### THEORY OF ODE

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ABSTRACT. We shall learn some theory related to ordinary differential equations.

# Part 1. Preliminaries

#### 1. Preliminaries from Real Analysis

**Definition 1** (Pointwise convergence). Let I be any interval in  $\mathbb{R}$ . Let  $f_n: I \to \mathbb{R}, n = 1, 2, \ldots$ , be a sequence of functions. We say that the  $\{f_n\}$  converges pointwise to a function  $f: I \subset \mathbb{R} \to \mathbb{R}$  if the sequence  $\{f_n(x)\}$  converges to f(x) for every  $x \in I$ .

**Remark 1.** Uniform convergence preserves continuity, interchange of limit and integral.

**Theorem 1.** Uniform limit of the sequence of continuous function is continuous.

Remark 2. Converse of the above theorem is not always true.

**Theorem 2** (Cauchy Criterion). Let  $\{f_n\}_{n\geq 1}$  a sequence of function defined on a metric space  $(X, d_X)$  with values in a complete metric space  $(Y, d_Y)$ . Then there exists a function  $f: X \to Y$  such that

$$f_n \to f$$
 uniformly on X

if and only if the following condition is satisfied: For every  $\varepsilon > 0$ , there exists an integer  $n_0$  such that

$$m, n \ge n_0 \text{ implies } d_Y(f_m(x), f_n(x)) < \varepsilon$$

for every  $x \in X$ .

*Proof.* ( $\Longrightarrow$ ) Assume that the sequence  $\{f_n\}$  converges uniformly on X. For given  $\varepsilon > 0$  there exists  $n_0 \in \mathbb{N}$  such that  $\forall m > n_0$  and for all  $x \in X$ 

$$d_Y(f_m, f) \le \frac{\varepsilon}{2}$$

and there exists  $n_1 \in \mathbb{N}$  such that  $\forall n \geq n_1$ 

$$d_Y(f_n, f) \le \frac{\varepsilon}{2}$$

Take  $n_3 = \min\{n_0, n_1\}$  then for all  $m, n \ge n_3$ 

$$d_Y(f_m, f_n) \le d_Y(f_m, f) + d_Y(f, f_n)$$

$$\le \frac{\varepsilon}{2} + \frac{\varepsilon}{2}$$

$$\le \varepsilon \text{ for all } x \in X$$

( $\iff$ ) Conversely, suppose that  $m, n \geq n_0$  implies that  $d_Y(f_m, f_n) < \varepsilon$  for all  $x \in X$ . Then for each  $x \in X$ , the sequence  $\{f_n(x)\}_{n\geq 1}$  is cauchy in a complete space Y and therefore converges in Y. Let  $f(x) = \lim_{n \to \infty} f_n(x)$  for each  $x \in X$ . For k > 0 then

$$d_Y(f_n(x), f_{n+k}(x)) < \frac{\varepsilon}{2}$$

for every  $k = 0, 1, \ldots$  and every  $x \in X$ .

$$d_Y(f_n(x), f(x)) = \lim_{n \to \infty} d_Y(f_n(x), f_{n+k}(x)) \le \frac{\varepsilon}{2} < \varepsilon$$
$$d_Y(f_n(x), f(x)) < \varepsilon$$

for every  $x \in X$  and  $f_n \to f$  uniformly over X.

**Theorem 3** (Weierstrass M-test). Let  $f_1, f_2, ...$  be a sequence of real valued functions defined on a set X and suppose that

$$|f_n(x)| \le M_n$$

for all  $x \in X$  and all n = 1, 2, ... If  $\sum_{n=1}^{\infty} M_n$  converges, then  $\sum_{n=1}^{\infty} f_n$  converges uniformly.

*Proof.* If  $\sum_{n=1}^{\infty} M_n$  converges, then for given  $\varepsilon > 0$ ,

$$\left| \sum_{k=n}^{m} f_k(x) \right| \le \sum_{k=n}^{m} |f_k(x)| \le \sum_{k=n}^{m} M_k < \varepsilon$$

for every  $x \in X$  and provided that m, n are sufficiently large. Now, by cauchy criterion, uniform convergence follows.

**Remark 3.** For Uniform boundedness and equicontinuity, the underline space is assumed to be compact hence equicontinuity here is same as uniform equicontinuity.

**Remark 4.** Some authors define pointwise and uniform equicontinuity separately and then on compact space they prove the equivalence.

**Definition 2** (Uniform Boundedness). A sequence of functions  $\{f_n\}$  defined on I(compact) is said to be uniformly bounded if there exists a constant M > 0 such that  $|f_k(x)| \leq M$ , for all  $x \in I$ , for all  $k \in \mathbb{N}$ .

