i)$. By (13.1), one reduces to $dy_j \circ df_p = \sum_{i=1}^n \frac{\partial y_j}{\partial x_i} dx_i$. Using (10.5) of (1.10) we find $f^*(\sum_{j=1}^m a_j dy_j) = \sum_{j=1}^m a_j (\sum_{i=1}^n \frac{\partial y_j}{\partial x_i} dx_i)$. Thus we get $f^*dy_j = \sum_{i=1}^n \frac{\partial y_j}{\partial x_i} dx_i$. Note: 1. Using (10.9) of 1.10, one find from above theorem 13.5) $f^*(\sum_{j=1}^m a_j dy_j) = \sum_{j=1}^m a_j (\sum_{i=1}^n \frac{\partial y_j}{\partial x_i} dx_i)$. We can also write it as 13.6) $f^*(\sum_{j=1}^m a_j dy_j) = \sum_{j=1}^m a_j (\sum_{i=1}^n \frac{\partial y_j}{\partial x_i} dx_i)$.

58 Where a_{ij} 's are unknowns to be determined or $\{a_{ij}\}$. Using 10.5 of 1.10 we find that $f^*(\sum_{j=1}^m a_j dy_j) = \sum_{j=1}^m a_j (\sum_{i=1}^n a_{ij} dx_i)$. By (13.1), one reduces to $dy_j \circ df_p = \sum_{i=1}^n \frac{\partial y_j}{\partial x_i} dx_i$. Using (10.5) of (1.10) we find $f^*(\sum_{j=1}^m a_j dy_j) = \sum_{j=1}^m a_j (\sum_{i=1}^n \frac{\partial y_j}{\partial x_i} dx_i)$. Thus we get $f^*dy_j = \sum_{i=1}^n \frac{\partial y_j}{\partial x_i} dx_i$. Note: 1. Using (10.9) of 1.10, one find from above theorem 13.5) $f^*(\sum_{j=1}^m a_j dy_j) = \sum_{j=1}^m a_j (\sum_{i=1}^n \frac{\partial y_j}{\partial x_i} dx_i)$. We can also write it as 13.6) $f^*(\sum_{j=1}^m a_j dy_j) = \sum_{j=1}^m a_j (\sum_{i=1}^n \frac{\partial y_j}{\partial x_i} dx_i)$.

59 2. If ω is a 1-form, then, its pull-back 1-form $f^*\omega$ is given by 13.7) $f^*\omega = \sum_{j=1}^m a_j (df_p)^{-1} dy_j$, where a_j are the components of ω . (Prove it.) Exercises: 1. If $f: M \rightarrow N$ be such that $f(x, y) = (\cos x, \sin x)$ where $x \in \mathbb{R}$. Compute $f^*(\cos y)$ and $f^*(\sin y)$. 2. If $f: M \rightarrow N$ be such that $f(x, y) = (\sin x, \cos x)$. Compute $f^*(\cos y)$ and $f^*(\sin y)$. 3. Let ω be the 1-form in \mathbb{R}^2 by $\omega = y dx + x dy$. Let U be the set in the plane (x, y) given by $U = \{(x, y) \in \mathbb{R}^2 : x^2 + y^2 = 1\}$. Let $f: U \rightarrow \mathbb{R}^2$ be the map $f(x, y) = (\cos x, \sin x)$. Compute $f^*\omega$. Let us now suppose that ω be a r -form on N . In the same manner, as defined earlier, we define an r -form on M , called the pull-back r -form on M , denoted by $f^*\omega$, as follows: 13.8) $f^*\omega = \omega \circ df_p$.

60 We also write it as 13.9) $f^*(\sum_{j_1, \dots, j_r} a_{j_1, \dots, j_r} dy_{j_1} \wedge \dots \wedge dy_{j_r}) = \sum_{j_1, \dots, j_r} a_{j_1, \dots, j_r} (df_p)^{-1} dy_{j_1} \wedge \dots \wedge dy_{j_r}$. Proposition: 1. Let $f: M \rightarrow N$ be a map, ω and η be r -forms on N and g be a 0-form on N . Then a) $f^*(\omega + \eta) = f^*\omega + f^*\eta$ b) $f^*(g\omega) = g f^*\omega$. Proof: a) As ω and η are r -forms on N , $(\omega + \eta)$ is also so. Hence

$$f^*(\omega + \eta) = (\omega + \eta) \circ df_p = \omega \circ df_p + \eta \circ df_p = f^*\omega + f^*\eta$$

$$f^*(g\omega) = (g\omega) \circ df_p = g \circ df_p \wedge \omega \circ df_p = g f^*\omega$$

Note that if ω is a r -form and g is a 0-form, then $g\omega$ is again a r -form. Using (13.8) one gets $f^*(g\omega) = g f^*\omega$.

61 or $f g f g p f f p f p^* ()^* () () () () ? ? d i d i ?$ or $f g f g p f p p^* * * () () () , ? ? b g b g b g ? ? p$ Hence $f g f g f^* * * () () () . ? ? ?$ Exercises : 4. Show that $f f f^* * * () ? ? ? ? ? ? ?$ 5. Prove that $() ()^* * * f h h f ? ? ? ? ?$ Note : From Theorem 1 of 1.11, we see that, any r -form ω can be expressed as $\sum g_i dx_1 \wedge \dots \wedge dx_n$ where g_i are differentiable functions on N . Then $f \omega = \sum (f g_i + \sum_j g_j dx_j \wedge dx_i) dx_1 \wedge \dots \wedge dx_n$ by the Proposition 1(b) and Exercise 4 above $\omega = \sum g_i dx_1 \wedge \dots \wedge dx_n$ Using 13.5) of 1.13 we see that $f \omega = \sum (f g_i + \sum_j g_j dx_j \wedge dx_i) dx_1 \wedge \dots \wedge dx_n$ Exercise : 7. Let M be a circle and \mathbb{R}^2 so that $f : M \rightarrow \mathbb{R}^2$ be defined by $x = \cos t, y = \sin t$

62 If $\omega = a dx + b dy$ and $\eta = c dx + d dy$, find $f^*(\omega \wedge \eta)$ Solution : In this case, $f^*(a dx + b dy) \wedge f^*(c dx + d dy) = (a dx + b dy) \wedge (c dx + d dy) = (ac - bd) dx \wedge dy$ Using Exercise 5, one finds that $f^*(\omega \wedge \eta) = (a dx + b dy) \wedge (c dx + d dy) = (ac - bd) dx \wedge dy$ where $d f^*(\omega \wedge \eta) = (a dx + b dy) \wedge (c dx + d dy) = (ac - bd) dx \wedge dy$

63 where the symbols have their usual meanings. Proof : We shall consider the following cases. i) ω is a 0-form ii) ω is a 1-form Case i) : In this case, let $\omega = h$, where h is a differentiable function Then $f^*(dh) = d(f^*h)$ by (10.4) of 1.10 = $d(h \circ f)$ by (7.3) of 1.7 = $d(h \circ f)$ by (10.4) of 1.10 = $d(f^*h)$ by (10.4) of 1.10 or $f^*(dh) = d(f^*h)$ The result is true

in this case. Case ii) : In this case, we assume that the result is true for 1-form. Without any loss of generality, we may take an r -form ω as $\sum g_i dx_1 \wedge \dots \wedge dx_r$

$f^*(\omega) = \sum g_i dx_1 \wedge \dots \wedge dx_r$ or $f^*(\omega) = \sum g_i dx_1 \wedge \dots \wedge dx_r$ Using (12.1) of 1.12 we find that $d(f^*\omega) = \sum (df g_i + \sum_j g_j dx_j \wedge dx_i) dx_1 \wedge \dots \wedge dx_r$

64 $f^*(d\omega) = \sum (d g_i + \sum_j g_j dx_j \wedge dx_i) dx_1 \wedge \dots \wedge dx_r$ Note that $df g_i$ is a 1-form and hence the theorem is true in this case. Thus $d f^*(\omega) = f^*(d\omega)$ by (12.1) of 1.12

Hence $d f^*(\omega) = f^*(d\omega)$ as the result is true for 1-form ω by known result Thus $d f^*(\omega) = f^*(d\omega)$ and hence the result is true for r -form also. Combining we claim that $d f^*(\omega) = f^*(d\omega)$ i.e. d and f^* commute each other..

