

Numerical Methods

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Errors

- Computations involve errors because of noisy data, floating repr, etc.
- **Absolute error**: $\epsilon_x = |x^* - x| \Rightarrow x^* = x \pm \epsilon_x$ $\leftarrow x^*$ is our approx.
- **Relative error**: $\eta_x = \Delta x / |x| \Rightarrow x^* = x(1 + \eta_x)$
- For $x^* \pm y^*$, $\epsilon_{x \pm y} = \epsilon_x + \epsilon_y$. No quadrature because these are not distributions.
 - $\hookrightarrow$ be careful if subtraction results in near-zero because $\eta \uparrow$ a lot
- The error in $f(x^*)$ is $\epsilon_{f(x)} \approx |f'(x^*)| \epsilon_x$.
- If $f(x^*)$ is a **truncated Taylor series**, we have the **Lagrange error bound**:
$$\epsilon_{f(x)} \leq \frac{M(\epsilon_x)^{n+1}}{(n+1)!}$$
where M is the upper bound on $|f^{(n+1)}(k)|$, for $k \in (x, x^*)$ or $k \in (x^*, x)$.
- **Forward error analysis** propagates perturbations in the input
 - $\hookrightarrow$ sometimes overestimates because we propagate worst case but really there will be some cancellation.
- **Backward error analysis** instead finds the input error associated with a given output data and verifies it is within the uncertainty of the data
 - e.g. if $f(x) = x^2$ then $f^*(x) = (x^2)^* = x^2(1 + \eta_{x^2})$.
We can then find some $\tilde{\eta}$ such that $(1 + \tilde{\eta})^2 = 1 + \eta$
 $\Rightarrow f^*(x) = x^2(1 + \tilde{\eta})^2 = f(x(1 + \tilde{\eta}))$. $\tilde{\eta}$ can be compared with the unc. of the data.

Numerical differentiation

- For analytic functions, $f(x) = \sum_{n=0}^{\infty} \frac{f^{(n)}(a)}{n!} (x-a)^n \leftarrow$ Taylor series
- This allows the first approximation that $f'(x_0) = \frac{f(x) - f(x_0)}{x - x_0}$
- Alternatively, we can utilise the **discretisation parameter** h ; the **forward approx** is then $D_{f'}^+(x) \equiv \frac{f(x+h) - f(x)}{h}$
 - $\hookrightarrow$ by expanding $f(x+h)$ and rearranging for $f'(x)$, we can show that the **truncation error** is $f'(x) - D_{f'}^+(x) = -h f''(x)/2 + O(h^2)$

• The symmetric approx is $P_{f,1}^0(x) \equiv \frac{f(x+h) - f(x-h)}{2h}$

↳ truncation error is $h^2 f'''(x)/3! + O(h^3)$

↳ quadratic in h i.e. second-order method

• Generally, an n th order method has truncation error $O(h^n)$

↳ N.B.: not the same as order of convergence

Numerical integration

• A Riemann sum approximates $\int_a^b f(x) dx \approx \sum_{i=0}^{n-1} f(t_i)(x_{i+1} - x_i)$

↳ $t_i \in [x_i, x_{i+1}]$ is the point within the interval at which height is evaluated.

• The trapezoidal rule approximates $f(x)$ as piecewise linear

$$\int_a^b f(x) dx \approx \sum_{i=0}^{n-1} (f(x_{i+1}) + f(x_i)) \frac{b-a}{2n}, \quad x_i = a + (b-a)i/n$$

• Simpson's rule treats $f(x)$ as piecewise quadratic (fits $\cap$ to 3 points)

$$\int_a^b f(x) dx \approx \frac{b-a}{3n} \left(f(x_0) + 4f(x_1) + \underbrace{2f(x_2) + 4f(x_3) + \dots + 4f(x_{n-1})}_{\text{alternates}} + f(x_n) \right)$$

↳ n must be even

↳ fourth-order method

• Higher order polynomial interpolation is not always better.

• Monte Carlo integration can be useful for high dimension integrals

- randomly sample a domain of known area

- count number of points that fall within a target region

- area of target is approx'd by the fraction of points that are within target.