**Definition 3** (Equicontinuity). A sequence of functions  $\{f_k\}$  defined on I(compact) is said to be equicontinuity on I, if for every  $\varepsilon > 0$ , there exists  $\delta > 0$  such that  $|f_k(x) - f_k(y)| < \varepsilon$  whenever  $x, y \in I$  and  $|x - y| < \delta$ , and for all k.

**Remark 5.** If the family of functions is equicontinuous, then each member in the family is uniformly continuous. Converse may not be true always. For e.g.  $f_k(x) = x^k, 0 \le x \le 1, k = 1, 2, ...$  is not an equicontinuous family; however each member is uniformly continuous function (continuous on closed and bounded interval).

**Example 1.** Each finite set of functions defined on compact set is equicontinuous. If the set is singleton then its trivial. For two element set, take minimum of the  $\delta$ 's and then for finite set take minimum of all  $\delta$ 's.

**Theorem 4.** For the sequence of functions  $f_n$  and f defined on compact set E and  $f_n \to f$  uniformly then the family of functions  $A = \{f_n : n \in \mathbb{N}\}$  is equicontinuous.

*Proof.* Since  $f_n \to f$  uniformly, by cauchy criterion of uniform continuity  $\exists n_0 \in \mathbb{N}$  such that  $\forall m, n \geq n_0$ 

$$|f_m(x) - f_n(x)| < \frac{\varepsilon}{3}$$

Look at the family of functions  $B = \{f_1, f_2, \dots, f_{n_0}\}$ , since it is finite, it is equicontinuous. For  $k \in \mathbb{N}$ 

$$|f_k(x) - f_k(y)| < \frac{\varepsilon}{3}$$

whenever  $|x-y|<\delta$  for all  $f_k\in B$  and for all  $x,y\in E$ . Now, for each  $f_n\in A$ 

$$|f_n(x) - f_n(y)| \le |f_n(x) - f_{n_0}(x)| + |f_{n_0}(x) - f_{n_0}(y)| + |f_{n_0}(y) - f_n(y)|$$

$$< \frac{\varepsilon}{3} + \frac{\varepsilon}{3} + \frac{\varepsilon}{3}$$

$$< \varepsilon \text{ whenever } |x - y| < \delta$$

**Theorem 5** (Arzela-Ascoli). Let  $\{f_k\}$  be a sequence of functions in C[a,b] which is uniformly bounded and equicontinuous. Then, there exists a subsequence  $\{f_{k_n}\}$  of  $\{f_k\}$  such that  $\{f_{k_n}\}$  converges uniformly to a function  $f \in C[a,b]$ .

Proof. TODO

**Definition 4** (Lipschitz continuity). A function  $f: D \subset \mathbb{R} \to \mathbb{R}$ , is said to be *locally* Lipschitz in D if for any  $x_0 \in D$ , there exists a neighbourhood  $N_{x_0}$  of  $x_0$  and an  $\alpha = \alpha(x_0) > 0$  such that

$$|f(x)-f(y)| \leq \alpha(x_0)|x-y|$$
, for all  $x,y \in N_{x_0}$ 

The function  $f: D \subset \mathbb{R} \to \mathbb{R}$  is said to be *globally* Lipschitz in D if there exists  $\alpha > 0$  such that

$$|f(x) - f(y)| \le \alpha |x - y|$$
, for all  $x, y \in D$ 

**Remark 6.** The smallest  $\alpha$  satisfying the above equation is called *Lipschitz constant* of f. It should be finite.

**Remark 7.** If f is globally Lipschitz, then it is uniformly continuous. (Take  $\delta = \frac{\varepsilon}{\alpha}$ ).

**Theorem 6** (Sufficient condition for Lipschitz continuity). Suppose D is an open interval in  $\mathbb{R}$  and  $f: D \to \mathbb{R}$  is differentiable on D and  $\alpha = \sup_{x \in D} |f'(x)| < \infty$ . Then, f is Lipschitz with a Lipschitz constant less than or equal to  $\alpha$ .

*Proof.* Use mean value theorem:

$$\frac{|f(b) - f(a)|}{|b - a|} = f'(c) \le \alpha$$

And Lipschitz constant is by definition the smallest  $\alpha$  satisfying above inequality. Therefore, Lipschitz constant  $\leq \alpha$ .

**Example 2.** Example of Lipschitz continuous functions: polynomials, polynomials of sine and cosine, exponential functions etc.

**Definition 5** (Lipschitz continuity for vector valued map). A function  $\mathbf{f}(t, \mathbf{y}) : (a, b) \times D \to \mathbb{R}^n$  is said to be Lipschitz continuous (globally) with respect to  $\mathbf{y}$  if there exists  $\alpha > 0$  such that

$$|\mathbf{f}(t, \mathbf{y_1}) - \mathbf{f}(t, \mathbf{y_1})| \le \alpha |\mathbf{y_1} - \mathbf{y_2}|$$

for all  $(t, \mathbf{y_1})$  and  $(t, \mathbf{y_2})$  in  $(a, b) \times D$ .  $\alpha$  should be finite.

