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JTMS-MAT-13: Numerical Methods

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Notes for numerical methods course 2024. This can be downloaded from: https://djps.github.io/courses/numericalmethods24/notes

Note that proofs for theorems marked with * were presented in class.

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Recommended Reading

- R. L. Burden and J. D. Faires "Numerical Analysis", Brooks/Cole $9^{\rm th}$ Edition (2011).

1 Taylor Series

The Taylor series, or the Taylor expansion of a function, is defined as

Definition 1.1 (Taylor Series). For a function $f : \mathbb{R} \to \mathbb{R}$ which is infinitely differentiable at a point c, the Taylor series of f(c) is given by

$$\sum_{k=0}^{\infty} \frac{f^{(k)}(c)}{k!} (x-c)^k.$$

This is a power series, which is convergent for some radius.

Theorem 1 (Taylor's Theorem). For a function $f \in C^{n+1}([a,b])$, i.e. f is (n+1)-times continuously differentiable in the interval [a,b], then for some c in the interval, the function can be written as

$$f(x) = \sum_{k=0}^{n} \frac{f^{(k)}(c)}{k!} (x - c)^{k} + \frac{f^{(n+1)}(\xi)}{(n+1)!} (x - c)^{n+1}$$

for some value $\xi \in [a,b]$ where

$$\lim_{\xi \to c} \frac{f^{(n+1)}(\xi)}{(n+1)!} (x-c)^{n+1} = 0.$$

2 **Errors**

Definition 2.1 (Absolute and Relative Errors).

Let \tilde{a} be an approximation to a, then the **absolute error** is given by

$$|\tilde{a} - a|$$

If $|a| \neq 0$, the **relative error** may be given by

$$\left|\frac{\tilde{a}-a}{a}\right|$$

The error bound is the magnitude of the admissible error.

Theorem 2. For both addition and subtraction the bounds for the absolute error are added. In division and multiplication the bounds for the relative errors are added.

Definition 2.2 (Linear Sensitivity to Uncertainties).

If y(x) is a smooth function, i.e. is differentiable, then |y'| can be interpreted as the **linear sensitivity** of y(x) to uncertainties in x.

For functions of several variables, i.e. $f: \mathbb{R}^n \to \mathbb{R}$, then

$$|\Delta y| \le \sum_{i=1}^{n} \left| \frac{\partial y}{\partial x_i} \right| |\Delta x_i|$$

 $|\Delta y| \leq \sum_{i=1}^{n} \left| \frac{\partial y}{\partial x_i} \right| |\Delta x_i|$ where $|\Delta x_i| = |\tilde{x}_i - x_i|$ for an approximation \tilde{x}_i , $|\Delta y_i| = |\tilde{y}_i - y_i| = |f(\tilde{x}_i) - (x_i)|$. thus

3 Number Representations

Definition 3.1 (Base Representation).

Let $b \in \mathbb{N} \setminus \{1\}$. Every number $x \in \mathbb{N}_0$ can be written as a unique expansion with respect to base b as

$$(x)_b = a_0 b^0 + a_1 b^1 + \ldots + a_n b^n = \sum_{i=0}^n a_i b^i$$

A number can be written in a nested form:

$$(x)_b = a_0 b^0 + a_1 b^1 + \dots + a_n b^n$$

= $a_0 + b (a_1 + b (a_2 + b (a_3 + \dots + b a_n) \dots)$

with $a_i < \mathbb{N}_0$ and $a_i < b$, i.e. $a_i \in \{0, \dots, b-1\}$. For a real number, $x \in \mathbb{R}$, write

$$x = \sum_{i=0}^{n} a_i b^i + \sum_{i=1}^{\infty} \alpha_i b^{-i}$$
$$= a_n \dots a_0 \cdot \alpha_1 \alpha_2 \dots$$

Algorithm (Euclid).

Euclid's algorithm can convert number x in base 10, i.e. $(x)_{10}$ into another base, b, i.e. $(x)_b$.

- 1. Input $(x)_{10}$
- 2. Determine the smallest integer n such that $x < b^{n+1}$
- 3. Let y = x. Then for $i = n, \dots, 0$

$$\begin{array}{rcl}
a_i & = & y \operatorname{div} b^i \\
y & = & y \operatorname{mod} b^i
\end{array}$$

which at each steps provides an a_i and updates y.

4. Output as $(x)_b = a_n a_{n-1} \cdots a_0$

Algorithm (Horner).

- 1. Input $(x)_{10}$
- 2. Set i = 0

3. Let y = x. Then while y > 0

$$a_i = y \operatorname{div} b$$

$$y = y \operatorname{mod} b$$

$$i = i + 1$$

which at each steps provides an a_i and updates y.

4. Output as $(x)_b = a_n a_{n-1} \cdots a_0$

Definition 3.2 (Normalized Floating Point Representations).

Normalized floating point representations with respect to some base b, store a number x as

$$x = 0 \cdot a_1 \dots a_k \times b^n$$

where the $a_i \in \{0, 1, \dots b-1\}$ are called the **digits**, k is the **precision** and n is the **exponent**. The set a_1, \dots, a_k is called the **mantissa**. Impose that $a_1 \neq 0$, it makes the representation unique.

Theorem 3. Let x and y be two normalized floating point numbers with x > y > 0 and base b = 2. If there exists integers p and $q \in \mathbb{N}_0$ such that

$$2^{-p} \le 1 - \frac{y}{x} \le 2^{-q}$$

then, at most p and at least q significant bits (i.e. significant figures written in base 2) are lost during subtraction.

