

Numerical Methods

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Errors

- Computations involve errors because of noisy data, floating repr., etc.
- Absolute error: $E_x = |x^* - x| \Rightarrow x^* = x + E_x$ ← x^* is our approx.
- Relative error: $\eta_x = \frac{\Delta x}{|x|} \Rightarrow x^* = x(1 + \eta_x)$
- For $x^* \pm y^*$, $E_{x+y} = E_x + E_y$. No quadrature because these are not distributions.
↳ be careful if subtraction results in near-zero because $\eta \uparrow$ a lot
- The error in $f(x^*)$ is $E_{f(x)} \approx |f'(x^*)| E_x$.
- If $f(x^*)$ is a truncated Taylor series, we have the Lagrange error bound:

$$E_{f(x)} \leq \frac{M(E_x)^{n+1}}{(n+1)!}$$
 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
↳ 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, $f(x) = \sum_{n=0}^{\infty} \frac{f^{(n)}(a)}{n!} (x-a)^n$ ← 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 apprx is then $D_{f'}^+(x) \equiv \frac{f(x+h) - f(x)}{h}$
↳ 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 $D_{f_1}^0(x) = \frac{f(x+h) - f(x-h)}{2h}$

↳ 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 \curvearrowright to 3 points)

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

alternates

↳ 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.
- k th order convergence $\Rightarrow \epsilon_n = M \epsilon_{n-1}^K$ for $M > 0$ as $n \rightarrow \infty$.
 - ↳ linear convergence is slow
 - ↳ 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 $\forall x \in I$.

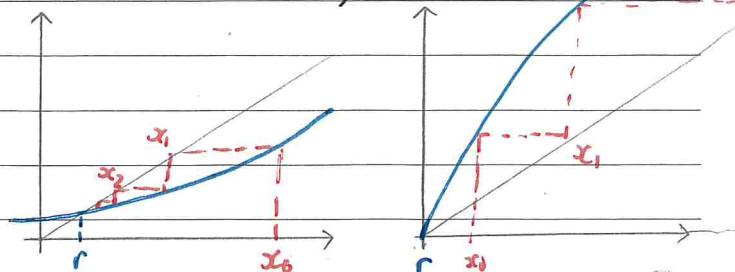
- This is because the error depends on the gradient:

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

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

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

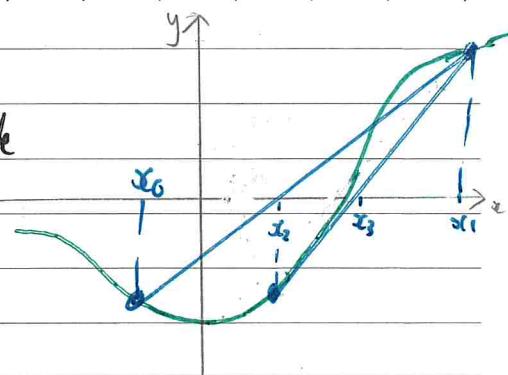
Using $\frac{f(x_n)}{f'(x_n)} = \varepsilon_{n+1} - \varepsilon_n$ from NR:
$$|\varepsilon_{n+1}| \leq \frac{M \varepsilon_n^2}{2 |f'(x_n)|}$$

- In practice, x_0 may be found as the root of a crude linear approx to $F(x)$.
- If the $f'(x)$ 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: $x_{n+1} = x_n - m \frac{f(x_n)}{f'(x_n)}$ ← m is multiplicity.

Secant method

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

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



- The convergence is superlinear:

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

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

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

- 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

$$\tilde{x}_{n+1} = \tilde{x}_n + -\gamma_n \nabla f(\tilde{x}_n)$$

\nwarrow learning rate

- We then repeat this until $\|\nabla f(\tilde{x}_n)\| < \epsilon$
- This is a fixed point iteration.
- γ_n can be chosen by an adaptive scheme: choose γ_n that minimizes the projection of f onto the plane given by $Df(\tilde{x}_n)$
i.e. $\gamma = \arg \min_{\gamma} f(\tilde{x}_n - \gamma \nabla f(\tilde{x}_n))$
- 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 $P(f(c_n), f(c_{n+1}), T)$

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

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 ~~that~~ 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_{ij} to zero by $\underline{a}_i = \underline{a}_i - \frac{a_{ij}}{a_{jj}} \underline{a}_j$

- 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
 - ↳ rows can be interchanged, so we should choose the one with the largest leading value.
 - ↳ columns can be permuted but then rows in Δ need to be as well.

LU factorization

- If we can write A as the product of lower and upper triangular forms, we can easily solve for x :
 - solve $L\bar{u} = b$ using forward sub
 - then solve $Ux = \bar{u}$ 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 $R_2 \rightarrow R_2 - 3R_1$ 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, these algorithm will fail even though a factorization exists
 - ↳ first permute the rows, corresponding to $PA = LU$.
 - ↳ any square matrix has an LUP factorization.
- LU can be used to determine if A has an inverse: A^{-1} exists iff all L/U diagonals are non zero
 - ↳ $A^{-1} = U^{-1}L^{-1}$, with inverse of Δ matrix easy to compute.

