NUMERICAL SOLUTION OF PARTIAL DIFFERENTIAL EQUATIONS
MA 3243 LECTURE NOTES

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This course is designed to respond to the needs of the aeronautical engineering curricula by providing an applications oriented introduction to the finite difference method of solving partial differential equations arising from various physical phenomenon. This course will emphasize design, coding, and debugging programs written by the students in order to fix ideas presented in the lectures. In addition, the course will serve as an introduction to a course on analytical solutions of PDE's. Elementary techniques including separation of variables, and the method of characteristics will be used to solve highly idealized problems for the purpose of gaining physical insight into the physical processes involved, as well as to serve as a theoretical basis for the numerical work which follows.

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1 Introduction and Applications

This section is devoted to basic concepts in partial differential equations. We start the chapter with definitions so that we are all clear when a term like linear partial differential equation (PDE) or second order PDE is mentioned. After that we give a list of physical problems that can be modelled as PDEs. An example of each class (parabolic, hyperbolic and elliptic) will be derived in some detail. Several possible boundary conditions are discussed.

1.1 Basic Concepts and Definitions

Definition 1. A partial differential equation (PDE) is an equation containing partial derivatives of the dependent variable.

For example, the following are PDEs

\[ u_t + cu_x = 0 \]  \hspace{1cm} (1.1.1)

\[ u_{xx} + u_{yy} = f(x, y) \]  \hspace{1cm} (1.1.2)

\[ \alpha(x, y)u_{xx} + 2u_{xy} + 3x^2u_{yy} = 4e^x \]  \hspace{1cm} (1.1.3)

\[ u_x u_{xx} + (u_y)^2 = 0 \]  \hspace{1cm} (1.1.4)

\[ (u_{xx})^2 + u_{yy} + a(x, y)u_x + b(x, y)u = 0 \]  \hspace{1cm} (1.1.5)

Note: We use subscript to mean differentiation with respect to the variables given, e.g. \( u_t = \frac{\partial u}{\partial t} \). In general we may write a PDE as

\[ F(x, y, \cdots, u, u_x, u_y, \cdots, u_{xx}, u_{xy}, \cdots) = 0 \]  \hspace{1cm} (1.1.6)

where \( x, y, \cdots \) are the independent variables and \( u \) is the unknown function of these variables. Of course, we are interested in solving the problem in a certain domain \( D \). A solution is a function \( u \) satisfying (1.1.6). From these many solutions we will select the one satisfying certain conditions on the boundary of the domain \( D \). For example, the functions

\[ u(x, t) = e^{x-ct} \]
\[ u(x, t) = \cos(x - ct) \]

are solutions of (1.1.1), as can be easily verified. We will see later (section 7.1) that the general solution of (1.1.1) is any function of \( x - ct \).

Definition 2. The order of a PDE is the order of the highest order derivative in the equation. For example (1.1.1) is of first order and (1.1.2) - (1.1.5) are of second order.

Definition 3. A PDE is linear if it is linear in the unknown function and all its derivatives with coefficients depending only on the independent variables.
For example (1.1.1) - (1.1.3) are linear PDEs.

**Definition 4.** A PDE is nonlinear if it is not linear. A special class of nonlinear PDEs will be discussed in this book. These are called quasilinear.

**Definition 5.** A PDE is quasilinear if it is linear in the highest order derivatives with coefficients depending on the independent variables, the unknown function and its derivatives of order lower than the order of the equation.

For example (1.1.4) is a quasilinear second order PDE, but (1.1.5) is not.

We shall primarily be concerned with linear second order PDEs which have the general form

\[ A(x, y)u_{xx} + B(x, y)u_{xy} + C(x, y)u_{yy} + D(x, y)u_x + E(x, y)u_y + F(x, y)u = G(x, y). \]  
(1.1.7)

**Definition 6.** A PDE is called homogeneous if the equation does not contain a term independent of the unknown function and its derivatives.

For example, in (1.1.7) if \( G(x, y) \equiv 0 \), the equation is homogeneous. Otherwise, the PDE is called inhomogeneous.

Partial differential equations are more complicated than ordinary differential ones. Recall that in ODEs, we find a particular solution from the general one by finding the values of arbitrary constants. For PDEs, selecting a particular solution satisfying the supplementary conditions may be as difficult as finding the general solution. This is because the general solution of a PDE involves an arbitrary function as can be seen in the next example. Also, for linear homogeneous ODEs of order \( n \), a linear combination of \( n \) linearly independent solutions is the general solution. This is not true for PDEs, since one has an infinite number of linearly independent solutions.

**Example**

Solve the linear second order PDE

\[ u_{\xi \eta}(\xi, \eta) = 0 \]  
(1.1.8)

If we integrate this equation with respect to \( \eta \), keeping \( \xi \) fixed, we have

\[ u_{\xi} = f(\xi) \]

(Since \( \xi \) is kept fixed, the integration constant may depend on \( \xi \).)

A second integration yields (upon keeping \( \eta \) fixed)

\[ u(\xi, \eta) = \int f(\xi)d\xi + G(\eta) \]

Note that the integral is a function of \( \xi \), so the solution of (1.1.8) is

\[ u(\xi, \eta) = F(\xi) + G(\eta). \]  
(1.1.9)

To obtain a particular solution satisfying some boundary conditions will require the determination of the two functions \( F \) and \( G \). In ODEs, on the other hand, one requires two constants. We will see later that (1.1.8) is the one dimensional wave equation describing the vibration of strings.
Problems

1. Give the order of each of the following PDEs
   a. \( u_{xx} + u_{yy} = 0 \)
   b. \( u_{xxx} + u_{xy} + a(x)u_y + \log u = f(x, y) \)
   c. \( u_{xxx} + u_{xyy} + a(x)u_{xxy} + u^2 = f(x, y) \)
   d. \( u u_{xx} + u^2_{yy} + e^u = 0 \)
   e. \( u_x + cu_y = d \)

2. Show that
   \( u(x, t) = \cos(x - ct) \)
   is a solution of
   \( u_t + cu_x = 0 \)

3. Which of the following PDEs is linear? quasilinear? nonlinear? If it is linear, state whether it is homogeneous or not.
   a. \( u_{xx} + u_{yy} - 2u = x^2 \)
   b. \( u_{xy} = u \)
   c. \( u u_x + x u_y = 0 \)
   d. \( u_x^2 + \log u = 2xy \)
   e. \( u_{xx} - 2u_{xy} + u_{yy} = \cos x \)
   f. \( u_x(1 + u_y) = u_{xx} \)
   g. \( (\sin u_x)u_x + u_y = e^x \)
   h. \( 2u_{xx} - 4u_{xy} + 2u_{yy} + 3u = 0 \)
   i. \( u_x + u_x u_y - u_{xy} = 0 \)

4. Find the general solution of
   \( u_{xy} + u_y = 0 \)
   (Hint: Let \( v = u_y \))

5. Show that
   \( u = F(xy) + x G\left(\frac{y}{x}\right) \)
   is the general solution of
   \( x^2 u_{xx} - y^2 u_{yy} = 0 \)
1.2 Applications

In this section we list several physical applications and the PDE used to model them. See, for example, Fletcher (1988), Haltiner and Williams (1980), and Pedlosky (1986).

For the heat equation (parabolic, see definition 7 later).

\[ u_t = ku_{xx} \quad \text{(in one dimension)} \quad (1.2.1) \]

the following applications

1. Conduction of heat in bars and solids
2. Diffusion of concentration of liquid or gaseous substance in physical chemistry
3. Diffusion of neutrons in atomic piles
4. Diffusion of vorticity in viscous fluid flow
5. Telegraphic transmission in cables of low inductance or capacitance
6. Equilization of charge in electromagnetic theory.
7. Long wavelength electromagnetic waves in a highly conducting medium
8. Slow motion in hydrodynamics

Laplace’s equation (elliptic)

\[ u_{xx} + u_{yy} = 0 \quad \text{(in two dimensions)} \quad (1.2.2) \]

or Poisson’s equation

\[ u_{xx} + u_{yy} = S(x, y) \quad (1.2.3) \]

is found in the following examples

1. Steady state temperature
2. Steady state electric field (voltage)
3. Inviscid fluid flow
4. Gravitational field.

Wave equation (hyperbolic)

\[ u_{tt} - c^2u_{xx} = 0 \quad \text{(in one dimension)} \quad (1.2.4) \]

appears in the following applications
1. Linearized supersonic airflow
2. Sound waves in a tube or a pipe
3. Longitudinal vibrations of a bar
4. Torsional oscillations of a rod
5. Vibration of a flexible string
6. Transmission of electricity along an insulated low-resistance cable
7. Long water waves in a straight canal.

Remark: For the rest of this book when we discuss the parabolic PDE

\[ u_t = k \nabla^2 u \]  \hspace{1cm} (1.2.5)

we always refer to \( u \) as temperature and the equation as the heat equation. The hyperbolic PDE

\[ u_{tt} - c^2 \nabla^2 u = 0 \]  \hspace{1cm} (1.2.6)

will be referred to as the wave equation with \( u \) being the displacement from rest. The elliptic PDE

\[ \nabla^2 u = Q \]  \hspace{1cm} (1.2.7)

will be referred to as Laplace’s equation (if \( Q = 0 \)) and as Poisson’s equation (if \( Q \neq 0 \)).

In the following sections we give details of several applications. The first example leads to a parabolic one dimensional equation. Here we model the heat conduction in a wire (or a rod) having a constant cross section. The boundary conditions and their physical meaning will also be discussed. The second example is a hyperbolic one dimensional wave equation modelling the vibrations of a string. We close with a three dimensional advection diffusion equation describing the dissolution of a substance into a liquid or gas. A special case (steady state diffusion) leads to Laplace’s equation.

1.3 Conduction of Heat in a Rod

Consider a rod of constant cross section \( A \) and length \( L \) (see Figure 1) oriented in the \( x \) direction.

Let \( \epsilon(x, t) \) denote the thermal energy density or the amount of thermal energy per unit volume. Suppose that the lateral surface of the rod is perfectly insulated. Then there is no thermal energy loss through the lateral surface. The thermal energy may depend on \( x \) and \( t \) if the bar is not uniformly heated. Consider a slice of thickness \( \Delta x \) between \( x \) and \( x + \Delta x \).
If the slice is small enough then the total energy in the slice is the product of thermal energy density and the volume, i.e.

\[ e(x, t)A\Delta x . \] (1.3.1)

The rate of change of heat energy is given by

\[ \frac{\partial}{\partial t}[e(x, t)A\Delta x] . \] (1.3.2)

Using the conservation law of heat energy, we have that this rate of change per unit time is equal to the sum of the heat energy generated inside per unit time and the heat energy flowing across the boundaries per unit time. Let \( \varphi(x, t) \) be the heat flux (amount of thermal energy per unit time flowing to the right per unit surface area). Let \( S(x, t) \) be the heat energy per unit volume generated per unit time. Then the conservation law can be written as follows

\[ \frac{\partial}{\partial t}[e(x, t)A\Delta x] = \varphi(x, t)A - \varphi(x + \Delta x, t)A + S(x, t)A\Delta x . \] (1.3.3)

This equation is only an approximation but it is exact at the limit when the thickness of the slice \( \Delta x \to 0 \). Divide by \( A\Delta x \) and let \( \Delta x \to 0 \), we have

\[ \frac{\partial}{\partial t}e(x, t) = - \lim_{\Delta x \to 0} \frac{\varphi(x + \Delta x, t) - \varphi(x, t)}{\Delta x} + S(x, t) . \] (1.3.4)

We now rewrite the equation using the temperature \( u(x, t) \). The thermal energy density \( e(x, t) \) is given by

\[ e(x, t) = c(x)\rho(x)u(x, t) \] (1.3.5)

where \( c(x) \) is the specific heat (heat energy to be supplied to a unit mass to raise its temperature by one degree) and \( \rho(x) \) is the mass density. The heat flux is related to the temperature via Fourier’s law

\[ \varphi(x, t) = -K \frac{\partial u(x, t)}{\partial x} \] (1.3.6)

where \( K \) is called the thermal conductivity. Substituting (1.3.5) - (1.3.6) in (1.3.4) we obtain

\[ e(x)\rho(x)\frac{\partial u}{\partial t} = \frac{\partial}{\partial x} \left( K \frac{\partial u}{\partial x} \right) + S . \] (1.3.7)

For the special case that \( e, \rho, K \) are constants we get

\[ u_t = ku_{xx} + Q \] (1.3.8)
where
\[ k = \frac{K}{c\rho} \]  
and
\[ Q = \frac{S}{c\rho} \]  

(1.3.9) (1.3.10)

1.4 Boundary Conditions

In solving the above model, we have to specify two boundary conditions and an initial condition. The initial condition will be the distribution of temperature at time \( t = 0 \), i.e.
\[ u(x, 0) = f(x) . \]

The boundary conditions could be of several types.

1. Prescribed temperature (Dirichlet b.c.)
\[ u(0, t) = p(t) \]
or
\[ u(L, t) = q(t) . \]

2. Insulated boundary (Neumann b.c.)
\[ \frac{\partial u(0, t)}{\partial n} = 0 \]

where \( \frac{\partial}{\partial n} \) is the derivative in the direction of the outward normal. Thus at \( x = 0 \)
\[ \frac{\partial}{\partial n} = -\frac{\partial}{\partial x} \]

and at \( x = L \)
\[ \frac{\partial}{\partial n} = \frac{\partial}{\partial x} \]  

(see Figure 2).

\[ \text{Figure 2: Outward normal vector at the boundary} \]

This condition means that there is no heat flowing out of the rod at that boundary.
3. Newton's law of cooling

When a one dimensional wire is in contact at a boundary with a moving fluid or gas, then there is a heat exchange. This is specified by Newton's law of cooling

\[-K(0) \frac{\partial u(0, t)}{\partial x} = -H \{u(0, t) - v(t)\}\]

where \(H\) is the heat transfer (convection) coefficient and \(v(t)\) is the temperature of the surroundings. We may have to solve a problem with a combination of such boundary conditions. For example, one end is insulated and the other end is in a fluid to cool it.

4. Periodic boundary conditions

We may be interested in solving the heat equation on a thin circular ring (see figure 3).

If the endpoints of the wire are tightly connected then the temperatures and heat fluxes at both ends are equal, i.e.

\[u(0, t) = u(L, t)\]
\[u_x(0, t) = u_x(L, t)\]
Problems

1. Suppose the initial temperature of the rod was

\[ u(x, 0) = \begin{cases} 2x & 0 \leq x \leq 1/2 \\ 2(1 - x) & 1/2 \leq x \leq 1 \end{cases} \]

and the boundary conditions were

\[ u(0, t) = u(1, t) = 0 , \]

what would be the behavior of the rod’s temperature for later time?

2. Suppose the rod has a constant internal heat source, so that the equation describing the heat conduction is

\[ u_t = ku_{xx} + Q , \quad 0 < x < 1 . \]

Suppose we fix the temperature at the boundaries

\[ u(0, t) = 0 \]
\[ u(1, t) = 1 . \]

What is the steady state temperature of the rod? (Hint: set \( u_t = 0 \).)

3. Derive the heat equation for a rod with thermal conductivity \( K(x) \).

4. Transform the equation

\[ u_t = k(u_{xx} + u_{yy}) \]

to polar coordinates and specialize the resulting equation to the case where the function \( u \) does NOT depend on \( \theta \). (Hint: \( r = \sqrt{x^2 + y^2} \), \( \tan \theta = y/x \))

5. Determine the steady state temperature for a one-dimensional rod with constant thermal properties and

a. \( Q = 0 \), \( u(0) = 1 \), \( u(L) = 0 \)
b. \( Q = 0 \), \( u_x(0) = 0 \), \( u(L) = 1 \)
c. \( Q = 0 \), \( u(0) = 1 \), \( u_x(L) = \varphi \)
d. \( \frac{Q}{k} = x^2 \), \( u(0) = 1 \), \( u_x(L) = 0 \)
e. \( Q = 0 \), \( u(0) = 1 \), \( u_x(L) + u(L) = 0 \)
1.5 A Vibrating String

Suppose we have a tightly stretched string of length $L$. We imagine that the ends are tied down in some way (see next section). We describe the motion of the string as a result of disturbing it from equilibrium at time $t = 0$, see Figure 4.

![Figure 4](image)

**Figure 4:** A string of length $L$

We assume that the slope of the string is small and thus the horizontal displacement can be neglected. Consider a small segment of the string between $x$ and $x + \Delta x$. The forces acting on this segment are along the string (tension) and vertical (gravity). Let $T(x, t)$ be the tension at the point $x$ at time $t$, if we assume the string is flexible then the tension is in the direction tangent to the string, see Figure 5.

![Figure 5](image)

**Figure 5:** The forces acting on a segment of the string

The slope of the string is given by

$$\tan \theta = \lim_{\Delta x \to 0} \frac{u(x + \Delta x, t) - u(x, t)}{\Delta x} = \frac{\partial u}{\partial x}. \quad (1.5.1)$$

Thus the sum of all vertical forces is:

$$T(x + \Delta x, t) \sin \theta(x + \Delta x, t) - T(x, t) \sin \theta(x, t) + \rho_0(x) \Delta x Q(x, t) \quad (1.5.2)$$

where $Q(x, t)$ is the vertical component of the body force per unit mass and $\rho_0(x)$ is the density. Using Newton’s law

$$F = ma = \rho_0(x) \Delta x \frac{\partial^2 u}{\partial t^2}. \quad (1.5.3)$$

Thus

$$\rho_0(x) u_{tt} = \frac{\partial}{\partial x} [T(x, t) \sin \theta(x, t)] + \rho_0(x) Q(x, t) \quad (1.5.4)$$

For small angles $\theta$,

$$\sin \theta \approx \tan \theta \quad (1.5.5)$$

Combining (1.5.1) and (1.5.5) with (1.5.4) we obtain

$$\rho_0(x) u_{tt} = (T(x, t) u_x)_x + \rho_0(x) Q(x, t) \quad (1.5.6)$$

For perfectly elastic strings \( T(x, t) \equiv T_0 \). If the only body force is the gravity then

\[
Q(x, t) = -g
\]  

(1.5.7)

Thus the equation becomes

\[
u_{tt} = c^2 u_{xx} - g
\]  

(1.5.8)

where \( c^2 = T_0 / \rho_0(x) \).

In many situations, the force of gravity is negligible relative to the tensile force and thus we end up with

\[
u_{tt} = c^2 u_{xx}.
\]  

(1.5.9)

### 1.6 Boundary Conditions

If an endpoint of the string is fixed, then the displacement is zero and this can be written as

\[
u(0, t) = 0
\]  

(1.6.1)

or

\[
u(L, t) = 0.
\]  

(1.6.2)

We may vary an endpoint in a prescribed way, e.g.

\[
u(0, t) = b(t).
\]  

(1.6.3)

A more interesting condition occurs if the end is attached to a dynamical system (see e.g. Haberman [4])

\[
T_0 \frac{\partial u(0, t)}{\partial x} = k (u(0, t) - u_E(t)).
\]  

(1.6.4)

This is known as an elastic boundary condition. If \( u_E(t) = 0 \), i.e. the equilibrium position of the system coincides with that of the string, then the condition is homogeneous.

As a special case, the free end boundary condition is

\[
\frac{\partial u}{\partial x} = 0.
\]  

(1.6.5)

Since the problem is second order in time, we need two initial conditions. One usually has

\[
u(x, 0) = f(x)
\]

\[
u_t(x, 0) = g(x)
\]

i.e. given the displacement and velocity of each segment of the string.
Problems

1. Derive the telegraph equation

\[ u_{tt} + au_t + bu = c^2 u_{xx} \]

by considering the vibration of a string under a damping force proportional to the velocity and a restoring force proportional to the displacement.