4 Linear Systems

Definition 4.1 (Systems of Linear Equations). A system of linear equations (or a linear system) is a collection of one or more linear equations involving the same variables. If there are m equations with n unknown variables to solve for, i.e.

$$a_{1,1}x_1 + a_{1,2}x_2 + \dots + a_{1,n}x_n = b_1$$

$$a_{2,1}x_1 + a_{2,2}x_2 + \dots + a_{2,n}x_n = b_2$$

$$\vdots$$

$$a_{m,1}x_1 + a_{m,2}x_2 + \dots + a_{m,n}x_n = b_m$$

then the system of linear equations can be written in matrix form Ax = b, where

$$A = \begin{pmatrix} a_{1,1} & a_{1,2} & \cdots & a_{1,n} \\ a_{2,1} & a_{2,2} & \cdots & a_{2,n} \\ \vdots & \vdots & \ddots & \vdots \\ a_{m,1} & a_{m,2} & \cdots & a_{m,n} \end{pmatrix}, \quad x = \begin{pmatrix} x_1 \\ x_2 \\ \vdots \\ x_n \end{pmatrix}, \quad \text{and} \quad b = \begin{pmatrix} b_1 \\ b_2 \\ \vdots \\ b_m \end{pmatrix},$$

so that $A \in \mathbb{R}^{m \times n}$, $x \in \mathbb{R}^n$ and $b \in \mathbb{R}^m$.

Definition 4.2 (Banded Systems). A banded matrix is a matrix whose non-zero entries are confined to a diagonal band, comprising the main diagonal and zero or more diagonals on either side.

Definition 4.3 (Symmetric & Hermitian Systems). A square matrix A is symmetric if $A = A^T$, that is, $a_{i,j} = a_{j,i}$ for all indices i and j.

A square matrix is said to be **Hermitian** if the matrix is equal to its conjugate transpose, i.e. $a_{i,j} = \overline{a_{j,i}}$ for all indices i and j. A Hermitian matrix is written as A^H .

Definition 4.4 (Positive Definite Matrices). A matrix, M, is said to be **positive definite** if it is symmetric (or Hermitian) and all its eigenvalues are real and positive.

Definition 4.5 (Nonsingular Matrices). A matrix is **non-singular** or **invertible** if there exists a matrix A^{-1} such that $A^{-1}A = AA^{-1} = I$, where I is the identity matrix.

Remarks (Properties of Nonsingular Matrices). For a nonsingular matrix, the following all hold:

- Nonsingular matrix has full rank
- A square matrix is nonsingular if and only if the determinant of the matrix is non-zero.
- If a matrix is singular, both versions of Gaussian elimination will fail due to division by zero, yielding a floating exception error.

Definition 4.6. If \tilde{x} is an approximate solution to the linear problem Ax = b, then the **residual** is defined as $r = A\tilde{x} - b$.

If the magnitude of the residual, |r|, is large due to rounding, the matrix is said to be **ill-conditioned**.

4.1 Direct Methods

Algorithm (Gaussian Elimination).

Gaussian elimination is a method to solve systems of linear equations based on forward elimination (a series of row-wise operations) to convert the matrix, A, to upper triangular form (echelon form), and then back-substitution to solve the system. The row operations are:

- row swapping
- row scaling, i.e. multiplying by a non-zero scalar
- row addition, i.e. adding a multiple of one row to another

```
1: procedure Forward Elimination
 2:
        for k = 1 to n - 1 do
            for i = k + 1 to n do
 3:
                for j = k to n do
 4:
                    a_{i,j} = a_{i,j} - \frac{a_{i,k}}{a_{k,k}} a_{k,j}
               end for b_i = b_i - \frac{a_{i,k}}{a_{k,k}} b_k
 6:
 7:
            end for
 8:
 9:
        end for
10: end procedure
11: procedure Back Substitution
        x_n = \frac{b_n}{a}
12:
        for i = n - 1 to 1 do
13:
14:
            y = b_i
            for j = n to i + 1 do
15:
16:
               y = y - a_{i,j} x_j
17:
18:
        end for
19:
20: end procedure
```

Algorithm (Gaussian Elimination with Scaled Partial Pivoting). A pivot element is the element of a matrix which is selected first to do certain calculations. Pivoting helps reduce errors due to rounding during forward elimination.

To use partial pivoting to produce a matrix in row-echelon form

```
maximal absolute
 1: Find
                                       values
                                                  vector
                                                                 with
                                                                          entries
   s_i = \max j = 1, \dots, n |a_{i,j}|
2: for k = 1 to n - 1 do
       for i = k to n do
3:
           Compute \left| \frac{a_{i,k}}{s_i} \right|
4:
       end for
5:
       Find row with largest relative pivot element, denote this as row j
6:
       Swap rows k and j in the matrix A
7:
       Swap entries k and j in the vector s
8:
       Do forward elimination on row k
10: end for
```

Definition 4.7 (Upper and Lower Triangular Matrices). A square matrix is said to be a **lower triangular matrix** if all the elements above the main diagonal are zero and an **upper triangular** if all the entries below the main diagonal are zero.

Theorem 4 (LU-Decomposition). Let $A \in \mathbb{R}^{n \times n}$ be invertible. Then there exists a decomposition of A such that A = LU, where L is a lower triangular matrix and U is an upper triangular matrix, and

$$L = U_1^{-1} U_2^{-1} \cdots U_{n-1}^{-1}$$

where each matrix U_i is a matrix which describes the i^{th} step in forward elimination part of Gaussian elimination. The upper triangular matrix U is given by

$$U = U_{n-1} \cdots U_2 U_1 A.$$

Definition 4.8 (Cholesky-Decomposition). A symmetric, positive definite matrix can be decomposed as $A = \tilde{L}\tilde{L}^T$, where $\tilde{L} = LD^{1/2}$, where D is a diagonal matrix whose elements d_i are all positive, so that $D^{1/2}$ has elements $\sqrt{d_i}$.