**Theorem 7** (Sufficient condition for Lipschitz continuity of  $\mathbf{f}(t, \mathbf{y})$ ). Let  $\mathbf{f}: (a, b) \times D \to \mathbb{R}^n$  be a  $C^1$  vector valued function, where D is a convex domain in  $\mathbb{R}^n$  such that

$$\sup_{(t,\mathbf{y})\in(a,b)\times D} \left| \frac{\partial f_i}{\partial y_j}(t,\mathbf{y}) \right| = \alpha < \infty,$$

for i, j = 1, 2, ..., n. Then,  $\mathbf{f}(t, \mathbf{y})$  is Lipschitz continuous on  $(a, b) \times D$  with respect to  $\mathbf{y}$  having a Lipschitz constant less than or equal to a multiple of  $\alpha$ .

Remark 8 (Convex Domain). A set is called a convex domain if it contains all the line segments between any two points of the set.

**Remark 9.** Lipschitz continuity is a smoothness property stronger than continuity, but weaker than differentiability, locally. There can be functions which are not differentiable but still they can be lipschitz continuous.

**Theorem 8** (Calculus Lemma). Let (a, b) be a finite or infinite interval and  $h: (a, b) \to \mathbb{R}$  satisfy either

- (i) h is bounded above and non-decreasing or
- (ii) h is bounded below and non-increasing,

then,  $\lim_{t\to b} h(t)$  exists.

*Proof.* Let  $\alpha = \sup_{t \in (a,b)} h(t)$ . For some  $\varepsilon > 0$ ,  $\alpha - \varepsilon$  is not a supremum.

Hence there exists  $t_0$  such that  $\alpha - \varepsilon < h(t_0) < \alpha$ . For some  $t \in (a, b)$  such that  $t \ge t_0$  then

$$\alpha - \varepsilon < h(t_0) \le h(t) < \alpha$$

Then  $\alpha - h(t) < \varepsilon$ . Also above can be written as  $\alpha < h(t_0) + \varepsilon \le h(t) + \varepsilon < \alpha + \varepsilon$ . Then  $h(t) + \varepsilon < \alpha + \varepsilon \implies h(t) < \alpha + \varepsilon \implies h(t) - \alpha < \varepsilon$ . Therefore

$$|h(t) - \alpha| < \varepsilon \ \forall t \ge t_0 \implies \lim_{t \to b} h(t) = \alpha$$

**Theorem 9** (Change of Variable Formula). Let  $g: [c, d] \to \mathbb{R}$  be a  $C^1$  function and let [a, b] be any interval containing the image of g, that is  $g[c, d] \subset g[a, b]$ . If  $f: [a, b] \to \mathbb{R}$  is a continuous function, then

$$\int_{c}^{d} f(g(t))g'(t)dt = \int_{g(c)}^{g(d)} f(x)dx.$$

**Theorem 10** (Generalized Leibnitz Formula). Let  $\alpha, \beta \colon [a,b] \to \mathbb{R}$  be differentiable functions and c,d be real numbers satisfying

$$c \le \alpha(t), \beta(t) \le d$$
, for all  $t \in [a, b]$ .

Let  $f:[a,b]\times[c,d]\to\mathbb{R}$  be a continuous function such that  $\frac{\partial f}{\partial t}(t,s)$  is also continuous. Define

$$F(t) = \int_{\alpha(t)}^{\beta(t)} f(t, s) ds.$$

Then, F is differentiable and

$$\frac{dF}{dt} = \int_{\alpha(t)}^{\beta(t)} \frac{\partial f}{\partial t}(t, s) ds + f(t, \beta(t)) \frac{d\beta}{dt} - f(t, \alpha(t)) \frac{d\alpha}{dt}.$$

**Definition 6** (Banach space). A normed linear space which is a complete metric space (metric space by the norm) is called a Banach space.

**Example 3** (Examples of Banach space).  $\mathbb{R}^n$  under supremum norm and p-norm. Function space C[a,b] with supremum norm. However, it is *not* a complete space with repect to  $\|\cdot\|_1$  norm.

**Remark 10.** A sequence  $\{f_n\} \subset C[a,b]$ ,  $f_n \to f$  in supremum norm is equivalent to saying that  $f_n$  converges uniformly to f ( $\exists \delta > 0$  which works in supremum norm for f, will work for all f in the sequence.)

