

# Vectors

- Combining a vector with its additive inverse gives the **zero vector**, with length 0 and undefined direction.
- A scalar product projects one vector onto another.
- We can resolve  $\underline{a}$  into  $\parallel$  and  $\perp$  vectors w.r.t some  $\hat{n}$

$$\underline{a} = \underbrace{\underline{a} - (\underline{a} \cdot \hat{n}) \hat{n}}_{\perp} + \underbrace{(\underline{a} \cdot \hat{n}) \hat{n}}_{\parallel}$$

- Distributive property of dot product can be proved diagrammatically.
- Derive cosine rule with  $|\underline{c}| = |\underline{a} + \underline{b}| \Rightarrow |\underline{c}|^2 = (\underline{a} + \underline{b}) \cdot (\underline{a} + \underline{b})$ .

## Vector product

- $\underline{a} \times \underline{b} = |\underline{a}| |\underline{b}| \sin \theta \hat{n}$  ← only unique in 3D.
- $\underline{a} \times \underline{b} = -\underline{b} \times \underline{a}$  (**anticommutative**)
- $\underline{a} \times \underline{b} = 0 \rightarrow \underline{a} \parallel \underline{b}$  OR  $\underline{a}$  or  $\underline{b} = 0$ .
- $|\underline{a} \times \underline{b}|$  is the area of a parallelogram.
- Non-associative, i.e.  $\underline{a} \times (\underline{b} \times \underline{c}) \neq (\underline{a} \times \underline{b}) \times \underline{c}$ .

## Vector area

- Vector area  $\underline{S}$  of a finite plane surface is defined such that  $|\underline{S}| = \text{area}$ , with  $\underline{S}$  pointing normal to surface.
- The area of a projection (e.g. onto xy plane) is  $\underline{S} \cdot \hat{z}$
- We can define a total vector area for a composite surface as the sum of vector area elements,  $\underline{S} = \int d\underline{S}$   
↳  $\sum \underline{S}$  for a closed surface  $= 0$ .

## Triple products

- Scalar triple product:  $[\underline{a}, \underline{b}, \underline{c}] \equiv \underline{a} \cdot (\underline{b} \times \underline{c})$   
↳ invariant under cyclic permutation, i.e.  $\underline{a} \cdot (\underline{b} \times \underline{c}) = \underline{c} \cdot (\underline{a} \times \underline{b}) = \underline{b} \cdot (\underline{c} \times \underline{a})$   
↳ gives the volume of a parallelepiped

• If scalar triple product is zero, vectors are coplanar.

• The vector triple product is  $\underline{a} \times (\underline{b} \times \underline{c})$ , which can be evaluated with the BAC-CAB rule:

$$\cancel{\underline{a} \times (\underline{b} \times \underline{c})} = \underline{b}(\underline{a} \cdot \underline{c}) - \underline{c}(\underline{a} \cdot \underline{b}).$$

$\hookrightarrow \underline{a} \times (\underline{b} \times \underline{c})$  lies in the plane of  $\underline{b}$  and  $\underline{c}$ .

### Lines and planes

• A line is parameterised by  $\lambda$ :  $\underline{L} = \underline{a} + \lambda \underline{l}$

• Because  $(\underline{r} - \underline{a}) \parallel \underline{l}$ , we can also write:  $\underline{L} \times \underline{l} = \underline{a} \times \underline{l}$

• For a plane:  $\underline{L} = \underline{a} + \lambda \underline{f} + \mu \underline{g}$

$$\Rightarrow \underline{L} \cdot \underline{n} = \underline{a} \cdot \underline{n} = d$$

$\hookrightarrow$  shortest distance to the origin is  $|d|$ .

### Orthogonal basis

• In 3D, any 3 non-coplanar vectors constitute a basis.

- basis spans the space, i.e  $\underline{L} = \lambda \underline{a} + \mu \underline{b} + \nu \underline{c}$  where the components  $\{\lambda, \mu, \nu\}$  are unique.

- basis vectors will have linear independence.

• Components can be extracted using the reciprocal basis

cyclic order  
preserved.  $\left\{ \begin{array}{l} \underline{A} = \frac{\underline{b} \times \underline{c}}{[\underline{a}, \underline{b}, \underline{c}]} \quad \underline{B} = \frac{\underline{c} \times \underline{a}}{[\underline{a}, \underline{b}, \underline{c}]} \quad \underline{C} = \frac{\underline{a} \times \underline{b}}{[\underline{a}, \underline{b}, \underline{c}]} \end{array} \right.$

$\hookrightarrow$  the component is just dot product of  $\underline{L}$  with the appropriate reciprocal basis vector:

$$\lambda = \underline{A} \cdot \underline{L} \quad \mu = \underline{B} \cdot \underline{L} \quad \nu = \underline{C} \cdot \underline{L}$$

• A basis is orthonormal if all basis vectors are  $\perp$  and have unit length

• Right-handed if  $[\underline{a}, \underline{b}, \underline{c}] > 0$

• Direction cosines are cosines of angles between  $\underline{a}$  and coordinate axes, i.e  $\underline{a} = |\underline{a}|(\cos \theta_x, \cos \theta_y, \cos \theta_z)$  in Cartesian.