Algorithm (Cholesky-Decomposition). Given a matrix A, the lower triangular matrix \tilde{L} can be constructed via

```
1: for i=1 to n do
        for j = 1 to i - 1 do
 2:
 3:
            y = a_{i,j}
            for k = 1 to j-1 do
 4:
 5:
                y = y - l_{i,k} l_{j,k}
            end for
 6:
            l_{i,j} = y/l_{j,j}
 7:
 8:
9:
        y = a_{i,i}
        for k=1 to i-1 do
10:
11:
            y = y - l_{i,k} l_{i,k}
        end for
12:
        if y \leq 0 then
13:
            there is no solution
14:
15:
16:
            l_{i,i} = \sqrt{y}
        end if
17:
18: end for
```

4.2 Indirect Methods

For a non-singular matrix A, consider the iterative scheme $Qx_{k+1} = (Q - A)x_k + b$. This is equivalent to

$$x_{k+1} = (I - Q^{-1}A) x_k + Q^{-1}b.$$

Definition 4.9 (Spectral Radius of a Matrix). The **spectral radius** of a matrix A is defined as

$$\rho(A) = \max\{|\lambda_1|, |\lambda_2|, \dots |\lambda_n|\}$$

where the λ_i are the eigenvalues of the matrix.

Theorem 5 (Convergence). The iterative scheme converges if and only if the spectral radius of the matrix $I-Q^{-1}A$ is less than one, i.e. $\rho\left(I-Q^{-1}A\right)<1$

Definition 4.10. If a sequence x_n converges to r as $n \to \infty$, then it is said to **converge linearly** if there exists a $\mu \in (0,1)$ such that

$$\lim_{n \to \infty} \frac{|x_{n+1} - r|}{|x_n - r|} = \mu.$$

The sequences converges super-linearly if

$$\lim_{n \to \infty} \frac{|x_{n+1} - r|}{|x_n - r|} = 0$$

and sub-linearly if

$$\lim_{n \to \infty} \frac{|x_{n+1} - r|}{|x_n - r|} = 1.$$

More generally, a sequence converges with order q if there exists a $\mu>0$ such that

$$\lim_{n\to\infty} \frac{|x_{n+1}-r|}{|x_n-r|^q} = \mu.$$

Thus a sequence is said to converge quadratically when q=2 and exhibit cubic convergence when q=3.

Definition 4.11 (Richardson Iteration). Let Q = I. Then **Richardson iteration** computes the sequence of vectors

$$x_{k+1} = (I - A) x_k + b.$$

This may converge, depending on A.

The modified Richardson iteration scales $Q = \omega I$, so that

$$x_{k+1} = x_k + \omega \left(b - Ax_k \right).$$

Definition 4.12 (Jacobi Iteration). The **Jacobi iteration** has Q = D, so

$$x_{k+1} = (I - D^{-1}A) x_k + D^{-1}b.$$

Definition 4.13 (Diagonally Dominant Matrices). A matrix $A \in \mathbb{R}^{n \times n}$ is said to be **diagonally dominant** if, for every row, the absolute value of the diagonal element is greater or equal to the sum of the magnitudes of all other elements, i.e.

$$||a_{i,i}|| \ge \sum_{j=1, j \ne i}^{n} ||a_{i,j}||$$
 for all $i \in (1, n)$.

Theorem 6 (Convergence of Jacobi Scheme). If a matrix A is diagonally dominant, then the Jacobi scheme converges for any initial guess x_0 .

Definition 4.14 (Gauss-Seidel Scheme). Let Q = L + D, then the **Gauss-Seidel** scheme is given by

$$(D+L)x^{(n+1)} = -Ux^{(n)} + b.$$

Theorem 7 (Convergence of Gauss-Seidel).

If a matrix A is diagonally dominant, then the Gauss-Seidel scheme converges for any initial guess x_0 .

Definition 4.15 (Successive Over Relaxation). The scheme uses $Q = L + \frac{1}{\omega}D$, thus

$$(D + \omega L)x^{(n+1)} = -((\omega - 1) D + \omega U) x^{(n)} + \omega b.$$

Theorem 8 (Convergence of Successive Over Relaxation).

Let A be a symmetric matrix with positive entries on the diagonal and let $\omega \in (0,2)$. Then, if and only if A is positive definite will the method of successive over relaxation converge.

5 Nonlinear Solvers

5.1 Bisection Method

Definition 5.1 (Bisection Method).

The bisection method, when applied in the interval [a,b] to a function $f \in C^0([a,b])$ with f(a)f(b) < 0

Bisect the interval into two subintervals [a, c] and [c, b] such that a < c < b.

- If f(c) = 0 or is sufficiently close, then c is a root,
- else, if f(c)f(a) < 0 continue in the interval [a, c],
- else, if f(c)f(b) < 0 continue in the interval [c, b].

Theorem 9 (Bisection Method).

The bisection method, when applied in the interval [a, b] to a function $f \in C^0([a, b])$ with f(a)f(b) < 0 will compute, after n steps, an approximation c_n of the root r with error

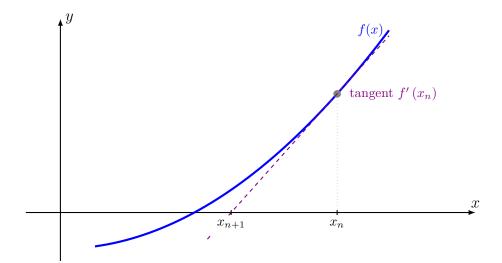
$$|r - c_n| < \frac{b - a}{2n}.$$

5.2 Newton's Method

Definition 5.2.

Let a function $f \in C^1([a,b])$, then for an initial guess x_0 , Newton's method is

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



Theorem 10.

When Newton's method converges, it converges to a root, r, of f(x), i.e. $f\left(r\right)=0.$

Theorem 11.