**Theorem 11** (Banach Fixed Point theorem). Suppose (X, d) is a complete metric space and  $T: X \to X$  is a contraction, that is, there exists an  $\alpha \in (0,1)$  such that

$$d(Tx, Ty) \le \alpha d(x, y)$$

for all  $x, y \in X$ . Then, T has a unique fixed point  $x^* \in X$ . Further, the sequence  $\{x_k\}$  defined by  $x_k = Tx_{k-1}, x_0 \in X$  is arbitrary and  $k = 1, 2, \ldots$ , converges to  $x^*$ .

**Remark 11.** If we omit the contraction condition then f(x) = x + 1 show there does not exists any fixed point. If completion condition is droped then  $f: (0,1) \to (0,1)$  defined by f(x) = mx, 0 < m < 1 is a contraction with  $\alpha = m$  but no fixed point.

*Proof.* Choose any  $x_0 \in X$  and define the sequence  $x_1 = Tx_0, x_2 = T^2x_0, \dots, x_k = T^kx_0, \dots$ . We need to show that sequence  $\{x_k\}$  is a cauchy sequence. To see this

$$d(x_{n+1}, x_n) = d(Tx_n, Tx_{n-1}) \le \alpha d(x_n, x_{n-1})$$

By induction,

$$d(x_{n+1}, x_n) \le \alpha^n d(x_1, x_0)$$

Consider, for m < n and using triangle inequality

$$d(x_n, x_m) \le d(x_n, x_{n-1}) + \dots + d(x_{m+1}, x_m) \le \alpha^m \frac{1 - \alpha^{n-m-1}}{1 - \alpha} d(x_1, x_0)$$

and the rhs  $\to 0$  as  $n, m \to \infty$ . Hence the sequence is cauchy and by completeness, there exists  $x^* \in X$  such that  $x_k \to x^*$ . By continuity of T, we get  $Tx_k \to Tx^*$ . By continuity of T, we get  $Tx_k \to Tx^*$ . But  $Tx_k = x_{k+1} \to x^*$ . Thus,  $Tx^* = x^*$ . If there exists  $y^*$  with same properties as  $x^*$  then by uniqueness of limits  $x^* = y^*$ .

Remark 12. Here, fixed point can be constructed with any desired accuracy as limiting value is fixed point. And secondly, any point can be taken as an initial guess.

**Corollary 1.** Let  $T: X \to X$  where X is a complete metric space, be such that  $T^k$  is a contraction for some  $k \geq 1$ . Then, T has a unique fixed point.

Proof. Since,  $T^k$  is contraction map hence it has a unique fixed point (by Banach fixed point theorem). Let  $x^*$  be that point i.e.  $T^kx^*=x^*$ . Applying T, we get  $T^k(Tx^*)=Tx^*$  implies  $Tx^*$  is the fixed point for  $T^k \implies Tx^*=x^*$ . Therefore,  $x^*$  is the fixed point of T. For uniqueness, assume  $x_1$  is another fixed point for  $T \implies Tx_1 = x_1$ . Applying T repeatedly, we get  $T^kx_1 = x_1$  but fixed point for  $T^k$  is  $x^*(unique) \implies x^* = x_1$ .

#### 2. Preliminaries from Linear Algebra

**Definition 7** (Normed Linear Space). A norm, denoted by  $\|\cdot\|$  on a vector space or a linear space X is a mapping from  $X \to \mathbb{R}$  that satisfies

- $||x|| \ge 0$ ; ||x|| = 0 if and only if x = 0,
- ||ax|| = |a|||x||,
- (Triangle inequality)  $||x + y|| \le ||x|| + ||y||$ ,

for all  $x, y \in X$  and scalar a.

Example 4.  $\mathbb{R}^n$ ,  $\mathbb{C}^n$ 

**Remark 13.** Every normed linear space is a metric space with metric induced by the norm. A metric space can be equipped with different norms which are fundamentally different. C[0,1] with sup norm is complete while with the integral norm  $\left(\int_0^1 |f(x)| \, dx\right)$  is not complete.

**Definition 8** (Vector Field). A vector field on a space(most commonly Euclidean space) is a function  $\vec{F}$  that assigns a vector to each point of the space.

**Definition 9** (Matrix norm). |A| should satisfy the following criterion

- $|A| \ge 0$ ; |A| = 0 if and only if A = 0,
- $\bullet ||aA| = |a||A|,$
- (Triangle inequality)  $|A + B| \le |A| + |B|$ ,
- $\bullet |AB| \leq |A||B|$

for all  $A, B \in M_n(\mathbb{R})$  and scalars a.