REFERENCES 1. W.M.Boothby : An Introduction to Differentiable Manifolds and Riemannian Geometry. 2. Kobayashi & Nomizu : Foundations of Differentiable Geometry, Volume I 3. N. J. Hicks : Differentiable Manifold 4. Y. Matsushima : Differentiable Manifold 65 UNIT - 2 . 2.1 Lie group, Left translation, Right translation : Let G be a differentiable manifold. If G is a group and if the map $(g, h) \rightarrow gh$ from $G \times G$ to G and the map $g \rightarrow g^{-1}$ from G to G are both differentiable, then G is called a Lie group. Exmample : Let $GL(n, \mathbb{R})$ denote the set of all nonsingular $n \times n$ matrices over real numbers. $GL(n, \mathbb{R})$ is a group under matrix multiplication. Define $(g, h) \rightarrow gh$

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$GL(n, \mathbb{R})$ is a mapping of class C^∞ . Hence $GL(n, \mathbb{R})$ is a Lie group. Note : Lie groups are the fundamental building blocks for gauge theories. For every $a \in G$, a mapping $L_a : G \rightarrow G$ defined by $L_a(x) = ax$ is called a Left translation on G . Similarly, a mapping $R_a : G \rightarrow G$ defined by $R_a(x) = xa$ is called a right translation on G .

66 Note that $L L x L b x a b a ? ? ? ()$ and $L x a b ? ? ? L a L a = L a b R R x a b ? R b x x b a a () ?$ and $R x x b a a b ? ? ? R R R a b b a ? L R x L x b a x b a b a ? ? ? ()$ and $R L x R a x a b b a ? ? ? () ? L R R L a b b a ?$ Thus 2.3) $L L L a b a b ? , R R R a b b a ? , L R R L a b b a ?$ Again $L L x L a x b a x a b x L L x b a b a b ? ? ? ? ? ()$, Thus 2.4) $L L L L b a a b ?$, unless G is commutative Taking $b a ? ? 1$ in 2.3) we find $L L L a a a a ? ? ? 1 1$ by 2.3) $? L e$ Thus 2.5) $L L a a ? ? ? 1 1 ()$ It is evident that, for every $a ? G$, each $L a$ and $R a$ are diffeomorphism on G . Exercise : 1 Show that the set of all left (right) translation on G form a group. 2. Let $? : G \rightarrow G$ be a homeomorphism of a Lie group G to another Lie group G . Show that i) $? ? ? ? ? L L a a ? ()$ ii) $? ? ? ? L R b b ? () , ? a b$, in G .

67 3. Let $? be a 1-1 non-identity map from G to G. If ? ? ? ? L L g g ? is satisfied for all g ? G, then there is a h ? G such that ? ? R h . Solution : 2. From the definition of group homeomorphism of a Lie group G 1 to another Lie group G 2 , ? ? ? () () () , a b a b ? ? a b, in G 1 i) () () () () ? ? ? ? ? ? L x L x a x a x a a ? ? ? ? ? L x L x a a ? ? ? ? () () () , ? ? x in G 1 ? ? ? ? ? ? L L a a ? () Similarly ii) can be proved. 3. As G is a group, e ? G (identity). Further ? is a 1-1 map from G to G, so for e ? G, there is h in G such that ?(e) = h? Note that ?(e) = e because, ? is not an identity map. Now for g ? G, g e ? ? ? ? ? () () g e ? ? () L e g ? () () ? ? L e g ? () () , L e g ? ? as given ? L g e ? () b g = L g h = g h = R h g ? ? ? ? R g h ,$

68 . 2.2. Invariant Vector Field : We have already defined a vector field to be invariant under a transformation in 1.8. Note that, in a Lie group G , for every a, b in G , each $L a, R b$ is a transformation on G . Thus we can define invariant vector field under $L a, R b$.

A vector field X on a Lie group G is called a left invariant

vector field on G if 2.6) $() , * () L X X a p L p a ? ? p ? G$, where $() * L a$ is the differential of $L a$.

Thus from 1.7) $() * () () L X X a p L p L p a a d i ?$ We write it as 2.7) $() * L X X a ?$ Similarly for a right invariant vector field, write 2.8) $() * R X X a ?$ From 1.7) we know that $() () * L X g X g L a p p a d i ? ?$ or $() () * () L X g X g L a p L p p a a d i ? ?$ If $L p q a () ?$ then $p L q L q a q a a ? ? ? ? ? ? () 1 1 1$ Thus the above relation reduces to 2.9) $() () * L X g X g L a q a q a b g ? ? 1 ?$ Let g be

the set of all left invariant vector field on G . If $X, Y, ? g, a, b ? R$, then 2.10) $() () * L a X b Y p ? ? ? a L X b L Y p p () * * ? ? a X b Y , () * L p$ being linear explained in Unit 1. 2.11) $() [,] () , () * * * L X Y L X L Y p a p ? , see 1.7 = [X, Y]$

69 Thus $a X b Y g ? ?$ and $[,] . X Y g ?$ Consequently g is a vector space over R and also a Lie- algebra. The Lie algebra formed by

the set of all left invariant vector fields on G is called the Lie algebra of the Lie group G .

Note that every left invariant vector field is a vector field i.e. $g G ? ? ()$ where $? () G$ denotes the set of all vector field on G . The converse is not necessarily true. The converse will be true if a condition is satisfied by a vector field. The following theorem states such condition. Theorem 1 :

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A vector field X on a Lie group G is left invariant if			

and only if for every $f \in F(G) ? () 2.12) () () X f L X f L a a ? ? ?$

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Proof : Let X be a left invariant vector field on a Lie group G . Then for every $f \in F(G) ? ()$, we have from (2.6) $() * () L X f X f$			

Let \mathfrak{g} be the Lie algebra of a Lie group G . Conversely let (2.12) be true. Reversing the steps one gets the desired result. Note : i) The behaviour of a Lie group is determined largely by its behaviour in the neighbourhood of the identity element e of G . The behaviour can be represented by an algebraic structure on the tangent space of e , called the Lie algebra of the group. ii) Note that, two vector spaces U and V are said to be isomorphic, if a mapping $f : U \rightarrow V$ is i) linear and ii) has an inverse $f^{-1} : V \rightarrow U$. Theorem 2 : As a vector space, the Lie subalgebra \mathfrak{g} of the Lie group G is isomorphic to the tangent space $T_e(G)$ at the identity element $e \in G$.