In Cartesian,  $\underline{a} \cdot \underline{b}$  is invariant under rotation.

$$\underline{a} \times \underline{b} = \begin{vmatrix} \mathbf{i} & \mathbf{j} & \mathbf{k} \\ a_x & a_y & a_z \\ b_x & b_y & b_z \end{vmatrix}$$

$$\underline{a} \cdot (\underline{b} \times \underline{c}) = \begin{vmatrix} a_x & a_y & a_z \\ b_x & b_y & b_z \\ c_x & c_y & c_z \end{vmatrix} \leftarrow \text{i.e transformed volume of a unit cube.}$$

### Polar coordinates

Point specified by  $(r, \phi)$

$$x = r\cos\phi \quad y = r\sin\phi$$

$$r = \sqrt{x^2 + y^2} \quad \phi = \tan^{-1}(y/x).$$

Circle described by  $r=a$

Straight line at angle  $\alpha$  to  $y$ -axis with shortest dist  $|d|$ :

$$r\cos(\phi - \alpha) = d.$$

We can use the following orthonormal basis:

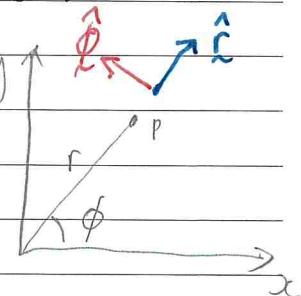
$$\hat{\mathbf{i}} = \cos\phi \hat{\mathbf{i}} + \sin\phi \hat{\mathbf{j}}$$

$$\hat{\mathbf{j}} = -\sin\phi \hat{\mathbf{i}} + \cos\phi \hat{\mathbf{j}}$$

We can evaluate  $\hat{\mathbf{r}}$ :

$$\hat{\mathbf{r}} = \hat{\mathbf{r}} \hat{\mathbf{i}} + r\phi \hat{\mathbf{j}}$$

The area element will be  $rdrd\phi$ .



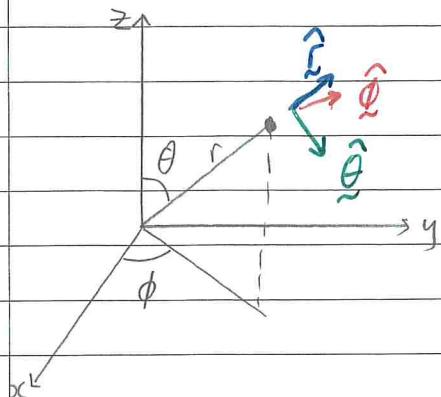
### Cylindrical coordinates

Extension of plane polar coordinates to include  $z$ .

$$x = r\cos\phi \quad y = r\sin\phi \quad z = z$$

Volume element is:  $dV = rdrd\phi dz$ .

## Spherical coordinates



Points described by radius, polar angle, azimuthal angle (i.e  $r, \theta, \phi$ ).

$$x = r \sin \theta \cos \phi \quad y = r \sin \theta \sin \phi \quad z = r \cos \theta$$

$$r = \sqrt{x^2 + y^2 + z^2} \quad \theta = \cos^{-1}\left(\frac{z}{r}\right) \quad \phi = \tan^{-1}\left(\frac{y}{x}\right)$$

We can find the orthonormal basis vectors using:

$$\hat{i} = \frac{\partial \mathbf{r}}{\partial r} / \left\| \frac{\partial \mathbf{r}}{\partial r} \right\| \quad \hat{j} = \frac{\partial \mathbf{r}}{\partial \theta} / \left\| \frac{\partial \mathbf{r}}{\partial \theta} \right\| \quad \hat{k} = \frac{\partial \mathbf{r}}{\partial \phi} / \left\| \frac{\partial \mathbf{r}}{\partial \phi} \right\|$$

$$\Rightarrow \hat{i} = \sin \theta \cos \phi \hat{i} + \sin \theta \sin \phi \hat{j} + \cos \theta \hat{k}$$

$$\hat{j} = \cos \theta \cos \phi \hat{i} + \cos \theta \sin \phi \hat{j} - \sin \theta \hat{k}$$

$$\hat{k} = -\sin \phi \hat{i} + \cos \phi \hat{j}$$

$$dV = (dr)(rd\theta)(rsin\theta d\phi) = r^2 \sin \theta dr d\theta d\phi.$$

# Complex Numbers

- Complex numbers are a closed field  $\rightarrow$  all operations return  $z$
- Complex conjugate  $z^* \equiv a - ib$  for  $z = a + ib$

$$\hookrightarrow z z^* = a^2 + b^2 > 0$$

$$\hookrightarrow z + z^* = 2 \operatorname{Re}(z)$$

$$\hookrightarrow z - z^* = 2i \operatorname{Im}(z).$$

$$\hookrightarrow \frac{1}{z} = \frac{z^*}{|z|^2}$$

- Multiplying corresponds to scaling and rotation.

- De Moivre's theorem:  $(\cos \theta + i \sin \theta)^n = \cos n\theta + i \sin n\theta$ .

$\hookrightarrow$  can be used to derive trig identities

$$\text{e.g. } \cos 5\theta = \operatorname{Re}(\cos 5\theta + i \sin 5\theta) = \operatorname{Re}((\cos \theta + i \sin \theta)^5).$$

$$\text{e.g. } \cos \theta = \frac{1}{2}(z + z^{-1}) \Rightarrow \cos 5\theta = \frac{1}{2^5}(z + z^{-1})^5 \text{ etc.}$$

- Euler's formula:  $e^{i\theta} = \cos \theta + i \sin \theta$ .

- The  $n$ th roots of unity are the solutions to  $z^n = 1$  for positive  $n$ .

$$e^{in\theta} = 1 \Rightarrow \theta = \frac{2\pi k}{n}, \quad k=0, 1, 2, \dots, n-1$$

$\therefore$  roots are  $1, \omega, \omega^2, \dots, \omega^{n-1}$  with  $\omega \equiv e^{2\pi i/n}$

- We define the complex logarithm as:

$$\ln z = \ln(r e^{i\theta}) = \ln r + i(\theta + 2\pi n) \quad n=0, \pm 1, \pm 2, \dots$$

$\hookrightarrow$  the principal value is  $\ln r + i\theta$  for  $\theta \in [0, 2\pi]$ .

- Likewise, general powers will be multi-valued

$$z_1^{z_2} = e^{z_2 \ln z_1}$$

- The fundamental theorem of algebra states:

$$a_n z^n + a_{n-1} z^{n-1} + \dots + a_1 z + a_0$$

has  $n$  complex roots for all possible complex coefficients.