Let $f \in C^1([a, b])$, with

- 1. f(a)f(b) < 0,
- 2. $f'(x) \neq 0$ for all $x \in (a, b)$,
- 3. f''(x) exists, is continuous and either f''(x) > 0 or f''(x) < 0 for all $x \in (a,b)$.

Then f(x) = 0 has exactly one root, r, in the interval and the sequence generated by Newton iterations converges to the root when the initial guess is chosen according to

• if f(a) < 0 and f''(a) < 0 or f(a) > 0 and f''(a) > 0 then $x \in [a, r]$

or

• if f(a) < 0 and f''(a) > 0 or f(a) > 0 and f''(a) < 0 then $x \in [r, b]$.

The iterate in the sequence satisfies

$$|x_n - r| < \frac{f(x_n)}{\min_{x \in [a,b]} |f'(x)|}.$$

Theorem 12.

- Theorem 12. Let $f \in C^1([a,b])$, with 1. f(a)f(b) < 0, 2. $f'(x) \neq 0$ for all $x \in (a,b)$, 3. f''(x) exists and is continuous, i.e. $f(x) \in C^2([a,b])$.

Then, if x_0 is close enough to the root r, Newton's method converges quad-

Secant Methods 5.3

Definition 5.3.

The secant method is defined as

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

Theorem 13.

Let $f \in C^{2}([a,b])$, and $r \in (a,b)$ such that f(r) = 0 and $f'(r) \neq 0$. Fur-

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

Then there exists a $\delta > 0$ such that when $|r - x_0| < \delta$ and $|r - x_1| < \delta$, then the following holds:

1.
$$\lim_{n \to \infty} |r - x_n| = 0 \Leftrightarrow \lim_{n \to \infty} x_n = r$$
,

1.
$$\lim_{n \to \infty} |r - x_n| = 0 \Leftrightarrow \lim_{n \to \infty} x_n = r,$$

2. $|r - x_{n+1}| \le \mu |r - x_n|^{\alpha}$ with $\alpha = \frac{1 + \sqrt{5}}{2}$.

Convergence 5.4

Method	Regularity	Proximity to r	Init. points	Func. calls	Convergence
Bisection	C^0	No	2	1	Linear
Newton	\mathcal{C}^2	Yes	1	2	Quadratic
Secant	\mathcal{C}^2	Yes	2	1	Superlinear

5.5 Systems of Nonlinear Equation

Definition 5.4 (Multi-Dimensional Newton Method). For a vector-valued function $f: \mathbb{R}^n \to \mathbb{R}^n$, which takes as an argument the vector

$$x = (x_1, x_2 \dots x_n) \in \mathbb{R}^n,$$

the **Jacobian** matrix is defined as

$$J = \begin{pmatrix} \frac{\partial f_1}{\partial x_1} & \frac{\partial f_1}{\partial x_2} & \cdots & \frac{\partial f_1}{\partial x_n} \\ \frac{\partial f_2}{\partial x_1} & & & \vdots \\ \vdots & & & \vdots \\ \frac{\partial f_n}{\partial x_1} & \frac{\partial f_n}{\partial x_2} & \cdots & \frac{\partial f_n}{\partial x_n} \end{pmatrix} \in \mathbb{R}^{n \times n}$$

where f_1 is the first component of the vector-valued function f.

If the derivatives are evaluated at the vector x, the Jacobian matrix can be parameterised as J(x). Newton's method can then be written as a vector equation,

$$x_{m+1} = x_m - J^{-1}(x_m) f(x_m)$$

where $J^{-1}(x_m)$ is the inverse of the Jacobian matrix evaluated at the m-iterate of the approximation vector which is denoted by x_m .

In practice, as matrix inversion can be computationally expensive, the system

$$J(x_n)(x_{n+1} - x_n) = -f(x_n)$$

is solved for the unknown vector $x_{m+1} - x_m$, and then x_{m+1} is found.

6 Interpolation

Definition 6.1. Given as set of points $p_0, \ldots, p_n \in \mathbb{R}$ and corresponding nodes $u_0, \ldots, u_n \in \mathbb{R}$, a function $f : \mathbb{R} \to \mathbb{R}$ with $f(u_i) = p_i$ is an **interpolating function**.

This can be generalised to higher dimensions, i.e. $f: \mathbb{R} \to \mathbb{R}^N$.

Definition 6.2. If the interpolating function is a polynomial, it can be written as

$$p(u) = \sum_{i=0}^{n} \alpha_i \varphi_i(u).$$

So that for every j, the polynomial satisfies $p(u_j) = \sum_{i=0}^{n} \alpha_i \varphi_i(u_j)$, for weights α_j . Thus, solving for values of α which fit the interpolating function to the data, leads to a linear system of the form

$$\Phi \alpha = p$$

where p is the vector defined the polynomial evaluated at the node points, i.e. $p = p(u_j)$ and Φ is the **collocation matrix**, given by

$$\Phi = \begin{pmatrix} \varphi_0(u_0) & \varphi_1(u_1) & \cdots & \varphi_n(u_n) \\ \vdots & & & \vdots \\ \varphi_0(u_n) & \cdots & \cdots & \varphi_n(u_n) \end{pmatrix}$$

Thus $\alpha = \Phi^{-1}p$.

The collocation matrix is invertible if and only if the set of functions φ are linearly independent.

When the $\varphi_i(u) = u^i$, then Φ is the **Vandermonde matrix**.

Definition 6.3 (Lagrange Polynomials). The Lagrange form of an interpolating polynomial is given by

$$p(x) = \sum_{i=0}^{n} \alpha_i l_i(x)$$

where $l_i \in \mathbb{P}_n$ are such that $l_i(x_j) = \delta_{ij}$. The polynomials $l_i(x) \in \mathbb{P}_n$ for $i = 0, \ldots, n$, are called **characteristic polynomials** and are given by

$$l_{i}(x) = \prod_{j=0, j \neq i}^{n} \frac{x - x_{j}}{x_{i} - x_{j}}.$$

Algorithm (Aitken's Algorithm).