70

Proof : Let us define a mapping $\phi : T_e(G) \rightarrow \mathfrak{g}$ by $\phi(X) = X_e$. Note that, for every $X, Y \in \mathfrak{g}$, $[X, Y]_e = [X_e, Y_e]$ by (2.12). Also for $b \in \mathbb{R}$, $[bX]_e = b[X]_e$ by (2.12). Thus ϕ is linear. We choose $X, Y \in \mathfrak{g}$ such that $\phi^{-1}(X_e) = X$, $\phi^{-1}(Y_e) = Y$. Then $[X, Y]_e = [X_e, Y_e]$ from above. $[X, Y]_e = [X_e, Y_e]$ by (2.3). $[X, Y]_e = [X_e, Y_e]$ as chosen or $[X, Y]_e = [X_e, Y_e]$ by (2.3). We define $\psi : \mathfrak{g} \rightarrow T_e(G)$ by $\psi(X) = X_e$. Then $\psi(X) = X_e$ where X_e is the identity differential on G . or $\psi(X) = X_e$. Further, $\psi([X, Y]) = [X_e, Y_e]$ by (2.12). Thus an inverse mapping exists and we claim that $\psi : \mathfrak{g} \rightarrow T_e(G)$ is an isomorphism. Exercises : 1. If X, Y are left invariant vector fields, show that $[X, Y]$ is also so. 2. If $c_{ijk} \in \mathbb{R}$ are structure constants on a Lie group G with respect to the basis X_1, \dots, X_n of \mathfrak{g} , show that i) $c_{ijk} = -c_{ikj}$ ii) $c_{ijk} + c_{kji} + c_{jki} = 0$. Solution : 1. From Q 1.7), we see that $[X, Y]_e = [X_e, Y_e]$ from the definition of Lie Bracket. $[X, Y]_e = [X_e, Y_e]$ by (2.7). Using (2.7), we see that $[X, Y]$ is a left invariant vector field. 2. Using problem 1 above, we see that every $[X, Y]_e = [X_e, Y_e]$ as $X, Y \in \mathfrak{g}$. Since X_1, \dots, X_n is a basis of \mathfrak{g} , every $[X, Y]_e$ can be expressed uniquely as $\sum c_{ijk} X_{ijk}$ where $c_{ijk} \in \mathbb{R}$. i) Note that if $[X, Y]_e = 0$. So, let $[X, Y]_e = 0$. Then from a known result, $[X, Y]_e = 0$. Using 1) we find that $c_{ijk} = -c_{ikj}$. As the set X_1, \dots, X_n is a basis of \mathfrak{g} and hence linearly independent, we must have $c_{ijk} = -c_{ikj}$. ii) Using Jacobi Identity, we find that $[X, [Y, Z]] + [Y, [Z, X]] + [Z, [X, Y]] = 0$. Hence from 1) $c_{ijk} + c_{kji} + c_{jki} = 0$. Again applying 1), we find that $c_{ijk} + c_{kji} + c_{jki} = 0$. As X_1, \dots, X_n is a basis and hence linearly independent, we must have $c_{ijk} + c_{kji} + c_{jki} = 0$.

71 iii) $[X, Y]_e = [X_e, Y_e]$ where X_e is the identity differential on G . or $[X, Y]_e = [X_e, Y_e]$. Further, $[X, Y]_e = [X_e, Y_e]$ by (2.12). Thus an inverse mapping exists and we claim that $\psi : \mathfrak{g} \rightarrow T_e(G)$ is an isomorphism. Exercises : 1. If X, Y are left invariant vector fields, show that $[X, Y]$ is also so. 2. If $c_{ijk} \in \mathbb{R}$ are structure constants on a Lie group G with respect to the basis X_1, \dots, X_n of \mathfrak{g} , show that i) $c_{ijk} = -c_{ikj}$ ii) $c_{ijk} + c_{kji} + c_{jki} = 0$. Solution : 1. From Q 1.7), we see that $[X, Y]_e = [X_e, Y_e]$ from the definition of Lie Bracket. $[X, Y]_e = [X_e, Y_e]$ by (2.7). Using (2.7), we see that $[X, Y]$ is a left invariant vector field. 2. Using problem 1 above, we see that every $[X, Y]_e = [X_e, Y_e]$ as $X, Y \in \mathfrak{g}$. Since X_1, \dots, X_n is a basis of \mathfrak{g} , every $[X, Y]_e$ can be expressed uniquely as $\sum c_{ijk} X_{ijk}$ where $c_{ijk} \in \mathbb{R}$. i) Note that if $[X, Y]_e = 0$. So, let $[X, Y]_e = 0$. Then from a known result, $[X, Y]_e = 0$. Using 1) we find that $c_{ijk} = -c_{ikj}$. As the set X_1, \dots, X_n is a basis of \mathfrak{g} and hence linearly independent, we must have $c_{ijk} = -c_{ikj}$. ii) Using Jacobi Identity, we find that $[X, [Y, Z]] + [Y, [Z, X]] + [Z, [X, Y]] = 0$. Hence from 1) $c_{ijk} + c_{kji} + c_{jki} = 0$. Again applying 1), we find that $c_{ijk} + c_{kji} + c_{jki} = 0$. As X_1, \dots, X_n is a basis and hence linearly independent, we must have $c_{ijk} + c_{kji} + c_{jki} = 0$.

72 ? () [,] [,] * L X Y X Y f a ? ? Using (2.7), we see that $[X, Y]$ is a left invariant vector field. 2. Using problem 1 above, we see that every $[X, Y]_e = [X_e, Y_e]$ as $X, Y \in \mathfrak{g}$. Since X_1, \dots, X_n is a basis of \mathfrak{g} , every $[X, Y]_e$ can be expressed uniquely as $\sum c_{ijk} X_{ijk}$ where $c_{ijk} \in \mathbb{R}$. i) Note that if $[X, Y]_e = 0$. So, let $[X, Y]_e = 0$. Then from a known result, $[X, Y]_e = 0$. Using 1) we find that $c_{ijk} = -c_{ikj}$. As the set X_1, \dots, X_n is a basis of \mathfrak{g} and hence linearly independent, we must have $c_{ijk} = -c_{ikj}$. ii) Using Jacobi Identity, we find that $[X, [Y, Z]] + [Y, [Z, X]] + [Z, [X, Y]] = 0$. Hence from 1) $c_{ijk} + c_{kji} + c_{jki} = 0$. Again applying 1), we find that $c_{ijk} + c_{kji} + c_{jki} = 0$. As X_1, \dots, X_n is a basis and hence linearly independent, we must have $c_{ijk} + c_{kji} + c_{jki} = 0$.

73 . 2.3 Invariant Differential Form : A differential form ω on a Lie group G is said to be left invariant if (2.13)

Let $L_p \omega$ be the pull-back of ω by L_p . We write it as (2.14) $L_p \omega = \omega_p$ and call $L_p \omega$, the pull-back differential form of ω . Similarly, a differential form ω on a Lie group G is said to be right invariant if (2.15)

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A differential form, which is both left and right invariant, is called a biinvariant differential form. Exercises : 1. If ω is a left invariant differential form, show that $d\omega$ is also so. 2. Prove that a differential 1-form ω on a Lie group is left invariant if and only if for every left invariant vector field X on G , $\omega(X)$ is a constant function on G . 3. Let ω be a left invariant 1-form on G . Show that ω is left invariant if and only if $\omega \circ \text{Ad}_g$ is right invariant. 4. Prove that the set of all left invariant forms on G is an algebra over \mathbb{R} . Such a set is denoted by \mathcal{A} , say. 5. If g^* denotes the dual space of g , then, prove that $\mathcal{A} \subset g^*$ where \mathcal{A} is the set already defined in Exercise 4 above. Solution : 1. From Q 1.13, we see that $L_{g^{-1}}^* \omega = \omega$ where $L_{g^{-1}}$ is the pull-back of ω . Using on (2.14) on the right hand side of the above equation, we see that $L_{g^{-1}}^* \omega = \omega$.