## Hyperbolic functions

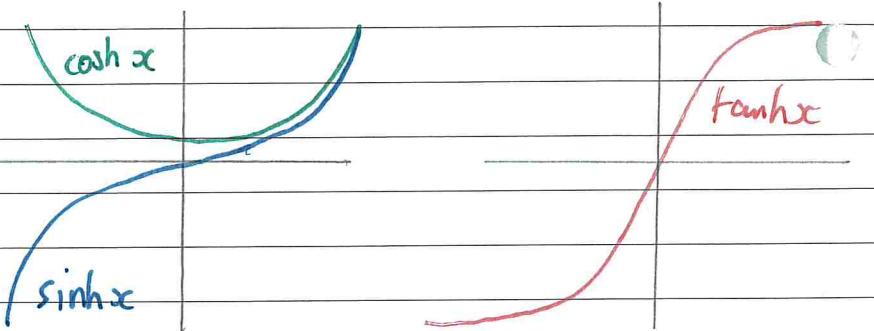
- Define:  $\cos z \equiv \frac{1}{2}(e^{iz} + e^{-iz})$  and  $\sin z \equiv \frac{1}{2i}(e^{iz} - e^{-iz})$ .
- The hyperbolic functions are these functions evaluated on the imaginary axis.

$$\cosh y \equiv \cos(iy) = \frac{1}{2}(e^y + e^{-y})$$

$$\sinh y \equiv \frac{1}{2i}\sin(iy) = \frac{1}{2i}(e^y - e^{-y})$$

$$\tanh y = \frac{1}{i} \sin(iy) = \frac{1}{2}(e^y - e^{-y}).$$

We can then define  
 $\tanh$ ,  $\text{sech}$ ,  $\text{cosech}$  etc.



- We can generate identities by substituting  $iy$  in and using  $\cos iy = \cosh y$ ,  $\sin(iy) = i \sinh y$ .
  - ↪  $\cosh^2 y - \sinh^2 y = 1$
  - ↪  $\cosh(A+B) = \cosh A \cosh B + \sinh A \sinh B$
- Inverse hyperbolic functions can be expressed as elementary functions.

# Calculus and Analysis

## Limits

- Intuitively,  $\lim_{x \rightarrow x_0} f(x) = k$  means  $f(x)$  can be made arbitrarily close to  $k$  by making  $x$  close enough to  $x_0$ .
- The  $\epsilon-\delta$  definition: For real  $f(x)$  defined on some open interval containing  $x_0$  (but not necessarily at  $x_0$ ),  $\lim_{x \rightarrow x_0} f(x) = k$  means for any  $\epsilon > 0$ ,  $\exists \delta > 0$  such that:  
 $|f(x) - k| < \epsilon$  for all  $0 < |x - x_0| < \delta$
- ↳ i.e if you give me an  $\epsilon$ , I can find  $\delta$  to stay within  $\epsilon$  of  $k$ .
- ↳ in practice, we guess the limit then prove with  $\epsilon-\delta$ .
- Limits at infinity:  $|f(x) - k| < \epsilon$  for all  $x > X$ .
- Limits can be manipulated by addition and multiplication.
- If a quotient is indeterminate (top and bottom both 0 or  $\pm\infty$ ), we can use L'Hôpital's rule:

$$\lim_{x \rightarrow a} \frac{f(x)}{g(x)} = \lim_{x \rightarrow a} \frac{f'(x)}{g'(x)}.$$

## Continuity and differentiability

- A real function  $f(x)$  is continuous at  $x=a$  iff:
  - $f(a)$  exists
  - $\lim_{x \rightarrow a} f(x)$  exists and equals  $f(a)$ . ~~exists~~
- A function  $f(x)$  is differentiable at  $x=a$  iff:
  - it is continuous at  $x=a$
  - $f'(a)$  exists i.e  $\lim_{\Delta x \rightarrow 0} \frac{f(x+a+\Delta x) - f(x)}{\Delta x}$  exists.

## Leibniz formula

- Used to find  $n$ th derivative of a product of functions  
(just like Binomial theorem):

$$\frac{d^n(fg)}{dx^n} = \sum_{m=0}^n \binom{n}{m} f^{(n-m)} g^{(m)}$$

$$= f^{(n)} g + n f^{(n-1)} g' + \frac{n(n-1)}{2} f^{(n-2)} g'' + \dots + f g^{(n)}$$

- Can be proved by induction.

## Infinite Series

- Given a sequence of terms  $u_0, u_1, u_2, \dots$  the  $n$ th partial sum is  $S_n \equiv \sum_{k=0}^n u_k$

- If the partial sums have a finite limit as  $n \rightarrow \infty$ , the infinite series is **convergent**.

↳ if it doesn't converge, it either diverges or oscillates.

- If  $\sum_{k=0}^{\infty} |u_k|$  converges, the series is **absolutely convergent**  
(which also implies  $\sum_{k=0}^{\infty} u_k$  converges)

↳ otherwise if  $\sum_{k=0}^{\infty} u_k$  converges but  $|u_k|$  doesn't, series is **conditionally convergent**.

↳ For absolutely convergent series we can rearrange terms.

## Geometric progressions

$$S_n = \sum_{k=0}^n r^k = \frac{1-r^{n+1}}{1-r} \quad \text{because } r S_n = r + r^2 + \dots + r^{n+1} = r^{n+1} + S_n - 1.$$

- Series is absolutely convergent for  $|r| < 1$

$$\therefore \sum_{k=0}^{\infty} r^k = \frac{1}{1-r}$$

- If  $|r| \geq 1$ , series cannot converge.

## Convergence tests

1.  $U_k \rightarrow 0$  as  $k \rightarrow \infty$  is a necessary condition for convergence (but insufficient, e.g. harmonic series).