Iterative methods

- Iterative methods all have an iterative step and some termination criterion, e.g. error small enough, time, change between iterations.
- Error can be analysed by substituting our guess back into the eq or finding $|x_n - x_{n-1}|$, though the latter is unreliable.
- Convergence may not be guaranteed, or can be slow.
- **kth order convergence** $\Rightarrow \epsilon_n = M \epsilon_{n-1}^k$ for $M > 0$ as $n \rightarrow \infty$.
 - $\hookrightarrow$ linear convergence is slow
 - $\hookrightarrow$ 2nd order (quadratic) convergence means that num sig. digits doubles at each iteration (i.e. much faster).
- Most common application of iterative methods is root finding.

Bisection

1. Choose a, b such that $\text{sign}(f(a)) \neq \text{sign}(f(b))$
(this is actually quite difficult to do quickly)
 2. Find midpoint $c = (a+b)/2$
 3. IF $|f(c)| < \text{desired_accuracy}$ stop.
 4. IF $\text{sign}(f(c)) = \text{sign}(f(a))$ then $a=c$ else $b=c$
 5. GOTO 2
- Similar to a binary search for the root
 - First order convergence

Fixed point iteration

- Rewrite $f(x)=0$ as $g(x)=x$ for some g . Iteratively compute $x_{n+1} = g(x_n)$. If this converges to r , r is a **fixed point** of g .
- Sufficient condition for convergence
 - $g: I \rightarrow I$ for some interval I
 - $|g'(x)| < 1$ for $\forall x \in I$.

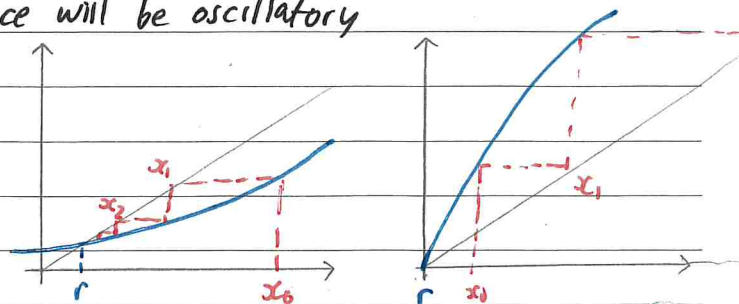
• This is because the error depends on the gradient:

$$x_n = r + \varepsilon_n \Rightarrow r + \varepsilon_{n+1} = g(r + \varepsilon_n) \approx g(r) + \varepsilon_n g'(r) \\ \Rightarrow \varepsilon_{n+1} \approx \varepsilon_n g'(r)$$

- if $0 < g'(x) < 1$, convergence will be monotonic

- if $-1 < g'(x) < 0$, convergence will be oscillatory

• Convergence/divergence can be analysed graphically by bounding with $y=x$



• Sufficient condition for divergence:

if $|g'(r)| > 1$ for all fixed points r , then they are *repelling fixed points* and all (x_n) will diverge (unless $x_0 = r$).

• Thus we must carefully choose $g(x)$ and x_0 to fall under the convergence criterion.

• Fixed point theory can be used to prove that there exist solutions.

Newton-Raphson

• NR is a special case of fixed point iteration: $x_{n+1} = x_n - \frac{f(x_n)}{f'(x_n)}$

↳ error can be analysed by considering the

1st order Taylor expansion of $f(r + \varepsilon_n)$, with $\varepsilon_n = r - x_n$

$$|f(r) - f(x_n)| = |f(r) - f(x_n) - \varepsilon_n f'(x_n)| \leq \frac{M \varepsilon_n^2}{2}$$

Using $\frac{f(x_n)}{f'(x_n)} = \varepsilon_{n+1} - \varepsilon_n$ from NR:

$$|\varepsilon_{n+1}| \leq \frac{M \varepsilon_n^2}{2 |f'(x_n)|}$$

• In practice, x_0 may be found as the root of a crude linear approx to $f(x)$.