74 Consequently, $d\omega$ is a left invariant differential form. It can be proved easily that ω is a left invariant differential form. 2. Let us consider a differential 1-form ω . Then for every $a \in G$, $L_a^* \omega$ will be defined as the pull-back differential 1-form. Consequently from the definition of pull-back. $L_a^* \omega(X) = \omega(L_a X)$, $d(L_a^* \omega)(X) = L_a^* d\omega(X)$. Let us consider X to be left invariant. Then on using (2.6) on the right hand side of the above equation, we get $L_a^* d\omega(X) = d\omega(L_a X)$. Let us now consider ω to be left invariant 1-form. Then by (2.13), we get from $L_a^* \omega(X) = \omega(L_a X)$. Taking $a = e$, we see that $\omega(X) = \omega(L_e X)$. Consequently, $\omega(X)$ is a constant function on G . Conversely, if $\omega(X)$ is a constant function on G , then $L_a^* \omega(X) = \omega(L_a X)$. Hence $L_a^* \omega(X) = \omega(L_a X)$ which is (2.13) Thus ω is a left invariant differential form. This completes the proof.

75 Theorem 1 : If g is a Lie subalgebra of a Lie group G and g^* denotes the set of all left invariant form on G , then $dX \wedge Y - X \wedge dY = [X, Y]$ where $X, Y \in g$. Note : Such an equation is called Maurer-Carter Equation. Proof : From theorem 1 of 1.12, we know that $dX \wedge Y - X \wedge dY = [X, Y]$ for every vector field X, Y . If X, Y are in g then by Exercise 2, X, Y are constant functions on G . Hence by Exercise 2 of 1.4), $X, Y = 0$. Thus the above equation reduces to $dX \wedge Y - X \wedge dY = [X, Y]$. Exercise : 6. Show that $d(c_1 X + c_2 Y) = c_1 dX + c_2 dY$. Solution : If X_1, X_2, \dots, X_n is a basis of g and $\omega_1, \dots, \omega_n$ is the dual basis of g^* , then $dX_i = \sum_j c_{ij} \omega_j$. Hence from theorem 1 above $d(X_i \wedge X_j) = dX_i \wedge X_j - X_i \wedge dX_j = [X_i, X_j]$.

76 Again from 1.11 $d(X_i \wedge X_j) = dX_i \wedge X_j - X_i \wedge dX_j = [X_i, X_j]$ of Exercise 1 of 2.2 $d(X_i \wedge X_j) = dX_i \wedge X_j - X_i \wedge dX_j = [X_i, X_j]$. Thus $d(X_i \wedge X_j) = dX_i \wedge X_j - X_i \wedge dX_j = [X_i, X_j]$. Take $i, j, k = 1, 2, 3$, then $d(X_i \wedge X_j) = dX_i \wedge X_j - X_i \wedge dX_j = [X_i, X_j]$ as $d(X_i \wedge X_j) = dX_i \wedge X_j - X_i \wedge dX_j = [X_i, X_j]$.

77 Thus, we write $d(X_i \wedge X_j) = dX_i \wedge X_j - X_i \wedge dX_j = [X_i, X_j]$. Hence $d(X_i \wedge X_j) = dX_i \wedge X_j - X_i \wedge dX_j = [X_i, X_j]$. 2.4 Automorphism : A mapping, denoted by α for every $a \in G$, $\alpha : G \rightarrow G$ defined by $\alpha(x) = axa^{-1}$ is said to be an inner automorphism if i) α is injective ii) α is surjective such α is written as Ad_a . Exercise : Show that if G is a Lie group, $h \in G$, then the map $\alpha : G \rightarrow G$ defined by $\alpha(x) = hxh^{-1}$ is an automorphism. An inner automorphism of a Lie group G is defined by 2.16) $\alpha(x) = axa^{-1}$. Now, $\alpha(x) = axa^{-1}$. Using 2.3) we get

78 2.17) $\alpha(L_a X) = L_a X$. Note that α is a diffeomorphism. Theorem 1 : Every inner automorphism of a Lie group G induces an automorphism of the Lie algebra g of G . Proof : For every $a \in G$ let us denote the inner automorphism on G by $\alpha(x) = axa^{-1}$. Now for every $X \in g$ and from 1.7 such $\alpha : G \rightarrow G$ induces a differential mapping $(\alpha)_* : g \rightarrow g$. Such a mapping is a linear mapping and by Theorem 2 of 2.2, the Lie subalgebra g of a Lie group G is such that $g = T_e(G)$. Thus to show every α induces an automorphism of the Lie algebra g of G we are to show ii) $(\alpha)_*$ is a mapping from g to g iii) $(\alpha)_*$ is a homomorphism i.e. $(\alpha)_*(X \wedge Y) = (\alpha)_* X \wedge (\alpha)_* Y$ iv) $(\alpha)_*$ is injective v) $(\alpha)_*$ is surjective ii) Let $Y \in g$. Then on using 2.17) we get $(\alpha)_*(Y) = Y$. Thus iii) $(\alpha)_*(X \wedge Y) = (\alpha)_* X \wedge (\alpha)_* Y = X \wedge Y$. iv) Let $Y \in g$. Then on using 2.17) we get $(\alpha)_*(Y) = Y$. Thus v) $(\alpha)_*(X \wedge Y) = (\alpha)_* X \wedge (\alpha)_* Y = X \wedge Y$. Thus $(\alpha)_*$ is an automorphism of g .

$X^i Y^j Z^k + Z^l$ ii) $Y^m (Y^p + Z^q) X^r X^s$??? iii) X^t
 $f^u X^v Y^w$
 $f^x Y^y Z^z$ iv) $X^i X^j (f^k Y^l) (X^m f^n Y^o) , , , , X^p, Y^q,$
 Z^r
 $M), F($
 $M) ? f$ The vector field ?

$X^i Y^j$ is called the covariant derivative of Y in the direction of X
 with respect to the connection If P is a tensor field of type (o, s) we define $v^i X^j P = X^j P^i$, if $s = 0$ vi) $X^i X^j P^k = X^i X^j P^k + P^k X^i X^j$
 Exercise 1: Let $M = R^n$ and $X, Y, Z \in \mathfrak{X}(M)$ be such that $Y = b^i X^i$
 where $b^i = 1$ for $i = 1, \dots, n$ Show that Z determines a linear connection on M .