2. Use Kirchoff's law to show that the current and potential in a wire satisfy

\[
\begin{align*}
i_x + C v_t + G v &= 0 \\
v_x + L i_t + R i &= 0
\end{align*}
\]

where \( i = \) current, \( v = L = \) inductance potential, \( C = \) capacitance, \( G = \) leakage conductance, \( R = \) resistance,

b. Show how to get the one dimensional wave equations for \( i \) and \( v \) from the above.
1.7 Diffusion in Three Dimensions

Diffusion problems lead to partial differential equations that are similar to those of heat conduction. Suppose \( C(x, y, z, t) \) denotes the concentration of a substance, i.e. the mass per unit volume, which is dissolving into a liquid or a gas. For example, pollution in a lake. The amount of a substance (pollutant) in the given domain \( V \) with boundary \( \Gamma \) is given by

\[
\int_V C(x, y, z, t) dV.
\]

The law of conservation of mass states that the time rate of change of mass in \( V \) is equal to the rate at which mass flows into \( V \) minus the rate at which mass flows out of \( V \) plus the rate at which mass is produced due to sources in \( V \). Let’s assume that there are no internal sources. Let \( \vec{q} \) be the mass flux vector, then \( \vec{q} \cdot \vec{n} \) gives the mass per unit area per unit time crossing a surface element with outward unit normal vector \( \vec{n} \).

\[
\frac{d}{dt} \int_V C dV = \int_V \frac{\partial C}{\partial t} dV = - \int_{\Gamma} \vec{q} \cdot \vec{n} dS.
\]

Use Gauss divergence theorem to replace the integral on the boundary

\[
\int_{\Gamma} \vec{q} \cdot \vec{n} dS = \int_V \text{div} \vec{q} dV.
\]

Therefore

\[
\frac{\partial C}{\partial t} = -\text{div} \vec{q}.
\]

Fick’s law of diffusion relates the flux vector \( \vec{q} \) to the concentration \( C \) by

\[
\vec{q} = -D \text{grad} C + C \vec{v}
\]

where \( \vec{v} \) is the velocity of the liquid or gas, and \( D \) is the diffusion coefficient which may depend on \( C \). Combining (1.7.4) and (1.7.5) yields

\[
\frac{\partial C}{\partial t} = \text{div} (D \text{grad} C) - \text{div}(C \vec{v}).
\]

If \( D \) is constant then

\[
\frac{\partial C}{\partial t} = D \nabla^2 C - \nabla \cdot (C \vec{v}).
\]

If \( \vec{v} \) is negligible or zero then

\[
\frac{\partial C}{\partial t} = D \nabla^2 C
\]

which is the same as (1.3.8).

If \( D \) is relatively negligible then one has a first order PDE

\[
\frac{\partial C}{\partial t} + \vec{v} \cdot \nabla C + C \text{div} \vec{v} = 0.
\]
At steady state (\( t \) large enough) the concentration \( C \) will no longer depend on \( t \). Equation (1.7.6) becomes
\[
\nabla \cdot (D \nabla C) - \nabla \cdot (C \bar{v}) = 0
\]
(1.7.10)
and if \( \bar{v} \) is negligible or zero then
\[
\nabla \cdot (D \nabla C) = 0
\]
(1.7.11)
which is Laplace’s equation.
2 Separation of Variables-Homogeneous Equations

In this chapter we show that the process of separation of variables solves the one dimensional heat equation subject to various homogeneous boundary conditions and solves Laplace’s equation. All problems in this chapter are homogeneous. We will not be able to give the solution without the knowledge of Fourier series. Therefore these problems will not be fully solved until Chapter 6 after we discuss Fourier series.

2.1 Parabolic equation in one dimension

In this section we show how separation of variables is applied to solve a simple problem of heat conduction in a bar whose ends are held at zero temperature.

\[ u_t = ku_{xx}, \quad (2.1.1) \]
\[ u(0, t) = 0, \quad \text{zero temperature on the left}, \quad (2.1.2) \]
\[ u(L, t) = 0, \quad \text{zero temperature on the right}, \quad (2.1.3) \]
\[ u(x, 0) = f(x), \quad \text{given initial distribution of temperature}. \quad (2.1.4) \]

Note that the equation must be linear and for the time being also homogeneous (no heat sources or sinks). The boundary conditions must also be linear and homogeneous. In Chapter 8 we will show how inhomogeneous boundary conditions can be transferred to a source/sink and then how to solve inhomogeneous partial differential equations. The method there requires the knowledge of eigenfunctions which are the solutions of the spatial parts of the homogeneous problems with homogeneous boundary conditions.

The idea of separation of variables is to assume a solution of the form

\[ u(x, t) = X(x)T(t), \quad (2.1.5) \]

that is the solution can be written as a product of a function of \( x \) and a function of \( t \). Differentiate (2.1.5) and substitute in (2.1.1) to obtain

\[ X(x)\ddot{T}(t) = kX''(x)T(t), \quad (2.1.6) \]

where prime denotes differentiation with respect to \( x \) and dot denotes time derivative. In order to separate the variables, we divide the equation by \( kX(x)T(t) \),

\[ \frac{\ddot{T}(t)}{kT(t)} = \frac{X''(x)}{X(x)}. \quad (2.1.7) \]

The left hand side depends only on \( t \) and the right hand side only on \( x \). If we fix one variable, say \( t \), and vary the other, then the left hand side cannot change (\( t \) is fixed) therefore the right hand side cannot change. This means that each side is really a constant. We denote that so called separation constant by \(-\lambda\). Now we have two ordinary differential equations

\[ X''(x) = -\lambda X(x), \quad (2.1.8) \]
\[ \dot{T}(t) = -k\lambda T(t). \] 

Remark: This does NOT mean that the separation constant is negative.

The homogeneous boundary conditions can be used to provide boundary conditions for (2.1.8). These are

\[ X(0)T(t) = 0, \]
\[ X(L)T(t) = 0. \]

Since \( T(t) \) cannot be zero (otherwise the solution \( u(x, t) = X(x)T(t) \) is zero), then

\[ X(0) = 0, \] 
\[ X(L) = 0. \] 

First we solve (2.1.8) subject to (2.1.10)-(2.1.11). This can be done by analyzing the following 3 cases. (We will see later that the separation constant \( \lambda \) is real.)

**case 1:** \( \lambda < 0 \).

The solution of (2.1.8) is

\[ X(x) = Ae^{\sqrt{\mu}x} + Be^{-\sqrt{\mu}x}, \]

where \( \mu = -\lambda > 0 \).

Recall that one should try \( e^{rx} \) which leads to the characteristic equation \( r^2 = \mu \). Using the boundary conditions, we have two equations for the parameters \( A, B \)

\[ A + B = 0, \]
\[ Ae^{\sqrt{\mu}L} + Be^{-\sqrt{\mu}L} = 0. \]

Solve (2.1.13) for \( B \) and substitute in (2.1.14)

\[ B = -A \]
\[ A \left( e^{\sqrt{\mu}L} - e^{-\sqrt{\mu}L} \right) = 0. \]

Note that

\[ e^{\sqrt{\mu}L} - e^{-\sqrt{\mu}L} = 2 \sinh \sqrt{\mu}L \neq 0 \]

Therefore \( A = 0 \) which implies \( B = 0 \) and thus the solution is trivial (the zero solution).

Later we will see the use of writing the solution of (2.1.12) in one of the following four forms

\[ X(x) = Ae^{\sqrt{\mu}x} + Be^{-\sqrt{\mu}x} = C \cosh \sqrt{\mu}x + D \sinh \sqrt{\mu}x = E \cosh \left( \sqrt{\mu}x + F \right) = G \sinh \left( \sqrt{\mu}x + H \right). \]

In figure 6 we have plotted the hyperbolic functions \( \sinh x \) and \( \cosh x \), so one can see that the hyperbolic sine vanishes only at one point and the hyperbolic cosine never vanishes.

**case 2:** \( \lambda = 0 \).
This leads to

\[ X''(x) = 0, \quad (2.1.16) \]
\[ X(0) = 0, \]
\[ X(L) = 0. \]

The ODE has a solution

\[ X(x) = Ax + B. \quad (2.1.17) \]

Using the boundary conditions

\[ A \cdot 0 + B = 0, \]
\[ A \cdot L + B = 0, \]

we have

\[ B = 0, \]
\[ A = 0, \]

and thus

\[ X(x) = 0, \]

which is the trivial solution (leads to \( u(x, t) = 0 \)) and thus of no interest.

\underline{case 3: \( \lambda > 0 \).}

The solution in this case is

\[ X(x) = A \cos \sqrt{\lambda} x + B \sin \sqrt{\lambda} x. \quad (2.1.18) \]
The first boundary condition leads to

\[ X(0) = A \cdot 1 + B \cdot 0 = 0 \]

which implies

\[ A = 0. \]

Therefore, the second boundary condition (with \( A = 0 \)) becomes

\[ B \sin \sqrt{\lambda}L = 0. \]

Clearly \( B \neq 0 \) (otherwise the solution is trivial again), therefore

\[ \sin \sqrt{\lambda}L = 0, \]

and thus

\[ \sqrt{\lambda}L = n\pi, \quad n = 1, 2, \ldots \quad \text{(since } \lambda > 0, \text{ then } n \geq 1) \]

and

\[ \lambda_n = \left( \frac{n\pi}{L} \right)^2, \quad n = 1, 2, \ldots \]

(2.1.20)

These are called the eigenvalues. The solution (2.1.18) becomes

\[ X_n(x) = B_n \sin \frac{n\pi}{L}x, \quad n = 1, 2, \ldots \]

(2.1.21)

The functions \( X_n \) are called eigenfunctions or modes. There is no need to carry the constants \( B_n \), since the eigenfunctions are unique only to a multiplicative scalar (i.e. if \( X_n \) is an eigenfunction then \( KX_n \) is also an eigenfunction).

The eigenvalues \( \lambda_n \) will be substituted in (2.1.9) before it is solved, therefore

\[ \dot{T}_n(t) = -k \left( \frac{n\pi}{L} \right)^2 T_n. \]

(2.1.22)

The solution is

\[ T_n(t) = e^{-k \left( \frac{n\pi}{L} \right)^2 t}, \quad n = 1, 2, \ldots \]

(2.1.23)

Combine (2.1.21) and (2.1.23) with (2.1.5)

\[ u_n(x, t) = e^{-k \left( \frac{n\pi}{L} \right)^2 t} \sin \frac{n\pi}{L}x, \quad n = 1, 2, \ldots \]

(2.1.24)

Since the PDE is linear, the linear combination of all the solutions \( u_n(x, t) \) is also a solution

\[ u(x, t) = \sum_{n=1}^{\infty} b_n e^{-k \left( \frac{n\pi}{L} \right)^2 t} \sin \frac{n\pi}{L}x. \]

(2.1.25)

This is known as the principle of superposition. As in power series solution of ODEs, we have to prove that the infinite series converges (see section 3.5). This solution satisfies the PDE and the boundary conditions. To find \( b_n \), we must use the initial condition and this will be done after we learn Fourier series.
2.2 Other Homogeneous Boundary Conditions

If one has to solve the heat equation subject to one of the following sets of boundary conditions
1. 
   \[ u(0, t) = 0, \quad u_x(L, t) = 0. \]  

2. 
   \[ u_x(0, t) = 0, \quad u(L, t) = 0. \]  

3. 
   \[ u_x(0, t) = 0, \quad u_x(L, t) = 0. \]  

4. 
   \[ u(0, t) = u(L, t), \quad u_x(0, t) = u_x(L, t). \]  

the procedure will be similar. In fact, (2.1.8) and (2.1.9) are unaffected. In the first case, (2.2.1)-(2.2.2) will be
\[ X(0) = 0, \quad X'(L) = 0. \]  

It is left as an exercise to show that
\[ \lambda_n = \left( n - \frac{1}{2} \right)^2 \frac{\pi^2}{L^2}, \quad n = 1, 2, \ldots \]  
\[ X_n = \sin \left( n - \frac{1}{2} \right) \frac{\pi}{L} x, \quad n = 1, 2, \ldots \]  
The boundary conditions (2.2.3)-(2.2.4) lead to
\[ X'(0) = 0, \quad X(L) = 0. \]  

and the eigenpairs are
\[ \lambda_n = \left( n - \frac{1}{2} \right)^2 \frac{\pi^2}{L^2}, \quad n = 1, 2, \ldots \]  
\[ X_n = \cos \left( n - \frac{1}{2} \right) \frac{\pi}{L} x, \quad n = 1, 2, \ldots \]  
The third case leads to
\[ X'(0) = 0, \]
\begin{equation}
X'(L) = 0.
\end{equation}

Here the eigenpairs are
\begin{align}
\lambda_0 &= 0, & X_0 &= 1, \quad (2.2.19) \\
\lambda_n &= \left( \frac{n\pi}{L} \right)^2, & n &= 1, 2, \ldots \quad (2.2.21) \\
X_n &= \cos \frac{n\pi}{L} x, & n &= 1, 2, \ldots \quad (2.2.22)
\end{align}

The case of periodic boundary conditions require detailed solution.

\underline{case 1:} \(\lambda < 0\).

The solution is given by (2.1.12)
\[ X(x) = Ae^{\sqrt{\mu}x} + Be^{-\sqrt{\mu}x}, \quad \mu = -\lambda > 0. \]

The boundary conditions (2.2.7)-(2.2.8) imply
\begin{align}
A + B &= Ae^{\sqrt{\mu}L} + Be^{-\sqrt{\mu}L}, \quad (2.2.23) \\
A\sqrt{\mu}B - B\sqrt{\mu}A &= A\sqrt{\mu}e^{\sqrt{\mu}L} - B\sqrt{\mu}e^{-\sqrt{\mu}L}. \quad (2.2.24)
\end{align}

This system can be written as
\begin{align}
A \left( 1 - e^{\sqrt{\mu}L} \right) + B \left( 1 - e^{-\sqrt{\mu}L} \right) &= 0, \quad (2.2.25) \\
\sqrt{\mu}A \left( 1 - e^{\sqrt{\mu}L} \right) + \sqrt{\mu}B \left( -1 + e^{-\sqrt{\mu}L} \right) &= 0. \quad (2.2.26)
\end{align}

This homogeneous system can have a solution only if the determinant of the coefficient matrix is zero, i.e.
\[ \begin{vmatrix}
1 - e^{\sqrt{\mu}L} & 1 - e^{-\sqrt{\mu}L} \\
\left( 1 - e^{\sqrt{\mu}L} \right) \sqrt{\mu} & \left( -1 + e^{-\sqrt{\mu}L} \right) \sqrt{\mu}
\end{vmatrix} = 0. \]

Evaluating the determinant, we get
\[ 2\sqrt{\mu} \left( e^{\sqrt{\mu}L} + e^{-\sqrt{\mu}L} - 2 \right) = 0, \]

which is not possible for \(\mu > 0\).

\underline{case 2:} \(\lambda = 0\).

The solution is given by (2.1.17). To use the boundary conditions, we have to differentiate \(X(x)\),
\[ X'(x) = A, \quad (2.2.27) \]

The conditions (2.2.8) and (2.2.7) correspondingly imply
\[ A = A, \]
Thus for the eigenvalue
\[ \lambda_0 = 0, \]  
(2.2.28)
the eigenfunction is
\[ X_0(x) = 1. \]  
(2.2.29)

\textbf{case 3: } \( \lambda > 0. \) 
The solution is given by
\[ X(x) = A \cos \sqrt{\lambda} x + B \sin \sqrt{\lambda} x. \]  
(2.2.30)
The boundary conditions give the following equations for \( A, B, \)
\[ A = A \cos \sqrt{\lambda} L + B \sin \sqrt{\lambda} L, \]
\[ \sqrt{\lambda} B = -\sqrt{\lambda} A \sin \sqrt{\lambda} L + \sqrt{\lambda} B \cos \sqrt{\lambda} L, \]
or
\[ A \left( 1 - \cos \sqrt{\lambda} L \right) - B \sin \sqrt{\lambda} L = 0, \]  
(2.2.31)
\[ A \sqrt{\lambda} \sin \sqrt{\lambda} L + B \sqrt{\lambda} \left( 1 - \cos \sqrt{\lambda} L \right) = 0. \]  
(2.2.32)
The determinant of the coefficient matrix
\[ \begin{vmatrix} 1 - \cos \sqrt{\lambda} L & -\sin \sqrt{\lambda} L \\ \sqrt{\lambda} \sin \sqrt{\lambda} L & \sqrt{\lambda} \left( 1 - \cos \sqrt{\lambda} L \right) \end{vmatrix} = 0, \]
or
\[ \sqrt{\lambda} \left( 1 - \cos \sqrt{\lambda} L \right)^2 + \sqrt{\lambda} \sin^2 \sqrt{\lambda} L = 0. \]
Expanding and using some trigonometric identities,
\[ 2 \sqrt{\lambda} \left( 1 - \cos \sqrt{\lambda} L \right) = 0, \]
or
\[ 1 - \cos \sqrt{\lambda} L = 0. \]  
(2.2.33)
Thus (2.2.31)-(2.2.32) become
\[ -B \sin \sqrt{\lambda} L = 0, \]
\[ A \sqrt{\lambda} \sin \sqrt{\lambda} L = 0, \]
which imply
\[ \sin \sqrt{\lambda} L = 0. \]  
(2.2.34)
Thus the eigenvalues \( \lambda_n \) must satisfy (2.2.33) and (2.2.34), that is
\[ \lambda_n = \left( \frac{2n\pi}{L} \right)^2, \quad n = 1, 2, \ldots \]  
(2.2.35)
Condition (2.2.34) causes the system to be true for any $A, B$, therefore the eigenfunctions are

$$X_n(x) = \begin{cases} 
\cos \frac{2n\pi}{L} x & n = 1, 2, \ldots \\
\sin \frac{2n\pi}{L} x & n = 1, 2, \ldots 
\end{cases}$$  \hspace{1cm} (2.2.36)

In summary, for periodic boundary conditions

$$\lambda_0 = 0,$$  \hspace{1cm} (2.2.37)

$$X_0(x) = 1,$$  \hspace{1cm} (2.2.38)

$$\lambda_n = \left(\frac{2n\pi}{L}\right)^2, \quad n = 1, 2, \ldots$$  \hspace{1cm} (2.2.39)

$$X_n(x) = \begin{cases} 
\cos \frac{2n\pi}{L} x & n = 1, 2, \ldots \\
\sin \frac{2n\pi}{L} x & n = 1, 2, \ldots 
\end{cases}$$  \hspace{1cm} (2.2.40)

Remark: The ODE for $X$ is the same even when we separate the variables for the wave equation. For Laplace’s equation, we treat either the $x$ or the $y$ as the marching variable (depending on the boundary conditions given).