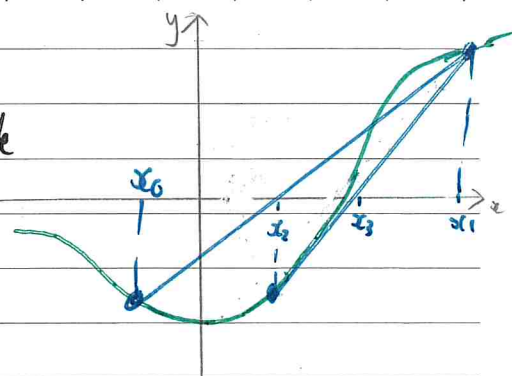
• If $f'(x)$ is badly behaved near $x=r$, e.g. $f(x) = |x|^{0.5}$, then NR may overshoot and thus diverge/oscillate.

• If the ~~real~~ multiplicity of roots > 1 , we must use *modified NR* to preserve quadratic convergence: $x_{n+1} = x_n - m \frac{f(x_n)}{f'(x_n)}$ $\leftarrow m$ is multiplicity.

Secant method

- NR requires $f'(x_n)$. We might approximate this by using a secant between any two points.

$$\therefore x_{n+1} = x_n - f(x_n) \frac{x_n - x_{n-1}}{f(x_n) - f(x_{n-1})}$$



- The convergence is superlinear:

$$\epsilon_{n+1} = \epsilon_n - f(x_n) \frac{\epsilon_n - \epsilon_{n-1}}{f(x_n) - f(x_{n-1})} \approx \epsilon_n \epsilon_{n-1}$$

$$\therefore \epsilon_{n-1}^{p^2-p-1} = M. \quad \epsilon_{n-1} \rightarrow 0 \text{ for convergence}$$

$$\therefore p^2-p-1=0$$

- Thus the order of convergence is ϕ .
- Faster than bisection, but doesn't guarantee convergence unless we ensure that we only construct secants between points of opposite sign (can be done with caching).

Gradient descent

- To find the minimum, we look in the direction of fastest descent

$$\underline{x}_{n+1} = \underline{x}_n + \underset{\substack{\uparrow \\ \text{learning rate}}}{-\gamma_n} \nabla f(\underline{x}_n)$$

- We then repeat this until $\|\nabla f(\underline{x}_n)\| < \epsilon$
- This is a fixed point iteration.
- γ_n can be chosen by an adaptive scheme: choose γ_n that minimises the projection of f onto the plane given by $\nabla f(\underline{x}_n)$
i.e. $\gamma = \arg\min_{\gamma} f(\underline{x}_n - \gamma \nabla f(\underline{x}_n))$
- Gradient descent only finds local minima (unless f is convex)
- Slow convergence when going from side to side down a valley.



Simulated annealing

- A Monte Carlo method for finding the global minimum.
↳ in each iteration, a random candidate is jumped to with some probability.

Start with x_0 and some temperature T_0

1. Randomly select $C_{n+1} = C(x_n, T_n)$, e.g. $C(x, T) = x + \text{uniform}(-1, 1)e^T$
2. $x_{n+1} = C_{n+1}$ with probability $P(f(x_n), f(C_{n+1}), T)$

$$\text{e.g. } P(y, y', T) = \frac{1}{1 + e^{\frac{y' - y}{T}}}$$

3. Decrease temp, e.g. $T_{n+1} = A(T_n)$ with $A(T) = 0.999T$.

- As $T \rightarrow 0$, $P(y, y', T) \rightarrow 0$ if $y' > y$ i.e. over time we are less likely to jump to a higher value.

Linear Systems

- All about solving matrix equations of the form $Ax = b$
- The simplest way is to find A^{-1} , but ~~but~~ this is unstable for large matrices.

Gaussian elimination

- Uses the principle that multiples of rows can be added to achieve an upper-triangular form matrix

for each row $\underline{a}_i$:

for j in range $(0, i)$:

set $a_{i,j}$ to zero by $\underline{a}_i = \underline{a}_i - \frac{a_{i,j}}{a_{j,j}} \underline{a}_j$

- Then back-substitute to find unknowns, which is $O(n^2)$
- $O(n^3)$ in total

- If the **pivot element** is too small, the coefficients may become huge
 $\rightarrow$ rows can be interchanged, so we should choose the one with the largest leading value.
 $\rightarrow$ columns can be permuted but then rows in x need to be as well.