### 2. Comparison test:

- Compare with a series of known convergence,  $V_k$
- If all terms  $\leq$  ~~less than~~  $V_k$  for all  $k > K$ ,  $S_n$  converges
- If all terms  $>$   $V_k$  for divergent  $V$ ,  $S_n$  diverges.
- Try to compare with geometric series or harmonic series
- p-series test:  $\sum_{k=1}^{\infty} \frac{1}{k^p}$  converges for  $p > 1$  by comparison with geometric series  
 $\hookrightarrow$  diverges for  $p \leq 1$  by comp. with harmonic

### 3. Ratio test

If  $\lim_{K \rightarrow \infty} \frac{U_{K+1}}{U_K} < 1$ ,  $S_n$  converges

If  $\lim_{K \rightarrow \infty} \frac{U_{K+1}}{U_K} > 1$ ,  $S_n$  diverges

If ratio = 1, test indeterminate.

### 4. Alternating series:

- Use the Leibniz criterion:

$\sum_{K=0}^{\infty} (-1)^{K+1} a_k$  with  $a_k > 0$  converges if  $a_k$  is monotonic decreasing for large enough  $k$  and  $\lim_{K \rightarrow \infty} a_k = 0$ .

### 5. Integral test:

- If  $f(n)$  is continuous, positive, and decreasing on  $[1, \infty)$ :

$\sum_{n=1}^{\infty} f(n)$  converges / diverges as  $\int_1^{\infty} f(x) dx$ .

## Power series

Series of the form  $f(x) = \sum_{k=0}^{\infty} a_k x^k = a_0 + a_1 x + a_2 x^2 + \dots$

Either:

- converges for  $x=0$  only

- converges for all finite  $x$

- converges for  $|x| < R$ , diverges for  $|x| > R$ .

Using ratio test and  $L \equiv \lim_{k \rightarrow \infty} \frac{a_{k+1}}{a_k}$ :

convergent for  $|x| < 1/L$ , divergent for  $|x| > 1/L$

For a complex power series, this will define a circle of convergence

## Taylor series

$$f(x) = f(a) + (x-a) f'(a) + \frac{(x-a)^2}{2!} f''(a) + \frac{(x-a)^3}{3!} f'''(a) + \dots$$

Or Maclaurin series when  $a=0$ :  $f(x) = f(0) + xf'(0) + \frac{x^2}{2!} f''(0) + \dots$

We can truncate the Taylor series and add a remainder term:

$$f(x) = f(0) + xf'(0) + \frac{x^2}{2!} f''(0) + \dots + \frac{x^k}{k!} f^{(k)}(0) + R_n$$

$$\text{with } R_n = \frac{1}{k!} \int_0^x (x-t)^k f^{(k+1)}(t) dt.$$

↳ derived by  $f(x) = f(0) + \int_0^x f'(t) dt$  (FTC) then IBP.

$$e^x = 1 + x + \frac{x^2}{2!} + \frac{x^3}{3!}$$

$$\sin x = x - \frac{x^3}{3!} + \frac{x^5}{5!} - \dots$$

$$\cos x = 1 - \frac{x^2}{2!} + \frac{x^4}{4!}$$

$$\ln(1+x) = x - \frac{x^2}{2} + \frac{x^3}{3} - \frac{x^4}{4}$$

## Newton-Raphson

- Helps us find  $x^*$  such that  $f(x^*) = 0$
- If we have an initial guess  $x_0$ , we need  $h$  such that  $f(x_0 + h) = 0$ .

$$0 = f(x_0 + h) \approx f(x_0) + hf'(x_0).$$

$$\Rightarrow x_1 = x_0 - \frac{f(x_0)}{f'(x_0)}$$

- Then we can iterate this to converge on  $x^*$ .

- If  $\epsilon_i$  is the error in  $x_i$ :

$$x_{i+1} = x_i - \frac{f(x_i)}{f'(x_i)} \Rightarrow \epsilon_{i+1} = \epsilon_i - \frac{f(x^* + \epsilon_i)}{f'(x^* + \epsilon_i)}.$$

↪ approximating the last term with a Taylor expansion:

$$\epsilon_{i+1} \approx \epsilon_i^2 \frac{f''(x^*)}{2f'(x^*)} \text{ i.e rapid quadratic convergence}$$

- If there is a turning point between the root and  $x_i$ , it may not converge.

# Integration

• Formally:  $\int_a^b f(x) dx = \lim_{N \rightarrow \infty} \sum_{i=0}^{N-1} f(x_i) (x_{i+1} - x_i)$ .  $\leftarrow$  area under curve.

• Hyperbolic substitutions:

$$\sqrt{x^2 + a^2} \leftarrow x = a \sinh y$$

$$\sqrt{x^2 - a^2} \leftarrow x = a \cosh y$$

$$a^2 - x^2 \leftarrow x = a \tanh y$$

• Integrate using complex numbers, e.g.  $\int \cos x e^{\alpha x} dx = \operatorname{Re} \left( \int e^{(\alpha+i)x} dx \right)$

• If  $I(\alpha) = \int_{\alpha(a)}^{b(\alpha)} f(x; \alpha) dx$

$$I'(\alpha) = \int_{\alpha(a)}^{b(\alpha)} \frac{\partial f}{\partial \alpha} dx + \frac{db}{d\alpha} f(b; \alpha) - \frac{da}{d\alpha} f(a; \alpha)$$

## Stirling's approximation

$$\ln n! = \sum_{k=1}^n \ln k. \quad \text{But} \quad \int_1^n \ln x dx \leq \sum_{k=1}^n \ln k \leq \int_1^{n+1} \ln x dx$$

$$\therefore \ln n! \sim n \ln n - n \quad \text{for large } n.$$

## Get Cauchy-Schwarz inequality

$$|\langle a, b \rangle|^2 \leq \|a\| \|b\| \quad \text{where } \langle \cdot, \cdot \rangle \text{ is the inner product.}$$

• For an  $N$ -dimensional vectors

$$\left( \sum_{i=1}^N a_i b_i \right)^2 \leq \left( \sum_{i=1}^N a_i^2 \right) \left( \sum_{i=1}^N b_i^2 \right)$$

• Taking  $N \rightarrow \infty$ , we get Schwarz's inequality.

$$\left( \int_a^b f(x) g(x) dx \right)^2 \leq \int_a^b (f^2(x))^2 dx \int_a^b (g(x))^2 dx.$$

## Multiple integrals

$$\iiint f(x) dV = \lim_{\delta V \rightarrow 0} \sum f(r) dV.$$

- Cartesian:  $dV = dx dy dz$
- Cylindrical:  $dV = r dr d\theta dz$
- Spherical:  $dV = r^2 \sin\theta dr d\theta d\phi$
- We can do the integrals in any order.
- If limits are independent, we can factor the integral out.