**Example.**

$$u_{xx} + u_{yy} = 0 \quad 0 \leq x, y \leq 1$$  \hspace{1cm} (2.2.41)

$$u(x, 0) = u_0 = \text{constant}$$  \hspace{1cm} (2.2.42)

$$u(x, 1) = 0$$  \hspace{1cm} (2.2.43)

$$u(0, y) = u(1, y) = 0.$$  \hspace{1cm} (2.2.44)

This leads to

$$X'' + \lambda X = 0$$  \hspace{1cm} (2.2.45)

$$X(0) = X(1) = 0$$  \hspace{1cm} (2.2.46)

and

$$Y'' - \lambda Y = 0$$  \hspace{1cm} (2.2.47)

$$Y(1) = 0.$$  \hspace{1cm} (2.2.48)

The eigenvalues and eigenfunctions are

$$X_n = \sin n\pi x, \quad n = 1, 2, \ldots$$  \hspace{1cm} (2.2.49)

$$\lambda_n = (n\pi)^2, \quad n = 1, 2, \ldots$$  \hspace{1cm} (2.2.50)

The solution for the $y$ equation is then

$$Y_n = \sinh n\pi (y - 1)$$  \hspace{1cm} (2.2.51)
and the solution of the problem is

\[ u(x, y) = \sum_{n=1}^{\infty} \alpha_n \sin n\pi x \sinh n\pi (y - 1) \]

(2.2.52)

and the parameters \( \alpha_n \) can be obtained from the Fourier expansion of the nonzero boundary condition, i.e.

\[ \alpha_n = \frac{2u_0}{n\pi} \frac{(-1)^n - 1}{\sinh n\pi}. \]

(2.2.53)
Problems

1. Consider the differential equation

\[ X''(x) + \lambda X(x) = 0 \]

Determine the eigenvalues \( \lambda \) (assumed real) subject to

a. \( X(0) = X(\pi) = 0 \)

b. \( X'(0) = X'(L) = 0 \)

c. \( X(0) = X'(L) = 0 \)

d. \( X'(0) = X(L) = 0 \)

e. \( X(0) = 0 \) and \( X'(L) + X(L) = 0 \)

Analyze the cases \( \lambda > 0 \), \( \lambda = 0 \) and \( \lambda < 0 \).
2.3 Eigenvalues and Eigenfunctions

As we have seen in the previous sections, the solution of the $X$-equation on a finite interval subject to homogeneous boundary conditions, results in a sequence of eigenvalues and corresponding eigenfunctions. Eigenfunctions are said to describe natural vibrations and standing waves. $X_1$ is the fundamental and $X_i$, $i > 1$ are the harmonics. The eigenvalues are the natural frequencies of vibration. These frequencies do not depend on the initial conditions. This means that the frequencies of the natural vibrations are independent of the method to excite them. They characterize the properties of the vibrating system itself and are determined by the material constants of the system, geometrical factors and the conditions on the boundary.

The eigenfunction $X_n$ specifies the profile of the standing wave. The points at which an eigenfunction vanishes are called “nodal points” (nodal lines in two dimensions). The nodal lines are the curves along which the membrane at rest during eigenvibration. For a square membrane of side $\pi$ the eigenfunction (as can be found in Chapter 4) are $\sin nx \sin my$ and the nodal lines are lines parallel to the coordinate axes. However, in the case of multiple eigenvalues, many other nodal lines occur.

Some boundary conditions may not be exclusive enough to result in a unique solution (up to a multiplicative constant) for each eigenvalue. In case of a double eigenvalue, any pair of independent solutions can be used to express the most general eigenfunction for this eigenvalue. Usually, it is best to choose the two solutions so they are orthogonal to each other. This is necessary for the completeness property of the eigenfunctions. This can be done by adding certain symmetry requirement over and above the boundary conditions, which pick either one or the other. For example, in the case of periodic boundary conditions, each positive eigenvalue has two eigenfunctions, one is even and the other is odd. Thus the symmetry allows us to choose. If symmetry is not imposed then both functions must be taken.

The eigenfunctions, as we proved in Chapter 6 of Neta, form a complete set which is the basis for the method of eigenfunction expansion described in Chapter 5 for the solution of inhomogeneous problems (inhomogeneity in the equation or the boundary conditions).
SUMMARY

\[ X'' + \lambda X = 0 \]

Boundary conditions | Eigenvalues \( \lambda_n \) | Eigenfunctions \( X_n \)
---|---|---
\( X(0) = X(L) = 0 \) | \( \left( \frac{n\pi}{L} \right)^2 \) | \( \sin \frac{n\pi}{L} x \) \hspace{1cm} \( n = 1, 2, \ldots \)
\( X(0) = X'(L) = 0 \) | \( \left[ \frac{(n-\frac{1}{2})\pi}{L} \right]^2 \) | \( \sin \frac{(n-\frac{1}{2})\pi}{L} x \) \hspace{1cm} \( n = 1, 2, \ldots \)
\( X'(0) = X(L) = 0 \) | \( \left[ \frac{(n-\frac{1}{2})\pi}{L} \right]^2 \) | \( \cos \frac{(n-\frac{1}{2})\pi}{L} x \) \hspace{1cm} \( n = 1, 2, \ldots \)
\( X'(0) = X'(L) = 0 \) | \( \left( \frac{n\pi}{L} \right)^2 \) | \( \cos \frac{n\pi}{L} x \) \hspace{1cm} \( n = 0, 1, 2, \ldots \)
\( X(0) = X(L), X'(0) = X'(L) \) | \( \left( \frac{2n\pi}{L} \right)^2 \) | \( \sin \frac{2n\pi}{L} x \) \hspace{1cm} \( n = 1, 2, \ldots \)
| \( \left( \frac{2n\pi}{L} \right)^2 \) | \( \cos \frac{2n\pi}{L} x \) \hspace{1cm} \( n = 0, 1, 2, \ldots \)
3 Fourier Series

In this chapter we discuss Fourier series and the application to the solution of PDEs by the method of separation of variables. In the last section, we return to the solution of the problems in Chapter 4 and also show how to solve Laplace’s equation. We discuss the eigenvalues and eigenfunctions of the Laplacian. The application of these eigenpairs to the solution of the heat and wave equations in bounded domains will follow in Chapter 7 (for higher dimensions and a variety of coordinate systems) and Chapter 8 (for nonhomogeneous problems.)

3.1 Introduction

As we have seen in the previous chapter, the method of separation of variables requires the ability of presenting the initial condition in a Fourier series. Later we will find that generalized Fourier series are necessary. In this chapter we will discuss the Fourier series expansion of $f(x)$, i.e.

$$f(x) \sim \frac{a_0}{2} + \sum_{n=1}^{\infty} \left( a_n \cos \frac{n\pi}{L} x + b_n \sin \frac{n\pi}{L} x \right). \quad (3.1.1)$$

We will discuss how the coefficients are computed, the conditions for convergence of the series, and the conditions under which the series can be differentiated or integrated term by term.

**Definition 11.** A function $f(x)$ is piecewise continuous in $[a, b]$ if there exists a finite number of points $a = x_1 < x_2 < \ldots < x_n = b$, such that $f$ is continuous in each open interval $(x_j, x_{j+1})$ and the one sided limits $f(x_{j+})$ and $f(x_{j+1-})$ exist for all $j \leq n - 1$.

**Examples**

1. $f(x) = x^2$ is continuous on $[a, b]$.

2. $f(x) = \begin{cases} x & 0 \leq x < 1 \\ x^2 - x & 1 \leq x \leq 2 \end{cases}$

The function is piecewise continuous but not continuous because of the point $x = 1$.

3. $f(x) = \frac{1}{x}$ for $-1 \leq x \leq 1$. The function is not piecewise continuous because the one sided limit at $x = 0$ does not exist.

**Definition 12.** A function $f(x)$ is piecewise smooth if $f(x)$ and $f'(x)$ are piecewise continuous.

**Definition 13.** A function $f(x)$ is periodic if $f(x)$ is piecewise continuous and $f(x+p) = f(x)$ for some real positive number $p$ and all $x$. The number $p$ is called a period. The smallest period is called the fundamental period.

**Examples**

1. $f(x) = \sin x$ is periodic of period $2\pi$. 

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2. \( f(x) = \cos x \) is periodic of period \( 2\pi \).

Note: If \( f_i(x), \quad i = 1, 2, \cdots, n \) are all periodic of the same period \( p \) then the linear combination of these functions

\[ \sum_{i=1}^{n} c_i f_i(x) \]

is also periodic of period \( p \).

3.2 Orthogonality

Recall that two vectors \( \vec{a} \) and \( \vec{b} \) in \( \mathbb{R}^n \) are called orthogonal vectors if

\[ \vec{a} \cdot \vec{b} = \sum_{i=1}^{n} a_i b_i = 0. \]

We would like to extend this definition to functions. Let \( f(x) \) and \( g(x) \) be two functions defined on the interval \([\alpha, \beta]\). If we sample the two functions at the same points \( x_i, \quad i = 1, 2, \cdots, n \) then the vectors \( \vec{F} \) and \( \vec{G} \), having components \( f(x_i) \) and \( g(x_i) \) correspondingly, are orthogonal if

\[ \sum_{i=1}^{n} f(x_i) g(x_i) = 0. \]

If we let \( n \) to increase to infinity then we get an infinite sum which is proportional to

\[ \int_{\alpha}^{\beta} f(x) g(x) dx. \]

Therefore, we define orthogonality as follows:

**Definition 14.** Two functions \( f(x) \) and \( g(x) \) are called orthogonal on the interval \((\alpha, \beta)\) with respect to the weight function \( w(x) > 0 \) if

\[ \int_{\alpha}^{\beta} w(x) f(x) g(x) dx = 0. \]

**Example 1**

The functions \( \sin x \) and \( \cos x \) are orthogonal on \([-\pi, \pi]\) with respect to \( w(x) = 1 \),

\[ \int_{-\pi}^{\pi} \sin x \cos x dx = \frac{1}{2} \int_{-\pi}^{\pi} \sin 2x dx = -\frac{1}{4} \cos 2x|_{-\pi}^{\pi} = -\frac{1}{4} + \frac{1}{4} = 0. \]

**Definition 15.** A set of functions \( \{\phi_n(x)\} \) is called orthogonal system with respect to \( w(x) \) on \([\alpha, \beta]\) if

\[ \int_{\alpha}^{\beta} \phi_n(x) \phi_m(x) w(x) dx = 0, \quad \text{for } m \neq n. \] (3.2.1)
**Definition 16.** The norm of a function $f(x)$ with respect to $w(x)$ on the interval $[\alpha, \beta]$ is defined by

$$
\| f \| = \left\{ \int_{\alpha}^{\beta} w(x) f^2(x) dx \right\}^{1/2}
$$

(3.2.2)

**Definition 17.** The set \{\phi_n(x)\} is called orthonormal system if it is an orthogonal system and if

$$
\| \phi_n \| = 1.
$$

(3.2.3)

**Examples**

1. \{\sin \frac{n\pi}{L} x\} is an orthogonal system with respect to $w(x) = 1$ on $[-L, L]$.

For $n \neq m$

$$
\int_{-L}^{L} \sin \frac{n\pi}{L} x \sin \frac{m\pi}{L} x dx
$$

$$
= \int_{-L}^{L} \left[ -\frac{1}{2} \cos \frac{(n+m)\pi}{L} x + \frac{1}{2} \cos \frac{(n-m)\pi}{L} x \right] dx
$$

$$
= \left\{ -\frac{1}{2} \frac{L}{(n+m)\pi} \sin \frac{(n+m)\pi}{L} x + \frac{1}{2} \frac{L}{(n-m)\pi} \sin \frac{(n-m)\pi}{L} x \right\}_{-L}^{L} = 0
$$

2. \{\cos \frac{n\pi}{L} x\} is also an orthogonal system on the same interval. It is easy to show that for $n \neq m$

$$
\int_{-L}^{L} \cos \frac{n\pi}{L} x \cos \frac{m\pi}{L} x dx
$$

$$
= \int_{-L}^{L} \left[ \frac{1}{2} \cos \frac{(n+m)\pi}{L} x + \frac{1}{2} \cos \frac{(n-m)\pi}{L} x \right] dx
$$

$$
= \left\{ \frac{1}{2} \frac{L}{(n+m)\pi} \sin \frac{(n+m)\pi}{L} x + \frac{1}{2} \frac{L}{(n-m)\pi} \sin \frac{(n-m)\pi}{L} x \right\}_{-L}^{L} = 0
$$

3. The set \{1, \cos x, \sin x, \cos 2x, \sin 2x, \cdots, \cos nx, \sin nx, \cdots\} is an orthogonal system on $[-\pi, \pi]$ with respect to the weight function $w(x) = 1$.

We have shown already that

$$
\int_{-\pi}^{\pi} \sin nx \sin mx dx = 0 \quad \text{for } n \neq m
$$

(3.2.4)

$$
\int_{-\pi}^{\pi} \cos nx \cos mx dx = 0 \quad \text{for } n \neq m.
$$

(3.2.5)
The only thing left to show is therefore
\[
\int_{-\pi}^{\pi} 1 \cdot \sin nx \, dx = 0 \quad (3.2.6)
\]
\[
\int_{-\pi}^{\pi} 1 \cdot \cos nx \, dx = 0 \quad (3.2.7)
\]
and
\[
\int_{-\pi}^{\pi} \sin nx \cos mx \, dx = 0 \quad \text{for any } n, m. \quad (3.2.8)
\]
Note that
\[
\int_{-\pi}^{\pi} \sin nx \, dx = -\frac{\cos nx}{n} \bigg|_{-\pi}^{\pi} = -\frac{1}{n} \left( \cos n\pi - \cos (-n\pi) \right) = 0
\]
since
\[
\cos n\pi = \cos (-n\pi) = (-1)^n. \quad (3.2.9)
\]
In a similar fashion we demonstrate (3.2.7). This time the antiderivative \( \frac{1}{n} \sin nx \) vanishes at both ends.

To show (3.2.8) we consider first the case \( n = m \). Thus
\[
\int_{-\pi}^{\pi} \sin nx \cos nx \, dx = \frac{1}{2} \int_{-\pi}^{\pi} \sin 2nx \, dx = -\frac{1}{4n} \cos 2nx \bigg|_{-\pi}^{\pi} = 0
\]
For \( n \neq m \), we can use the trigonometric identity
\[
\sin ax \cos bx = \frac{1}{2} \left[ \sin(a + b)x + \sin(a - b)x \right]. \quad (3.2.10)
\]
Integrating each of these terms gives zero as in (3.2.6). Therefore the system is orthogonal.

### 3.3 Computation of Coefficients

Suppose that \( f(x) \) can be expanded in Fourier series
\[
f(x) \sim a_0 \frac{2}{2} + \sum_{k=1}^{\infty} \left( a_k \cos \frac{k\pi}{L} x + b_k \sin \frac{k\pi}{L} x \right). \quad (3.3.1)
\]
The infinite series may or may not converge. Even if the series converges, it may not give the value of \( f(x) \) at some points. The question of convergence will be left for later. In this section we just give the formulae used to compute the coefficients \( a_k, b_k \).

\[
a_0 = \frac{1}{L} \int_{-L}^{L} f(x) \, dx, \quad (3.3.2)
\]
\[
a_k = \frac{1}{L} \int_{-L}^{L} f(x) \cos \frac{k\pi}{L} x \, dx \quad \text{for } k = 1, 2, \ldots \quad (3.3.3)
\]
\[ b_k = \frac{1}{L} \int_{-L}^{L} f(x) \sin \frac{k\pi}{L} x \, dx \quad \text{for } k = 1, 2, \ldots \] (3.3.4)

Notice that for \( k = 0 \) (3.3.3) gives the same value as \( a_0 \) in (3.3.2). This is the case only if one takes \( \frac{a_0}{2} \) as the first term in (3.3.1), otherwise the constant term is

\[ \frac{1}{2L} \int_{-L}^{L} f(x) \, dx. \] (3.3.5)

The factor \( L \) in (3.3.3)-(3.3.4) is exactly the square of the norm of the functions \( \sin \frac{k\pi}{L} x \) and \( \cos \frac{k\pi}{L} x \). In general, one should write the coefficients as follows:

\[ a_k = \frac{\int_{-L}^{L} f(x) \cos \frac{k\pi}{L} x \, dx}{\int_{-L}^{L} \cos^2 \frac{k\pi}{L} x \, dx}, \quad \text{for } k = 1, 2, \ldots \] (3.3.6)

\[ b_k = \frac{\int_{-L}^{L} f(x) \sin \frac{k\pi}{L} x \, dx}{\int_{-L}^{L} \sin^2 \frac{k\pi}{L} x \, dx}, \quad \text{for } k = 1, 2, \ldots \] (3.3.7)

These two formulae will be very helpful when we discuss generalized Fourier series.

Example 2
Find the Fourier series expansion of

\[ f(x) = x \quad \text{on } [-L, L] \]

\[ a_k = \frac{1}{L} \int_{-L}^{L} x \cos \frac{k\pi}{L} x \, dx \]

\[ = \frac{1}{L} \left[ \frac{L}{k\pi} x \sin \frac{k\pi}{L} x + \left( \frac{L}{k\pi} \right)^2 \cos \frac{k\pi}{L} x \right]_{-L}^{L}. \]

The first term vanishes at both ends and we have

\[ = \frac{1}{L} \left( \frac{L}{k\pi} \right)^2 [\cos k\pi - \cos(-k\pi)] = 0. \]

\[ b_k = \frac{1}{L} \int_{-L}^{L} x \sin \frac{k\pi}{L} x \, dx \]

\[ = \frac{1}{L} \left[ -\frac{L}{k\pi} x \cos \frac{k\pi}{L} x + \left( \frac{L}{k\pi} \right)^2 \sin \frac{k\pi}{L} x \right]_{-L}^{L}. \]

Now the second term vanishes at both ends and thus

\[ b_k = -\frac{1}{k\pi} [L \cos k\pi - (-L) \cos(-k\pi)] = -\frac{2L}{k\pi} \cos k\pi = \frac{2L}{k\pi} (-1)^k = \frac{2L}{k\pi} (-1)^{k+1}. \]
Therefore the Fourier series is

\[ x \sim \sum_{k=1}^{\infty} \frac{2L}{k\pi} (-1)^{k+1} \sin \frac{k\pi}{L} x. \] (3.3.8)

In Figure 7 we graphed the function \( f(x) = x \) and the \( N^{th} \) partial sum for \( N = 1, 5, 10, 20 \). Notice that the partial sums converge to \( f(x) \) except at the endpoints where we observe the well known Gibbs phenomenon. (The discontinuity produces spurious oscillations in the solution).

Example 3
Find the Fourier coefficients of the expansion of

\[ f(x) = \begin{cases} -1 & \text{for } -L < x < 0 \\ 1 & \text{for } 0 < x < L \end{cases} \] (3.3.9)

\[ a_k = \frac{1}{L} \int_{-L}^{0} (-1) \cos \frac{k\pi}{L} x dx + \frac{1}{L} \int_{0}^{L} 1 \cdot \cos \frac{k\pi}{L} x dx \]

\[ = \frac{1}{L} \frac{L}{k\pi} \sin \frac{k\pi}{L} x\Big|_{-L}^{0} + \frac{1}{L} \frac{L}{k\pi} \sin \frac{k\pi}{L} x\Big|_{0}^{L} = 0, \]

\[ a_0 = \frac{1}{L} \int_{-L}^{0} (-1) dx + \frac{1}{L} \int_{0}^{L} 1 dx \]

\[ = -\frac{1}{L} x\bigg|_{-L}^{0} + \frac{1}{L} x\bigg|_{0}^{L} = \frac{1}{L} (-L) + \frac{1}{L} \cdot L = 0, \]
Figure 8: Graph of $f(x)$ given in Example 3 and the $N^{th}$ partial sums for $N = 1, 5, 10, 20$

$$b_k = \frac{1}{L} \int_{-L}^{0} (-1) \sin \frac{k\pi}{L} x \, dx + \frac{1}{L} \int_{0}^{L} 1 \cdot \sin \frac{k\pi}{L} x \, dx$$

$$= \frac{1}{L} \left(-1 \right) \left(-\frac{L}{k\pi}\right) \cos \frac{k\pi}{L} x \bigg|_{-L}^{0} + \frac{1}{L} \left(-\frac{L}{k\pi}\right) \cos \frac{k\pi}{L} x \bigg|_{0}^{L}$$

$$= \frac{1}{k\pi} \left[1 - \cos(-k\pi)\right] - \frac{1}{k\pi} \left[\cos k\pi - 1\right]$$

$$= \frac{2}{k\pi} \left[1 - (-1)^k\right].$$

Therefore the Fourier series is

$$f(x) \sim \sum_{k=1}^{\infty} \frac{2}{k\pi} \left[1 - (-1)^k\right] \sin \frac{k\pi}{L} x. \quad (3.3.10)$$

The graphs of $f(x)$ and the $N^{th}$ partial sums (for various values of $N$) are given in figure 8.