93 Solution: Let $X = a^i X^i, Z = c^i X^i$ with $a^i, c^i \in C^\infty(M)$, $i = 1, \dots, n$ Then i) $X^j Y^k = X^j (b^i X^i) = X^j b^i X^i + b^i X^j X^i$
 as defined $X^j X^i = X^i X^j + X^i X^j$ Similarly it can be shown that $(Y+Z) X^i = Y X^i + Z X^i$
 Again, $X^j (f^i X^i) = X^j f^i X^i + f^i X^j X^i$ as $(X^j f^i) X^i + f^i X^j X^i$
 $f^i X^j X^i = Y X^i + Z X^i$ and $X^j X^i = X^i X^j + X^i X^j$ as $X^j (X^i) = X^j X^i + X^i X^j$
 $f^i X^j X^i = Y X^i + Z X^i$

Thus Z determines a linear connection on M . Let (x^1, \dots, x^n) be a system of co-ordinates in a neighbourhood U of p of M . We define $\Gamma^k_{ij} = \Gamma^k_{ij}$ where $\Gamma^k_{ij} \in C^\infty(M)$ Such are called the christoffel symbols or the connection coefficients or the components

of the connection. Hence if $X = X^i X^i, Y = Y^j X^j$ where each $X^i, Y^j \in C^\infty(M)$, $i, j = 1, \dots, n$, we see that $X^j Y^k = X^j (Y^i X^i) = X^j Y^i X^i + Y^i X^j X^i$
 $Y^i X^j X^i = Y^i X^i X^j + X^i X^j Y^i$ by iii) $X^j Y^k = X^j Y^i X^i + Y^i X^i X^j + X^i X^j Y^i$
 Exercise 2: Let Γ^k_{ij} be the connection co-efficients of the linear connection with respect to the local coordinate system (x^1, \dots, x^n) and (y^1, \dots, y^n) respectively. Show that in the intersection of the two coordinate neighbourhoods $\Gamma^k_{ij} = \Gamma^k_{ij}$ Solution: In the intersection of the two coordinates $x^i = x^i(y^j)$ or $y^j = y^j(x^i)$
 Again, from 3.1) we see that $\Gamma^k_{ij} = \Gamma^k_{ij}$ from above

95 $\Gamma^k_{ij} = \Gamma^k_{ij}$
 $X^j Y^k = X^j (Y^i X^i) = X^j Y^i X^i + Y^i X^j X^i$
 $Y^i X^j X^i = Y^i X^i X^j + X^i X^j Y^i$ from above
 $X^j Y^k = X^j Y^i X^i + Y^i X^i X^j + X^i X^j Y^i$
 $\Gamma^k_{ij} = \Gamma^k_{ij}$ by 3.1) $\Gamma^k_{ij} = \Gamma^k_{ij}$
 $X^j Y^k = X^j Y^i X^i + Y^i X^i X^j + X^i X^j Y^i$

Changing s, r, t $\Gamma^k_{ij} = \Gamma^k_{ij}$
 $X^j Y^k = X^j Y^i X^i + Y^i X^i X^j + X^i X^j Y^i$
 $\Gamma^k_{ij} = \Gamma^k_{ij}$ from above $\Gamma^k_{ij} = \Gamma^k_{ij}$
 $X^j Y^k = X^j Y^i X^i + Y^i X^i X^j + X^i X^j Y^i$ Since $\Gamma^k_{ij} = \Gamma^k_{ij}$
 $\Gamma^k_{ij} = \Gamma^k_{ij}$ $R^k_{st} = R^k_{st}$ is a basis of the tangent space and hence linearly independent and the result follows immediately. 3.2 Torsion tensor field and curvature tensor field on a linear connection we define a mapping $T: \mathfrak{X}(M) \times \mathfrak{X}(M) \rightarrow \mathfrak{X}(M)$
 $T(X, Y) = X(Y) - Y(X) - \Gamma^k_{ij} X^i Y^j + \Gamma^k_{ji} Y^i X^j$ and another $R: \mathfrak{X}(M) \times \mathfrak{X}(M) \times \mathfrak{X}(M) \rightarrow \mathfrak{X}(M)$
 $R(X, Y, Z) = X(YZ) - Y(XZ) - \Gamma^k_{ij} X^i Y^j Z^k + \Gamma^k_{ji} Y^i X^j Z^k - X^i Y^j Z^k \Gamma^k_{ij} + X^i Y^j Z^k \Gamma^k_{ji}$

96 3.3) $R^k_{st} = R^k_{st}$
 $X^j Y^k = X^j Y^i X^i + Y^i X^i X^j + X^i X^j Y^i$
 Then T is a tensor field of type $(1, 2)$
 and is called the torsion tensor field and R is a tensor field of type $(1, 3)$, called the curvature tensor field of M . A linear connection is said to be symmetric if 3.4) $T(X, Y) = 0$ In such case 3.5)
 $X^j Y^k = X^j Y^i X^i + Y^i X^i X^j + X^i X^j Y^i$ Exercise: 1. Verify that i) $T(X, Y) = -T(Y, X)$; ii) $T(X, Z) = T(Y, Z) + X(YZ) - Y(XZ) - \Gamma^k_{ij} X^i Y^j Z^k + \Gamma^k_{ji} Y^i X^j Z^k$
 $Z(XY) - Y(XZ) - \Gamma^k_{ij} X^i Y^j Z^k + \Gamma^k_{ji} Y^i X^j Z^k$; iii) $T(fX, gY) = fg T(X, Y)$. 2. If $T(X, Y) = 0$
 $X^j Y^k = X^j Y^i X^i + Y^i X^i X^j + X^i X^j Y^i$ show that Γ^k_{ij} is a linear connection and $T(X, Y) = 0$ 3. Show that i) $T(X, Y, Z) = X(YZ) - Y(XZ) - \Gamma^k_{ij} X^i Y^j Z^k + \Gamma^k_{ji} Y^i X^j Z^k - X^i Y^j Z^k \Gamma^k_{ij} + X^i Y^j Z^k \Gamma^k_{ji}$
 $X^j Y^k = X^j Y^i X^i + Y^i X^i X^j + X^i X^j Y^i$ ii) $R(X, Y, Z) = X(YZ) - Y(XZ) - \Gamma^k_{ij} X^i Y^j Z^k + \Gamma^k_{ji} Y^i X^j Z^k - X^i Y^j Z^k \Gamma^k_{ij} + X^i Y^j Z^k \Gamma^k_{ji}$
 $X^j Y^k = X^j Y^i X^i + Y^i X^i X^j + X^i X^j Y^i$ iii) $R(X, Y, Z) = X(YZ) - Y(XZ) - \Gamma^k_{ij} X^i Y^j Z^k + \Gamma^k_{ji} Y^i X^j Z^k - X^i Y^j Z^k \Gamma^k_{ij} + X^i Y^j Z^k \Gamma^k_{ji}$

$b g c h d i d i c h \dots iv)$

R
X, $fY Z R fY, Y Z R X, Y fZ f R X, Y$
Z

$b g b g b g b g \dots$ Hence Show that $R fX, gY hZ fgh R X, Y Z b g b g \dots$ 4. Exercise 3 : Prove Ricci Identity a) for a 1-form $w : \dots F H I K \dots$

$X Y Y X X, Y Z W R X, Y Z \dots b g c h$

97 b) for a 2-form $W : \dots X Y W Y X W X, Y W Z, P W R X, Y Z, P W Z, R X, Y$

P
 $e j b g b g c h b g c h$ 5. If x^1, \dots, x^n is a local coordinate system and $T x^i, x^j, R x^i, x^j, y x^i x^j, i j k k i j k i j k h h \dots$
 $\dots F H G I K J \dots F H G I K J \dots$ Show that i) T and $i j k j i k j i k j k \dots$ for a symmetric linear connection ii) $R x^i x^j m k i j m k j i m k j m t i k i m t j t k \dots$ Solution : 1 i) From the definition

$Y X$
 $T($
 $Y,$
 $X)$

$X Y [Y, X] \dots Y X X Y [X, Y] \dots X Y Y X [X,$
 $Y] \dots$

$T X Y (,)$

Thus T

is skew-symmetric ii) $fX gY Z T($

$fX gY, Z) Z (fX gY) [fX gY,$
 $Z] \dots X Y Z Z f Z g Z (fX) (gY) [fX, Z] [gY, Z] \dots g Y Z Z g Y [,] ()$