**Remark 14.**  $M_n(\mathbb{R})$  is a complete metric space.

2.1. Matrix Exponential  $e^A$ . Let  $A \in \mathbb{M}_n\mathbb{R}$ , define the sequence of matrices

$$S_k = I + A + \frac{A^2}{2!} + \dots + \frac{A^k}{k!}.$$

For k > l,

$$|S_k - S_l| \le \sum_{j=l+1}^k \frac{|A|^j}{j!} \to 0 \text{ as } l, k \to \infty.$$

Thus  $\{S_k\}$  is a Cauchy sequence and converges to some  $S \in M_n(\mathbb{R})$ .

**Definition 10.** Given  $A \in M_n(\mathbb{R})$ ,  $e^A$  is defined as

$$e^A = S$$

where  $S = \lim_{k \to \infty} \sum_{j=0}^{k} \frac{A^j}{j!}$ .

**Remark 15.**  $e^A \in M_n(\mathbb{R})$ . Also,  $|e^A| \leq e^{|A|}$  (substitute the values to see this.)

**Remark 16.** For diagonal matrix  $A = diag(\lambda_1, ..., \lambda_n) \implies A^j = diag(\lambda_1^j, ..., \lambda_n^j)$  and  $e^A = \sum_{j=0}^{\infty} \frac{A^j}{j!}$  then

$$e^A = \sum_{j=0}^{\infty} \frac{1}{j!} diag(\lambda_1^j, \dots, \lambda_n^j) \implies e^A = diag(\sum_{j=0}^{\infty} \frac{\lambda_1^j}{j!}, \dots, \sum_{j=0}^{\infty} \frac{\lambda_n^j}{j!}),$$

therefore,

$$e^A = diag(e^{\lambda_1}, \dots, e^{\lambda_n}).$$

Here are important observations:

- (1) If  $A \sim B \Longrightarrow \exists$  non-singular matrix P such that  $B = PAP^{-1}$ . Then  $e^A \sim e^B$  (:  $B^j = PA^jP^{-1} \Longrightarrow e^B = P\left(\sum_{j=0}^{\infty} \frac{A^j}{j!}\right)P^{-1} \Longrightarrow e^B = Pe^AP^{-1}$ ).
- (2) If A represents a block diagonal matrix  $A = diag(A_1, \ldots, A_k)$  where all A's are square matrix on the diagonal (their size maybe different) then

$$e^{A} = diag(e^{A_1}, e^{A_2}, \dots, e^{A_k}).$$

**Remark 17.** (a)  $|A| \le \max\{|A_1|, \dots, |A_k|\} \implies$ 

- (b)  $|e^A| \leq \max\{e^{|A_1|}, \dots, e^{|A_k|}\}$ . ((a)suspicious identity to me!!)
- (c) Equality in above cases holds for Euclidean norm.

Computing  $e^A$ . If A is diagonalizable then we look for matrix P so that  $PAP^{-1} = B$  is diagonal matrix hence computing  $e^A$  is easy.

If A is not diagonalizable then we look for P so that  $PAP^{-1}$  is a block diagonal, with easily computable  $e^{A_i}$ .

If T is a linear transformation and it is invariant on all coordinate axes then A is diagonalizable with respect to standard basis.

If usual coordinate axes are invariant under T then we look for n distinct directions, if possible, which are invariant under T and then take these directions as new bases. The matrix with respect to this basis will be diagonal matrix.

**Remark 18.** The set of all eigenvalues of A is known as *spectrum* of A. It is denoted by  $\sigma(A)$ .

The eigenvalues of A are the roots of the *characteristics polynomials*  $\det(\lambda I - A)$  which is a real polynomial in  $\lambda$  of degree n. The roots maybe real or complex. If the eigenvalue is real then there exists a corresponding real eigenvector and if it complex then there exists a corresponding complex eigenvector.

#### 3. Doubts

**Ques 1.** How to prove theorem 7? Is there a mean value theorem for vector valued functions that I can apply here?

Ques 2. Proof for change of varible formula in Theorem 9?

Ques 2. Proof for generalized Leibnitz formula in Theorem 10?

### Part 2.

## First and Second order linear equations

- 4. First Order equations
- 4.1. **General form of IVP and BVP.** General form of IVP of first order ODE is

$$\begin{cases}
\dot{y} = f(t, y) \\
y(t_0) = y_0
\end{cases}$$

General form of IVP of second order ODE is

$$\begin{cases}
\ddot{y} = f(t, y, \dot{y}) \\
y(t_0) = y_0 \\
\dot{y}(t_0) = y_1
\end{cases}$$

General form of BVP

$$\begin{cases}
\ddot{y} = f(t, y, \dot{y}) \text{ for } t \in (a, b) \\
\alpha_1 y(a) + \beta_1 \dot{y}(a) = \gamma_1 \\
\alpha_2 y(b) + \beta_2 \dot{y}(b) = \gamma_2
\end{cases}$$

**Definition 11** (Solution of ODE). Let f be defined on the rectangle in  $R := (a, b) \times (c, d)$  containing the initial data  $(t_0, y_0)$ . A solution to the IVP is a function  $y : (\bar{a}, \bar{b}) \to \mathbb{R}$  which is differentiable and satisfies the IVP together with initial condition.