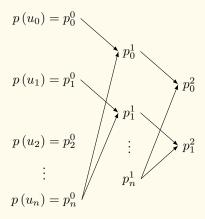
Aitken's algorithm is an iterative process for evaluating Lagrange interpolation polynomials at an arbitrary point, u^* , without explicitly constructing them. If the interpolating polynomial is given by p, and is derived from n data points (u_i, y_i) for $i = 0, \ldots, n$

$$p(u) = \sum_{i=0}^{n} p_i^n l_i^n(u).$$

The interpolation is achieved by constructing a series of polynomials, evaluated at the $u = u^*$, where $p_i^k(u)$ is given by

$$p_{i}^{k+1}(u) = p_{i}^{k}(u) \left(\frac{u - u_{n-k}}{u_{i} - u_{n-k}}\right) + p_{n-k}^{k}(u) \left(1 - \frac{u - u_{n-k}}{u_{i} - u_{n-k}}\right)$$

with initial values $p_i^0 = y_i$, i.e.



where the coefficients are evaluated from left to right, until p_0^n is evaluated.

6.1 Piecewise Polynomial Interpolation

Definition 6.4. A function s(u) is called a **spline** of degree k on the domain [a,b] if $s \in C^{k-1}([a,b])$ and there exists nodes $a = u_0 < u_1 < \ldots < u_m = b$ such that s is a polynomial of degree k for $i = 0, \ldots m-1$.

Definition 6.5 (B-Splines). A spline is said to be a **b-spline** if it is of the form

$$s\left(u\right) = \sum_{i=0}^{m} \alpha_{i} \mathcal{N}_{i}^{n}\left(u\right)$$

where \mathcal{N}^n are the **basis spline functions** of degree n with minimal support. (That is they are positive in the domain and zero outside). The functions are defined recursively. Let u_i be the set of nodes u_0, u_1, \ldots, u_m , then

$$\mathcal{N}_{i}^{0}(u) = \begin{cases} 1 & \text{for } u_{i} \leq u \leq u_{i+1} \\ 0 & \text{else.} \end{cases}$$

$$\mathcal{N}_{i}^{n}(u) = \alpha_{i}^{n-1}(u)\mathcal{N}_{i}^{n-1}(u) + \left(1 - \alpha_{i+1}^{n-1}(u)\right)\mathcal{N}_{i+1}^{n-1}(u)$$

$$\alpha_{i}^{n-1}(u) = \frac{u - u_{i}}{u_{i+n} - u_{i}}$$

$$\alpha_i^{n-1}(u) = \frac{u - u_i}{u_{i+n} - u_i}$$

Given data with nodes u_i and values p_i , to interpolate with splines, of order n, requires solving

Find
$$s = \sum_{i=0}^{m} \alpha_i \mathcal{N}_i^n(u)$$
 such that $s(u_i) = p_i$ for $i = 0, \dots, m$

which is matrix form is $\Phi \alpha = p$, where the collocation matrix, $\Phi \in \mathbb{R}^{(m+1)\times (m+1)}$ is given by

$$\Phi = \begin{pmatrix} \mathcal{N}_0^n (u_0) & \cdots & \mathcal{N}_m^n (u_0) \\ \vdots & & \vdots \\ \mathcal{N}_0^n (u_m) & \cdots & \mathcal{N}_m^n (u_m) \end{pmatrix}.$$

6.2**Least-Squares Approximation**

Definition 6.6 (Least-Squares Approximation). Given a set of points $y = (y_0, y_1, \dots y_n)$ at nodes x_i , seek a continuous function of x, with a given form characterized by m parameters $\beta = (\beta_0, \beta_1, \dots, \beta_m)$, i.e. $f(x, \beta)$, which approximates the points while minimizing the error, defined by the sum of

$$E = \sum_{i=0}^{n} (y - f(x_i, \beta))^{2}.$$

The minimum is found when

$$\frac{\partial E}{\partial \beta_j} = 0$$
 for all $j = 1, \dots m$

$$-2\sum_{i=0}^{n} (y_i - f(x_i, \beta_j)) \frac{\partial f(x_i, \beta)}{\partial \beta_j} = 0 \quad \text{for all} \quad j = 1, \dots m.$$

Definition 6.7 (Linear Least-Squares Approximation). If the function f is a function of the form

$$y = \sum_{j=1}^{m} \beta_j \varphi_j(x)$$

then the least squares problem can be expressed as

$$\frac{\partial E}{\partial \beta_{j}} = \sum_{j=1}^{m} \left(\sum_{i=1}^{n} \varphi_{j} (x_{i}) \varphi_{k} (x_{i}) \right).$$

Thus, the weights β can be determined by solving the linear system,

$$\Phi\Phi^T\beta = \Phi y$$

i.e. $\beta=\left(\Phi\Phi^T\right)^{-1}\Phi y$, where Φ is the collocation matrix.

7 Numerical Differentiation

Definition 7.1 (Finite-Difference Quotients). Consider the approximations to the first-order derivative:

1. Forward Difference Quotient:

$$D_j^+ u = \frac{u_{j+1} - u_j}{h}$$

2. Backwards Difference Quotient:

$$D_j^- u = \frac{u_j - u_{j-1}}{h}$$

3. Central Difference Quotient:

$$D_j^0 u = \frac{u_{j+1} - u_{j-1}}{2h}$$

The forward and backwards difference schemes are first order approximations to the derivative. The central difference scheme is a second order accurate approximation.