In the last two examples, we have seen that $a_k = 0$. Next, we give an example where all the coefficients are nonzero.

**Example 4**

$$f(x) = \begin{cases} 
\frac{1}{L} x + 1 & -L < x < 0 \\
\frac{1}{x} & 0 < x < L 
\end{cases} \quad (3.3.11)$$
Figure 9: Graph of \( f(x) \) given in Example 4

\[
a_0 = \frac{1}{L} \int_{-L}^{0} \left( \frac{1}{L} x + 1 \right) \, dx + \frac{1}{L} \int_{0}^{L} x \, dx
\]

\[
= \frac{1}{L^2} \left[ \frac{1}{2} x^2 \right]_{-L}^{0} + \frac{1}{L} \left[ \frac{1}{2} x \right]_{0}^{L} + \frac{1}{L} \left[ \frac{1}{2} x^2 \right]_{0}^{L}
\]

\[
= -\frac{1}{2} + 1 + \frac{L}{2} = \frac{L + 1}{2},
\]

\[
a_k = \frac{1}{L} \int_{-L}^{0} \left( \frac{1}{L} x + 1 \right) \cos \frac{k\pi}{L} x \, dx + \frac{1}{L} \int_{0}^{L} x \cos \frac{k\pi}{L} x \, dx
\]

\[
= \frac{1}{L^2} \left[ \frac{L}{k\pi} x \sin \frac{k\pi}{L} x + \left( \frac{L}{k\pi} \right)^2 \cos \frac{k\pi}{L} x \right]_{-L}^{0}
\]

\[
+ \frac{1}{L} \left[ \frac{L}{k\pi} x \sin \frac{k\pi}{L} x + \left( \frac{L}{k\pi} \right)^2 \cos \frac{k\pi}{L} x \right]_{0}^{L}
\]

\[
= \frac{1}{L^2} \left( \frac{L}{k\pi} \right)^2 - \frac{1}{L^2} \left( \frac{L}{k\pi} \right)^2 \cos k\pi + \frac{1}{L} \left( \frac{L}{k\pi} \right)^2 \cos k\pi - \frac{1}{L} \left( \frac{L}{k\pi} \right)^2
\]

\[
= \frac{1}{\left( \frac{L}{k\pi} \right)^2} - \frac{1}{\left( \frac{L}{k\pi} \right)^2} (-1)^k = \frac{1}{\left( \frac{L}{k\pi} \right)^2} \left( 1 - (-1)^k \right),
\]
Figure 10: Graph of \( f(x) \) given by example 4 \((L = 1)\) and the \( N^{th} \) partial sums for \( N = 1, 5, 10, 20 \). Notice that for \( L = 1 \) all cosine terms and odd sine terms vanish, thus the first term is the constant \( \frac{1}{2} \).

\[
\begin{align*}
b_k & = \frac{1}{L} \int_{-L}^{0} \left( \frac{1}{L} x + 1 \right) \sin \frac{k\pi}{L} x dx + \frac{1}{L} \int_{0}^{L} x \sin \frac{k\pi}{L} x dx \\
& = \frac{1}{L^2} \left[ -\frac{L}{k\pi} \cos \frac{k\pi}{L} x + \left( \frac{L}{k\pi} \right)^2 \sin \frac{k\pi}{L} x \right]_{-L}^{0} \\
& \quad + \frac{1}{L} \frac{L}{k\pi} \left( -\cos \frac{k\pi}{L} x \right)_{-L}^{0} + \frac{1}{L} \left[ -\frac{L}{k\pi} \cos \frac{k\pi}{L} x + \left( \frac{L}{k\pi} \right)^2 \sin \frac{k\pi}{L} x \right]_{0}^{L} \\
& = \frac{1}{L^2} \frac{L}{k\pi} \left( -1 \right) \cos \frac{k\pi}{L} - \frac{1}{k\pi} \cos k\pi + \frac{1}{k\pi} \cos k\pi - \frac{L}{k\pi} \cos k\pi \\
& = -\frac{1}{k\pi} \left( 1 + (-1)^k L \right),
\end{align*}
\]

therefore the Fourier series is

\[
f(x) = \frac{L + 1}{4} + \sum_{k=1}^{\infty} \left\{ \frac{1 - L}{(k\pi)^2} \left[ 1 - (-1)^k \right] \cos \frac{k\pi}{L} x - \frac{1}{k\pi} \left[ 1 + (-1)^k L \right] \sin \frac{k\pi}{L} x \right\}
\]

The sketches of \( f(x) \) and the \( N^{th} \) partial sums are given in figures 10-12 for various values of \( L \).
Figure 11: Graph of $f(x)$ given by example 4 ($L = 1/2$) and the $N^{th}$ partial sums for $N = 1, 5, 10, 20$

Figure 12: Graph of $f(x)$ given by example 4 ($L = 2$) and the $N^{th}$ partial sums for $N = 1, 5, 10, 20$
Problems

1. For the following functions, sketch the Fourier series of \( f(x) \) on the interval \([-L, L]\). Compare \( f(x) \) to its Fourier series
   
a. \( f(x) = 1 \)
   b. \( f(x) = x^2 \)
   c. \( f(x) = e^x \)
   d. \( f(x) = \begin{cases} \frac{1}{2}x & x < 0 \\ 3x & x > 0 \end{cases} \)
   e. \( f(x) = \begin{cases} 0 & x < \frac{L}{2} \\ x^2 & x > \frac{L}{2} \end{cases} \)

2. Sketch the Fourier series of \( f(x) \) on the interval \([-L, L]\) and evaluate the Fourier coefficients for each
   
a. \( f(x) = x \)
   b. \( f(x) = \sin \frac{\pi}{L} x \)
   c. \( f(x) = \begin{cases} 1 & |x| < \frac{L}{2} \\ 0 & |x| > \frac{L}{2} \end{cases} \)

3. Show that the Fourier series operation is linear, i.e. the Fourier series of \( \alpha f(x) + \beta g(x) \) is the sum of the Fourier series of \( f(x) \) and \( g(x) \) multiplied by the corresponding constant.
3.4 Relationship to Least Squares

It can be shown that the Fourier series expansion of $f(x)$ gives the best approximation of $f(x)$ in the sense of least squares. That is, if one minimizes the squares of differences between $f(x)$ and the $n^{th}$ partial sum of the series

$$\frac{a_0}{2} + \sum_{k=1}^{\infty} \left( a_k \cos \frac{k\pi}{L} x + b_k \sin \frac{k\pi}{L} x \right)$$

then the coefficients $a_0$, $a_k$ and $b_k$ are exactly the Fourier coefficients given by (3.3.6)-(3.3.7).

3.5 Convergence

If $f(x)$ is piecewise smooth on $[-L, L]$ then the series converges to either the periodic extension of $f(x)$, where the periodic extension is continuous, or to the average of the two limits, where the periodic extension has a jump discontinuity.

3.6 Fourier Cosine and Sine Series

In the examples in the last section we have seen Fourier series for which all $a_k$ are zero. In such a case the Fourier series includes only sine functions. Such a series is called a Fourier sine series. The problems discussed in the previous chapter led to Fourier sine series or Fourier cosine series depending on the boundary conditions.

Let us now recall the definition of odd and even functions. A function $f(x)$ is called odd if

$$f(-x) = -f(x)$$

and even, if

$$f(-x) = f(x).$$

Since $\sin kx$ is an odd function, the sum is also an odd function, therefore a function $f(x)$ having a Fourier sine series expansion is odd. Similarly, an even function will have a Fourier cosine series expansion.

Example 5

$$f(x) = x,$$

on $[-L, L].$ (3.6.3)

The function is odd and thus the Fourier series expansion will have only sine terms, i.e. all $a_k = 0$. In fact we have found in one of the examples in the previous section that

$$f(x) \sim \sum_{k=1}^{\infty} \frac{2L}{k\pi} (-1)^{k+1} \sin \frac{k\pi}{L} x$$

(3.6.4)

Example 6

$$f(x) = x^2$$

on $[-L, L].$ (3.6.5)
The function is even and thus all $b_k$ must be zero.

$$a_0 = \frac{1}{L} \int_{-L}^{L} x^2 \, dx = \frac{2}{L} \int_{0}^{L} x^2 \, dx = \frac{2}{3} \left. x^3 \right|_{0}^{L} = \frac{2L^2}{3}. \quad (3.6.6)$$

$$a_k = \frac{1}{L} \int_{-L}^{L} x^2 \cos \frac{k\pi}{L} \, dx =$$

Use table of integrals

$$a_k = \frac{1}{L} \left[ 2x \left( \frac{L}{k\pi} \right)^2 \cos \frac{k\pi}{L} x \bigg|_{-L}^{L} + \left( \frac{k\pi}{L} \right)^2 x^2 - 2 \right] \left( \frac{L}{k\pi} \right)^3 \sin \frac{k\pi}{L} \bigg|_{-L}^{L} \right].$$

The sine terms vanish at both ends and we have

$$a_k = \frac{1}{L} \left[ 4L \left( \frac{L}{k\pi} \right)^2 \cos k\pi = 4 \left( \frac{L}{k\pi} \right)^2 (-1)^k. \quad (3.6.7)$$

Notice that the coefficients of the Fourier sine series can be written as

$$b_k = \frac{2}{L} \int_{0}^{L} f(x) \sin \frac{k\pi}{L} x \, dx, \quad (3.6.8)$$

that is the integration is only on half the interval and the result is doubled. Similarly for the Fourier cosine series

$$a_k = \frac{2}{L} \int_{0}^{L} f(x) \cos \frac{k\pi}{L} x \, dx. \quad (3.6.9)$$

If we go back to the examples in the previous chapter, we notice that the partial differential equation is solved on the interval $[0, L]$. If we end up with Fourier sine series, this means that the initial solution $f(x)$ was extended as an odd function to $[-L, 0]$. It is the odd extension that we expand in Fourier series.

**Example 7**

Give a Fourier cosine series of

$$f(x) = x \quad \text{for} \quad 0 \leq x \leq L. \quad (3.6.10)$$

This means that $f(x)$ is extended as an even function, i.e.

$$f(x) = \begin{cases} -x & -L \leq x \leq 0 \\ x & 0 \leq x \leq L \end{cases} \quad (3.6.11)$$

or

$$f(x) = |x| \quad \text{on} \quad [-L, L]. \quad (3.6.12)$$

The Fourier cosine series will have the following coefficients

$$a_0 = \frac{2}{L} \int_{0}^{L} x \, dx = \frac{2}{L} \left. \frac{1}{2} x^2 \right|_{0}^{L} = L, \quad (3.6.13)$$

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Figure 13: Graph of \( f(x) = x^2 \) and the \( N^{th} \) partial sums for \( N = 1, 5, 10, 20 \)

\[
a_k = \frac{2}{L} \int_{0}^{L} x \cos \frac{k\pi}{L} x \, dx = \frac{2}{L} \left[ \frac{L}{k\pi} x \sin \frac{k\pi}{L} x + \left( \frac{L}{k\pi} \right)^2 \cos \frac{k\pi}{L} x \right]_{0}^{L} = \frac{2}{L} \left[ 0 + \left( \frac{L}{k\pi} \right)^2 \cos k\pi - 0 - \left( \frac{L}{k\pi} \right)^2 \right] = \frac{2}{L} \left( \frac{L}{k\pi} \right)^2 \left[ (-1)^k - 1 \right]. \tag{3.6.14}
\]

Therefore the series is

\[
|x| \sim \frac{L}{2} + \sum_{k=1}^{\infty} \frac{2L}{(k\pi)^2} \left[ (-1)^k - 1 \right] \cos \frac{k\pi}{L} x. \tag{3.6.15}
\]

In the next four figures we have sketched \( f(x) = |x| \) and the \( N^{th} \) partial sums for various values of \( N \).

To sketch the Fourier cosine series of \( f(x) \), we first sketch \( f(x) \) on \([0, L]\), then extend the sketch to \([-L, L]\) as an even function, then extend as a periodic function of period \(2L\). At points of discontinuity, take the average.

To sketch the Fourier sine series of \( f(x) \) we follow the same steps except that we take the odd extension.

**Example 8**

\[
f(x) = \begin{cases} 
\sin \frac{\pi}{L} x, & -L < x < 0 \\
x, & 0 < x < \frac{L}{2} \\
L - x, & \frac{L}{2} < x < L
\end{cases} \tag{3.6.16}
\]
Figure 14: Graph of $f(x) = |x|$ and the $N^{th}$ partial sums for $N = 1, 5, 10, 20$

The Fourier cosine series and the Fourier sine series will ignore the definition on the interval $[-L, 0]$ and take only the definition on $[0, L]$. The sketches follow on figures 15-17:

Figure 15: Sketch of $f(x)$ given in Example 8

Notes:

1. The Fourier series of a piecewise smooth function $f(x)$ is continuous if and only if $f(x)$ is continuous and $f(-L) = f(L)$.
2. The Fourier cosine series of a piecewise smooth function $f(x)$ is continuous if and only if $f(x)$ is continuous. (The condition $f(-L) = f(L)$ is automatically satisfied.)
3. The Fourier sine series of a piecewise smooth function $f(x)$ is continuous if and only if $f(x)$ is continuous and $f(0) = f(L)$. 
Figure 16: Sketch of the Fourier sine series and the periodic odd extension

Figure 17: Sketch of the Fourier cosine series and the periodic even extension

Example 9
The previous example was for a function satisfying this condition. Suppose we have the following $f(x)$

$$f(x) = \begin{cases} 
0 & -L < x < 0 \\
x & 0 < x < L
\end{cases} \quad (3.6.17)$$

The sketches of $f(x)$, its odd extension and its Fourier sine series are given in figures 18-20 correspondingly.
Figure 19: Sketch of the odd extension of $f(x)$

Figure 20: Sketch of the Fourier sine series is not continuous since $f(0) \neq f(L)$
Problems

1. For each of the following functions
   i. Sketch \( f(x) \)
   ii. Sketch the Fourier series of \( f(x) \)
   iii. Sketch the Fourier sine series of \( f(x) \)
   iv. Sketch the Fourier cosine series of \( f(x) \)
   a. \( f(x) = \begin{cases} 
   x & x < 0 \\
   1 + x & x > 0 
   \end{cases} \)
   b. \( f(x) = e^x \)
   c. \( f(x) = 1 + x^2 \)
   d. \( f(x) = \begin{cases} 
   \frac{1}{2}x + 1 & -2 < x < 0 \\
   x & 0 < x < 2 
   \end{cases} \)

2. Sketch the Fourier sine series of \( f(x) = \cos \frac{\pi}{L}x \).

   Roughly sketch the sum of the first three terms of the Fourier sine series.

3. Sketch the Fourier cosine series and evaluate its coefficients for

   \[
   f(x) = \begin{cases} 
   1 & x < \frac{L}{6} \\
   3 & \frac{L}{6} < x < \frac{L}{2} \\
   0 & \frac{L}{2} < x 
   \end{cases}
   \]

4. Fourier series can be defined on other intervals besides \([-L, L]\). Suppose \( g(y) \) is defined on \([a, b]\) and periodic with period \( b - a \). Evaluate the coefficients of the Fourier series.

5. Expand \( f(x) = \begin{cases} 
   1 & 0 < x < \frac{\pi}{2} \\
   0 & \frac{\pi}{2} < x < \pi 
   \end{cases} \)

   in a series of \( \sin nx \).

   a. Evaluate the coefficients explicitly.
   b. Graph the function to which the series converges to over \(-2\pi < x < 2\pi\).
3.7 Full solution of Several Problems

In this section we give the Fourier coefficients for each of the solutions in the previous chapter.

Example 10

\[ u_t = ku_{xx}, \]  
(3.7.1)
\[ u(0, t) = 0, \]  
(3.7.2)
\[ u(L, t) = 0, \]  
(3.7.3)
\[ u(x, 0) = f(x). \]  
(3.7.4)

The solution given in the previous chapter is

\[ u(x, t) = \sum_{n=1}^{\infty} b_n e^{-k(\frac{\pi}{L})^2 t} \sin \frac{n\pi}{L} x. \]  
(3.7.5)

Upon substituting \( t = 0 \) in (3.7.5) and using (3.7.4) we find that

\[ f(x) = \sum_{n=1}^{\infty} b_n \sin \frac{n\pi}{L} x, \]  
(3.7.6)

that is \( b_n \) are the coefficients of the expansion of \( f(x) \) into Fourier sine series. Therefore

\[ b_n = \frac{2}{L} \int_0^L f(x) \sin \frac{n\pi}{L} x \, dx. \]  
(3.7.7)

Example 11

\[ u_t = ku_{xx}, \]  
(3.7.8)
\[ u(0, t) = u(L, t), \]  
(3.7.9)
\[ u_x(0, t) = u_x(L, t), \]  
(3.7.10)
\[ u(x, 0) = f(x). \]  
(3.7.11)

The solution found in the previous chapter is

\[ u(x, t) = \frac{a_0}{2} + \sum_{n=1}^{\infty} (a_n \cos \frac{2n\pi}{L} x + b_n \sin \frac{2n\pi}{L} x) e^{-k(\frac{2n\pi}{L})^2 t}, \]  
(3.7.12)

As in the previous example, we take \( t = 0 \) in (3.7.12) and compare with (3.7.11) we find that

\[ f(x) = \frac{a_0}{2} + \sum_{n=1}^{\infty} (a_n \cos \frac{2n\pi}{L} x + b_n \sin \frac{2n\pi}{L} x). \]  
(3.7.13)

Therefore (notice that the period is \( L \))

\[ a_n = \frac{2}{L} \int_0^L f(x) \cos \frac{2n\pi}{L} x \, dx, \quad n = 0, 1, 2, \ldots \]  
(3.7.14)
\[ b_n = \frac{2}{L} \int_0^L f(x) \sin \frac{2n\pi}{L} x \, dx, \quad n = 1, 2, \ldots \] (3.7.15)

(Note that \[ \int_0^L \sin^2 \frac{2n\pi}{L} x \, dx = \frac{L}{2} \])

Example 12
Solve Laplace’s equation inside a rectangle:
\[ u_{xx} + u_{yy} = 0, \quad 0 \leq x \leq L, \quad 0 \leq y \leq H, \] (3.7.16)
subject to the boundary conditions:
\[ u(0, y) = g_1(y), \] (3.7.17)
\[ u(L, y) = g_2(y), \] (3.7.18)
\[ u(x, 0) = f_1(x), \] (3.7.19)
\[ u(x, H) = f_2(x). \] (3.7.20)

Note that this is the first problem for which the boundary conditions are inhomogeneous. We will show that \( u(x, y) \) can be computed by summing up the solutions of the following four problems each having 3 homogeneous boundary conditions:

Problem 1:
\[ u_{xx}^1 + u_{yy}^1 = 0, \quad 0 \leq x \leq L, \quad 0 \leq y \leq H, \] (3.7.21)
subject to the boundary conditions:
\[ u^1(0, y) = g_1(y), \] (3.7.22)
\[ u^1(L, y) = 0, \] (3.7.23)
\[ u^1(x, 0) = 0, \] (3.7.24)
\[ u^1(x, H) = 0. \] (3.7.25)

Problem 2:
\[ u_{xx}^2 + u_{yy}^2 = 0, \quad 0 \leq x \leq L, \quad 0 \leq y \leq H, \] (3.7.26)
subject to the boundary conditions:
\[ u^2(0, y) = 0, \] (3.7.27)
\[ u^2(L, y) = g_2(y), \] (3.7.28)
\[ u^2(x, 0) = 0, \] (3.7.29)
\[ u^2(x, H) = 0. \] (3.7.30)

Problem 3:
\[ u^3_{xx} + u^3_{yy} = 0, \quad 0 \leq x \leq L, \quad 0 \leq y \leq H, \quad (3.7.31) \]

subject to the boundary conditions:

\[ u^3(0, y) = 0, \quad (3.7.32) \]
\[ u^3(L, y) = 0, \quad (3.7.33) \]
\[ u^3(x, 0) = f_1(x), \quad (3.7.34) \]
\[ u^3(x, H) = 0. \quad (3.7.35) \]

Problem 4:

\[ u^4_{xx} + u^4_{yy} = 0, \quad 0 \leq x \leq L, \quad 0 \leq y \leq H, \quad (3.7.36) \]

subject to the boundary conditions:

\[ u^4(0, y) = 0, \quad (3.7.37) \]
\[ u^4(L, y) = 0, \quad (3.7.38) \]
\[ u^4(x, 0) = 0, \quad (3.7.39) \]
\[ u^4(x, H) = f_2(x). \quad (3.7.40) \]

It is clear that since \( u^1, u^2, u^3, \) and \( u^4 \) all satisfy Laplace’s equation, then

\[ u = u^1 + u^2 + u^3 + u^4 \]

also satisfies that same PDE (the equation is linear and the result follows from the principle of superposition.) It is also as straightforward to show that \( u \) satisfies the inhomogeneous boundary conditions (3.7.17)-(3.7.20).