LU factorisation

- If we can write A as the product of lower and upper triangular forms, we can easily solve for x :
 - solve $Lw = b$ using forward sub
 - then solve $Ux = w$ with back sub.
- **Doolittle algorithm**: do a normal Gaussian elimination on A to get U , but keep track of the steps in a new matrix
 - start with $L = \begin{pmatrix} 1 & 0 & 0 \\ 0 & 1 & 0 \\ 0 & 0 & 1 \end{pmatrix}$.
 - e.g. if $R2 \rightarrow R2 - 3R1$ then $L = \begin{pmatrix} 1 & 0 & 0 \\ 3 & 1 & 0 \\ 0 & 0 & 1 \end{pmatrix}$.
 - $O(n^3)$, small overhead on Gauss
- Unlike Gaussian elimination, b is not part of the decomposition. Thus we can factorise $A=LU$ then apply it to a b in $O(n^2)$.
- If $A_{1,1}$ is 0, then algorithm will fail even though a factorisation exists
 $\rightarrow$ first permute the rows, corresponding to **$PA=LU$** .
 $\rightarrow$ any square matrix has an LUP factorisation.

- LU can be used to determine if A has an inverse: A^{-1} exists iff all L/U diagonals are non zero
 $\rightarrow A^{-1} = U^{-1}L^{-1}$, with inverse of Δ matrix easy to compute.

Cholesky factorisation

- The matrix A is **positive definite** if $v^T A v > 0$ for all v :
 - eigenvalues are positive
 - restrictions how far an input vector v can be rotated
 - alternatively, positive semi-definite (≥ 0), negative definite (< 0) ...
- If A is symmetric and positive definite, LU fact. can be optimised by setting $U = L^T$

$$LL^T = \begin{pmatrix} L_{00}^2 & L_{10}^2 + L_{11}^2 & L_{20}^2 + L_{21}^2 + L_{22}^2 \\ L_{10}L_{00} & L_{10}L_{10} + L_{11}L_{11} & L_{10}L_{20} + L_{11}L_{21} + L_{12}L_{22} \\ L_{20}L_{00} & L_{20}L_{10} + L_{21}L_{11} & L_{20}L_{20} + L_{21}L_{21} + L_{22}L_{22} \end{pmatrix} = A$$

SYMMETRIC

- We can then forward substitute to find L_{ij}
- The Cholesky decomp can be used to test for positive definiteness, by seeing if the decomp exists.
- We can find A^{-1} using $A^{-1} = A^T(AA^T)^{-1}$, where AA^T is symmetric

III - conditionedness

- It is important to know how a small input perturbation affects the output.
- The **condition number** is the maximum ratio of the relative change in output to a relative change in input

↳ problem is **ill-conditioned** if it has a high condition number

- Conditionedness is a property of the problem, not algorithm

$$K(x) = \lim_{\epsilon \rightarrow 0} \left| \frac{f(x+\epsilon) - f(x)}{f(x)} \right| / \left| \frac{\epsilon}{x} \right| = \left| \frac{xf'(x)}{f(x)} \right|$$

- For linear systems, we want to know how a vector perturbation of b affects the value of x .

$$K(A) = \max_{x, \epsilon \neq 0} \frac{\|A^T \epsilon\|}{\|x\|} / \frac{\|\epsilon\|}{\|Ax\|} = \max_{\epsilon \neq 0} \frac{\|A^T \epsilon\|}{\|\epsilon\|} \times \max_{x \neq 0} \frac{\|Ax\|}{\|x\|}$$

- When using the Cholesky decomposition, we may end up with negative, due to rounding errors → problem in sqrt. A solution is to add a diagonal matrix:

$$A = LDL^T = \begin{pmatrix} 1 & 0 & 0 \\ L & 1 & 0 \\ L & L & 1 \end{pmatrix} \begin{pmatrix} D & 0 & 0 \\ 0 & D & 0 \\ 0 & 0 & D \end{pmatrix} \begin{pmatrix} 1 & L & L \\ 0 & 1 & L \\ 0 & 0 & 1 \end{pmatrix}$$

↳ LDL^T decomp has the same $O(n^3)$ complexity as LL^T .

