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# CTMS-MAT-13: Numerical Methods

### Summary from 29 April 2025

Notes for numerical methods course 2025. This can be downloaded from:  $\verb|https://djps.github.io/docs/numericalmethods/notes|$ 

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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

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

 $\sum_{k=0}^{\infty}\frac{f^{(k)}\left(c\right)}{k!}\left(x-c\right)^k$  where  $f^{(k)}=\frac{\mathrm{d}^kf}{\mathrm{d}x^k}$  is the  $k^{\mathrm{th}}$  derivative.

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

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

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

 $\left|\tilde{a}-a\right|.$  If  $\left|a\right|\neq0,$  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 errors 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

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

where  $|\Delta x_i| = |\tilde{x}_i - x_i|$  for an approximation  $\tilde{x}_i$ .

# 3 Number Representations

**Definition 3.1** (Base Representation). Every number  $x \in \mathbb{N}_0$  can be written as a unique expansion with respect to base  $b \in \mathbb{N} \setminus \{1\}$  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 an integer 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  $i = n, \ldots, 0$

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

which at each steps provides an  $a_i$  and updates y.

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

where div is integer division, and mod is the remainder operator.

There are two issues: finding n maybe difficult and for large values of  $b^i$  division maybe computationally costly. Horner's algorithm seeks to overcome these issues.

### Algorithm (Horner).

Horner's algorithm is:

- 1. Input  $(x)_{10}$
- 2. Set i = 0
- 3. Let y = x. Then while y > 0

$$\begin{array}{rcl} a_i & = & y \bmod b \\ y & = & y \operatorname{div} b \\ i & = & i+1 \end{array}$$

which at each steps provides an  $a_i$  and updates y.

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

### **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

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

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

**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.

An equivalent definition is that for any non-zero real column vector z, then  $z^T M z$  is always 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 (i.e. with and without pivoting) will fail due to division by zero, yielding a floating exception error.

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

#### 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
a_{i,j} = a_{i,j} - \frac{a_{i,k}}{a_{k,k}} a_{k,j}
 4:
                  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_{n,n}} for i = n - 1 to 1 do
12:
13:
14:
             y = b_i
             for j = n to i + 1 do
15:
             y = y - a_{i,j}x_j
end for
x_i = \frac{y}{a_{i,i}}
16:
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

2: **for** k = 1 to n - 1 **do** for i = k to n do

1: Find maximal absolute values vector  $\vec{s}$  with entries

$$s_i = \max_{j=1,...,n} |a_{i,j}|$$
 $= 1 \text{ to } n-1 \text{ do}$ 
 $: i = k \text{ to } n \text{ do}$ 
Compute relative pivot elements  $\left| \frac{a_{i,k}}{s_i} \right|$ 

5:

3:

4:

- Find row with largest relative pivot element, denote this as row j6:
- 7: Swap rows k and j in the matrix A
- Swap entries k and j in the vector  $\vec{s}$
- Do forward elimination on row k
- 10: end for

It is important to remember that while the elements of the augmented matrix are rescaled during forward elimination, the elements of the vector  $\vec{s}$  are reordered.

**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.

Note that

- The product of two upper triangular matrices is also upper triangular.
- The inverse of an upper triangular matrix is also upper triangular.

**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.

A single pass of Gaussian forward elimination can be written as a matrix operation, where the elementary row operations are expressed as a Gaussian transformation matrix, M. This is a lower triangular matrix. It can be shown that

$$L = M_1^{-1} M_2^{-1} \cdots M_{n-1}^{-1}$$

where each matrix  $M_i$  is a matrix which describes the  $i^{\rm th}$  pass of forward elimination part of Gaussian elimination. Equivalently, the upper triangular matrix U is given by

$$U = M_{n-1} \cdots M_2 M_1 A.$$

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

#### 4.2 Indirect Methods

For a non-singular matrix A, consider an iterative scheme of the form

$$Qx_{k+1} = (Q - A)x_k + b.$$

This is equivalent to

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

**Definition 4.9** (Convergent Methods). A numerical method is said to be **convergent** if  $x_k \to x^*$  as  $k \to \infty$  where  $x^*$  is the exact solution.

**Definition 4.10** (Spectral Radius of a Matrix). The **spectral radius** of a square matrix A is defined as

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

where the  $\lambda_i$  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.

$$\rho\left(I - Q^{-1}A\right) < 1.$$

**Corollary 4.11.** A sequence of vectors  $e_k = M^k e_0$  convergences to the zero vector if and only if the spectral radius of the matrix M is less than one.