## Gaussian distribution Integrals

•  $I = \int_{-\infty}^{\infty} e^{-x^2} dx$  is a common improper integral.

• Evaluate with polar coordinates:

$$I^2 = \left( \int_{-\infty}^{\infty} e^{-x^2} dx \right) \left( \int_{-\infty}^{\infty} e^{-y^2} dy \right) = \iint e^{-(x^2+y^2)} dx dy.$$

• Technically should use a in limits then  $\lim_{a \rightarrow \infty}$

# Probability

- Outcomes  $w_i$  are mutually exclusive
- The sample space is the set of all possible outcomes:  $\Omega = \{w_i\}$
- An event is a subset of  $\Omega$
- $P(A|B) = \frac{P(A \cap B)}{P(B)}$   $\Rightarrow P(A|B) = \frac{P(B|A)P(A)}{P(B)}$   
*(Bayes' Theorem)*
- Law of total probability:  $P(A) = \sum_i P(A|B_i) P(B_i)$

## Random variables

- Map sample states to an allowed value of the random variable such that the subsets partition the space.
- Assign a prob. distribution  $P(x)$ .
- Poisson distribution :  $P(X=x) = e^{-\lambda} \frac{\lambda^x}{x!}$   
↳ can be shown that it is the limit of a binomial dist as  $n \rightarrow \infty$ , with  $np = \lambda$ .
- For continuous random variables, the prob. density function is  $f(x) dx = P(x - \frac{dx}{2} \leq X \leq x + \frac{dx}{2})$ .  
 $P(a \leq X \leq b) = \int_a^b f(x) dx$  with  $\int_{-\infty}^{\infty} f(x) dx = 1$ .  
 $F(a) = \int_{-\infty}^a f(x) dx$ .
- Median is  $a$  such that  $F(a) = 1/2$
- Variance of a distribution is the same even when conditioned

# Ordinary Differential Equations

- A first-order ODE has the form  $F(y', y, x) = 0$ .
- An  $n$ th-order ODE:  $F(y^{(n)}, y^{(n-1)}, \dots, y', y, x) = 0$ .
- A **separable** 1st order ODE:

$$\frac{dy}{dx} = \frac{F(x)}{g(y)} \Rightarrow \int g(y) dy = \int f(x) dx.$$

The general solution (including a constant) can be fixed by an initial/boundary condition.

- A **linear 1st order ODE**:  $\frac{dy}{dx} + p(x)y = f(x)$  if  $f(x) = 0$ , it is **homogeneous**, and **separable**.

↳  $y$  and  $\frac{dy}{dx}$  appear linearly  
 ↳ can be solved with an **integrating factor**,  $\mu(x)$ , such that  $\mu(x) \cdot \text{LHS}$  is the derivative of something w.r.t  $x$ .

$$\Rightarrow \mu(x) = e^{\int p(x) dx}$$

↳  $\therefore \frac{d}{dx}(\mu(x)y) = \mu(x)f(x)$  which is easy to solve.

- Substitutions may be required to make an ODE linear/separable.
- Homogeneous ODE**:  $\frac{dy}{dx} = f\left(\frac{y}{x}\right)$  ← i.e.  $f$  invariant when  $x$  and  $y$  scaled.

↳ solve by sub  $u = y/x$

$$\Rightarrow y = u(x)x \Rightarrow x \frac{du}{dx} + u = f(u) \leftarrow \text{separable.}$$

- Bernoulli ODE**:  $\frac{dy}{dx} + p(x)y = q(x)y^n$

$$\hookrightarrow \text{sub } z = y^{1-n} \Rightarrow \frac{dz}{dx} = (1-n)y^{-n} \frac{dy}{dx}$$

$$\therefore \frac{dz}{dx} = (1-n)[-p(x)z + q(x)] \leftarrow \text{linear.}$$

## Second-order equations

- A linear 2<sup>nd</sup> order ODE:  $\frac{d^2y}{dx^2} + p(x)\frac{dy}{dx} + q(x)y = f(x)$ .

↳ with the linear differential operator  $L$ , we can rewrite

$$L = \frac{d}{dx^2} + p(x)\frac{d}{dx} + q(x) \Rightarrow Ly = f(x)$$

↳  $L(\alpha u) = \alpha L(u)$  if  $\alpha$  constant

↳  $L(u+v) = L(u) + L(v)$

- ↳ For a homogeneous 2<sup>nd</sup> order ODE ( $Ly = 0$ ), any linear combination of solutions is a solution  $\Leftarrow$  principle of superposition.

- For inhomogeneous case, i.e.  $Ly = f(x)$ :

- a particular integral is any solution of  $Ly = f(x)$

- the complementary function  $y_c$  is the general solution of  $Ly = 0$

- the general solution is the sum:  $y(x) = y_c(x) + y_p(x)$ .

• 2<sup>nd</sup> order ODEs are generally hard to solve unless constant coefficients.

• Consider homogeneous 2<sup>nd</sup> order linear ODE:

$$\frac{d^2y}{dx^2} + 2a\frac{dy}{dx} + by = 0.$$

↳ sub  $y = e^{\lambda x}$  as a trial gives the auxiliary equation  $\lambda^2 + 2a\lambda + b = 0$

↳ if roots are negative complex, we have oscillatory behavior.

↳ if  $\lambda_1 = \lambda_2$ , we have critical damping:  $y = (C_1 + C_2x)e^{-\alpha x}$ .

• For linear 2<sup>nd</sup> order inhomogeneous ODEs with constant coefficients:

-  $y_c$  can be found as above.