98 $\dots X Z Y Z f Z X [X, Z] g Z Y [Y, Z] \dots f T X Z g T Y Z (,) (,)$

Again,

using the definition, given in 3.1 and also from 1.5 we get Thus T is a bilinear mapping. 2. To prove that ∇ is a linear connection, we have to prove i), ii), iii), iv) of 3.1. Now

X
 $X (Y Z) (Y Z) T(X, Y Z) \dots$ as defined $X X Y Z T(X, Y) T(X, Z) \dots X X Y Z, \dots$

as defined similarly, other results can be proved and hence ∇ is a linear connection. Now,,
 $X Y T(X, Y) Y X [X, Y], \dots$ by definition $X Y Y T(X, Y) X T(Y, X) [X, Y], \dots$ as defined $\dots T X Y T X Y$

$T X$
 $Y (,) (,) (,)$ by Ex 1 (i) above $\dots T X Y (,) \dots T T \dots$

99 3. (iv) From the definition $R(X, fY) Z X fY fY X [$
 $X, fY] Z Z \dots X Y Y X f[X, Y] (XfY) (fZ) fZ Z \dots Y X Y Y X [X, Y] Y (Xf) Z f Z f Z f Z (Xf) Z \dots$
 $\dots X Y Y X [X, Y] f Z Z \dots fR(X, Y) Z \dots$

by definition. 5. From the given condition $i j i j i j x x T, \dots x x x x x \dots$
 \dots Using 3.1) we find $k k i j j i k k 0 x x \dots$ or, $\dots k k k j i j j i k k T, x x \dots$ as defined Since $k : k 1,$
 $, n x \dots$ is a basis and hence linearly independent and thus i) $k k k i j j i j \dots$ If the linear connection is
symmetric, then $T = 0$. consequently, the above equation reduces to $k k i j j i \dots$

100 ii) From the definition, we see that $i j j i i j j m m m,$

x
 $x x x x R, x x x x x \dots$

$i j k k j m i m k k x x x x \dots$ as $i j, 0 x x \dots$

k
 $k t k k j m j m i k i m i k t j k t x x x x x \dots$

Changing the dummy indices $t k, k t$ in the 2nd and 4th term we get $k k t k k t k i j m j m i t j m i m j t k i k j k k R x x x x$
 $x x x \dots$ Since $k : k 1, n x \dots$ is a basis and hence linearly
independent, we get from above $k k k t k t k i j m j m i m j m i t j i j R x x \dots$ 3.2 Covariant Differential

of
A Tensor Field of type (o, s) The covariant differential of a tensor field of type $(0,$

s)
 is a tensor field of type $(0, s + 1)$ and is defined as $\sum_{i_1, \dots, i_s} S_{i_1, \dots, i_s}(X) P_{i_1, \dots, i_s}(X)$
 $X_{i_1} \dots X_{i_s}$? ? ? ? ? ? ? ?
 Exercise : 1 Let i^j be the components of a vector field Y with respect to a local coordinate system $\{x^1, \dots, x^n\}$ i.e. $Y = i^j \frac{\partial}{\partial x^j}$
 If i^j be the components of the covariant differential $\nabla_X Y$ so that $\nabla_X Y = x^k \frac{\partial i^j}{\partial x^k} \frac{\partial}{\partial x^j} + \Gamma_{ik}^j i^k \frac{\partial}{\partial x^j}$ then, show that $\nabla_X Y = i^j \frac{\partial}{\partial x^j} + \Gamma_{ik}^j i^k \frac{\partial}{\partial x^j}$
 2. Let ω be a 1 form and $l = dx^i$ If we write $\omega = i^k x^k$? ? ? ? ? ? ? ? ? ?
 $101 \nabla_X \omega = i^k \frac{\partial}{\partial x^i} dx^k = \frac{\partial i^k}{\partial x^i} dx^k + i^k \frac{\partial dx^k}{\partial x^i} = \frac{\partial i^k}{\partial x^i} dx^k$
 102
 103 UNIT - 4 ?
 104 Theorem 1 : Every Riemannian manifold (M, g) admits a unique Riemannian Connection. Proof : To prove the existence of such a connection, let us define a mapping $\nabla : (M) \times (M) \rightarrow (M)$ denoted by $X \nabla Y = (X \cdot Y) - Y(X)$ as follows 4.3) Clearly, $X \nabla Y = X(Y) - Y(X)$
 $X \nabla (Y \cdot Z) = X(Y \cdot Z) - Y \cdot X(Z) = X(Y) \cdot Z + Y \cdot X(Z) - Y \cdot X(Z) = X(Y) \cdot Z$
 $Y \cdot X(Z) = (Y \cdot X) \cdot Z + X \cdot Y(Z) = g(X, Y) \cdot Z + X \cdot Y(Z)$
 $g(X, [Y, Z]) = X \cdot g(Y, Z) - Y \cdot g(X, Z) + g(X, [Y, Z]) = X \cdot g(Y, Z) - Y \cdot g(X, Z) + g(X, [Y, Z])$
 $g(Y, [X, Z]) = Y \cdot g(X, Z) - X \cdot g(Y, Z) + g(Y, [X, Z]) = Y \cdot g(X, Z) - X \cdot g(Y, Z) + g(Y, [X, Z])$
 $g(X, [Y, Z]) - g(Y, [X, Z]) = X \cdot g(Y, Z) - Y \cdot g(X, Z) - Y \cdot g(X, Z) + X \cdot g(Y, Z) = 2(X \cdot g(Y, Z) - Y \cdot g(X, Z))$
 is linear Whence $X \nabla (Y \cdot Z) = X(Y) \cdot Z + Y \cdot X(Z)$
 Similarly it can be shown that $X \nabla Y = X(Y) - Y(X)$
 $Z \cdot X = X \cdot Z + [X, Z]$
 $X \nabla (fY) = (Xf)Y + f \nabla_X Y$
 Thus such a mapping determines a linear connection on M . Also, from (4.3) it can be shown that
 $X \nabla (Y \cdot Z) = X(Y) \cdot Z + Y \cdot X(Z)$
 $X \nabla (Y \cdot Z) = X(Y) \cdot Z + Y \cdot X(Z)$
 $X \nabla (X \cdot X) = X(X) \cdot X + X \cdot X(X) = 2X(X) \cdot X = 2X(X) \cdot X$
 $X \nabla (X \cdot X) = 2X(X) \cdot X = 2X(X) \cdot X$
 by v) of 3.1 or, $X \nabla (X \cdot X) = 2X(X) \cdot X$
 Thus such a linear connection admits a metric connection. Further, it can be shown that $X \nabla Y = X(Y) - Y(X)$ Hence such a metric connection admits a Riemannian connection To prove the uniqueness, let ∇ be another such connection. Then we must have
 $X \nabla Y = X(Y) - Y(X)$
 $X \nabla (Y \cdot Z) = X(Y) \cdot Z + Y \cdot X(Z)$
 $g(Y, Z) = 0$? ? ? ? ? and
 $X \nabla Y = X(Y) - Y(X) = 0$? ? ? ? ? $X \nabla (Y \cdot Z) = X(Y) \cdot Z + Y \cdot X(Z) = 0$? ? ? ? ? and $X \nabla Y = X(Y) - Y(X) = 0$? ? ? ? ? Subtracting, $X \nabla (Y \cdot Z) = X(Y) \cdot Z + Y \cdot X(Z) = 0$
 $X \nabla (Y \cdot Z) = X(Y) \cdot Z + Y \cdot X(Z) = 0$
 $Z \cdot X = X \cdot Z + [X, Z]$
 and $X \nabla Y = X(Y) - Y(X)$
 $X \nabla Y = X(Y) - Y(X)$
 $Y \cdot X = X \cdot Y + [X, Y]$
 $X \nabla Y = X(Y) - Y(X) = 0$? ? ? ? ? where form, we get $X \nabla Y = 0$? ? ? ? ? $X \nabla X = Y$
 $Y \cdot X = X \cdot Y + [X, Y]$
 $Y \cdot X = X \cdot Y + [X, Y]$
 Thus uniqueness is established. This completes the proof
 Exercise : 1 In terms of a local coordinate system $\{x^1, \dots, x^n\}$ in a neighbourhood U of p of a Riemannian Manifold (M, g) show that i) the components Γ_{ij}^k defined in UNIT 3 is symmetric and ii) the Riemannian metric is covariantly constant. 2.
 Let ∇ be a metric connection of a Riemannian manifold (M, g) and $\tilde{\nabla}$ be another linear connection given by $X \nabla Y = T(X, Y) + \nabla_X Y$ where T is the torsion tensor of M . Show that the following conditions are equivalent i) $\nabla g = 0$ and ii) $\nabla(T(X, Y), Z) = 0$