**Definition 4.12** (Order of Convergence). 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 4.13** (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 4.14** (Jacobi Iteration). The **Jacobi iteration** has Q = D, where D is the diagonal matrix of A, so

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

**Definition 4.15** (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 4.16** (Gauss-Seidel Scheme). Let Q = L + D, where L is the lower triangular matrix of A and D is the diagonal matrix of A, then the **Gauss-Seidel** scheme is given by

$$(D+L)x_{k+1} = -Ux_k + b.$$

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 4.17** (Successive Over Relaxation). The scheme uses  $Q = L + \frac{1}{\omega}D$ , where L is the lower triangular matrix of A and D is the diagonal matrix of A, thus

$$(D + \omega L)x_{k+1} = -((\omega - 1)D + \omega U)x_k + \omega b.$$

### 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.

### 4.3 Fundamental Theorem of Numerical Analysis

**Definition 4.18** (Stability). A numerical method is said to be **stable** if and only if any initial error  $e_0$  is damped during the iterations, i.e.  $||e_k|| < ||e_0||$ .

Note that ||x|| is a *norm* of a vector, such as  $||x||_2 = \sqrt{x_0^2 + x_1^2 + \dots x_n^2}$ .

**Definition 4.19** (Consistency). A numerical method is said to be **consistent** if any fixed point  $x^*$  of the iteration is a solution to the problem being solved.

For a linear system, a fixed point,  $x^*$ , fulfils

$$x^* = (I - Q^{-1}A) x^* + Q^{-1}b \Leftrightarrow Ax^* = b.$$

Thus, a fixed point of the iterative scheme is a solution of the linear system. If the method is stable then

$$e_k = (I - Q^{-1}A)^k e_0,$$

so then  $||I - Q^{-1}A|| < 1$  for  $||e_k|| < ||e_0||$ . Note that  $||I - Q^{-1}A||$  is the norm of a matrix, which is induced by a vector norm.

**Theorem 9 (Fundamental Theorem of Numerical Analysis).** A numerical method is convergent if and only if it is consistent and stable.

### 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 10.

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

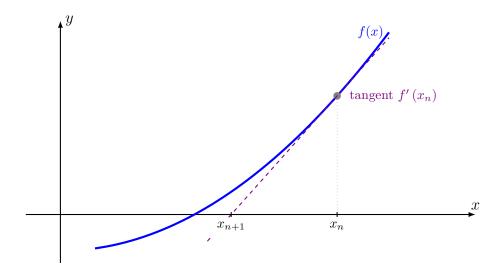
$$|r - c_n| < \frac{b - a}{2n}.$$

#### 5.2 Newton's Method

#### **Definition 5.2** (Newton's Method).

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 11.

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

#### 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, 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 (a,b) 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 iterates in the sequence will satisfy

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

#### Theorem 13.

- Theorem 13. Let  $f \in C^1([a,b])$ , with  $1. \ f(a)f(b) < 0,$   $2. \ f'(x) \neq 0 \text{ for all } x \in (a,b),$   $3. \ f''(x) \text{ exists and is continuous, i.e. } f(x) \in C^2([a,b]).$

If  $x_0$  is close enough to the root r, then Newton's method converges quad-

#### 5.3 Secant Method

### **Definition 5.3** (Secant Method).

The secant method is defined as

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

where the derivative of f is approximated via a Taylor expansion.

#### Theorem 14.

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  $\alpha = \frac{1 + \sqrt{5}}{2}$ .

The properties of the three nonlinear solvers can be summarised as follows:

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

### 5.4 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  $J^{-1}\left(x_{m}\right)$  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

$$J\left(x_{m}\right)\left(x_{m+1}-x_{m}\right)=-f\left(x_{m}\right)$$

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

## 6 Interpolation

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

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

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

where  $\varphi_{i}\left(u\right)$  are polynomials. Thus for every nodal value j, the polynomial exactly satisfies

$$p(u_j) = \sum_{i=0}^{n} \alpha_i \varphi_i(u_j),$$

for the weights  $\alpha_i$ . Thus, solving for all the values of  $\alpha$  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  $\Phi$  is the **collocation matrix**. If there are n data points, the collocation matrix is given by

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

Thus the weights which define the interpolating polynomial are found as  $\alpha = \Phi^{-1}p$ .