Definition 7.2 (Richardson Extrapolation). This is a method for deriving higher order approximations for derivatives from lower order approximations. Consider a first-order approximation to the derivative, $\varphi(h)$, such as backwards or forwards differencing, then

$$f'(x) = \varphi(h) + a_2h^2 + a_3h^3 + \dots$$

Now evaluate the derivative at h = h/2, so that

$$f'(x) = \varphi\left(\frac{h}{2}\right) + a_2\left(\frac{h}{2}\right)^2 + a_3\left(\frac{h}{2}\right)^3 + \dots$$

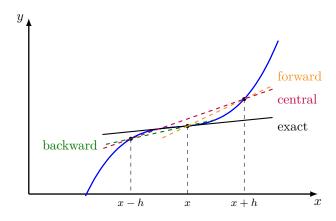
Combining the two terms so that the low order term cancel, i.e. via f'(x) - 4f'(x), then a better approximation can be found as

$$f'(x) = \varphi(h) - 4\varphi\left(\frac{h}{2}\right) + \mathcal{O}(h^3).$$

The process can also be applied to second order accurate schemes, such as central differencing, to produce more accurate approximations, as well as to higher order derivatives.

Definition 7.3 (Higher Order Derivatives). From the Taylor expansion of f(x+h) and f(x-h), the second order derivative can be expressed as

$$f''(x) = \frac{f(x-h) - 2f(x) + f(x+h)}{h^2} + \mathcal{O}\left(h^2\right).$$



8 Numerical Integration

Definition 8.1 (Riemann Sum). Create a partition, p, of the domain of integration: define n+1 nodes $a=x_0 < x_1 < \ldots < x_n = b$, so that there are n sub-intervals $[x_i, x_{i+1}]$. Then approximate the area under the curve by summing the areas in each subinterval defined as

$$\int_{a}^{b} f(x) dx \approx \sum_{i=0}^{n-1} (x_{i+1} - x_i) f(x^*), \quad x^* \in [x_i, x_{i+1}].$$

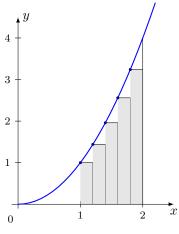
If f is continuous, the value of x_i^* may be chosen arbitrarily in the interval $[x_i, x_{i+1}]$. Then the **lower** and **upper Riemann sums** are given by

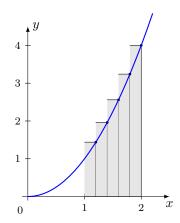
$$L(f, p) = \sum_{i=0}^{n-1} (x_{i+1} - x_i) m_i$$
 where $m_i = \min_{x \in [x_i, x_{i+1}]} f(x)$

$$U(f,p) = \sum_{i=0}^{n-1} (x_{i+1} - x_i) M_i \text{ where } M_i = \max_{x \in [x_i, x_{i+1}]} f(x)$$

so that bounds for the value of the quadrature can be made

$$L(f,p) \le \int_{a}^{b} f(x) dx \le U(f,p).$$





(a) Lower Riemann Sum

(b) Upper Riemann Sum

Additionally, the **left** and **right Riemann sums** are given by

$$\sum_{i=0}^{n-1} (x_{i+1} - x_i) f(x_{i-1}),$$

$$\sum_{i=0}^{n-1} (x_{i+1} - x_i) f(x_i).$$

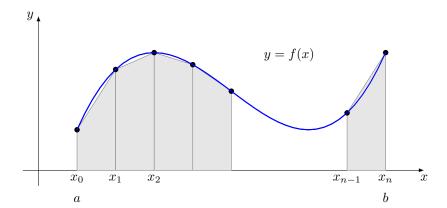
Definition 8.2 (Trapezoidal Rule). Rather than rectangles, use trapezoids to approximate the integral in a sub domain

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

If the nodes of the partition are equally spaced, so that $h = x_{i+1} - x_i$, then the formula can be given by

$$T(f,p) = \frac{h}{2} (f(x_0) + f(x_n)) + h \sum_{i=1}^{n-2} f(x_i).$$

Theorem 14 (Error for Trapezoidal Rule). Let $f \in C^2([a,b])$ and p be equidistant partition of [a,b], with $h=x_{i+1}-x_i$. The error can be shown



to have the form:

$$\left| \int_{a}^{b} f(x) \, dx - T(f, p) \right| = a_{2}h^{2} + a_{4}h^{4} + \dots$$

that is, the error terms are even powers of the discretization. The error for the trapezium rule is

$$\left| \int_{a}^{b} f(x) dx - T(f, p) \right| = \frac{1}{12} \left| (b - a) h^{2} f''(\xi) \right|$$

for a $\xi \in (a, b)$.

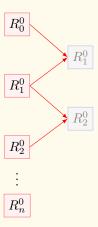
Definition 8.3 (Simpson's Rule). The integral is approximated as

$$\int_{a}^{b} f(x) dx \approx \frac{b-a}{6} \left(f(a) + 4f\left(\frac{a+b}{2}\right) + f(b) \right).$$

Simpson's rule uses quadratic interpolation. It can be applied in a composite manner, i.e. on many subdomains. It has an asymptotic error of $\mathcal{O}(h^4)$.

Algorithm (Romberg Algorithm). Romberg's method uses the Trapezoidal Rule 8.2 and then Richardson Extrapolation 7.2 to estimate integrals. First consider a sequence of partitions, p_i , of equal spacing given by $h_i = \frac{b-a}{2^i}$ for $i=0,\ldots,n$, which yield a sequence of integrals $R_i^0 = T_i(f,p_i)$. Refinements

of the integrals can then be produced by Richardson Extrapolation



Thus, consider the two integrals

$$\int_{a}^{b} f(x) dx = R_{i-1}^{0} + a_{2}h^{2} + a_{4}h^{4} + \dots$$

$$\int_{a}^{b} f(x) dx = R_{i}^{0} + a_{2}\left(\frac{h}{2}\right)^{2} + a_{4}\left(\frac{h}{2}\right)^{4} + \dots$$