$Z, Y) g(Z, X) \dots 3.4.3$

Einstein Manifold : Let $\{e_1, e_2, \dots, e_n\}$ be an orthonormal basis of $T_p(M)$. Then the Ricci tensor field, denoted by S , is the covariant tensor field of degree 2 and is defined by $S(X, Y) = R(e_i, Y)(e_i)$. We write it as 4.12 $S(X, Y) = R(e_i, Y)(e_i)$. Such a tensor field $S(X, Y)$ is also called the Ricci Curvature of M . If there is a constant r such that 4.13 $S(X, Y) = r g(X, Y)$ then M is called an Einstein Manifold. The function r on M , defined by $r = \frac{1}{n-1} \text{tr}(S)$ is called the scalar curvature of M . We write it as 4.14 $r = \frac{1}{n-1} \text{tr}(S)$. Exercise : 1. Show that the Ricci tensor field is symmetric. At any $p \in M$, π we denote by π a plane section i.e., a two dimensional subspace of $T_p(M)$. The sectional curvature of π denoted by $K(\pi)$ with orthonormal basis X, Y is defined as 4.15 $K(\pi) = g(R(X, Y)Y, X) = R(X, Y, Y, X)$. If $K(\pi)$ is constant for all plane section and for all points of $p \in M$, 114 Then (M, g) is called a manifold of constant curvature. For such a manifold 4.16)

$R(X, Y)Z = k\{g(Y, Z)X - g(X, Z)Y\}$ where $k(\cdot)$ say Example : Euclidean space is of Constant Curvature

Exercise : 1, Show that a Riemannian manifold of constant curvature is an Einstein Manifold. 2.

If M is a 3-dimensional Einstein Manifold, then, it is a manifold of constant curvature

Solution : Let $\{X_1, X_2, X_3\}$ be an orthonormal basis of $T_p(M)$. Then, the sectional curvature with orthonormal basis X_1, X_2 denoted by $K(X_1, X_2)$ is given by 12 $K(X_1, X_2) = R(X_1, X_2, X_2, X_1)$

$R(X_1, X_2, X_2, X_1) = 2K(X_1, X_2) = 2K(X_1, X_2)$. Thus, $R(X_1, X_2, X_2, X_1) = 2K(X_1, X_2)$. Again from 4.12 $R(X_1, X_2, X_2, X_1) = R(X_1, X_2, X_1, X_2) + R(X_1, X_2, X_1, X_2) + R(X_1, X_2, X_1, X_2) + R(X_1, X_2, X_1, X_2)$

$R(X_1, X_2, X_1, X_2) = R(X_1, X_2, X_1, X_2) + R(X_1, X_2, X_1, X_2) + R(X_1, X_2, X_1, X_2) + R(X_1, X_2, X_1, X_2)$

$R(X_1, X_2, X_1, X_2) = 0$

$R(X_1, X_2, X_1, X_2) = 0$

$R(X_1, X_2, X_1, X_2) = 0$

$K(X_1, X_2) = 0$. As it is a 3-dimensional Einstein manifold, so from 4.13 $S(X_1, X_2) = r g(X_1, X_2)$

$S(X_1, X_2) = r g(X_1, X_2) = r g(X_1, X_2)$

$S(X_1, X_2) = r g(X_1, X_2) = r g(X_1, X_2)$

115

116

using

the above result

in 4.8)

we get

$X_1 Y_1 Y_1 \{T(X_1, Y_1)(X_1)Y_1 - g(X_1, Y_1)X_1\} = 0$. Again using 4.17, one gets $X_1 Y_1 Y_1 (Y_1)X_1 = g(X_1, Y_1)X_1$

Exercise 1. If ∇

and $\tilde{\nabla}$ correspond to a semi-symmetric connection and the Levi-Civita connection respectively, then for any 1-form ω $X(X(\omega)) - \omega(X(X)) = (\tilde{\nabla}_X \omega)(X) - \omega(\tilde{\nabla}_X X)$ where $\tilde{\nabla}_X \omega = X(\omega) - \omega(X)$. Let $\tilde{\nabla}$ be the Levi-Civita Connection and ∇ be another linear connection such that $X(X(Y)) - \nabla_X(Y) = \omega(X)$ where ω is a 1-form. Show that ∇ is a semi-symmetric connection for which $X(g(X, Y)) = 0$. Hints : 1. Note that $X(X(Y)) - \nabla_X(Y) = \omega(X)$. Use Theorem 1 in the second term on the right hand side, one gets the desired result. 2. Note

that

$X(Y) - \nabla_X(Y) = \omega(X)$

X, Y

$Y(Y) - \nabla_Y(Y) = \omega(Y)$

$Y(X) - \nabla_Y(X) = \omega(X)$

$X(X)$

Y, X

on using the hypothesis $(Y)X = 0$

$X(Y) - \nabla_X(Y) = \omega(X)$. Again, $X(X) - \nabla_X(X) = \omega(X)$. $X(g(Y, Z)) = X(g(Y, Z)) - \nabla_X(g(Y, Z)) = 0$

$R(X, Y)Z - (AX, Y)Z - (X, AY)Z = (X, Y)Z$ $n^2(n-1)(n-2) \dots$ i.e. $1^r R(X, Y) = (AX, Y) - (X, AY)$
 $n^2(n-1)(n-2) \dots$ Theorem 3: If in a conformally flat manifold, for a symmetric linear transformation A , $R(X, Y)A = A \cdot R(X, Y)$ then $2^r A \cdot X \cdot X \cdot 0 \dots$ Proof: Note that $R(X, Y) = -R(Y, X)$ As A is symmetric, so by Exercise 1 of this article $A \cdot R(X, Y) = R(X, Y) \cdot A$ is skew-symmetric. Thus $R(Z, W)A$ is a skew-symmetric linear transformation and from 4.24) we can write
 $g((R(Z, W)A)X, X) = -g(X, (R(Z, W)A)X)$ or $g(R(Z, W)AX, X) = -g(X, R(Z, W)AX) = -g(R(Z, W)AX, X)$, as g is symmetric. $g(R(Z, W)AX, X) = 0$ Using 4.7) one gets $g(R(AX, X)Z, W) = 0$
 Whence $R(AX, X)Z = 0$ i.e., $R(AX, X) = 0$ Again $(AX, AX)Z = 0$ i.e., $AX, AX = 0$ for every Z . Using Theorem 2, one gets $2^r R(X, AX) = (AX, AX) - (X, AX) = 0$
 122 AS $R(AX, X) = -R(X, AX)$ and $R(AX, X) = 0$, we get from above, $2^r X \cdot A \cdot X \cdot X \cdot AX = 0$ Note that $X \cdot Y$ is skew-symmetric and thus $2^r A \cdot X \cdot X \cdot AX = 0$
 Definition: A curve $x(t)$, $a \cdot t \cdot b$ is called a geodesic on M with a linear connection ∇ if 4.25) $X \cdot X = 0$ Where X is the vector tangent to the integral curve γ at $x(t)$. Note that the integral curves of a left invariant vector fields are geodesic.
 4.7 Biinvariant Riemannian metric on a Lie group: A Riemannian metric g on a Lie group