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

When the polynomials are given by  $\varphi_i(u) = u^i$ , then  $\Phi$  is called 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

 $i = 0, \dots, n$ , are called **characteristic polynomials** and are given by

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

**Definition 6.4** (Newton Interpolation). Newton interpolation interpolates a set of points  $(x_i, y_i)$  as

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

using a linear combination of Newton basis polynomials, which are defined as

$$n_0(x) = 1,$$
  $n_i(x) = (x - x_0)(x - x_1) \cdots (x - x_{i-1})$   
=  $\prod_{j=0}^{j-1} (x - x_j).$ 

By construction,  $\alpha_0 = y_0$ , and subsequent terms must be solved by evaluating the interpolating polynomial at increasing orders, leading to the formula

$$\alpha_{i+1} = \frac{y_{i+1} - p_i(x_{i+1})}{n_i(x_{i+1})}.$$

#### Algorithm (Aitken's Algorithm).

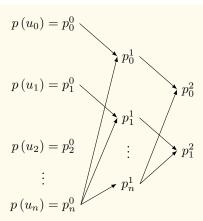
**Aitken's algorithm** is an iterative process for evaluating Lagrange interpolation polynomials at an arbitrary point,  $u^*$ , without explicitly constructing them. If the interpolating polynomial is given by p, and is derived from n data points  $(u_i, y_i)$  for  $i = 0, \ldots, n$ 

$$p\left(u\right) = \sum_{i=0}^{n} p_{i}^{n} l_{i}^{n}\left(u\right).$$

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

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

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.

### Piecewise Polynomial Interpolation

**Definition 6.5** (Spline Functions). 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 6.6** (B-Splines). A spline is said to be a **b-spline** if it is of the

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

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

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

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

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

where

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

is a local parameter.

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 6.7** (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 the squares

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

The minimum is found when

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

i.e.

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

**Definition 6.8** (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  $\beta$  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  $\Phi^T$  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

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

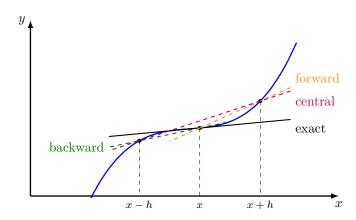
Combining the two terms so that the low order term cancel, i.e. via f'(x) - 4f'(x), then a better approximation can be found as

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

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

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



# 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  $m_i = \min_{x \in [x_i, x_{i+1}]} f(x)$ 

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

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

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

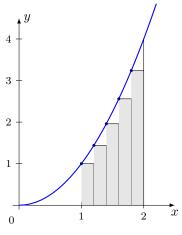
$$\sum_{i=0}^{n-1} (x_{i+1} - x_i) f(x_{i-1}),$$

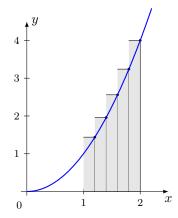
$$\sum_{i=0}^{n-1} (x_{i+1} - x_i) f(x_i).$$

**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





(a) Lower Riemann Sum

(b) Upper Riemann Sum

the formula can be given by

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

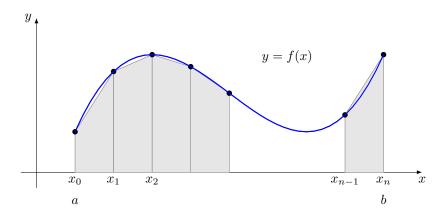
**Theorem 15 (Error for Trapezoidal Rule).** Let  $f \in C^2([a,b])$  and p be equidistant partition of [a,b], with  $h=x_{i+1}-x_i$ . The error can be shown to have the form:

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

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



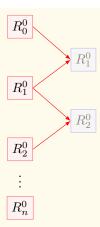
**Definition 8.3** (Simpson's Rule). The integral is approximated as

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

Simpson's rule uses interpolation with quadratic polynomials. 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 and then Richardson Extrapolation 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 16 (Gaussian Quadrature).** Let q(x) 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,  $P_1=x,\,P_2=\left(3x^2-1\right)/2$  etc. 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  $x=\frac{b-a}{2}t+\frac{a+b}{2}, \text{ so that}$   $\int_a^b f(x)\,\mathrm{d}x = \frac{b-a}{2}\int_{-1}^1 f\left(\frac{b-a}{2}t+\frac{a+b}{2}\right)\,\mathrm{d}t.$ 

$$\int_a^b f(x) \, \mathrm{d}x = \frac{b-a}{2} \int_{-1}^1 f\left(\frac{b-a}{2}t + \frac{a+b}{2}\right) \, \mathrm{d}t$$

## 9 Differential Equations

### 9.1 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 or 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), with  $t \ge a$ .