-  $y_p$  can be found with trial solutions

- if  $f(x)$  is a polynomial, try  $y_p = \text{polynomial of same degree}$

- if  $f(x) = ce^{kx}$ , try  $y_p = de^{kx}$

- if  $f(x) = c_1 \cos kx + c_2 \sin kx$ , try  $y_p = d_1 \cos kx + d_2 \sin kx$ .

- but if scalar multiples of these trial solutions are already solutions of the homogeneous eq, we may need to multiply by  $x$  or  $x^2$  and try again

• Alternatively, since it is linear and differential operators commute,  
we can factorise:  $(\frac{d}{dx} - \lambda_1)(\frac{d}{dx} - \lambda_2) = f(x)$

↳ let  $z(x) = (\frac{d}{dx} - \lambda_2)y \Rightarrow (\frac{d}{dx} - \lambda_1)z = f(x)$ .

↳ solve for  $z$  then for  $y$ .

↳ this gives us a particular integral

# Multivariable calculus

- Mixed partial derivatives are always equal, and partial derivatives commute.  $\therefore f_{xy} = f_{yx}$
- Integrating w.r.t one variable, we can treat others as constant but then we will need to add an arbitrary function.
- For  $f(x, y)$ ,  $df = \frac{\partial f}{\partial x} dx + \frac{\partial f}{\partial y} dy$ .
- Taylor series becomes:  
$$f(x+h, y+k) = f(x, y) + f_x(x, y)h + f_y(x, y)k + \frac{1}{2} f_{xx} h^2 + f_{xy} hk + \frac{1}{2} f_{yy} k^2 + \dots$$

- Suppose  $f(x, y)$  where  $x = x(u, v)$   $y = y(u, v)$ . By an abuse of notation, we write  $f(x, y) = f(u, v)$  even though they are different functions:

$$\begin{aligned} \left(\frac{\partial f}{\partial u}\right)_v &= \left(\frac{\partial f}{\partial x}\right)_y \left(\frac{\partial x}{\partial u}\right)_v + \left(\frac{\partial f}{\partial y}\right)_x \left(\frac{\partial y}{\partial u}\right)_v \\ \left(\frac{\partial f}{\partial v}\right)_u &= \left(\frac{\partial f}{\partial x}\right)_y \left(\frac{\partial x}{\partial v}\right)_u + \left(\frac{\partial f}{\partial y}\right)_x \left(\frac{\partial y}{\partial v}\right)_u \end{aligned} \quad \left. \begin{array}{l} \\ \end{array} \right\} \text{multivariable chain rule.}$$

e.g.  $f(x, y) \rightarrow f(r, \phi)$  :  $x = r\cos\phi, y = r\sin\phi$   
 $\therefore \left(\frac{\partial f}{\partial r}\right)_\phi = \cos\phi \left(\frac{\partial f}{\partial x}\right)_y + \sin\phi \left(\frac{\partial f}{\partial y}\right)_x$  etc.

- If both  $x$  and  $y$  are functions of  $t$ :

$$\frac{df}{dt} = \left(\frac{\partial f}{\partial x}\right)_y \frac{dx}{dt} + \left(\frac{\partial f}{\partial y}\right)_x \frac{dy}{dt}$$

- If we have  $F(x, y, z) = 0$ , then the partial derivatives have reciprocity and are cyclic.

$$\text{i.e. } \left(\frac{\partial x}{\partial y}\right)_z = 1 / \left(\frac{\partial y}{\partial x}\right)_z \quad \text{and} \quad \left(\frac{\partial x}{\partial z}\right)_y \left(\frac{\partial z}{\partial y}\right)_x \left(\frac{\partial y}{\partial x}\right)_z = -1$$

## Exact differentials

- $\omega = P(x,y)dx + Q(x,y)dy$  is a differential form in  $x$  and  $y$ .
- $\omega$  is an exact differential if  $\exists f(x,y)$  such that  $df = Pdx + Qdy$ .
  - ↳ equivalently, exact iff  $\frac{\partial P}{\partial y} = \frac{\partial Q}{\partial x}$ .
  - ↳ if  $Pdx + Qdy$  is exact,  $f(x,y) = c$ .
- We can make an inexact differential form exact with an integrating factor:  $M(x,y)[Pdx + Qdy]$ 
  - ↳ this is very difficult to solve for  $M$ , so we instead try to find  $M(x)$  or  $M(y)$  only.
  - e.g.  $M(x)$ :  $M \frac{\partial P}{\partial y} = Q \frac{\partial u}{\partial x} + M \frac{\partial Q}{\partial x}$  if exact
  $\Rightarrow \frac{1}{M} \frac{\partial M}{\partial x} = \frac{1}{Q} \left( \frac{\partial P}{\partial y} - \frac{\partial Q}{\partial x} \right)$ .
  - ↳ likewise for  $M(y)$ :  $\frac{1}{M} \frac{\partial M}{\partial y} = - \frac{1}{P} \left( \frac{\partial P}{\partial y} - \frac{\partial Q}{\partial x} \right)$ .



## Maxwell's relations

- Any two of  $(P, V, T, S)$  can describe the state of a gas.
- Given a thermodynamic relation, we
- The fundamental thermodynamic relation is
 
$$dV = TdS - pdV$$
  - ↳ if we treat  $V$  as a function of  $(S, V)$ :
  - $dV = \left(\frac{\partial V}{\partial S}\right)_V dS + \left(\frac{\partial V}{\partial V}\right)_S dV \Rightarrow \left(\frac{\partial V}{\partial S}\right)_V = T \text{ and } \left(\frac{\partial V}{\partial V}\right)_S = -p$
  - $\therefore \left(\frac{\partial T}{\partial V}\right)_S = - \left(\frac{\partial p}{\partial S}\right)_V$  by mixed partials. This is one of Maxwell's relations.