**Remark 19.** Interval  $(\bar{a}, b)$  is referred as interval of existence of solution. For all  $t \in (\bar{a}, \bar{b}) \subset (a, b), y(t) \in (c, d)$ . If  $(\bar{a}, \bar{b}) = (a, b)$  then y is called the global solution otherwise it is a local solution.

For vector valued functions, definition of solution can be extended. Let  $\mathbf{f}:(a,b)\times\Omega\to\mathbb{R}^n$  be a vector valued continuous function so that  $\mathbf{f}=(f_1,\ldots,f_n)$  and each  $f_i$  is a real valued continuous (we assume this throughout the notes) function, where  $\Omega\subset\mathbb{R}^n$  is open domain. For a given initial value  $y_0\in\Omega$ , the IVP is given by

$$\left. egin{aligned} \dot{\mathbf{y}} &= \mathbf{f}(t,\mathbf{y}) \\ \mathbf{y}(t_0) &= \mathbf{y}_0 \end{aligned} 
ight. 
ight.$$

4.2. First Order linear equations. A general first order ODE can be written as

$$f(t,y,\dot{y}) = h(t)$$

where h(t) is the function of t only. From linear algebra it can be shown that f takes the form (TODO)

$$f(t, y, \dot{y}) = p_0(t)\dot{y} + p_1(t)y$$

Thus the linear differential operator is given by  $L = p_0(t) \frac{d}{dt} + p_1(t)$ . General linear and homogeneous ODE is given by

$$Ly = 0$$

and non-homogeneous ODE is given by

$$Ly = q(t)$$

where  $p_0, p_1$  and q are the given functions of t.

**Remark 20.** Due to the linear structure of the operator, it follows the superposition principle i.e. if  $y_1$  and  $y_2$  are the two solutions of the homogeneous equation then  $\alpha y_1 + \beta y_2$  is also a solution of the homogeneous equation.

**Definition 12** (Singular equations). The equations in which coefficient of the highest order term vanished at one or more points are called singular equations. E.g. Bessel's eqn, Lagrange eqn, Legendre eqn.

**Definition 13** (Regular equations). The coefficient of highest order term is never 0. Their general form is given by  $Ly := \dot{y} + p(t)y = q(t)$ .

**Remark 21.** A continuous function defined on a interval in  $\mathbb{R}$  whose modulus is a constant, then f itself is constant.

For non-homogeneous equation

$$\dot{y} + p(t)y = q(t)$$

we can find a function h(t) such that  $\dot{h}(t) = \dot{y} + p(t)y$  then we  $\dot{h}(t) = q(t)$  which can be integrated easily to find the solution. Let a differentiable function  $\mu(t)$  such that  $d(\mu y)/dt = \mu \dot{y} + \dot{\mu}y$ . Multiplying both sides by  $\mu(t)$  in main non-homogeneous equation and then comparing with the derivative of  $\mu(t)$ , we get  $\dot{\mu}(t) = \mu(t)p(t)$ . Here  $\mu(t)$  is called *integrating* 

$$\mu(t)\dot{y}(t) + \mu(t)p(t)y(t) = \mu(t)q(t). \tag{3.1.16}$$

If  $\mu$  is positive, then any solution of (3.1.15) is a solution of (3.1.16) and vice versa. The term on the left hand side of (3.1.16) can be written as  $\frac{d}{dt}(\mu y)$ , provided  $\mu$  satisfies  $\dot{\mu}(t) - p(t)\mu(t) = 0$ . Thus, (3.1.16) is exact. Note that the equation satisfied by  $\mu$  is a homogeneous linear DE in  $\mu$  and hence,  $\mu(t) = \exp\left(\int^t p(\tau)d\tau\right)$  is a solution and it is positive. Thus, (3.1.16) becomes

FIGURE 1. why  $\mu$  has to be positive? Check this.

factor associated with homogeneous part.

#### 4.3. Exact Differential equations.

**Definition 14** (Exact Differential equations). If the differential equation  $\dot{y} = f(t,y)$  can be written as  $\frac{d}{dt}\phi(t,y(t)) = 0$  for a two variable function  $\phi$  in a domain in the t-y plane, then the differential equation is said to be an exact differential equation.