We will solve only problem 3 and leave the other 3 problems as exercises.

Separation of variables method applied to (3.7.31)-(3.7.35) leads to the following two ODEs

\[ X'' + \lambda X = 0, \quad (3.7.41) \]
\[ X(0) = 0, \quad (3.7.42) \]
\[ X(L) = 0, \quad (3.7.43) \]
\[ Y'' - \lambda Y = 0, \quad (3.7.44) \]
\[ Y(H) = 0. \quad (3.7.45) \]

The solution of the first was obtained earlier, see (2.1.20)-(2.1.21)

\[ X_n = \sin \frac{n\pi}{L} x, \quad (3.7.46) \]
\[ \lambda_n = \left( \frac{n\pi}{L} \right)^2, \quad n = 1, 2, \ldots \quad (3.7.47) \]
Using these eigenvalues in (3.7.44) we have

$$Y''_n - \left( \frac{n\pi}{L} \right)^2 Y_n = 0$$  \hspace{1cm} (3.7.48)

which has a solution

$$Y_n = A_n \cosh \frac{n\pi}{L} y + B_n \sinh \frac{n\pi}{L} y.$$  \hspace{1cm} (3.7.49)

Because of the boundary condition and the fact that \( \sinh y \) vanishes at zero, we prefer to write the solution as a shifted hyperbolic sine (see (2.1.15)), i.e.

$$Y_n = A_n \sinh \frac{n\pi}{L} (y - H).$$  \hspace{1cm} (3.7.50)

Clearly, this vanishes at \( y = H \) and thus (3.7.45) is also satisfied. Therefore, we have

$$u^3(x, y) = \sum_{n=1}^{\infty} A_n \sinh \frac{n\pi}{L} (y - H) \sin \frac{n\pi}{L} x.$$  \hspace{1cm} (3.7.51)

In the exercises, the reader will have to show that

$$u^1(x, y) = \sum_{n=1}^{\infty} B_n \sinh \frac{n\pi}{H} (x - L) \sin \frac{n\pi}{H} y,$$  \hspace{1cm} (3.7.52)

$$u^2(x, y) = \sum_{n=1}^{\infty} C_n \sinh \frac{n\pi}{H} x \sin \frac{n\pi}{H} y,$$  \hspace{1cm} (3.7.53)

$$u^4(x, y) = \sum_{n=1}^{\infty} D_n \sinh \frac{n\pi}{L} y \sin \frac{n\pi}{L} x.$$  \hspace{1cm} (3.7.54)

To get \( A_n, B_n, C_n, \) and \( D_n \) we will use the inhomogeneous boundary condition in each problem:

$$A_n \sinh \frac{n\pi}{L} (-H) = \frac{2}{L} \int_{0}^{L} f_1(x) \sin \frac{n\pi}{L} x \, dx,$$  \hspace{1cm} (3.7.55)

$$B_n \sinh \frac{n\pi}{H} (-L) = \frac{2}{H} \int_{0}^{H} g_1(y) \sin \frac{n\pi}{H} y \, dy,$$  \hspace{1cm} (3.7.56)

$$C_n \sinh \frac{n\pi}{H} L = \frac{2}{H} \int_{0}^{H} g_2(y) \sin \frac{n\pi}{H} y \, dy,$$  \hspace{1cm} (3.7.57)

$$D_n \sinh \frac{n\pi}{L} H = \frac{2}{L} \int_{0}^{L} f_2(x) \sin \frac{n\pi}{L} x \, dx.$$  \hspace{1cm} (3.7.58)

Example 13

Solve Laplace’s equation inside a circle of radius \( a \),

$$\nabla^2 u = \frac{1}{r} \frac{\partial}{\partial r} \left( r \frac{\partial u}{\partial r} \right) + \frac{1}{r^2} \frac{\partial^2 u}{\partial \theta^2} = 0,$$  \hspace{1cm} (3.7.59)
subject to

\[ u(a, \theta) = f(\theta). \quad (3.7.60) \]

Let

\[ u(r, \theta) = R(r)\Theta(\theta), \quad (3.7.61) \]

then

\[ \Theta \frac{1}{r} (r R')' + \frac{1}{r^2} R\Theta'' = 0. \]

Multiply by \( \frac{r^2}{R\Theta} \)

\[ \frac{r (r R')'}{R} = - \frac{\Theta''}{\Theta} = \mu. \quad (3.7.62) \]

Thus the ODEs are

\[ \Theta'' + \mu \Theta = 0, \quad (3.7.63) \]

and

\[ r(r R')' - \mu R = 0. \quad (3.7.64) \]

The solution must be periodic in \( \theta \) since we have a complete disk. Thus the boundary conditions for \( \Theta \) are

\[ \Theta(0) = \Theta(2\pi), \quad (3.7.65) \]

\[ \Theta'(0) = \Theta'(2\pi). \quad (3.7.66) \]

The solution of the \( \Theta \) equation is given by

\[ \mu_0 = 0, \quad \Theta_0 = 1, \quad (3.7.67) \]

\[ \mu_n = n^2, \quad \Theta_n = \left\{ \begin{array}{ll}
\sin n\theta & n = 1, 2, \ldots \\
\cos n\theta & \end{array} \right. \quad (3.7.68) \]

The only boundary condition for \( R \) is the boundedness, i.e.

\[ |R(0)| < \infty. \quad (3.7.69) \]

The solution for the \( R \) equation is given by (see Euler’s equation in any ODE book)

\[ R_0 = C_0 \ln r + D_0, \quad (3.7.70) \]

\[ R_n = C_n r^{-n} + D_n r^n. \quad (3.7.71) \]

Since \( \ln r \) and \( r^{-n} \) are not finite at \( r = 0 \) (which is in the domain), we must have \( C_0 = C_n = 0 \). Therefore

\[ u(r, \theta) = \frac{1}{2} \alpha_0 + \sum_{n=1}^{\infty} r^n (\alpha_n \cos n\theta + \beta_n \sin n\theta). \quad (3.7.72) \]

Using the inhomogeneous boundary condition

\[ f(\theta) = u(a, \theta) = \frac{1}{2} \alpha_0 + \sum_{n=1}^{\infty} a^n (\alpha_n \cos n\theta + \beta_n \sin n\theta), \quad (3.7.73) \]
we have the coefficients (Fourier series expansion of $f(\theta)$)
\[
\alpha_0 = \frac{1}{\pi} \int_0^{2\pi} f(\theta) d\theta, \tag{3.7.74}
\]
\[
\alpha_n = \frac{1}{\pi a^n} \int_0^{2\pi} f(\theta) \cos n\theta d\theta, \tag{3.7.75}
\]
\[
\beta_n = \frac{1}{\pi a^n} \int_0^{2\pi} f(\theta) \sin n\theta d\theta. \tag{3.7.76}
\]
The boundedness condition at zero is necessary only if $r = 0$ is part of the domain.

In the next example, we show how to overcome the Gibbs phenomenon resulting from discontinuities in the boundary conditions.

Example 14

Solve Laplace’s equation inside a rectangular domain $(0, a) \times (0, b)$ with nonzero Dirichlet boundary conditions on each side, i.e.
\[
\nabla^2 u = 0 \tag{3.7.77}
\]
\[
u(x, 0) = g_1(x), \tag{3.7.78}
u(a, y) = g_2(y), \tag{3.7.79}
u(x, b) = g_3(x), \tag{3.7.80}
u(0, y) = g_4(y), \tag{3.7.81}
\]
assuming that $g_1(a) \neq g_2(0)$ and so forth at other corners of the rectangle. This discontinuity causes spurious oscillations in the solution, i.e. we have Gibbs phenomenon.

The way to overcome the problem is to decompose $u$ to a sum of two functions
\[
u = v + w \tag{3.7.82}
\]
where $w$ is bilinear function and thus satisfies $\nabla^2 w = 0$, and $v$ is harmonic with boundary conditions vanishing at the corners, i.e.
\[
\nabla^2 v = 0 \tag{3.7.83}
\]
\[
v = g - w, \text{ on the boundary.} \tag{3.7.84}
\]

In order to get zero boundary conditions on the corners, we must have the function $w$ be of the form
\[
w(x, y) = g(0, 0) \frac{(a-x)(b-y)}{ab} + g(a, 0) \frac{x(b-y)}{ab} + g(a, b) \frac{xy}{ab} + g(0, b) \frac{(a-x)y}{ab}, \tag{3.7.85}
\]
and
\[
g(x, 0) = g_1(x) \tag{3.7.86}
g(a, y) = g_2(y) \tag{3.7.87}
g(x, b) = g_3(x) \tag{3.7.88}
g(0, y) = g_4(y). \tag{3.7.89}
\]

It is easy to show that this $w$ satisfies Laplace's equation and that $v$ vanishes at the corners and therefore the discontinuities disappear.
Problems

1. Solve the heat equation

\[ u_t = ku_{xx}, \quad 0 < x < L, \quad t > 0, \]

subject to the boundary conditions

\[ u(0, t) = u(L, t) = 0. \]

Solve the problem subject to the initial value:

a. \( u(x, 0) = 6 \sin \frac{9}{L}x. \)

b. \( u(x, 0) = 2 \cos \frac{3}{L}x. \)

2. Solve the heat equation

\[ u_t = ku_{xx}, \quad 0 < x < L, \quad t > 0, \]

subject to

\[ u_x(0, t) = 0, \quad t > 0 \]
\[ u_x(L, t) = 0, \quad t > 0 \]

a. \( u(x, 0) = \begin{cases} 
0 & x < \frac{L}{2} \\
1 & x > \frac{L}{2} 
\end{cases} \)

b. \( u(x, 0) = 6 + 4 \cos \frac{3}{L}x. \)

3. Solve the eigenvalue problem

\[ \phi'' = -\lambda \phi \]

subject to

\[ \phi(0) = \phi(2\pi) \]
\[ \phi'(0) = \phi'(2\pi) \]

4. Solve Laplace’s equation inside a wedge of radius \( a \) and angle \( \alpha \),

\[ \nabla^2 u = \frac{1}{r} \frac{\partial}{\partial r} \left( r \frac{\partial u}{\partial r} \right) + \frac{1}{r^2} \frac{\partial^2 u}{\partial \theta^2} = 0 \]

subject to

\[ u(a, \theta) = f(\theta), \]
\[ u(r, 0) = u_\theta(r, \alpha) = 0. \]
5. Solve Laplace’s equation inside a rectangle $0 \leq x \leq L$, $0 \leq y \leq H$ subject to
   a. $u_x(0, y) = u_x(L, y) = u(x, 0) = 0$, $u(x, H) = f(x)$.
   b. $u(0, y) = g(y)$, $u(L, y) = u_y(x, 0) = u(x, H) = 0$.
   c. $u(0, y) = u(L, y) = 0$, $u(x, 0) - u_y(x, 0) = 0$, $u(x, H) = f(x)$.

6. Solve Laplace’s equation outside a circular disk of radius $a$, subject to
   a. $u(a, \theta) = \ln 2 + 4 \cos 3\theta$.
   b. $u(a, \theta) = f(\theta)$.

7. Solve Laplace’s equation inside the quarter circle of radius 1, subject to
   a. $u_\theta(r, 0) = u(r, \pi/2) = 0$, $u(1, \theta) = f(\theta)$.
   b. $u_\theta(r, 0) = u_\theta(r, \pi/2) = 0$, $u_r(1, \theta) = g(\theta)$.
   c. $u(r, 0) = u(r, \pi/2) = 0$, $u_r(1, \theta) = 1$.

8. Solve Laplace’s equation inside a circular annulus $(a < r < b)$, subject to
   a. $u(a, \theta) = f(\theta)$, $u(b, \theta) = g(\theta)$.
   b. $u_r(a, \theta) = f(\theta)$, $u_r(b, \theta) = g(\theta)$.

9. Solve Laplace’s equation inside a semi-infinite strip $(0 < x < \infty, \ 0 < y < H)$ subject to
   $u_y(x, 0) = 0$, $u_y(x, H) = 0$, $u(0, y) = f(y)$.

10. Consider the heat equation
    
    $u_t = u_{xx} + q(x, t)$,
    $0 < x < L,$

    subject to the boundary conditions
    
    $u(0, t) = u(L, t) = 0$.

    Assume that $q(x, t)$ is a piecewise smooth function of $x$ for each positive $t$. Also assume that $u$ and $u_x$ are continuous functions of $x$ and $u_{xx}$ and $u_t$ are piecewise smooth. Thus
    
    $u(x, t) = \sum_{n=1}^{\infty} b_n(t) \sin \frac{n\pi}{L} x$.

    Write the ordinary differential equation satisfied by $b_n(t)$.

11. Solve the following inhomogeneous problem
    
    $\frac{\partial u}{\partial t} = k \frac{\partial^2 u}{\partial x^2} + e^{-t} + e^{-2t} \cos \frac{3\pi}{L} x$. 

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subject to
\[
\frac{\partial u}{\partial x}(0, t) = \frac{\partial u}{\partial x}(L, t) = 0, \\
u(x, 0) = f(x).
\]

Hint: Look for a solution as a Fourier cosine series. Assume \( k \neq \frac{2\pi^2}{L^2} \).

12. Solve the wave equation by the method of separation of variables
\[
\begin{align*}
\frac{\partial^2 u}{\partial t^2} - c^2 \frac{\partial^2 u}{\partial x^2} &= 0, & 0 < x < L, \\
u(0, t) &= 0, \\
u(L, t) &= 0, \\
u(x, 0) &= f(x), \\
u_t(x, 0) &= g(x).
\end{align*}
\]

13. Solve the heat equation
\[
\frac{\partial u}{\partial t} = 2 \frac{\partial^2 u}{\partial x^2},
\]
subject to the boundary conditions
\[
u(0, t) = u_x(L, t) = 0,
\]
and the initial condition
\[
u(x, 0) = \sin \frac{3\pi}{2L} x.
\]

14. Solve the heat equation
\[
\frac{\partial u}{\partial t} = k \left( \frac{1}{r} \frac{\partial}{\partial r} \left( r \frac{\partial u}{\partial r} \right) + \frac{1}{r^2} \frac{\partial^2 u}{\partial \theta^2} \right)
\]
inside a disk of radius \( a \) subject to the boundary condition
\[
\frac{\partial u}{\partial r}(a, \theta, t) = 0,
\]
and the initial condition
\[
u(r, \theta, 0) = f(r, \theta)
\]
where \( f(r, \theta) \) is a given function.

15. Determine which of the following equations are separable:
16. (a) Solve the one dimensional heat equation in a bar

\[ u_t = ku_{xx} \quad 0 < x < L \]

which is insulated at either end, given the initial temperature distribution

\[ u(x,0) = f(x) \]

(b) What is the equilibrium temperature of the bar? and explain physically why your answer makes sense.

17. Solve the 1-D heat equation

\[ u_t = ku_{xx} \quad 0 < x < L \]

subject to the nonhomogeneous boundary conditions

\[ u(0) = 1 \quad u_x(L) = 1 \]

with an initial temperature distribution \( u(x,0) = 0 \). (Hint: First solve for the equilibrium temperature distribution \( v(x) \) which satisfies the steady state heat equation with the prescribed boundary conditions. Once \( v \) is found, write \( u(x,t) = v(x) + w(x,t) \) where \( w(x,t) \) is the transient response. Substitute this \( u \) back into the PDE to produce a new PDE for \( w \) which now has homogeneous boundary conditions.

18. Solve Laplace’s equation,

\[ \nabla^2 u = 0, \quad 0 \leq x \leq \pi, \quad 0 \leq y \leq \pi \]

subject to the boundary conditions

\[ u(x,0) = \sin x + 2 \sin 2x \]

\[ u(\pi,y) = 0 \]

\[ u(x,\pi) = 0 \]

\[ u(0,y) = 0 \]

19. Repeat the above problem with

\[ u(x,0) = -\pi^2 x^2 + 2\pi x^3 - x^4 \]
**SUMMARY**

Fourier Series

\[ f(x) \sim \frac{a_0}{2} + \sum_{n=1}^{\infty} \left( a_n \cos \frac{n\pi}{L} x + b_n \sin \frac{n\pi}{L} x \right) \]

\[ a_k = \frac{1}{L} \int_{-L}^{L} f(x) \cos \frac{k\pi}{L} x \, dx \quad \text{for } k = 0, 1, 2, \ldots \]

\[ b_k = \frac{1}{L} \int_{-L}^{L} f(x) \sin \frac{k\pi}{L} x \, dx \quad \text{for } k = 1, 2, \ldots \]

Solution of Euler’s equation

\[ r (r R')' - \lambda R = 0 \]

For \( \lambda_0 = 0 \) the solution is \( R_0 = C_1 \ln r + C_2 \)

For \( \lambda_n = n^2 \) the solution is \( R_n = D_1 r^n + D_2 r^{-n}, \quad n = 1, 2, \ldots \)
4  PDEs in Higher Dimensions

4.1  Introduction

In the previous chapters we discussed homogeneous time dependent one dimensional PDEs with homogeneous boundary conditions. Also Laplace's equation in two variables was solved in cartesian and polar coordinate systems. The eigenpairs of the Laplacian will be used here to solve time dependent PDEs with two or three spatial variables. We will also discuss the solution of Laplace's equation in cylindrical and spherical coordinate systems, thus allowing us to solve the heat and wave equations in those coordinate systems.

In the top part of the following table we list the various equations solved to this point. In the bottom part we list the equations to be solved in this chapter.