- We can derive the others using Legendre transformations
  - ↳  $F = U - TS \Rightarrow dF = -SdT - pdV$
  - ↳  $H = U + PV \Rightarrow dH = TdS + Vdp$
  - ↳  $G = H - TS \Rightarrow dG = -SdT + Vdp$

- We can also derive a different type of relation:

$$dV = TdS - PdV \text{ but let } U = U(T, S)$$

$$\therefore dU = TdS - P \left[ \left( \frac{\partial V}{\partial T} \right)_S dT + \left( \frac{\partial V}{\partial S} \right)_T dS \right]$$

then we take partial derivatives and equate.

### Stationary points

- Because  $f(\underline{x}) = f(\underline{x}_0) + \nabla f(\underline{x}_0) \cdot \delta \underline{x}$ , a point is stationary if  $\nabla f(\underline{x}_0) = \underline{0}$ .

- To find the character of the stationary points, we use the determinant of the Hessian:  $H = \begin{pmatrix} f_{xx} & f_{xy} \\ f_{yx} & f_{yy} \end{pmatrix}$ .

$\det H > 0$  and  $f_{xx} > 0 \Rightarrow$  minimum

$\det H > 0$  and  $f_{xx} < 0 \Rightarrow$  maximum

$\det H < 0 \Rightarrow$  saddle

$\det H = 0$  inconclusive.

- For more variables:

↳ if all eigenvalues  $> 0$ , min

↳ if all eigenvalues  $< 0$ , max

↳ else saddle.

### Conditional stationary values

- To optimise  $f(x, y)$  subject to  $g(x, y) = c$ , solve

$\nabla f = \lambda \nabla g$ , where  $\lambda$  is a Lagrange Multiplier

↳ consider some displacement  $d\underline{x}$

↳  $d\underline{x}$  must be tangent to  $g(x, y) = 0 \therefore (\nabla g) \cdot d\underline{x} = 0$

↳ Likewise,  $df = (\nabla f) \cdot d\underline{x} = 0$  by definition of a stationary point

↳  $\therefore \nabla f \parallel \nabla g$ .

- For more constraints:  $\nabla f = \lambda \nabla g + \mu \nabla h$

## Boltzmann distribution

Consider a system which has  $n$  possible discrete states, with which holds  $N_i$  particles whose energy is  $E_i$ :

$$\hookrightarrow \text{total number of particles is } N = \sum_{i=1}^n N_i$$

$$\hookrightarrow \text{total energy is } E = \sum_{i=1}^n N_i E_i$$

- A given distribution of particles can be achieved in  $W$  ways:

$$W = \frac{N!}{N_1! N_2! \dots N_n!}$$

- The most likely state maximises  $W$ , or  $\ln W$  equivalently

$$\ln W = \ln(N!) - \sum_{i=1}^n \ln(N_i!)$$

$\hookrightarrow$  in an isolated system,  $N = \hat{N}$  and  $E = \hat{E}$

$$\therefore L = \ln(N!) - \sum_{i=1}^n \ln(N_i!) - \alpha \left( \sum_{i=1}^n N_i - \hat{N} \right) - \beta \left( \sum_{i=1}^n N_i E_i - \hat{E} \right).$$

$\hookrightarrow N_i$  are the variables,  $\therefore$  need  $\frac{\partial L}{\partial N_i}$ .

$$\frac{\partial L}{\partial N_i} = \ln N - \ln N_i - \alpha - \beta E_i \quad \text{because} \quad \frac{\partial \ln N!}{\partial N_i} = \frac{\partial \ln N!}{\partial N} \frac{\partial N}{\partial N_i} = \ln N.$$

$$\hookrightarrow \text{then set } \frac{\partial L}{\partial N_i} = 0 \text{ and solve for } N_i \quad \therefore N_i = N e^{-\alpha} e^{-\beta E_i}$$

- This gives the Boltzmann dist.

- Different assumptions about particle states leads to different  $W$ .

# Vector calculus

- Let  $\phi(x, y, z)$  be a scalar field.

$$\text{grad } \phi = \nabla \phi = \left( \frac{\partial \phi}{\partial x}, \frac{\partial \phi}{\partial y}, \frac{\partial \phi}{\partial z} \right)$$

- The rate of change of  $\phi$  in direction  $t$  is the directional derivative

$$\frac{d\phi}{ds} = t \cdot \nabla \phi$$

- this implies that  $\nabla \phi$  is the direction of most rapid increase.
- Given a surface  $f(x, y, z) = c$ ,  $\nabla f$  must be normal to the surface because  $f$  is constant along the surface

$$\Rightarrow \hat{n} = \frac{\nabla f}{|\nabla f|}$$

## Line integrals

- Consider a curve parameterised by  $t$ :  $\underline{c} = (x(t), y(t), z(t))$

$$d\underline{c} = \frac{dx}{dt} dt$$

- For a scalar field parameterised by an arc length  $s$ :

$$\int_C \phi d\underline{c} = \int_{s_1}^{s_2} \phi(\underline{c}(s)) ds$$

- For a more general parameter  $t$ :

$$\int_C \phi d\underline{c} = \int_{t_1}^{t_2} \phi(\underline{c}(t)) \left| \frac{d\underline{c}}{dt} \right| dt$$

- For a vector field  $\underline{F}(t)$

$$\int_C \underline{F} \cdot d\underline{c} = \int_{t_1}^{t_2} \underline{F}(\underline{c}(t)) \frac{d\underline{c}}{dt} dt.$$

- The Gradient theorem:

$$\int_C (\nabla \phi) \cdot d\underline{c} = \int_c d\phi = \phi(\underline{c}_2) - \phi(\underline{c}_1)$$

## Conservative fields

- Q Line integral independent of the path.
  - $\tilde{F} = -\nabla \phi$  for some  $\phi(r)$
  - $\tilde{F} \cdot d\tilde{x}$  is exact
  - $\oint_C \tilde{F} \cdot d\tilde{x} = 0$  for all closed curves.
  - $\nabla \times \tilde{F} = 0$
- } each implies the other.