Note that there no odd terms in the error. Then, define the next set of refinements as

$$R_i^1 = \frac{1}{3} \left(4R_i^0 - R_{i-1}^0 \right).$$

which has an error $\mathcal{O}\left(h^4\right)$. The extrapolated values are equivalent to integrals approximated by Simpson's rule. The recurrence formula can be derived

$$R_i^m = \frac{1}{4^m-1} \left(4^m R_i^{m-1} - R_{i-1}^{m-1} \right).$$

8.1 Gauss Quadrature

Generalise the quadrature formula so that an integral is approximated as

$$I_n[f] = \sum_{i=0}^{n} \alpha_i f(x_i)$$

The above equation is a weighted sum of the values of f at the points x_i , for i = 0, ..., n. These points are said to be the *nodes* of the quadrature formula, while the $\alpha_i \in \mathbb{R}$ are its *coefficients* or *weights*. Both weights and nodes depend in general on n.

- Can the weights be chosen such that the error in an integral is minimized?
- Furthermore, can the nodes be chosen such that the integral can be improved?

Definition 8.4 (Orthogonal functions). Two real-valued functions f(x) and g(x) are said to be **orthogonal** if

$$\langle f, g \rangle = \int_a^b f(x)g(x) \, \mathrm{d}x = 0.$$

Theorem 15 (Gaussian Quadrature). Let $q\left(x\right)$ be a non-trivial polynomial of degree n+1 such that

- 1. it has n+1 distinct roots, denoted as x_i , in [a,b],
- 2. the polynomial satisfies

$$\int_{a}^{b} x^{k} q(x) dx = 0 \quad \text{for} \quad k = 0, \dots, n.$$

i.e. is orthogonal to x^k .

Then, denote the integral as

$$I[f] = \int_{a}^{b} f(x) dx = \sum_{i=0}^{n} A_{i} f(x_{i})$$

with $A_i = \int_a^b L_i(x) dx$ for all polynomials f(x) of degree less than or equal to 2n+1. The integral I[f] integrates all polynomials of degree 2n+1 exactly.

The degree of exactness of I[f] is 2n + 1.

Definition 8.5 (Gauss-Legendre Quadrature). The Legendre polynomials are a set of orthogonal polynomials where

$$\int_{-1}^{1} P_m(x) P_n(x) dx = 0 \quad \text{for} \quad n \neq m.$$

and $P_0 = 1$. Thus, $P_1 = x$, $P_2 = (3x^2 - 1)/2$. The Legendre polynomials obey a recursive formula:

$$P_n = \frac{2n-1}{n}xP_{n-1}(x) - \frac{n-1}{n}P_{n-2}(x), \text{ for } n \ge 2.$$

Gauss-Legendre quadrature uses the roots of the Legendre polynomials as the nodes for integration, and weights found by equating the quadrature

expressions with the exact integrals for $f = 1, x, x^2$.

The domain of integration can be scaled via the invertible transformation $x=\frac{b-a}{2}t+\frac{a+b}{2},$ so that

$$\int_{a}^{b} f(x) dx = \frac{b-a}{2} \int_{-1}^{1} f\left(\frac{b-a}{2}t + \frac{a+b}{2}\right) dt.$$

Differential Equations

Finite Difference Methods for Differential Equations

Solutions to differential equations are functions.

Definition 9.1 (Ordinary Differential Equations). An ordinary differential equation (ODE) is an equation that involves one of more derivatives of a function of a single variable.

For example, with only the first derivative y'(t) = f(y(t), t).

Definition 9.2 (Initial Value Problems). An initial value problem (IVP) is given by an ordinary differential equation of the form y'(t) = f(y(t), t)and initial value $y(a) = y_a$ for the unknown function y(t).

Often a = 0, and the initial condition is denoted by $y(0) = y_0$.

Definition 9.3 (One step methods). A numerical method for approximating the solution to a differential equation is called a **one step method** if the solution at time step i + 1 depends only on the previous one. The time step is $t_{i+1} = t_i + h$, for some time step $h = \Delta t$.

Forward Euler Method

Definition 9.4. Forward Euler approximates the derivative through a firstorder forward difference approximation of the first-order derivative, i.e. for u_n , the solution u at t_n , the computed solution to a differential equation $\dot{u} = f(u)$, evolves according to

$$u_{m+1} = u_m + h f_m$$

 $u_{n+1}=u_n+hf_n$ where $f_n=f\left(u_n,t_n\right)$ and $h=t_{n+1}-t_n.$ The error is $\mathcal{O}\left(h^2\right).$

Backward Euler Method

Definition 9.5. Backward Euler uses the backward finite difference approximation of the first-order derivative, so that the solution is computed

$$u_{n+1} = u_n + h f_{n+1}$$

where $f_{n+1} = f(u_{n+1}, t_{n+1})$.

Crank-Nicolson Method

Definition 9.6. The **Crank-Nicolson method** is given by

$$u_{n+1} = u_n + \frac{h}{2} (f_n + f_{n+1}).$$

Heun's Method

Definition 9.7. Heun's method is

$$u_{n+1} = u_n + \frac{h}{2} (f_n + f_{n+1}).$$

Alternatively, one-step methods can be considered to be integrators

$$y(t+h) = y(t) + \int_{t}^{t+h} f(y(\tau), \tau) d\tau.$$

Thus, the forward Euler method is the left Riemann sum, backward Euler is the right Riemann sum and the Crank-Nicolson is the trapezoidal rule.

9.2 Analysis of One-Step Methods

Definition 9.8. A numerical method is said to be **explicit** if an approximation y_{k+1} can be calculated directly from already computed values y_i , i < k. Otherwise, the method is said to be **implicit**.

Often, implicit methods require, at each step, the solution of a nonlinear equation for computing y_{k+1} .

Both the Forward Euler and Heun's method are explicit, whereas the Backward Euler and Crank-Nicolson methods are implicit.