<table>
<thead>
<tr>
<th>Equation</th>
<th>Type</th>
<th>Comments</th>
</tr>
</thead>
<tbody>
<tr>
<td>( u_t = ku_{xx} )</td>
<td>heat</td>
<td>1D constant coefficients</td>
</tr>
<tr>
<td>( c(x)\rho(x)u_t = (K(x)u_x)_x )</td>
<td>heat</td>
<td>1D</td>
</tr>
<tr>
<td>( u_{tt} - c^2 u_{xx} = 0 )</td>
<td>wave</td>
<td>1D constant coefficients</td>
</tr>
<tr>
<td>( \rho(x)u_{tt} - T_0(x)u_{xx} = 0 )</td>
<td>wave</td>
<td>1D</td>
</tr>
<tr>
<td>( u_{xx} + u_{yy} = 0 )</td>
<td>Laplace</td>
<td>2D constant coefficients</td>
</tr>
<tr>
<td>( u_t = k(u_{xx} + u_{yy}) )</td>
<td>heat</td>
<td>2D constant coefficients</td>
</tr>
<tr>
<td>( u_t = k(u_{xx} + u_{yy} + u_{zz}) )</td>
<td>heat</td>
<td>3D constant coefficients</td>
</tr>
<tr>
<td>( u_{tt} - c^2(u_{xx} + u_{yy}) = 0 )</td>
<td>wave</td>
<td>2D constant coefficients</td>
</tr>
<tr>
<td>( u_{tt} - c^2(u_{xx} + u_{yy} + u_{zz}) = 0 )</td>
<td>wave</td>
<td>3D constant coefficients</td>
</tr>
<tr>
<td>( u_{xx} + u_{yy} + u_{zz} = 0 )</td>
<td>Laplace</td>
<td>3D Cartesian</td>
</tr>
<tr>
<td>( \frac{1}{r}(ru_r)<em>r + \frac{1}{r^2}u</em>{\theta\theta} + u_{zz} = 0 )</td>
<td>Laplace</td>
<td>3D Cylindrical</td>
</tr>
<tr>
<td>( u_{rr} + \frac{2}{r}u_r + \frac{1}{r^2}u_{\theta\theta} + \frac{\cos \theta}{r^2}u_{\theta} + \frac{1}{r^2 \sin^2 \theta}u_{\phi\phi} = 0 )</td>
<td>Laplace</td>
<td>3D Spherical</td>
</tr>
</tbody>
</table>
4.2 Heat Flow in a Rectangular Domain

In this section we solve the heat equation in two spatial variables inside a rectangle \( L \) by \( H \). The equation is

\[
\frac{\partial u}{\partial t} = k (u_{xx} + u_{yy}), \quad 0 < x < L, \quad 0 < y < H, \tag{4.2.1}
\]

\[
u(0, y, t) = 0, \tag{4.2.2}
\]

\[
u(L, y, t) = 0, \tag{4.2.3}
\]

\[
u(x, 0, t) = 0, \tag{4.2.4}
\]

\[
u(x, H, t) = 0, \tag{4.2.5}
\]

\[
u(x, y, 0) = f(x, y). \tag{4.2.6}
\]

Notice that the term in parentheses in (4.2.1) is \( \nabla^2 u \). Note also that we took Dirichlet boundary conditions (i.e. specified temperature on the boundary). We can write this condition as

\[
u(x, y, t) = 0. \quad \text{on the boundary} \tag{4.2.7}
\]

Other possible boundary conditions are left to the reader.

The method of separation of variables will proceed as follows:

1. Let

\[
u(x, y, t) = T(t)\phi(x, y) \tag{4.2.8}
\]

2. Substitute in (4.2.1) and separate the variables

\[
\frac{\hat{T}}{kT} = \frac{\nabla^2 \phi}{\phi} = -\lambda
\]

3. Write the ODEs

\[
\hat{T}(t) + k\lambda T(t) = 0 \tag{4.2.9}
\]

\[
\nabla^2 \phi + \lambda \phi = 0 \tag{4.2.10}
\]

4. Use the homogeneous boundary condition (4.2.7) to get the boundary condition associated with (4.2.10)

\[
\phi(x, y) = 0. \quad \text{on the boundary} \tag{4.2.11}
\]

The only question left is how to get the solution of (4.2.10) - (4.2.11). This can be done in a similar fashion to solving Laplace’s equation.

Let

\[
\phi(x, y) = X(x)Y(y), \tag{4.2.12}
\]

then (4.2.10) - (4.2.11) yield 2 ODEs

\[
X'' + \mu X = 0, \tag{4.2.13}
\]

\[
X(0) = X(L) = 0, \tag{4.2.14}
\]

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\[ Y'' + (\lambda - \mu)Y = 0, \quad (4.2.15) \]
\[ Y(0) = Y(H) = 0. \quad (4.2.16) \]

The boundary conditions (4.2.14) and (4.2.16) result from (4.2.2) - (4.2.5). Equation (4.2.13) has a solution

\[ X_n = \sin \frac{n\pi}{L} x, \quad n = 1, 2, \ldots \quad (4.2.17) \]
\[ \mu_n = \left( \frac{n\pi}{L} \right)^2, \quad n = 1, 2, \ldots \quad (4.2.18) \]
as we have seen in Chapter 2. For each \( n \), equation (4.2.15) is solved the same way

\[ Y_{mn} = \sin \frac{m\pi}{H} y, \quad m = 1, 2, \ldots, n = 1, 2, \ldots \quad (4.2.19) \]
\[ \lambda_{mn} - \mu_n = \left( \frac{n\pi}{H} \right)^2, \quad m = 1, 2, \ldots, n = 1, 2, \ldots \quad (4.2.20) \]

Therefore by (4.2.12) and (4.2.17)-(4.2.20),

\[ \phi_{mn}(x, y) = \sin \frac{n\pi}{L} x \sin \frac{m\pi}{H} y, \quad (4.2.21) \]
\[ \lambda_{mn} = \left( \frac{n\pi}{L} \right)^2 + \left( \frac{m\pi}{H} \right)^2, \quad (4.2.22) \]
\[ n = 1, 2, \ldots, m = 1, 2, \ldots \]

Using (4.2.8) and the principle of superposition, we can write the solution of (4.2.1) as

\[ u(x, y, t) = \sum_{n=1}^{\infty} \sum_{m=1}^{\infty} A_{mn} e^{-k\lambda_{mn} t} \sin \frac{n\pi}{L} x \sin \frac{m\pi}{H} y, \quad (4.2.23) \]

where \( \lambda_{mn} \) is given by (4.2.22).

To find the coefficients \( A_{mn} \), we use the initial condition (4.2.6), that is for \( t = 0 \) in (4.2.23) we get:

\[ f(x, y) = \sum_{n=1}^{\infty} \sum_{m=1}^{\infty} A_{mn} \sin \frac{n\pi}{L} x \sin \frac{m\pi}{H} y, \quad (4.2.24) \]

\( A_{mn} \) are the generalized Fourier coefficients (double Fourier series in this case). We can compute \( A_{mn} \) by

\[ A_{mn} = \frac{\int_0^L \int_0^H f(x, y) \sin \frac{n\pi}{L} x \sin \frac{m\pi}{H} y dy dx}{\int_0^L \int_0^H \sin^2 \frac{n\pi}{L} x \sin^2 \frac{m\pi}{H} y dy dx}. \quad (4.2.25) \]

(See next section.)

Remarks:

i. Equation (4.2.10) is called Helmholtz equation.

ii. A more general form of the equation is

\[ \nabla \cdot (p(x, y)\nabla \phi(x, y)) + q(x, y)\phi(x, y) + \lambda \sigma(x, y)\phi(x, y) = 0 \quad (4.2.26) \]

iii. A more general boundary condition is

\[ \beta_1(x, y)\phi(x, y) + \beta_2(x, y)\nabla \phi \cdot \vec{n} = 0 \quad \text{on the boundary} \quad (4.2.27) \]

where \( \vec{n} \) is a unit normal vector pointing outward. The special case \( \beta_2 \equiv 0 \) yields (4.2.11).
Problems

1. Solve the heat equation

\[ u_t(x, y, t) = k \left( u_{xx}(x, y, t) + u_{yy}(x, y, t) \right), \]

on the rectangle \( 0 < x < L, 0 < y < H \) subject to the initial condition

\[ u(x, y, 0) = f(x, y), \]

and the boundary conditions

a. \[ u(0, y, t) = u_x(L, y, t) = 0, \]
\[ u(x, 0, t) = u(x, H, t) = 0. \]

b. \[ u_x(0, y, t) = u(L, y, t) = 0, \]
\[ u_y(x, 0, t) = u_y(x, H, t) = 0. \]

c. \[ u(0, y, t) = u(L, y, t) = 0, \]
\[ u(x, 0, t) = u_y(x, H, t) = 0. \]

2. Solve the heat equation on a rectangular box

\[ 0 < x < L, 0 < y < H, 0 < z < W; \]

\[ u_t(x, y, z, t) = k(u_{xx} + u_{yy} + u_{zz}), \]

subject to the boundary conditions

\[ u(0, y, z, t) = u(L, y, z, t) = 0, \]
\[ u(x, 0, z, t) = u(x, H, z, t) = 0, \]
\[ u(x, y, 0, t) = u(x, y, W, t) = 0, \]

and the initial condition

\[ u(x, y, z, 0) = f(x, y, z). \]
4.3 Vibrations of a rectangular Membrane

The method of separation of variables in this case will lead to the same Helmholtz equation. The only difference is in the T equation. the problem to solve is as follows:

\[ u_{tt} = c^2(u_{xx} + u_{yy}), \quad 0 < x < L, 0 < y < H, \]  
(4.3.1)

\[ u(0, y, t) = 0, \]  
(4.3.2)

\[ u(L, y, t) = 0, \]  
(4.3.3)

\[ u(x, 0, t) = 0, \]  
(4.3.4)

\[ u_y(x, H, t) = 0, \]  
(4.3.5)

\[ u(x, y, 0) = f(x, y), \]  
(4.3.6)

\[ u_t(x, y, 0) = g(x, y). \]  
(4.3.7)

Clearly there are two initial conditions, (4.3.6)-(4.3.7), since the PDE is second order in time. We have decided to use a Neumann boundary condition at the top \( y = H \), to show how the solution of Helmholtz equation is affected.

The steps to follow are: (the reader is advised to compare these equations to (4.2.8)-(4.2.25))

\[ u(x, y, t) = T(t)\phi(x, y), \]  
(4.3.8)

\[ \frac{\ddot{T}}{c^2 T} = \frac{\nabla^2 \phi}{\phi} = -\lambda \]  
(4.3.9)

\[ \ddot{T} + \lambda c^2 T = 0, \]  
(4.3.10)

\[ \nabla^2 \phi + \lambda \phi = 0, \]  
(4.3.11)

\[ \beta_1 \phi(x, y) + \beta_2 \phi_y(x, y) = 0, \]  
(4.3.11)

where either \( \beta_1 \) or \( \beta_2 \) is zero depending on which side of the rectangle we are on.

\[ \phi(x, y) = X(x)Y(y), \]  
(4.3.12)

\[ X'' + \mu X = 0, \]  
(4.3.13)

\[ X(0) = X(L) = 0, \]  
(4.3.14)

\[ Y'' + (\lambda - \mu)Y = 0, \]  
(4.3.15)

\[ Y(0) = Y'(H) = 0, \]  
(4.3.16)

\[ X_n = \sin \frac{n\pi}{L} x, \quad n = 1, 2, \ldots \]  
(4.3.17)

\[ \mu_n = \left( \frac{n\pi}{L} \right)^2, \quad n = 1, 2, \ldots \]  
(4.3.18)

\[ Y_{mn} = \sin \frac{(m - \frac{1}{2})\pi}{H} y, \quad m = 1, 2, \ldots \quad n = 1, 2, \ldots \]  
(4.3.19)
\[ \lambda_{mn} = \left( \frac{(m - \frac{1}{2})\pi}{H} \right)^2 + \left( \frac{n\pi}{L} \right)^2, \quad m = 1, 2, \ldots \quad n = 1, 2, \ldots \] (4.3.20)

Note the similarity of (4.3.1)-(4.3.20) to the corresponding equations of section 4.2.

The solution
\[ u(x, y, t) = \sum_{m=1}^{\infty} \sum_{n=1}^{\infty} \left( A_{mn} \cos \sqrt{\lambda_{mn}}ct + B_{mn} \sin \sqrt{\lambda_{mn}}ct \right) \sin \frac{n\pi}{L}x \sin \frac{(m - \frac{1}{2})\pi}{H}y. \] (4.3.21)

Since the \( T \) equation is of second order, we end up with two sets of parameters \( A_{mn} \) and \( B_{mn} \). These can be found by using the two initial conditions (4.3.6)-(4.3.7).

\[ f(x, y) = \sum_{n=1}^{\infty} \sum_{m=1}^{\infty} A_{mn} \sin \frac{n\pi}{L}x \sin \frac{(m - \frac{1}{2})\pi}{H}y, \] (4.3.22)

\[ g(x, y) = \sum_{n=1}^{\infty} \sum_{m=1}^{\infty} c \sqrt{\lambda_{mn}} B_{mn} \sin \frac{n\pi}{L}x \sin \frac{(m - \frac{1}{2})\pi}{H}y. \] (4.3.23)

To get (4.3.23) we need to evaluate \( u_t \) from (4.3.21) and then substitute \( t = 0 \). The coefficients are then

\[ A_{mn} = \frac{\int_{0}^{L} \int_{0}^{H} f(x, y) \sin \frac{m\pi}{L}x \sin \frac{(m - \frac{1}{2})\pi}{H}y \, dy \, dx}{\int_{0}^{L} \int_{0}^{H} \sin^2 \frac{m\pi}{L}x \sin^2 \frac{(m - \frac{1}{2})\pi}{H}y \, dy \, dx}, \] (4.3.24)

\[ c \sqrt{\lambda_{mn}} B_{mn} = \frac{\int_{0}^{L} \int_{0}^{H} g(x, y) \sin \frac{n\pi}{L}x \sin \frac{(m - \frac{1}{2})\pi}{H}y \, dy \, dx}{\int_{0}^{L} \int_{0}^{H} \sin^2 \frac{n\pi}{L}x \sin^2 \frac{(m - \frac{1}{2})\pi}{H}y \, dy \, dx}. \] (4.3.25)
Problems

1. Solve the wave equation

\[ u_{tt}(x, y, t) = c^2 \left( u_{xx}(x, y, t) + u_{yy}(x, y, t) \right), \]

on the rectangle \(0 < x < L, 0 < y < H\) subject to the initial conditions

\[ u(x, y, 0) = f(x, y), \]
\[ u_t(x, y, 0) = g(x, y), \]

and the boundary conditions

a. \[ u(0, y, t) = u_x(L, y, t) = 0, \]
\[ u(x, 0, t) = u(x, H, t) = 0. \]

b. \[ u(0, y, t) = u(L, y, t) = 0, \]
\[ u(x, 0, t) = u(x, H, t) = 0. \]

c. \[ u_x(0, y, t) = u(L, y, t) = 0, \]
\[ u_y(x, 0, t) = u_y(x, H, t) = 0. \]

2. Solve the wave equation on a rectangular box

\[ 0 < x < L, 0 < y < H, 0 < z < W, \]
\[ u_{tt}(x, y, z, t) = c^2 (u_{xx} + u_{yy} + u_{zz}), \]

subject to the boundary conditions

\[ u(0, y, z, t) = u(L, y, z, t) = 0, \]
\[ u(x, 0, z, t) = u(x, H, z, t) = 0, \]
\[ u(x, y, 0, t) = u(x, y, W, t) = 0, \]

and the initial conditions

\[ u(x, y, z, 0) = f(x, y, z), \]
\[ u_t(x, y, z, 0) = g(x, y, z). \]

3. Solve the wave equation on an isosceles right-angle triangle with side of length \(a\)

\[ u_{tt}(x, y, t) = c^2 (u_{xx} + u_{yy}), \]
subject to the boundary conditions

\[ u(x, 0, t) = u(0, y, t) = 0, \]
\[ u(x, y, t) = 0, \quad \text{on the line} \quad x + y = a \]

and the initial conditions

\[ u(x, y, 0) = f(x, y), \]
\[ u_t(x, y, 0) = g(x, y). \]
4.4 Helmholtz Equation

As we have seen in this chapter, the method of separation of variables in two independent variables leads to Helmholtz equation,

$$\nabla^2 \phi + \lambda \phi = 0$$

subject to the boundary conditions

$$\beta_1 \phi(x, y) + \beta_2 \phi_x(x, y) + \beta_3 \phi_y(x, y) = 0.$$

Here we state a result generalizing Sturm-Liouville’s from Chapter 6 of Neta.

Theorem:
1. All the eigenvalues are real.
2. There exists an infinite number of eigenvalues. There is a smallest one but no largest.
3. Corresponding to each eigenvalue, there may be many eigenfunctions.
4. The eigenfunctions $$\phi_i(x, y)$$ form a complete set, i.e. any function $$f(x, y)$$ can be represented by

$$\sum_i a_i \phi_i(x, y)$$ (4.4.1)

where the coefficients $$a_i$$ are given by,

$$a_i = \frac{\int \int \phi_i f(x, y) dx dy}{\int \int \phi_i^2 dx dy}$$ (4.4.2)

5. Eigenfunctions belonging to different eigenvalues are orthogonal.
6. An eigenvalue $$\lambda$$ can be related to the eigenfunction $$\phi(x, y)$$ by Rayleigh quotient:

$$\lambda = \frac{\int \int (\nabla \phi)^2 dx dy - \oint \phi \nabla \phi \cdot \hat{n} ds}{\int \int \phi^2 dx dy}$$ (4.4.3)

where $$\oint$$ symbolizes integration on the boundary. For example, the following Helmholtz problem (see 4.2.10-11)

$$\nabla^2 \phi + \lambda \phi = 0, \quad 0 \leq x \leq L, 0 \leq y \leq H,$$

$$\phi = 0, \quad \text{on the boundary},$$ (4.4.4) (4.4.5)

was solved and we found

$$\lambda_{mn} = \left( \frac{n\pi}{L} \right)^2 + \left( \frac{m\pi}{H} \right)^2, \quad n = 1, 2, \ldots, \quad m = 1, 2, \ldots$$ (4.4.6)

$$\phi_{mn}(x, y) = \sin \frac{n\pi}{L} x \sin \frac{m\pi}{H} y, \quad n = 1, 2, \ldots, \quad m = 1, 2, \ldots$$ (4.4.7)

Clearly all the eigenvalues are real. The smallest one is $$\lambda_{11} = \left( \frac{\pi}{L} \right)^2 + \left( \frac{\pi}{H} \right)^2$$, $$\lambda_{mn} \to \infty$$ as $$n$$ and $$m \to \infty$$. There may be multiple eigenfunctions in some cases. For example, if
\( L = 2H \) then \( \lambda_{11} = \lambda_{22} \) but the eigenfunctions \( \phi_{11} \) and \( \phi_{22} \) are different. The coefficients of expansion are

\[
a_{mn} = \frac{\int_0^L \int_0^H f(x,y) \phi_{mn} \, dx \, dy}{\int_0^L \int_0^H \phi_{mn}^2 \, dx \, dy}
\]  

(4.4.8)

as given by (4.2.25).
Problems

1. Solve

\[ \nabla^2 \phi + \lambda \phi = 0 \quad [0, 1] \times [0, 1/4] \]

subject to

\[
\begin{align*}
\phi(0, y) &= 0 \\
\phi_x(1, y) &= 0 \\
\phi(x, 0) &= 0 \\
\phi_y(x, 1/4) &= 0.
\end{align*}
\]

Show that the results of the theorem are true.

2. Solve Helmholtz equation on an isosceles right-angle triangle with side of length \(a\)

\[ u_{xx} + u_{yy} + \lambda u = 0, \]

subject to the boundary conditions

\[
\begin{align*}
u(x, 0, t) &= u(0, y, t) = 0, \\
u(x, y, t) &= 0, \quad \text{on the line} \quad x + y = a.
\end{align*}
\]
4.5 Vibrating Circular Membrane

In this section, we discuss the solution of the wave equation inside a circle. As we have seen in sections 4.2 and 4.3, there is a similarity between the solution of the heat and wave equations. Thus we will leave the solution of the heat equation to the exercises.