Necessary condition for DE to be exact. Consider the differential form Mdt + Ndy, this can be written as general first order equation as

$$M(t,y) + N(t,y)\dot{y} = 0$$

Above DE is exact if and only if there exists  $\phi$  such that  $\frac{d}{dt}\phi(t,y) = 0$ . Therefore,

$$\frac{\partial}{\partial t}\phi(t,y) + \frac{\partial}{\partial y}\phi(t,y)\dot{y} = M(t,y) + N(t,y)\dot{y}$$

and this implies

$$M = \frac{\partial \phi}{\partial t}$$
 and  $N = \frac{\partial \phi}{\partial y}$ .

Assuming  $\phi$  is double differentiable, we get

$$\frac{\partial M}{\partial y} = \frac{\partial N}{\partial t}$$

**Theorem 12.** Assume M, N are defined on a rectangle  $D = (a, b) \times (c, d)$  and  $M, N \in C^1(D)$ . Then, there exists a function  $\phi$  defined in D, such that  $M = \frac{\partial \phi}{\partial t}$  and  $N = \frac{\partial \phi}{\partial u}$  if and only if  $\frac{\partial M}{\partial u} = \frac{\partial N}{\partial t}$ .

*Proof.* ( $\Longrightarrow$ ) Assume there exists a function  $\phi$  following the needed conditions given in theorem. The first relation  $M = \frac{\partial \phi}{\partial t}$  suggests that  $\phi(t,y) = \int M dt + h(y)$  for some h. Hence,

$$N = \frac{\partial \phi}{\partial t} = \int \frac{\partial M}{\partial y} dt + \frac{\partial h(y)}{\partial y} \implies \frac{\partial h}{\partial y} = N - \int \frac{\partial M}{\partial y} dt$$

Since, LHS is a function of y alone hence RHS should be a function of y alone. Therefore,

$$\frac{\partial}{\partial t} \left( N - \int \frac{\partial M}{\partial y} dt \right) = 0 \implies \frac{\partial N}{\partial t} - \frac{\partial M}{\partial y} = 0$$

(  $\iff$  ) Assume M and N satisfies  $\frac{\partial M}{\partial y} = \frac{\partial N}{\partial t}$ . Let a point  $(t_0, y_0) \in D$ . Define  $\phi(t, y) = \int_{t_0}^t M(s, y) ds + h(y)$  where h has to be determined.

$$\frac{\partial \phi(t,y)}{\partial y} = \int_{t_0}^t \frac{\partial M(s,y)}{\partial y} ds + \frac{\partial h(y)}{\partial y}$$
$$= \int_{t_0}^t \frac{\partial N(s,y)}{\partial t} ds + \frac{\partial h(y)}{\partial y}$$
$$= N(t,y) - N(t_0,y) + \frac{\partial h(y)}{\partial y}$$

If we want  $\frac{\partial \phi}{\partial y}$  to be equal to N(t,y) then  $\frac{\partial h(y)}{\partial y} = N(t_0,y) \implies h(y) = \int_{y_0}^{y} N(t_0,\xi)d\xi$ . Therefore,

$$\phi(t,y) = \int_{t_0}^{t} M(s,y)ds + \int_{y_0}^{y} N(t_0,\xi)d\xi$$

 $\phi$  is determined upto a constant. Hence on changing the point  $(t_0, y_0)$  the constant term in  $\phi$  will change. We may discard all the constant in expression of  $\phi$ .

**Definition 15** (Alternate definition of Exact ODE). The DE,  $M(t, y) + N(t, y)\dot{y} = 0$  is said to be exact if  $\frac{\partial M}{\partial y} = \frac{\partial N}{\partial t}$ .

**Solution of exact DEs.** If a DE is exact then it can be written as  $\frac{\partial \phi}{\partial t}(t,y) = 0$  and the solution of this DE can be written as  $\phi = constant$ .

#### 5. SECOND ORDER LINEAR EQUATIONS

A second order linear differential equation is given by

$$Ly \equiv \ddot{y} + p(t)\dot{y} + q(t)y = r(t).$$

The operator  $Ly = L(t, y, \dot{y}, \ddot{y})$  is a multi-linear operator (linear in each of its variable). The IVP for second order linear ODE is given by

$$\begin{cases}
\ddot{y} + p(t)\dot{y} + q(t)y = r(t), t \in I \\
y(t_0) = y_0, \dot{y}(t_0) = y_1
\end{cases}$$

**Theorem 13.** Let p, q, r be continuous functions defined in a compact interval  $I(t_0)$  and  $y_0, y_1$  be any real numbers. Then, the IVP above has a unique solution y defined in  $I(t_0)$  satisfying  $y(t_0) = y_0$  and  $\dot{y}(t_0) = y_1$ 

**Lemma 1.** Wronskian  $W(z, w) \equiv 0$  if and only if z and w are dependent. Alternatively,  $W \not\equiv 0$  if and only if z and w are independent.