Definition 9.9. A function f is **Hölder continuous** if there exists real constants C > 0 and $\alpha \le 0$ such that

$$|f(x) - f(y)| \le C||x - y||^{\alpha}$$

for all x and y. If $\alpha = 1$ the function is said to be **Lipshitz continuous**.

Definition 9.10 (Stable). stable

Definition 9.11 (Consistent). Let y(t) be the solution to an initial value problem and $h_{\max} = \max_{x_k} h_k$ A method is said to be **consistent** if the func-

tion f satisfies the Lipshitz condition, with respect to y, in an interval then

$$\lim_{h_{\max}} \left(\max_{x_k \in I_h} \left\| f\left(t_k, y\left(t_k\right)\right) - \Phi\left(t_k, y\left(t_k\right), h_k\right) \right\| \right) = 0.$$

Definition 9.12 (Order). A one-step method is of order $p \in \mathbb{N}$, if for all $t \in [0, \tau]$, the solution satisfies the condition that $\tau(h) = \mathcal{O}(h^p)$ as $h \to 0$.

Definition 9.13 (Zero Stable Methods). A method of the form

$$u_{n+1} = u_n + h\Phi\left(t_n, u_n, f_n, h\right)$$

is called **zero-stable** if there exists both a maximal step size, h_0 and a constant, C, such that for all $h \in [0, h_{\max}]$ and for $\varepsilon > 0$, then the following holds:

If, for all time-steps $0 \le n \le N$, there exists a $\delta_n \le \varepsilon$ and

$$z_{n+1} = z_n + h\Phi(t_n, z_n, f_n(z_n, t_n), h) + \delta_{n+1}$$

and $z_0 = u_0 + \delta_0$, then

$$|z_n - u_n| \le C\varepsilon$$
 for $0 \le n \le N$.

Zero stability means that small perturbations in the computations lead to small perturbations in the approximations.

Theorem 16. If increment function is Lipshitz continuous for y_n for any h and t_n , then the one step method is zero-stable.

Theorem 17. If the increment function Φ is

(i) Lipshitz continuous for u_n for any h and t_{n+1}

and

(ii) the method is consistent

then

$$\lim_{h \to 0} |y_n - u_n| = 0.$$

Also, if the method is of order p and if $|y_0 - u_0| = \mathcal{O}(h^p)$ as $h \to 0$, the convergence is of order p.

Definition 9.14 (Absolute Stability). A numerical scheme for approximating the solution to the linear differential equation $y'(t) = \lambda y(t)$ with $\lambda \in \mathbb{C}$ and initial condition $y_0 = 1$ is said to be **absolutely stable** if $|u_n| \to 0$ as $n \to \infty$, when $\text{Re}(\lambda) < 0$, for a fixed value of h.

Runge-Kutta Schemes And Multi-Step Schemes If a ordinary differential equation is given by $\dot{y}=f\left(y,t\right)$, then a Runge-Kutta scheme takes the form

$$u_{n+1} = u_n + hF(t_n, u_n, h; f)$$

where F is an increment function given by

$$F(t_n, u_n, h; f) = \sum_{i=1}^{s} b_i k_i,$$

with k_i defined as

$$k_i = f\left(u_n + h\sum_{j=1}^{s} a_{i,j}k_j, t_n + c_i h\right)$$
 $i = 1, \dots, s$

where s is referred to as the number of stages of the method.

Thus, an s-stage scheme is characterised by coefficients b_i , c_i and $a_{i,j}$.

If the matrix defined by the elements $a_{i,j}$ is lower triangular, i.e. $a_{i,j} = 0$ for all $i \leq j$, then each k_i can be computed explicitly in terms of the previous coefficients k_1, \ldots, k_{i-1} . Such schemes are called **explicit**, otherwise they are said to be **implicit**.

The local truncation error is defined as

$$h\tau_{n+1}(h) = u(t_{n+1}) - u(t_n) - hF(t_n, u_n, h; f).$$

It can be shown that $\tau(h) = \max |\tau_{n+1}(h)| \to 0$ as $h \to 0$ if and only if $\sum_{i=1}^{s} b_i = 1$.

A Runge-Kutta method is of order $p \ge 1$ if $\tau(h) = \mathcal{O}(h^p)$ as $h \to 0$.

The order of an s-stage explicit Runge-Kutta method cannot be greater than s. Additionally, there does not exist a s-stage explicit Runge-Kutta method with order s if $s \ge 5$.

The order of an s-stage implicit Runge-Kutta method cannot be greater than 2s

The most common form of the Runge-Kutta method is the fourth order Runge-Kutta method (RK4). It takes the form:

$$u_{n+1} = u_n + \frac{h}{6} \left(k_1 + 2k_2 + 2k_3 + k_4 \right)$$

where

$$k_{1} = f(u_{n}, t_{n}),$$

$$k_{2} = f\left(u_{n} + \frac{h}{2}k_{1}, t_{n} + \frac{h}{2}\right),$$

$$k_{3} = f\left(u_{n} + \frac{h}{2}k_{2}, t_{n} + \frac{h}{2}\right),$$

$$k_{4} = f\left(u_{n} + hk_{3}, t_{n} + h\right).$$

9.3 Partial Differential Equations

Definition 9.15. A partial differential equation is relation involving a unknown function of several free variables and partial derivatives with respect to these variables.

A partial differential equation is said to be linear if it only contains linear terms of the unknown and its derivatives. For example, a second-order linear partial differential equation for a unknown function u(x,t) has the form

$$a_1u_{xx}+a_2u_{xt}+a_3u_{tt}+a_3u_x+a_4u_t+a_5u=f(x,t).$$
 where $u_{xx}=\frac{\partial^2 u}{\partial x^2}.$

For finite-difference schemes, all partial derivatives must be approximated by discretized operators.