The problem is:

\[ u_{tt}(r, \theta, t) = c^2 \nabla^2 u, \quad 0 \leq r \leq a, 0 \leq \theta \leq 2\pi, t > 0 \]  \hspace{1cm} (4.5.1)

subject to the boundary condition

\[ u(a, \theta, t) = 0, \quad \text{(clamped membrane)} \]  \hspace{1cm} (4.5.2)

and the initial conditions

\[ u(r, \theta, 0) = \alpha(r, \theta), \]  \hspace{1cm} (4.5.3)

\[ u_t(r, \theta, 0) = \beta(r, \theta). \]  \hspace{1cm} (4.5.4)

The method of separation of variables leads to the same set of differential equations

\[ \ddot{T}(t) + \lambda c^2 T = 0, \]  \hspace{1cm} (4.5.5)

\[ \nabla^2 \phi + \lambda \phi = 0, \]  \hspace{1cm} (4.5.6)

\[ \phi(a, \theta) = 0, \]  \hspace{1cm} (4.5.7)

Note that in polar coordinates

\[ \nabla^2 \phi = \frac{1}{r} \frac{\partial}{\partial r} \left( r \frac{\partial \phi}{\partial r} \right) + \frac{1}{r^2} \frac{\partial^2 \phi}{\partial \theta^2} \]  \hspace{1cm} (4.5.8)

Separating the variables in the Helmholtz equation (4.5.6) we have

\[ \phi(r, \theta) = R(r) \Theta(\theta), \]  \hspace{1cm} (4.5.9)

\[ \Theta'' + \mu \Theta = 0 \]  \hspace{1cm} (4.5.10)

\[ \frac{d}{dr} \left( r \frac{dR}{dr} \right) + \left( \lambda r - \frac{\mu}{r} \right) R = 0. \]  \hspace{1cm} (4.5.11)

The boundary equation (4.5.7) yields

\[ R(a) = 0. \]  \hspace{1cm} (4.5.12)

What are the other boundary conditions? Check the solution of Laplace’s equation inside a circle!

\[ \Theta(0) = \Theta(2\pi), \quad \text{(periodicity)} \]  \hspace{1cm} (4.5.13)

\[ \Theta'(0) = \Theta'(2\pi), \]  \hspace{1cm} (4.5.14)
\[ |R(0)| < \infty \quad \text{(boundedness)} \quad (4.5.15) \]

The equation for \( \Theta(\theta) \) can be solved (see Chapter 2)

\[
\mu_m = m^2 \quad m = 0, 1, 2, \ldots \quad (4.5.16)
\]

\[
\Theta_m = \begin{cases} 
\sin m\theta & m = 0, 1, 2, \ldots \\
\cos m\theta & m = 0, 1, 2, \ldots 
\end{cases} \quad (4.5.17)
\]

In the rest of this section, we discuss the solution of (4.5.11) subject to (4.5.12), (4.5.15). After substituting the eigenvalues \( \mu_m \) from (4.5.16), we have

\[
\frac{d}{dr} \left( r \frac{dR_m}{dr} \right) + \left( \lambda r - \frac{m^2}{r} \right) R_m = 0 \quad (4.5.18)
\]

\[ |R_m(0)| < \infty \quad (4.5.19) \]

\[ R_m(a) = 0. \quad (4.5.20) \]

Using Rayleigh quotient for this singular Sturm-Liouville problem, we can show that \( \lambda > 0 \), thus we can make the transformation

\[ \rho = \sqrt{\lambda r} \quad (4.5.21) \]

which will yield Bessel’s equation

\[
\rho^2 \frac{d^2 R(\rho)}{d\rho^2} + \rho \frac{dR(\rho)}{d\rho} + \left( \rho^2 - m^2 \right) R(\rho) = 0 \quad (4.5.22)
\]

Consulting a textbook on the solution of Ordinary Differential Equations, we find:

\[ R_m(\rho) = C_{1m} J_m(\rho) + C_{2m} Y_m(\rho) \quad (4.5.23) \]

where \( J_m, Y_m \) are Bessel functions of the first, second kind of order \( m \) respectively. Since we are interested in a solution satisfying (4.5.15), we should note that near \( \rho = 0 \)

\[
J_m(\rho) \sim \begin{cases} 
1 & m = 0 \\
\frac{1}{2^m m!} \rho^m & m > 0
\end{cases} \quad (4.5.24)
\]

\[
Y_m(\rho) \sim \begin{cases} 
\frac{2 \ln \rho}{\pi} & m = 0 \\
-\frac{2^m m!}{\pi} \frac{1}{\rho^m} & m > 0.
\end{cases} \quad (4.5.25)
\]

Thus \( C_{2m} = 0 \) is necessary to achieve boundedness. Thus

\[ R_m(\rho) = C_{1m} J_m(\sqrt{\lambda r}) \quad (4.5.26) \]

In figure 21 we have plotted the Bessel functions \( J_0 \) through \( J_5 \). Note that all \( J_n \) start at 0 except \( J_0 \) and all the functions cross the axis infinitely many times. In figure 22 we have plotted the Bessel functions (also called Neumann functions) \( Y_0 \) through \( Y_5 \). Note that the vertical axis is through \( x = 3 \) and so it is not so clear that \( Y_n \) tend to \(-\infty\) as \( x \to 0 \).
Figure 21: Bessel functions $J_n, n = 0, \ldots, 5$

To satisfy the boundary condition (4.5.20) we get an equation for the eigenvalues $\lambda$

$$J_m(\sqrt{\lambda}a) = 0. \quad (4.5.27)$$

There are infinitely many solutions of (4.5.27) for any $m$. We denote these solutions by

$$\xi_{mn} = \sqrt{\lambda_{mn}a} \quad m = 0, 1, 2, \ldots \quad n = 1, 2, \ldots \quad (4.5.28)$$

Thus

$$\lambda_{mn} = \left(\frac{\xi_{mn}}{a}\right)^2, \quad (4.5.29)$$

$$R_{mn}(r) = J_m\left(\frac{\xi_{mn}}{a}r\right). \quad (4.5.30)$$

We leave it as an exercise to show that the general solution to (4.5.1) - (4.5.2) is given by

$$u(r, \theta, t) = \sum_{m=0}^{\infty} \sum_{n=1}^{\infty} J_m \left(\frac{\xi_{mn}}{a}r\right) \left\{ a_{mn} \cos m\theta + b_{mn} \sin m\theta \right\} \left\{ A_{mn} \cos c\frac{\xi_{mn}}{a}t + B_{mn} \sin c\frac{\xi_{mn}}{a}t \right\} \quad (4.5.31)$$

We will find the coefficients by using the initial conditions (4.5.3)-(4.5.4)

$$\alpha(r, \theta) = \sum_{m=0}^{\infty} \sum_{n=1}^{\infty} J_m \left(\frac{\xi_{mn}}{a}r\right) A_{mn} \left\{ a_{mn} \cos m\theta + b_{mn} \sin m\theta \right\} \quad (4.5.32)$$
\[ \beta(r, \theta) = \sum_{m=0}^{\infty} \sum_{n=1}^{\infty} J_m \left( \frac{\xi_{mn} a}{r} \right) c \frac{\xi_{mn}}{a} B_{mn} \left\{ a_{mn} \cos m\theta + b_{mn} \sin m\theta \right\}. \] (4.533)

\[ A_{mn} a_{mn} = \frac{\int_0^{2\pi} \int_0^a \alpha(r, \theta) J_m \left( \frac{\xi_{mn} a}{r} \right) \cos m\theta r dr d\theta}{\int_0^{2\pi} \int_0^a J_m^2 \left( \frac{\xi_{mn} a}{r} \right) \cos^2 m\theta r dr d\theta}, \] (4.534)

\[ c \frac{\xi_{mn}}{a} B_{mn} a_{mn} = \frac{\int_0^{2\pi} \int_0^a \beta(r, \theta) J_m \left( \frac{\xi_{mn} a}{r} \right) \cos m\theta r dr d\theta}{\int_0^{2\pi} \int_0^a J_m^2 \left( \frac{\xi_{mn} a}{r} \right) \cos^2 m\theta r dr d\theta}. \] (4.535)

Replacing \( \cos m\theta \) by \( \sin m\theta \) we get \( A_{mn} b_{mn} \) and \( c \frac{\xi_{mn}}{a} B_{mn} b_{mn} \).

Remarks

1. Note the weight \( r \) in the integration. It comes from having \( \lambda \) multiplied by \( r \) in (4.5.18).

2. We are computing the four required combinations \( A_{mn} a_{mn}, A_{mn} b_{mn}, B_{mn} a_{mn} \), and \( B_{mn} b_{mn} \). We do not need to find \( A_{mn} \) or \( B_{mn} \) and so on.

Example:

Solve the circularly symmetric case

\[ u_{tt}(r, t) = \frac{c^2}{r} \frac{\partial}{\partial r} \left( r \frac{\partial u}{\partial r} \right), \] (4.5.36)

\[ u(a, t) = 0, \] (4.5.37)
The reader can easily show that the separation of variables give

\[ \ddot{T} + \lambda c^2 T = 0, \quad (4.5.40) \]

\[ \frac{d}{dr} \left( r \frac{dR}{dr} \right) + \lambda r R = 0, \quad (4.5.41) \]

\[ R(a) = 0, \quad (4.5.42) \]

\[ |R(0)| < \infty. \quad (4.5.43) \]

Since there is no dependence on \( \theta \), the \( r \) equation will have no \( \mu \), or which is the same \( m = 0 \). Thus

\[ R_0(r) = J_0(\sqrt{\lambda_n}r) \quad (4.5.44) \]

where the eigenvalues \( \lambda_n \) are computed from

\[ J_0(\sqrt{\lambda_n}a) = 0. \quad (4.5.45) \]

The general solution is

\[ u(r, t) = \sum_{n=1}^{\infty} a_n J_0(\sqrt{\lambda_n}r) \cos c \sqrt{\lambda_n} t + b_n J_0(\sqrt{\lambda_n}r) \sin c \sqrt{\lambda_n} t. \quad (4.5.46) \]

The coefficients \( a_n, b_n \) are given by

\[ a_n = \frac{\int_0^a J_0(\sqrt{\lambda_n}r) \alpha(r) r \, dr}{\int_0^a J_0^2(\sqrt{\lambda_n}r) r \, dr}, \quad (4.5.47) \]

\[ b_n = \frac{\int_0^a J_0(\sqrt{\lambda_n}r) \beta(r) r \, dr}{c \sqrt{\lambda_n} \int_0^a J_0^2(\sqrt{\lambda_n}r) r \, dr}. \quad (4.5.48) \]
Problems

1. Solve the heat equation
   \[ u_t(r, \theta, t) = k \nabla^2 u, \quad 0 \leq r < a, \ 0 < \theta < 2\pi, \ t > 0 \]
   subject to the boundary condition
   \[ u(a, \theta, t) = 0 \]
   (zero temperature on the boundary)
   and the initial condition
   \[ u(r, \theta, 0) = \alpha(r, \theta). \]

2. Solve the wave equation
   \[ u_{tt}(r, t) = c^2 (u_{rr} + \frac{1}{r} u_r), \]
   \[ u_r(a, t) = 0, \]
   \[ u(r, 0) = \alpha(r), \]
   \[ u_t(r, 0) = 0. \]
   Show the details.

3. Consult numerical analysis textbook to obtain the smallest eigenvalue of the above problem.

4. Solve the wave equation
   \[ u_{tt}(r, \theta, t) - c^2 \nabla^2 u = 0, \quad 0 \leq r < a, \ 0 < \theta < 2\pi, \ t > 0 \]
   subject to the boundary condition
   \[ u_r(a, \theta, t) = 0 \]
   and the initial conditions
   \[ u(r, \theta, 0) = 0, \]
   \[ u_t(r, \theta, 0) = \beta(r) \cos \theta. \]

5. Solve the wave equation
   \[ u_{tt}(r, \theta, t) - c^2 \nabla^2 u = 0, \quad 0 \leq r < a, \ 0 < \theta < \pi/2, \ t > 0 \]
   subject to the boundary conditions
   \[ u(a, \theta, t) = u(r, 0, t) = u(r, \pi/2, t) = 0 \]
   (zero displacement on the boundary)
   and the initial conditions
   \[ u(r, \theta, 0) = \alpha(r, \theta), \]
   \[ u_t(r, \theta, 0) = 0. \]
4.6 Laplace’s Equation in a Circular Cylinder

Laplace’s equation in cylindrical coordinates is given by:

\[ \frac{1}{r} (ru_r)_r + \frac{1}{r^2} u_{\theta\theta} + u_{zz} = 0, \quad 0 \leq r < a, 0 < z < H, 0 < \theta < 2\pi. \]  

(4.6.1)

The boundary conditions we discuss here are:

\[ u(r, \theta, 0) = \alpha(r, \theta), \quad \text{on bottom of cylinder}, \]  

(4.6.2)

\[ u(r, \theta, H) = \beta(r, \theta), \quad \text{on top of cylinder}, \]  

(4.6.3)

\[ u(a, \theta, z) = \gamma(\theta, z), \quad \text{on lateral surface of cylinder}. \]  

(4.6.4)

Similar methods can be employed if the boundary conditions are not of Dirichlet type (see exercises).

As we have done previously with Laplace’s equation, we use the principle of superposition to get two homogeneous boundary conditions. Thus we have the following three problems to solve, each differ from the others in the boundary conditions:

Problem 1:

\[ \frac{1}{r} (ru_r)_r + \frac{1}{r^2} u_{\theta\theta} + u_{zz} = 0, \]  

(4.6.5)

\[ u(r, \theta, 0) = 0, \]  

(4.6.6)

\[ u(r, \theta, H) = \beta(r, \theta), \]  

(4.6.7)

\[ u(a, \theta, z) = 0. \]  

(4.6.8)

Problem 2:

\[ \frac{1}{r} (ru_r)_r + \frac{1}{r^2} u_{\theta\theta} + u_{zz} = 0, \]  

(4.6.9)

\[ u(r, \theta, 0) = \alpha(r, \theta), \]  

(4.6.10)

\[ u(r, \theta, H) = 0, \]  

(4.6.11)

\[ u(a, \theta, z) = 0. \]  

(4.6.12)

Problem 3:

\[ \frac{1}{r} (ru_r)_r + \frac{1}{r^2} u_{\theta\theta} + u_{zz} = 0, \]  

(4.6.13)

\[ u(r, \theta, 0) = 0, \]  

(4.6.14)

\[ u(r, \theta, H) = 0, \]  

(4.6.15)

\[ u(a, \theta, z) = \gamma(\theta, z). \]  

(4.6.16)

Since the PDE is the same in all three problems, we get the same set of ODEs

\[ \Theta'' + \mu \Theta = 0, \]  

(4.6.17)
Recalling Laplace’s equation in polar coordinates, the boundary conditions associated with (4.6.17) are

\[
\begin{align*}
\Theta(0) &= \Theta(2\pi), \\
\Theta'(0) &= \Theta'(2\pi),
\end{align*}
\]

and one of the boundary conditions for (4.6.19) is

\[
|R(0)| < \infty.
\]

The other boundary conditions depend on which of the three we are solving. For problem 1, we have

\[
\begin{align*}
Z(0) &= 0, \\
R(a) &= 0.
\end{align*}
\]

Clearly, the equation for \( \Theta \) can be solved yielding

\[
\begin{align*}
\mu_m &= m^2, & m &= 0, 1, 2, \ldots \\
\Theta_m &= \begin{cases} 
\sin m\theta & \\
\cos m\theta &
\end{cases}
\end{align*}
\]

Now the \( R \) equation is solvable

\[
R(r) = J_m(\sqrt{\lambda_{mn}}r),
\]

where \( \lambda_{mn} \) are found from (4.6.24) or equivalently

\[
J_m(\sqrt{\lambda_{mn}}a) = 0, & n = 1, 2, 3, \ldots
\]

Since \( \lambda > 0 \) (related to the zeros of Bessel’s functions), then the \( Z \) equation has the solution

\[
Z(z) = \sinh \sqrt{\lambda_{mn}}z.
\]

Combining the solutions of the ODEs, we have for problem 1:

\[
u(r, \theta, z) = \sum_{m=0}^{\infty} \sum_{n=1}^{\infty} \sinh \sqrt{\lambda_{mn}}z J_m(\sqrt{\lambda_{mn}}r) (A_{mn} \cos m\theta + B_{mn} \sin m\theta),
\]
where $A_{mn}$ and $B_{mn}$ can be found from the generalized Fourier series of $\beta(r, \theta)$. The second problem follows the same pattern, replacing $(4.6.23)$ by

$$Z(H) = 0,$$  \hfill (4.6.31)

leading to

$$u(r, \theta, z) = \sum_{m=0}^{\infty} \sum_{n=1}^{\infty} \sinh \left( \sqrt{\lambda_{mn}} (z - H) \right) J_m(\sqrt{\lambda_{mn}} r) (C_{mn} \cos m\theta + D_{mn} \sin m\theta),$$  \hfill (4.6.32)

where $C_{mn}$ and $D_{mn}$ can be found from the generalized Fourier series of $\alpha(r, \theta)$. The third problem is slightly different. Since there is only one boundary condition for $R$, we must solve the $Z$ equation (4.6.18) before we solve the $R$ equation. The boundary conditions for the $Z$ equation are

$$Z(0) = Z(H) = 0,$$  \hfill (4.6.33)

which result from (4.6.14-4.6.15). The solution of (4.6.18), (4.6.33) is

$$Z_n = \sin \frac{n\pi}{H} z, \quad n = 1, 2, \ldots$$  \hfill (4.6.34)

The eigenvalues

$$\lambda_n = \left( \frac{n\pi}{H} \right)^2, \quad n = 1, 2, \ldots$$  \hfill (4.6.35)

should be substituted in the $R$ equation to yield

$$r (r R')' - \left[ \left( \frac{n\pi}{H} \right)^2 r^2 + m^2 \right] R = 0.$$  \hfill (4.6.36)

This equation looks like Bessel’s equation but with the wrong sign in front of $r^2$ term. It is called the modified Bessel’s equation and has a solution

$$R(r) = c_1 I_m \left( \frac{n\pi}{H} r \right) + c_2 K_m \left( \frac{n\pi}{H} r \right).$$  \hfill (4.6.37)

The modified Bessel functions of the first ($I_m$) and the second ($K_m$) kinds behave at zero and infinity similar to $J_m$ and $Y_m$ respectively. In figure 23 we have plotted the Bessel functions $I_0$ through $I_5$. In figure 24 we have plotted the Bessel functions $K_n, n = 0, 1, 2, 3$. Note that the vertical axis is through $x = .9$ and so it is not so clear that $K_n$ tend to $\infty$ as $x \to 0$.