## Surface integrals

- For a general curved surfaces  $S$  in space, the vector area element is defined by  $d\tilde{S} = \hat{n} dS$ . The total vector area is  $\int_S \hat{n} dS$
- The flux of  $E$  through  $S$  is defined by:

$$\int_S \tilde{E} \cdot d\tilde{S} = \int_S E \cdot \hat{n} dS$$

## Divergence

$$\cdot \operatorname{div} F \equiv \nabla \cdot F \equiv \frac{\partial F_x}{\partial x} + \frac{\partial F_y}{\partial y} + \frac{\partial F_z}{\partial z} \leftarrow \text{SCALAR}$$

• The divergence theorem:

$$\iiint_V (\nabla \cdot F) dV = \int_S F \cdot dS$$

• Can be used to define divergence:

$$\nabla \cdot F = \lim_{\delta V \rightarrow 0} \frac{1}{\delta V} \int_{\delta S} F \cdot dS$$

• If a surface is not closed, we can first construct a closed one then apply the divergence theorem.

• The Laplacian is the divergence of a gradient

$$\nabla^2 \phi = \nabla \cdot (\nabla \phi) = \frac{\partial^2 \phi}{\partial x^2} + \frac{\partial^2 \phi}{\partial y^2} + \frac{\partial^2 \phi}{\partial z^2}$$

↳ it is also a scalar

Curl

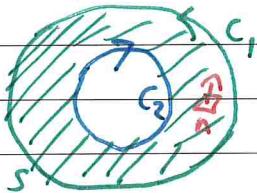
$$\cdot \nabla \times \underline{F} = \begin{vmatrix} \hat{i} & \hat{j} & \hat{k} \\ \frac{\partial}{\partial x} & \frac{\partial}{\partial y} & \frac{\partial}{\partial z} \\ F_x & F_y & F_z \end{vmatrix} \quad (\text{Vector})$$

• Stokes theorem:  $\int_S \nabla \times \underline{F} \cdot d\underline{s} = \int_C \underline{F} \cdot d\underline{r}$   
where  $C$  bounds  $S$ .

↳ use RH grip rule for direction.

- This leads to a geometric definition:  $\hat{n} \cdot (\nabla \times \underline{F}) = \lim_{S \rightarrow 0} \frac{1}{|S|} \int_S \underline{F} \cdot d\underline{s}$
- For any vector conservative field  $\underline{F} = -\nabla \phi$ ,  $\nabla \times \underline{F} = 0$ .
- Many different surfaces can be bounded by a closed curve, but only one volume is bounded by a closed surface
- A multiply connected surface may have multiple bounding curves

e.g. annulus



$$\int_S (\nabla \times \underline{F}) \cdot d\underline{s} = \int_{C_1} \underline{F} \cdot d\underline{r} + \int_{C_2} \underline{F} \cdot d\underline{r}$$

- For a planar surface, we can use Green's theorem, a special case.

# Fourier Series

- Functions are orthogonal on an interval if their inner product is zero:

$$\int_a^b f(x)g(x) dx = 0$$

- On the interval  $[-\pi, \pi]$ , all  $\cos nx$  and  $\sin mx$ ,  $\forall n, m \in \mathbb{Z}$  are mutually orthogonal (but not normalized):

$$\int_{-\pi}^{\pi} \cos mx \cos nx dx = \begin{cases} 2\pi, & m=n=0 \\ \pi, & m=n \neq 0 \\ 0, & m \neq n \end{cases}$$

$$\int_{-\pi}^{\pi} \sin mx \sin nx dx = \begin{cases} \pi, & m=n \neq 0 \\ 0, & \text{otherwise.} \end{cases}$$

$$\int_{-\pi}^{\pi} \cos mx \sin nx = 0$$

- We can change bounds to  $\pm L$  provided we scale  $mx \rightarrow \frac{m\pi x}{L}$   
 ↳ then this will work for any  $[a, b]$  such that  $2L = b - a$ .
- $\sin mx$  and  $\cos nx$  thus form a basis, such that almost any  $f(x)$  can be represented with a Fourier series:  $f(x) = \frac{a_0}{2} + \sum_{n=1}^{\infty} (a_n \cos \frac{n\pi x}{L} + b_n \sin \frac{n\pi x}{L})$

↳ Fourier coefficients can be found by integrating w.r.t after mult. with  $\cos(\frac{m\pi x}{L})$  or  $\sin(\frac{m\pi x}{L})$

$$a_n = \frac{1}{L} \int_{-L}^L f(x) \cos \frac{n\pi x}{L} dx \quad b_n = \frac{1}{L} \int_{-L}^L f(x) \sin \frac{n\pi x}{L} dx$$

↳ For even functions, all  $b_n = 0$ , so it is a cosine series

↳ For odd functions, all  $a_n = 0$ , so it is a sine series

- Fourier coefficients decline like  $\frac{1}{n^2}$ , so we can approx functions.
- We can observe how fast the coefficients decline to understand convergence.
- Around a discontinuity, the Fourier series will always overshoot, even in the limit, though the width of the overshoot ↓. Gibbs phenomenon.
- Differentiating always reduces smoothness:  
 ↳ Fourier coefficients drop less rapidly.

The mean-square value of a periodic function can be evaluated using Parseval's theorem:  $\frac{1}{2L} \int_{-L}^L f(x)^2 dx = \frac{a_0^2}{4} + \frac{1}{2} \sum_{n=1}^{\infty} (a_n^2 + b_n^2)$

The set of values, for different  $n$ , is the power spectrum and describes how power is distributed amongst the harmonics.

### Complex Fourier series

$$f(x) = \sum_{n=-\infty}^{\infty} c_n e^{inx/L}$$

$$c_n = \frac{1}{2L} \int_{-L}^L f(x) e^{-inx/L} dx$$

- $e^{inx/L}$  is used as a basis
- For complex functions,  $f(x)$  and  $g(x)$  are orthogonal if:

$$\int_a^b [f(x)]^* g(x) dx = 0.$$