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is said to be biinvariant if it is both left and right			

invariants. Exercise 1: If g is a left invariant covariant tensor field of order 2 on G and X, Y are left invariant vector fields on G , show that $g(X, Y)$ is a constant function. Theorem 1: If G is a Lie group admitting a biinvariant Riemannian metric g , then 4.26)
 $g([X, Y], Z) = g(X, [Y, Z]) - g(Y, [X, Z])$ 4.27) $1^r R(X, Y)Z = [[X, Y], Z] - [X, [Y, Z]] - [Y, [X, Z]]$ 4.28) $1^r g(R(X, Y)Z, W) = g([X, Y], [Z, W])$
 Proof: Since X, Y are left invariant vector fields, $X + Y$ is also so and hence from 4.25) $(X + Y) \cdot (X + Y) = 0$
 123 Using 4.25, we find from above i) $X \cdot Y \cdot X = 0$ since M admits a unique Riemannian connection, we must have
 $X \cdot Y = X \cdot [X, Y] - [X, Y] \cdot X$ or $X \cdot Y = [X, Y] - [X, Y] \cdot X$ from i) Now for a Riemannian Manifold $(g)(X, Z) = 0$ or, $Y \cdot Y \cdot g(X, Z) = g(X, Y)$
 $g(X, Z) = 0$
 Using Exercise 1 of this article and Exercise 2 of 1.4 we see that $Y \cdot g(X, Z) = 0$ Thus from ii) we find that $1^r g([Y, X], Z) = g(X, [Y, Z]) - g([X, Y], Z)$ or, $g([X, Y], Z) = g(X, [Y, Z]) - g([X, Y], Z)$
 Again from the definition $Z \cdot X \cdot Y = X \cdot [X, Y] - [X, Y] \cdot X$ or $1^r X \cdot [X, Y] = [X, Y] \cdot X$ or $1^r X \cdot [X, Y] = [X, Y] \cdot X$ by using ii) $1^r X \cdot [X, Y] = [X, Y] \cdot X$
 $1^r X \cdot [X, Y] = [X, Y] \cdot X$ or $1^r X \cdot [X, Y] = [X, Y] \cdot X$ by Jacobi Identity $1^r [X, Y], Z = [X, Y], Z$

$Z^4 \cdot 2 \cdot \dots \cdot 1 \cdot [X, Y], Z^4 \cdot \dots$

Again $\dots \cdot 1 \cdot R(X, Y)Z, W) \cdot g[X, Y], Z, W^4 \cdot \dots$ by 4.27) $\dots \cdot 1 \cdot g[X, Y], Z, [Z,$

$W]^4 \cdot \dots$ by 4.26)

This completes the proof.

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Theorem 2 : If G is a Lie group admitting a biinvariant Riemannian metric g and σ is a plane section in $p \in T(M)$ where σ is determined by orthonormal left invariant vector fields X, Y at p on G , then the sectional curvature at p is zero if and only if $[X, Y] = 0$. Proof : From 4.15) $K(\sigma) = g(K(\sigma))$

$R(X, Y)Y, X) \cdot \dots \cdot 1 \cdot g([X, Y], [Y, X]) \cdot \dots$ by 4.28) $1 \cdot g([X, Y], [X, Y]) \cdot \dots$

The result follows immediately as g is nonsingular. Theorem 3 : If G is a Lie group admitting a biinvariant Riemannian metric g , then for all left invariant vector fields, X, Y, Z, W, P . Proof : From Jacobi's identity $[W, [X, Y]] + [X, [Y, Z]] + [Y, [Z, W]] = 0$ Taking $P = [X, Y]$, we get $[W, [X, Y], Z] + [[X, Y], [Z, W]] + [Z, [W, [X, Y]]] = 0$ or $[W, [[X, Y], Z]] - [[X, Y], [W, Z]] = [[W, [X, Y]], Z] = [-[X, [Y, W]] - [Y, [W, X]], Z]$ by Jacobi Identity i) $[W, [[X, Y], Z]] - [[X, Y], [W, Z]] = [[X, [W, Y]], Z] + [[W, X], Y]$

$Z]$ Again from the definition $W \cdot W \cdot W \cdot W \cdot (R)(P, Z, X, Y) \cdot R(P, Z, X, Y) \cdot R(P, Z, X, Y) \cdot R(P, Z, X, Y) \cdot \dots \cdot W \cdot W \cdot R(P, Z, X, Y) \cdot R(P, Z, X, Y) \cdot \dots$

$P) \cdot \dots$

$W \cdot W \cdot W \cdot 0 \cdot R(X, Y, Z, P) \cdot R(X, Y, Z, P) \cdot R(X, Y, Z, P) \cdot \dots \cdot W \cdot P(X, Y, Z, P) \cdot \dots$

$P) \cdot \dots$

125

REFERENCES 1. W. B. Boothby : An Introduction to differentiable Manifold and Riemannian Geometry. Using 4.28), one gets $\dots \cdot W \cdot 1 \cdot 1 \cdot (R)(P,$

$Z, X, Y) \cdot g[X, Y], Z, [W, P] \cdot g[W, Z], P, [X, Y] \cdot \dots \cdot 1 \cdot g[W, X]Y[Z, P] \cdot g[X, W, Y], [Z, P] \cdot \dots$ Using 4.26) successively we get $\dots \cdot 1 \cdot g[X, Y], Z, W, P \cdot g[X, Y], [W, Z], P \cdot \dots \cdot g[W,$

$X], Z, P] \cdot g[X, W, Y], P) \cdot \dots \cdot 1 \cdot g[W, [X, Y], Z, P] \cdot g[X, Y], [W, Z], P \cdot \dots \cdot 1 \cdot g[X, W, Y], Z], P] \cdot g[[W, X], Y], Z,$

$P \cdot \dots = 0$

by i) for all left invariant vector fields X, Z, Y, W, P . This completes the proof.

126 NOTES

127 NOTES

128 NOTES

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$a_0/1, a_1 a_0 + 1/a_1, a_2(a_1 a_0 + 1) + a_0/a_2 a_1 + 1,$

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A vector field X on a Lie group G is left invariant if

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6/10 SUBMITTED TEXT 37 WORDS **46% MATCHING TEXT** 37 WORDS

Proof : Let X be a left invariant vector field on a Lie group G . Then for every $f \in F(G)$, we have from (2.6) $(L_X f)^* = L_X f$

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7/10 SUBMITTED TEXT 12 WORDS **100% MATCHING TEXT** 12 WORDS

on a Lie group G is said to be left invariant if (2.13)

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10/10	SUBMITTED TEXT	12 WORDS	100%	MATCHING TEXT	12 WORDS
<p>is said to be biinvariant if it is both left and right</p> <p>SA MA4K9 Project.pdf (D27014346)</p>					