Cholesky factorization

- The matrix A is **positive definite** if $v^T A v > 0$ for all v :
 - eigenvalues are positive
 - restricts 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 $U = L^T$

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

SYMMETRIC

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

Ill-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
 \hookrightarrow problem is **ill-conditioned** if it has a high condition number
- Conditionedness is a property of the problem, not algorithm

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

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

$$k(A) = \max_{x, \epsilon \neq 0} \frac{\|A^{-1}\epsilon\|}{\|x\|} / \frac{\|\epsilon\|}{\|Ax\|} = \max_{\epsilon \neq 0} \frac{\|A^{-1}\epsilon\|}{\|\epsilon\|} \times \max_{x \neq 0} \frac{\|Ax\|}{\|x\|}$$

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

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

\hookrightarrow LDL^T decompo 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 full
- To build a set of orthogonal vectors starting from vectors a_0 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 $q_0 = \frac{a_0}{\|a_0\|}$, $r_0 = \begin{pmatrix} \|a_0\| \\ \vdots \end{pmatrix}$

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

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

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

3. $w_n = a_n - \sum_{i=0}^{n-1} (q_i \cdot a_n)q_i$

$$q_n = \frac{w_n}{\|w_n\|}$$

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

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

- 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 minimize the sum of squared residuals

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

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

- This is linear least squares because it is linear in β , so the partial derivatives give linear equations (this is why we use squares)

↳ Generally: $y = \sum_{j=1}^p \beta_j \phi_j(x)$ e.g. $y = \beta_0 + \beta_1 x^2 + \beta_2 x^4 + \dots + \beta_p x^p$.

↳ 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:

$$\underset{\text{observations}}{\downarrow} \begin{pmatrix} y_1 \\ \vdots \\ y_n \end{pmatrix} = \begin{pmatrix} x_{11} & \dots & x_{1p} \\ \vdots & \ddots & \vdots \\ x_{n1} & \dots & x_{np} \end{pmatrix} \begin{pmatrix} \beta_1 \\ \vdots \\ \beta_p \end{pmatrix} + \begin{pmatrix} \epsilon_1 \\ \vdots \\ \epsilon_n \end{pmatrix}$$

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

$$S(\beta) = \|r\|^2 = (y - X\beta)^T (y - X\beta)$$

$$\underset{\text{grad w.r.t } \beta}{\rightarrow} \nabla S(\beta) = \nabla (y^T y - 2\beta^T X^T y + \beta^T X^T X \beta) \quad \begin{matrix} \text{think of } \beta^T X^T X \beta \text{ as } \beta^2 \\ \therefore 2X^T X \beta \end{matrix}$$

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

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

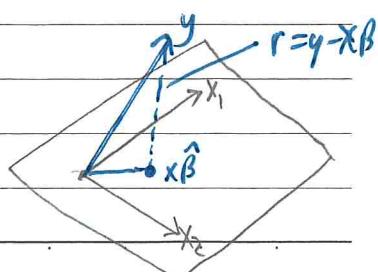
- This has a geometric interpretation:

- 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
 - ↳ use Cholesky to check. If it succeeds we can easily invert LL^T as $(LT)^{-1}L^{-1}$.
 - ↳ 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

 ↳ 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.
 - ↳ we can thus use the adjusted R^2 : $R_{adj}^2 = 1 - (1-R^2) \frac{n-1}{n-p-1}$

Eigenvalues and Eigen vectors

- Matrices are linear transformations: $A(\alpha x + \beta y) = \alpha Ax + \beta Ay$
- $A v$ will only span the space of v if A is nonsingular.
- Eigen vectors 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.
 - ↳ starting with any nonzero vector, we repeatedly apply the matrix A : this converges to the eigenvector $x^{(k+1)} = \frac{Ax^{(k)}}{\|Ax^{(k)}\|}$
 - ↳ 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, $\lim_{k \rightarrow \infty} \left(\frac{\lambda_j}{\lambda_1} \right)^k = 0$
- $$\therefore A^k(Qx) \rightarrow \lambda_1^k q_1 x_1 \text{ hence converges to eigenvector.}$$

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

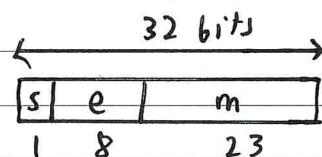
- A given 8-bit binary number can be interpreted as signed (Two's complement) or unsigned: $1000\ 1011 \rightarrow 2^7 + 2^3 + 2^1 + 2^0 = 139$
 $-2^7 + 2^3 + 2^1 + 2^0 = -117$
- Non-integer values could be represented by a **fixed point**, which reserves some bits for integer & fractional parts
 ↳ choosing 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 $(-1)^s \times 1.m \times \beta^{e-b}$

single precision (32 bits)

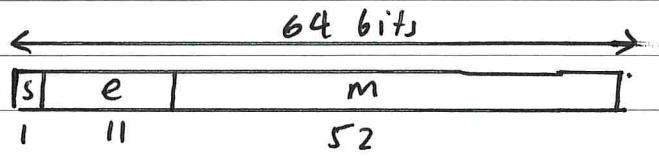
$$\hookrightarrow \text{bias} = 2^{8-1} - 1 = 127$$



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

double precision (64 bits)

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



- For a single prec number, the exponent can be used to 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.

$m=0 \quad m \neq 0$

$e=0$	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 ∞ 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.