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  $t_{n+1}$ , denoted by  $u_{n+1}$ , depends only on the previous one,  $t_n$ , where  $t_{n+1} = t_n + h$ , for some small increment  $h = \Delta t$ .

**Definition 9.4** (Forward Euler). This approximates the derivative through a first-order 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

$$u_{n+1} = u_n + hf_n$$

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

**Definition 9.5** (Backward Euler Method). This method uses the backward finite difference approximation of the first-order derivative, so that the solution is computed via

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

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

**Definition 9.6** (Crank-Nicolson Method). This method is given by

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

**Definition 9.7** (Heun's Method). This method is given by

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

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 equivalent to the left Riemann sum, backward Euler is equivalent to the right Riemann sum and the Crank-Nicolson is the trapezoidal rule.

**Definition 9.8** (Implicit and Explicit Schemes). A numerical method is said to be **explicit** if an approximation  $u_{n+1}$  can be calculated directly from already computed values  $u_i$ , i < n. Otherwise, the method is said to be **implicit**.

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

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

Huen's method can be interpreted as the Crank-Nicolson method with the approximation  $u_{n+1} \approx u_n + hf_n$  replacing the explicit  $f_{n+1}$  term, which depends on  $u_{n+1}$ .

#### 9.2 Analysis of One-Step Methods

Any explicit one-step method has the form

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

with  $\Phi$  the increment function.

**Definition 9.9** (Hölder & Lipshitz Continuity). A function f is **Hölder continuous** if there exists real constants C>0 and  $\alpha\geq 0$  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** (Consistent Schemes). For the exact solution to the differential equation,  $y(t_n) = y_n$ , the solution can be written as

$$y_{n+1} = y_n + h\Phi(t_n, y(t_n), f_n, h) + \varepsilon_n$$

so that

$$\tau_n = \frac{y_{n+1} - y_n}{h} - \Phi(t_n, y(t_n), f_n, h)$$

where  $\varepsilon_n = h\tau_n$  for a  $\tau_n = \tau_n (h)$  is defined as the **local truncation error** at step n.

The **consistency error** is given by  $\tau = \max_n |\tau_n|$ .

A method is said to be **consistent** if

$$\lim_{h \to 0} \Phi = f$$

This means the increment function is a good approximation to the differential equation as the step size tends to zero.

**Definition 9.11 (Order).** A one-step method is of order  $p \in \mathbb{N}$ , if for all  $t \in [0, T]$ , the solution satisfies the condition that  $\tau(h) = \mathcal{O}(h^p)$  as  $h \to 0$ .

#### **Definition 9.12** (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,  $h_{\max}$  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  $z_0 = y_0 + \delta_0$ , 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 17.** If the increment function is Lipshitz continuous for  $y_n$  for any h and  $t_n$ , then the one-step method is zero-stable.

#### **Theorem 18.** If the increment function $\Phi$ 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 9.13** (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.

**Definition 9.14** (Well-posed). A differential equation is said to be **well-posed** if

- a unique solution exists for any initial conditions and
- the solution's behaviour changes continuously with the initial conditions.

A differential equation which does not have these properties is said to be ill-posed.

**Theorem 19 (Lax Equivalence theorem).** The Lax Equivalence theorem or Lax–Richtmyer theorem is the equivalent form of the Fundamental Theorem of Numerical Analysis for differential equations, which states that for a *consistent* finite difference method for a well-posed linear initial value problem, the method is *convergent* if and only if it is *stable*.

#### Runge-Kutta Schemes And Multi-Step Schemes

**Definition 9.15** (Runge-Kutta Methods). If an ordinary differential equation is given by  $\dot{y} = f(y,t)$ , 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

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

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  $k_1, \ldots, k_{i-1}$ . Thus, 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 components of a Runge-Kutta scheme are expressed in a **Butcher array** 

$$\begin{array}{c|ccccc} c_1 & a_{1,1} & \dots & a_{1,s} \\ \vdots & \vdots & \ddots & \vdots \\ \hline c_s & a_{s,1} & \dots & a_{s,s} \\ \hline & b_1 & \dots & b_s \end{array}$$

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 9.16** (Partial Differential Equations). A partial differential equation is a relation involving an 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 an 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.