QR factorisation

- If we write $A = QR$, with Q being **orthogonal** and R being upper triangular form, we can easily solve. $y = Q^T b$ then $Rx = y$.
 - this works because $Q^{-1} = Q^T$
 - also works when A and R are tall
- To build a set of orthogonal vectors starting from vectors a_k in A , we use the **Gram-Schmidt algorithm**.
- Treat A, Q, R as their column vectors. Gram-Schmidt is like Gaussian elimination but we are subtracting to make orthogonal.

1. Initialise $q_0 = \frac{a_0}{\|a_0\|}$, $r_0 = \begin{pmatrix} \|a_0\| \\ 0 \\ \vdots \end{pmatrix}$

2. Using a_1 , construct a vector orthogonal to q_0 , called w_1 . Then normalise to get q_1 .

$$w_1 = a_1 - (q_0 \cdot a_1)q_0$$

$$\text{then } q_1 = \frac{w_1}{\|w_1\|} \quad r_1 = \begin{pmatrix} q_0 \cdot a_1 \\ \|w_1\| \\ 0 \\ \vdots \end{pmatrix}$$

3. $w_k = a_k - \sum_{i=0}^{k-1} (q_i \cdot a_k)q_i$

$$q_k = \frac{w_k}{\|w_k\|}$$

Then we have $Q = (q_1 | q_2 | \dots | q_m)$

$$R = \begin{pmatrix} a_1 \cdot q_1 & a_2 \cdot q_1 & \dots & a_n \cdot q_1 \\ 0 & a_2 \cdot q_2 & & a_n \cdot q_2 \\ 0 & 0 & \ddots & \vdots \\ 0 & 0 & 0 & a_n \cdot q_m \end{pmatrix}$$

- Small errors in the dot product can accumulate and hinder orthogonality.

Least Squares

- A linear model has the form $y = \beta_0 + \beta_1 x$
- The measurements are then $y_i = \beta_0 + \beta_1 x_i + \epsilon_i$, where ϵ_i are statistical errors (ie diff between measured and true value).
- The **residual** is the difference between the measurement and our model's prediction: $r_i = y_i - (\beta_0 + \beta_1 x_i)$
- We want to find $\hat{\beta}_0$ and $\hat{\beta}_1$ to minimise the **sum of squared residuals**

$$S(\beta_0, \beta_1) = \sum_i r_i^2 = \sum_i (y_i - \beta_0 - \beta_1 x_i)^2$$

$$\frac{\partial S}{\partial \beta_1} = 0 \text{ and } \frac{\partial S}{\partial \beta_0} = 0 \Rightarrow \hat{\beta}_1 = \frac{\sum x y - n \bar{x} \bar{y}}{\sum x^2 - n \bar{x}^2}, \hat{\beta}_0 = \bar{y} - \hat{\beta}_1 \bar{x}$$

- This is **linear least squares** because it is linear in β , so the partial derivatives give linear equations (this is why we use squares)
 $\hookrightarrow$ Generally: $y = \sum_{j=1}^p \beta_j \phi_j(x)$ e.g. $y = \beta_0 + \beta_1 x^2 + \beta_2 x^4 + \dots + \beta_p x^p$.

$\hookrightarrow$ We don't care what ϕ_j is. We can even use different independent variables $x_1, x_2, \dots, x_p$

- Can reformulate with matrices:

observations $\downarrow$

$$\begin{pmatrix} y_1 \\ \vdots \\ y_n \end{pmatrix} = \begin{pmatrix} x_{11} & \dots & x_{p1} \\ \vdots & \ddots & \vdots \\ x_{1n} & \dots & x_{pn} \end{pmatrix} \begin{pmatrix} \beta_1 \\ \vdots \\ \beta_p \end{pmatrix} + \begin{pmatrix} \epsilon_1 \\ \vdots \\ \epsilon_n \end{pmatrix}$$

$$\Rightarrow y = X\beta + \epsilon$$

grad w.r.t $\beta \rightarrow S(\beta) = \|r\|^2 = (y - X\beta)^T (y - X\beta)$

$$\therefore \nabla S(\beta) = \nabla (y^T y - 2\beta^T X^T y + \beta^T X^T X \beta)$$