*Proof.* Let  $t_0 \in I$  be a point in interval I on which z and w are defined. Then the Wronskian is given by

$$W = \begin{vmatrix} z(t_0) & w(t_0) \\ \dot{z}(t_0) & \dot{w}(t_0) \end{vmatrix}$$

This comes from the linear system (for  $y_0$  and  $y_1$  check next theorem)

$$\begin{cases} z(t_0)\alpha + w(t_0)\beta = y_0 \\ \dot{z}(t_0)\alpha + \dot{w}(t_0)\beta = y_1 \end{cases}$$

( $\iff$ ) If z = kw or w = kz then the  $W \equiv 0$ .

$$(\Longrightarrow)$$
 Assume that  $W \equiv 0 \implies z\dot{w} - \dot{z}w = 0$ 

$$z\frac{dw}{dt} = \frac{dz}{dt}w$$

$$\int \frac{dw}{wdt} \times dt = \int \frac{dz}{zdt} \times dt$$

$$\ln w = \ln z + c$$

$$w = e^{\ln z}e^{c}$$

$$w = ze^{c}$$

Therefore, w and z are dependent.

**Theorem 14.** Let z, w be two solutions of the linear homogeneous second order DE. Then, for any  $\alpha, \beta \in \mathbb{R}$ , the function  $y = \alpha z + \beta w$  is also a solution.

Further, if z and w are two linearly independent solution of the homogeneous DE in an interval  $I(t_0)$ , then every solution y can be written as  $\alpha z + \beta w$ .

*Proof.* First part of the theorem is trivial to verify. Let y be a solution of the homogeneous DE. For some  $t_0 \in I$ ,  $y(t_0) = y_0$  and  $\dot{y}(t_0) = y_1$ . We want to show that there exists  $\alpha$  and  $\beta$  such that  $y(t) = \alpha z(t) + \beta w(t)$ . Since z and w are independent,  $W \not\equiv 0$  then there exists a unique solution to the linear system

$$\begin{cases} z(t_0)\alpha + w(t_0)\beta = y_0 \\ \dot{z}(t_0)\alpha + \dot{w}(t_0)\beta = y_1 \end{cases}$$

Hence, there exists  $\alpha$  and  $\beta$ . Here, z and w are basis elements for the solution space of second order linear homogeneous DE. Therefore, all the solution for this equation can be written in this form.

**Theorem 15.**  $\dim(S) = 2$  where S is the solution space of second order linear homogeneous DE.

**Remark 22.** For *n*-th order linear homogeneous DE, the solution space is *n*-dimensional.

#### 5.1. Non-homogeneous equations.

$$Ly := \ddot{y} + p(t)\dot{y} + q(t)y = r(t) \tag{1}$$

**Theorem 16.** Assume that the functions p, q and r in (1) are continuous. Let S and  $\tilde{S}$  be respectively, the solution set of homogeneous equation and non homogeneous equation. Then, S is a linear space of dimension 2 and  $\tilde{S}$  is an affine space given by  $\tilde{S} = S + y_p$  for any particular solution  $y_p$  of (1).

*Proof.* The assumption of continuity is needed for existence of solution. We have already proved in theorem 15 that the dimension of S is 2. Since, any linear combination of solution is a solution and scalar times the solution is a solution, S is linear.

For some arbitrary  $y \in \tilde{S}$  and let some  $y_p$  be a particular solution for (1). Since Ly = r(t) and  $Ly_p = r(t)$  implies that  $L(y - y_p) = 0$  (due to linearity of operator L) then  $y - y_p$  is a solution of homogeneous equation. Hence, by theorem 14,  $y - y_p = \alpha z + \beta w$  where z and w are two independent solutions of homogeneous equation. Also, y is arbitrary

$$y = \alpha z + \beta w + y_p$$
$$\tilde{S} = S + y_p$$

Method of variation of parameters. This method can be used to find the particular solution, given the two independent solutions of corresponding homogeneous equation.

If  $y_1(t)$  and  $y_2(t)$  are the two independent solutions of homogeneous equation. We seek the solution to non homogeneous equation of the form

$$y(t) = c_1(t)y_1(t) + c_2(t)y_2(t)$$

where  $c_1(t)$  and  $c_2(t)$  are the functions to be determined. Compute Ly(t) = r(t), and then assume  $\dot{c_1}y_1 + \dot{c_2}y_2 = 0$  to be one equation. Another equation can be found by substituting above assumption in Ly(t) = r(t). Then by solving this system, we can find  $c_1(t)$  and  $c_2(t)$  and  $y_p(t) = y(t)$ .

#### 6. Doubts

**Ques 1.** In linear equations with constant coefficient, how is  $\dot{y} \propto y$ ? **Ques 2.** Why the assumption in method of variation is there?