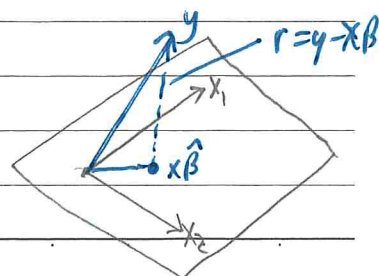
think of $\beta^T \beta$ as β^2
 $\therefore 2X^T X \beta$

$$= 0 - 2X^T y + 2X^T X \beta = 0$$

$$\therefore \hat{\beta} = (X^T X)^{-1} X^T y$$

- This has a geometric interpretation:

- the vectors of X span some subspace.
- we want to get as close to y as possible, so we project
- $r = y - X\beta$ is orthogonal to the subspace
 $\therefore X^T (y - X\beta) = 0$



- $X^T X$ isn't invertible $\Leftrightarrow$ columns linearly dependent
 $\hookrightarrow$ use Cholesky to check. If it succeeds we can easily invert LL^T as $(L^T)^{-1}L^{-1}$.
 $\hookrightarrow$ if singular, some of our variables are independent.
- Rather than explicitly computing $(X^T X)^{-1}$ to solve, we can just solve $(X^T X)\beta = X^T y$ using any of our techniques.

Goodness of fit

- The **coefficient of determination** (R^2) is defined by

$$R^2 = 1 - \frac{\sum_i (y_i - \hat{y}_i)^2}{\sum_i (y_i - \bar{y}_m)^2}$$
 $\leftarrow$ "unexplained" variance of residuals
 $\leftarrow$ total variance
 $\hookrightarrow$ hence R^2 is the explained variance
- Adding a new independent variable cannot decrease R^2 , in the worst case its β will be zero and R^2 unchanged.
 $\hookrightarrow$ we can thus use the **adjusted R^2** : $R^2_{adj} = 1 - (1 - R^2) \frac{n-1}{n-p-1}$

Eigenvalues and Eigenvectors

- Matrices are linear transformations: $A(ax + \beta y) = \alpha Ax + \beta Ay$
- Av will only span the space of v if A is nonsingular.
- Eigenvectors are only scaled (or rotated) when A is applied, $Av = \lambda v$.
- Solving the characteristic polynomial is very ill-conditioned, thus we want better numerical algorithms.
- **Power iteration** can be used to find the **dominant eigenvalue** for a diagonalisable matrix.
 $\hookrightarrow$ starting with any nonzero vector, we repeatedly apply the matrix A : this converges to the eigenvector

$$x^{(k+1)} = \frac{Ax^{(k)}}{\|Ax^{(k)}\|}$$
 $\hookrightarrow$ the corresponding eigenvalue can be found using the **Rayleigh quotient**: $\lambda = \frac{x^T A x}{x^T x}$

• If A is diagonalisable, it can be written as $Q\Lambda Q^{-1}$, with Q being a matrix of eigenvectors and $\Lambda = \text{diag}(\lambda_1, \lambda_2, \dots, \lambda_n)$

$$- A^k = (Q\Lambda Q^{-1})(Q\Lambda Q^{-1}) \dots (Q\Lambda Q^{-1}) = Q\Lambda^k Q^{-1}$$

- our initial random vector may be written as Qx for some x .

$$\therefore A^k(Qx) = Q\Lambda^k x = \sum q_j \lambda_j^k x_j$$

$$\therefore A^k(Qx) = \lambda_1^k \left(\sum q_j \left(\frac{\lambda_j}{\lambda_1}\right)^k x_j \right)$$

- because λ_1 (in our case) is the dominant value, $\lim_{k \rightarrow \infty} \left(\frac{\lambda_j}{\lambda_1}\right)^k = 0$

$\therefore A^k(Qx) \rightarrow \lambda_1^k q_{1,x}$ hence converges to eigenvector.

• Because power iteration only has linear convergence, it may not be suitable for huge matrices (e.g. PageRank)

• A fast iterative method uses QR factorisation instead.

$$A_k = Q_k R_k$$

$$A_{k+1} = Q_k^T A_k Q_k$$

$\hookrightarrow A_k$ converges to upper triangular, hence we can easily find all eigenvalues as the off-diagonal entries.

• For a real symmetric matrix: there is a set of n eigenvectors that form an orthonormal basis.

$\hookrightarrow$ similarly, if we have n unique eigenvalues, their eigenvectors will be orthogonal.

• **Principal component analysis (PCA)** is all about representing a dataset in a basis that maximises variance

- these PCs will be eigenvectors, and because $x^T x$ is real symmetric, they will be orthogonal

- it may be the case that only a few PCs are needed to show most of the variation.

Floating point representation

- A given 8-bit binary number can be interpreted as signed (Two's complement) or unsigned: $1000\ 1011 \rightarrow 2^7 + 2^3 + 2^1 + 2^0 = 139$
 $-2^7 + 2^3 + 2^1 + 2^0 = -117$

- Non-integer values could be represented by a **fixed point**, which reserves some bits for integer & fractional parts

$\rightarrow$ choosing position of decimal point is a tradeoff between prec. and range.

- Floating points** are written as $(-1)^s \times m \times \beta^e$

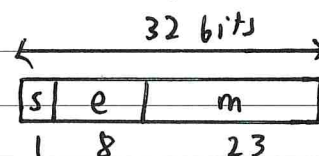
- s is the sign
- m is the **mantissa** (i.e. sig. digits)
- β is the base
- e is the exponent.

- The **IEEE 754** standard specifies constraints:

- because all mantissas start with 1, this is a **hidden bit** (not stored).
- in order for e to be stored unsigned (easier for radix sort), we include a **bias**, i.e. $(-1)^s \times 1.m \times \beta^{2^e - b}$

single precision (32 bits)

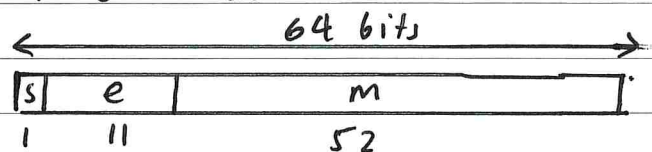
$\rightarrow \text{bias} = 2^{8-1} - 1 = 127$



$p \equiv \text{no. of mantissa bits}$

double precision (64 bits)

$\rightarrow \text{bias} = 2^{11-1} - 1 = 1023$



- For a single prec number, the exponent can be used to ~~represent~~ represent -126 to +127, i.e. e is from 1 to 254 inclusive. This is because $e=0$ and $e=255$ are reserved for special values:

- zero needs a special representation because the hidden bit is 1.
- **denormalised numbers** allow for extra small numbers by allowing m to have leading 0s.
- NaN allows us to recognise when a result isn't useful without aborting the computation.

	$m=0$	$m \neq 0$
$e=0$	zero	denormalised
$e=\text{MAX}$	∞	NaN

- Care must be taken because exact fractions in decimal may be recurring in binary, e.g. $0.1_{10} = 0.00011001100\dots_2$
- IEEE zeroes and infinities are signed, which can be useful to signify where a value started.

IEEE Arithmetic

- Treat operands as precise, then round the result to the nearest IEEE number. If there is a tie, choose the value with 0 LSB (even).
- This rounding (**unbiased**) is the default, but we can also choose rounding towards 0, $+\infty$, $-\infty$, using a global flag.
- For other functions, 'perfect accuracy' is considered 0.5 ulp (unit in the last place)
- Using floats, errors will come from **quantisation** and rounding errors.

Machine epsilon

- **ϵ_m** is defined as the difference between 1.0 and the smallest representable number greater than 1. $\epsilon_m = \beta^{-(p-1)}$
 $\hookrightarrow$ upper bound on the relative error caused by being off by 1 ulp.
- In IEEE, the diff between 1.0 and the greatest number smaller than 1.0 is $\epsilon_m/2$.

Range reduction

- Consider trying to evaluate $\sin x$ for large x using its Maclaurin series
- To avoid loss of significance, we can use modular arithmetic to bring x to $[-\pi, \pi)$. This is **range reduction**.
- However, these large numbers will have a ~~as~~ large absolute error (even though $\eta = \epsilon_m$), hence it will be difficult to know which cycle x falls into unless π is stored to very high precision.