Therefore the solution to the third problem is

$$u(r, \theta, z) = \sum_{m=0}^{\infty} \sum_{n=1}^{\infty} \sin \frac{n\pi}{H} z J_m \left( \frac{n\pi}{H} r \right) (E_{mn} \cos m\theta + F_{mn} \sin m\theta),$$  \hfill (4.6.38)

where $E_{mn}$ and $F_{mn}$ can be found from the generalized Fourier series of $\gamma(\theta, z)$. The solution of the original problem (4.6.1-4.6.4) is the sum of the solutions given by (4.6.30), (4.6.32) and (4.6.38).
Figure 23: Bessel functions $I_n, n = 0, \ldots, 4$

Figure 24: Bessel functions $K_n, n = 0, \ldots, 3$
Problems

1. Solve Laplace’s equation
   \[ \frac{1}{r} (ru_r)_r + \frac{1}{r^2} u_{\theta\theta} + u_{zz} = 0, \quad 0 \leq r < a, 0 < \theta < 2\pi, 0 < z < H \]
   subject to each of the boundary conditions
   a. \[ u(r, \theta, 0) = \alpha(r, \theta) \]
      \[ u(r, \theta, H) = u(a, \theta, z) = 0 \]
   b. \[ u(r, \theta, 0) = u(r, \theta, H) = 0 \]
      \[ u_r(a, \theta, z) = \gamma(\theta, z) \]
   c. \[ u_z(r, \theta, 0) = \alpha(r, \theta) \]
      \[ u(r, \theta, H) = u(a, \theta, z) = 0 \]
   d. \[ u(r, \theta, 0) = u_z(r, \theta, H) = 0 \]
      \[ u_r(a, \theta, z) = \gamma(z) \]

2. Solve Laplace’s equation
   \[ \frac{1}{r} (ru_r)_r + \frac{1}{r^2} u_{\theta\theta} + u_{zz} = 0, \quad 0 \leq r < a, 0 < \theta < \pi, 0 < z < H \]
   subject to the boundary conditions
   \[ u(r, \theta, 0) = 0, \]
   \[ u_z(r, \theta, H) = 0, \]
   \[ u(r, 0, z) = u(r, \pi, z) = 0, \]
   \[ u(a, \theta, z) = \beta(\theta, z). \]

3. Find the solution to the following steady state heat conduction problem in a box
   \[ \nabla^2 u = 0, \quad 0 \leq x < L, 0 < y < L, 0 < z < W, \]
   subject to the boundary conditions
   \[ \frac{\partial u}{\partial x} = 0, \quad x = 0, x = L, \]
\[ \frac{\partial u}{\partial y} = 0, \quad y = 0, y = L, \]
\[ u(x, y, W) = 0, \]
\[ u(x, y, 0) = 4 \cos \frac{3\pi}{L} x \cos \frac{4\pi}{L} y. \]

4. Find the solution to the following steady state heat conduction problem in a box
\[ \nabla^2 u = 0, \quad 0 \leq x < L, 0 < y < L, 0 < z < W, \]
subject to the boundary conditions
\[ \frac{\partial u}{\partial x} = 0, \quad x = 0, x = L, \]
\[ \frac{\partial u}{\partial y} = 0, \quad y = 0, y = L, \]
\[ u_z(x, y, W) = 0, \]
\[ u_z(x, y, 0) = 4 \cos \frac{3\pi}{L} x \cos \frac{4\pi}{L} y. \]

5. Solve the heat equation inside a cylinder
\[ \frac{\partial u}{\partial t} = \frac{1}{r} \frac{\partial}{\partial r} \left( r \frac{\partial u}{\partial r} \right) + \frac{1}{r^2} \frac{\partial^2 u}{\partial \theta^2} + \frac{\partial^2 u}{\partial z^2}, \quad 0 \leq r < a, 0 < \theta < 2\pi, 0 < z < H \]
subject to the boundary conditions
\[ u(r, \theta, 0, t) = u(r, \theta, H, t) = 0, \]
\[ u(a, \theta, z, t) = 0, \]
and the initial condition
\[ u(r, \theta, z, 0) = f(r, \theta, z). \]
4.7 Laplace’s equation in a sphere

Laplace’s equation in spherical coordinates is given in the form

\[ u_{rr} + \frac{2}{r} u_r + \frac{1}{r^2} u_{\theta\theta} + \cot \theta u_{\theta} + \frac{1}{r^2 \sin^2 \theta} u_{\varphi \varphi} = 0, \quad 0 \leq r < a, \ 0 < \theta < \pi, \ 0 < \varphi < 2\pi, \]

(4.7.1)

\( \varphi \) is the longitude and \( \frac{\pi}{2} - \theta \) is the latitude. Suppose the boundary condition is

\[ u(a, \theta, \varphi) = f(\theta, \varphi). \quad (4.7.2) \]

To solve by the method of separation of variables we assume a solution \( u(r, \theta, \varphi) \) in the form

\[ u(r, \theta, \varphi) = R(r) \Theta(\theta) \Phi(\varphi). \quad (4.7.3) \]

Substitution in Laplace’s equation yields

\[ \left( R'' + \frac{2}{r} R' \right) \Theta \Phi + \frac{1}{r^2} R \Theta'' \Phi + \frac{\cot \theta}{r^2} \Theta' R \Phi + \frac{1}{r^2 \sin^2 \theta} R \Theta \Phi'' = 0. \]

Multiplying by \( \frac{r^2 \sin^2 \theta}{R \Theta \Phi} \), we can separate the \( \varphi \) dependence:

\[ r^2 \sin^2 \theta \left[ \frac{R''}{R} + \frac{2}{r} \frac{R'}{R} + \frac{\cot \theta}{r^2} \frac{\Theta''}{\Theta} + \frac{\cot \theta}{r^2} \frac{\Theta'}{\Theta} \right] = \frac{-\Phi''}{\Phi} = \mu. \]

Now the ODE for \( \varphi \) is

\[ \Phi'' + \mu \Phi = 0 \quad (4.7.4) \]

and the equation for \( r, \theta \) can be separated by dividing through by \( \sin^2 \theta \)

\[ r^2 \frac{R''}{R} + 2r \frac{R'}{R} + \frac{\Theta''}{\Theta} + \cot \theta \frac{\Theta'}{\Theta} = \frac{\mu}{\sin^2 \theta}. \]

Keeping the first two terms on the left, we have

\[ r^2 \frac{R''}{R} + 2r \frac{R'}{R} = -\frac{\Theta''}{\Theta} - \cot \theta \frac{\Theta'}{\Theta} + \frac{\mu}{\sin^2 \theta} = \lambda. \]

Thus

\[ r^2 R'' + 2r R' - \lambda R = 0 \quad (4.7.5) \]

and

\[ \Theta'' + \cot \theta \Theta' - \frac{\mu}{\sin^2 \theta} \Theta + \lambda \Theta = 0. \]

The equation for \( \Theta \) can be written as follows

\[ \sin^2 \theta \Theta'' + \sin \theta \cos \theta \Theta' + (\lambda \sin^2 \theta - \mu) \Theta = 0. \quad (4.7.6) \]
What are the boundary conditions? Clearly, we have periodicity of \( \Phi \), i.e.

\[
\Phi(0) = \Phi(2\pi) \quad (4.7.7)
\]

\[
\Phi'(0) = \Phi'(2\pi). \quad (4.7.8)
\]

The solution \( R(r) \) must be finite at zero, i.e.

\[
|R(0)| < \infty \quad (4.7.9)
\]

as we have seen in other problems on a circular domain that include the pole, \( r = 0 \).

Thus we can solve the ODE (4.7.4) subject to the conditions (4.7.7) - (4.7.8). This yields the eigenvalues

\[
\mu_m = m^2 \quad m = 0, 1, 2, \ldots \quad (4.7.10)
\]

and eigenfunctions

\[
\Phi_m = \begin{cases} 
  \cos m\varphi & m = 1, 2, \ldots \\
  \sin m\varphi
\end{cases} \quad (4.7.11)
\]

and

\[
\Phi_0 = 1. \quad (4.7.12)
\]

We can solve (4.7.5) which is Euler’s equation, by trying

\[
R(r) = r^\alpha \quad (4.7.13)
\]

yielding a characteristic equation

\[
\alpha^2 + \alpha - \lambda = 0. \quad (4.7.14)
\]

The solutions of the characteristic equation are

\[
\alpha_{1,2} = \frac{-1 \pm \sqrt{1 + 4\lambda}}{2}. \quad (4.7.15)
\]

Thus if we take

\[
\alpha_1 = -1 + \sqrt{1 + 4\lambda} \quad (4.7.16)
\]

then

\[
\alpha_2 = -(1 + \alpha_1) \quad (4.7.17)
\]

and

\[
\lambda = \alpha_1(1 + \alpha_1). \quad (4.7.18)
\]

(Recall that the sum of the roots equals the negative of the coefficient of the linear term and the product of the roots equals the constant term.) Therefore the solution is

\[
R(r) = Cr^{\alpha_1} + Dr^{-(\alpha_1+1)} \quad (4.7.19)
\]
Using the boundedness condition (4.7.9) we must have $D = 0$ and the solution of (4.7.5) becomes

$$R(r) = Cr^{\alpha_1}. \quad (4.7.20)$$

Substituting $\lambda$ and $\mu$ from (4.7.18) and (4.7.10) into the third ODE (4.7.6), we have

$$\sin^2 \theta \Theta'' + \sin \theta \cos \theta \Theta' + \left(\alpha_1 (1 + \alpha_1) \sin^2 \theta - m^2\right) \Theta = 0. \quad (4.7.21)$$

Now, let’s make the transformation

$$\xi = \cos \theta \quad (4.7.22)$$

then

$$\frac{d\Theta}{d\theta} = \frac{d\Theta}{d\xi} \frac{d\xi}{d\theta} = -\sin \theta \frac{d\Theta}{d\xi} \quad (4.7.23)$$

and

$$\frac{d^2\Theta}{d\theta^2} = -\frac{d}{d\theta} \left(\sin \theta \frac{d\Theta}{d\xi}\right)$$

$$= -\cos \theta \frac{d\Theta}{d\xi} - \sin \theta \frac{d^2\Theta}{d\xi^2} \frac{d\xi}{d\theta}$$

$$= -\cos \theta \frac{d\Theta}{d\xi} + \sin^2 \theta \frac{d^2\Theta}{d\xi^2}. \quad (4.7.24)$$

Substitute (4.7.22) - (4.7.24) in (4.7.21) we have

$$\sin^4 \theta \frac{d^2\Theta}{d\xi^2} - \sin^2 \theta \cos \theta \frac{d\Theta}{d\xi} - \sin^2 \theta \cos \theta \frac{d\Theta}{d\xi} + \left(\alpha_1 (1 + \alpha_1) \sin^2 \theta - m^2\right) \Theta = 0.$$

Divide through by $\sin^2 \theta$ and use (4.7.22), we get

$$(1 - \xi^2)\Theta'' - 2\xi \Theta' + \left(\alpha_1 (1 + \alpha_1) - \frac{m^2}{1 - \xi^2}\right) \Theta = 0. \quad (4.7.25)$$

This is the so-called associated Legendre equation.

For $m = 0$, the equation is called Legendre’s equation. Using power series method of solution, one can show that Legendre’s equation (see e.g. Pinsky (1991))

$$(1 - \xi^2)\Theta'' - 2\xi \Theta' + \alpha_1 (1 + \alpha_1) \Theta = 0. \quad (4.7.26)$$

has a solution

$$\Theta(\xi) = \sum_{i=0}^{\infty} a_i \xi^i. \quad (4.7.27)$$

where

$$a_{i+2} = \frac{i(i + 1) - \alpha_1 (1 + \alpha_1)}{(i + 1)(i + 2)} a_i, \quad i = 0, 1, 2, \ldots. \quad (4.7.28)$$

and $a_0, a_1$ may be chosen arbitrarily.
If $a_1$ is an integer $n$, then the recurrence relation (4.7.28) shows that one of the solutions is a polynomial of degree $n$. (If $n$ is even, choose $a_1 = 0$, $a_0 \neq 0$ and if $n$ is odd, choose $a_0 = 0$, $a_1 \neq 0$.) This polynomial is denoted by $P_n(\xi)$. The first four Legendre polynomials are

\[
P_0 = 1
\]
\[
P_1 = \xi
\]
\[
P_2 = \frac{3}{2} \xi^2 - \frac{1}{2}
\]
\[
P_3 = \frac{5}{2} \xi^3 - \frac{3}{2} \xi
\]
\[
P_4 = \frac{35}{8} \xi^4 - \frac{30}{8} \xi^2 + \frac{3}{8}.
\]

In Figure 25, we have plotted the first 6 Legendre polynomials. The orthogonality of Legendre polynomials can be easily shown

\[
\int_{-1}^{1} P_n(\xi)P_\ell(\xi) d\xi = 0, \quad \text{for} \quad n \neq \ell
\]

or

\[
\int_{0}^{\pi} P_n(\cos \theta)P_\ell(\cos \theta) \sin \theta d\theta = 0, \quad \text{for} \quad n \neq \ell.
\]
The other solution is not a polynomial and denoted by $Q_n(\xi)$. In fact these functions can be written in terms of inverse hyperbolic tangent.

\[
Q_0 = \tanh^{-1} \xi \\
Q_1 = \xi \tanh^{-1} \xi - 1 \\
Q_2 = \frac{3\xi^2 - 1}{2} \tanh^{-1} \xi - \frac{3\xi}{2} \\
Q_3 = \frac{5\xi^3 - 3\xi}{2} \tanh^{-1} \xi - \frac{15\xi^2 - 4}{6}.
\] (4.7.32)

Now back to (4.7.25), differentiating (4.7.26) $m$ times with respect to $\theta$, one has (4.7.25). Therefore, one solution is

\[
P_n^m(\cos \theta) = \sin^m \theta \frac{d^m}{d\theta^m} P_n(\cos \theta), \quad \text{for} \quad m \leq n
\] (4.7.33)

or in terms of $\xi$

\[
P_n^m(\xi) = (1 - \xi^2)^{m/2} \frac{d^m}{d\xi^m} P_n(\xi), \quad \text{for} \quad m \leq n
\] (4.7.34)

which are the associated Legendre polynomials. The other solution is

\[
Q_n^m(\xi) = (1 - \xi^2)^{m/2} \frac{d^m}{d\xi^m} Q_n(\xi).
\] (4.7.35)
The general solution is then
\[ \Theta_{nm}(\theta) = A P_n^m(\cos \theta) + B Q_n^m(\cos \theta), \quad n = 0, 1, 2, \ldots \] (4.7.36)

Since \( Q_n^m \) has a logarithmic singularity at \( \theta = 0 \), we must have \( B = 0 \). Therefore, the solution becomes
\[ \Theta_{nm}(\theta) = A P_n^m(\cos \theta). \] (4.7.37)

Combining (4.7.11), (4.7.12), (4.7.19) and (4.7.37) we can write
\[
\begin{align*}
    u(r, \theta, \varphi) &= \sum_{n=0}^{\infty} A_n r^n P_n(\cos \theta) \\
    &+ \sum_{n=0}^{\infty} \sum_{m=1}^{n} r^n P_n^m(\cos \theta)(A_{nm} \cos m\varphi + B_{nm} \sin m\varphi). \\
\end{align*}
\] (4.7.38)

where \( P_n(\cos \theta) = P_n(\cos \theta) \) are Legendre polynomials. The boundary condition (4.7.2) implies
\[
\begin{align*}
    f(\theta, \varphi) &= \sum_{n=0}^{\infty} A_{n0} a^n P_n(\cos \theta) \\
    &+ \sum_{n=0}^{\infty} \sum_{m=1}^{n} a^n P_n^m(\cos \theta)(A_{nm} \cos m\varphi + B_{nm} \sin m\varphi). \\
\end{align*}
\] (4.7.39)

The coefficients \( A_{n0}, A_{nm}, B_{nm} \) can be obtained from
\[
\begin{align*}
    A_{n0} &= \frac{\int_0^{2\pi} \int_0^\pi f(\theta, \varphi) P_n(\cos \theta) \sin \theta d\theta d\varphi}{2\pi a^n I_0}, \\
    A_{nm} &= \frac{\int_0^{2\pi} \int_0^\pi f(\theta, \varphi) P_n^m(\cos \theta) \cos m\varphi \sin \theta d\theta d\varphi}{\pi a^n I_m}, \\
    B_{nm} &= \frac{\int_0^{2\pi} \int_0^\pi f(\theta, \varphi) P_n^m(\cos \theta) \sin m\varphi \sin \theta d\theta d\varphi}{\pi a^n I_m}, \\
\end{align*}
\] (4.7.40, 4.7.41, 4.7.42)

where
\[
I_m = \int_0^\pi [P_n^m(\cos \theta)]^2 \sin \theta d\theta = \frac{2(n+m)!}{(2n+1)(n-m)!}. \] (4.7.43)
Problems

1. Solve Laplace’s equation on the sphere

\[ u_{rr} + \frac{2}{r} u_r + \frac{1}{r^2} u_{\theta \theta} + \cot \theta \frac{1}{r^2} u_{\theta \theta} + \frac{1}{r^2 \sin^2 \theta} u_{\varphi \varphi} = 0, \quad 0 \leq r < a, \ 0 < \theta < \pi, \ 0 < \varphi < 2\pi, \]

subject to the boundary condition

\[ u_r(a, \theta, \varphi) = f(\theta). \]

2. Solve Laplace’s equation on the half sphere

\[ u_{rr} + \frac{2}{r} u_r + \frac{1}{r^2} u_{\theta \theta} + \cot \theta \frac{1}{r^2} u_{\theta \theta} + \frac{1}{r^2 \sin^2 \theta} u_{\varphi \varphi} = 0, \quad 0 \leq r < a, \ 0 < \theta < \pi, \ 0 < \varphi < \pi, \]

subject to the boundary conditions

\[ u(a, \theta, \varphi) = f(\theta, \varphi), \]
\[ u(r, \theta, 0) = u(r, \theta, \pi) = 0. \]

3. Solve Laplace’s equation on the surface of the sphere of radius \( a \).
SUMMARY

Heat Equation

\[ u_t = k (u_{xx} + u_{yy}) \]
\[ u_t = k (u_{xx} + u_{yy} + u_{zz}) \]
\[ u_t = k \left\{ \frac{1}{r} \frac{\partial}{\partial r} \left( r \frac{\partial u}{\partial r} \right) + \frac{1}{r^2} \frac{\partial^2 u}{\partial \theta^2} \right\} \]

Wave equation

\[ u_{tt} - c^2 (u_{xx} + u_{yy}) = 0 \]
\[ u_{tt} - c^2 (u_{xx} + u_{yy} + u_{zz}) = 0 \]
\[ u_{tt} = c^2 \left\{ \frac{1}{r} \frac{\partial}{\partial r} \left( r \frac{\partial u}{\partial r} \right) + \frac{1}{r^2} \frac{\partial^2 u}{\partial \theta^2} \right\} \]

Laplace’s Equation

\[ u_{xx} + u_{yy} + u_{zz} = 0 \]
\[ \frac{1}{r} (ru_r)_r + \frac{1}{r^2} u_{\theta\theta} + u_{zz} = 0 \]
\[ u_{rr} + \frac{2}{r} u_r + \frac{1}{r^2} u_{\theta\theta} + \cot \theta \frac{\partial u}{\partial \theta} + \frac{1}{r^2 \sin^2 \theta} u_{\phi\phi} = 0 \]

Bessel’s Equation (inside a circle)

\[ (r R'_m)' + \left( \lambda r - \frac{m^2}{r} \right) R_m = 0, \quad m = 0, 1, 2, \ldots \]
\[ |R_m(0)| < \infty \]
\[ R_m(a) = 0 \]
\[ R_m(r) = J_m \left( \sqrt{\lambda_{mn}} r \right) \quad \text{eigenfunctions} \]
\[ J_m \left( \sqrt{\lambda_{mn}} a \right) = 0 \quad \text{equation for eigenvalues.} \]

Bessel’s Equation (outside a circle)

\[ (r R'_m)' + \left( \lambda r - \frac{m^2}{r} \right) R_m = 0, \quad m = 0, 1, 2, \ldots \]
\[ R_m \to 0 \quad \text{as} \quad r \to \infty \]
\[ R_m(a) = 0 \]
\[ R_m(r) = Y_m \left( \sqrt{\lambda_{mn}} r \right) \quad \text{eigenfunctions} \]
\[ Y_m \left( \sqrt{\lambda_{mn}} a \right) = 0 \quad \text{equation for eigenvalues.} \]
Modified Bessel’s Equation

\[
(r R'_m) - \left( \lambda^2 r + \frac{m^2}{r} \right) R_m = 0, \quad m = 0, 1, 2, \ldots
\]

\[|R_m(0)| < \infty\]

\[R_m(r) = C_{1m} I_m(\lambda r) + C_{2m} K_m(\lambda r)\]

Legendre’s Equation

\[
(1 - \xi^2)\Theta'' - 2\xi \Theta' + \alpha(1 + \alpha)\Theta = 0
\]

\[\Theta(\xi) = C_1 P_n(\xi) + C_2 Q_n(\xi)\]

\[\alpha = n\]

Associated Legendre Equation

\[
(1 - \xi^2)\Theta'' - 2\xi \Theta' + \left( \alpha(1 + \alpha) - \frac{m^2}{1 - \xi^2} \right)\Theta = 0
\]

\[\Theta(\xi) = C_1 P^m_n(\xi) + C_2 Q^m_n(\xi)\]

\[\alpha = n\]
5 Separation of Variables-Nonhomogeneous Problems

In this chapter, we show how to solve nonhomogeneous problems via the separation of variables method. The first section will show how to deal with inhomogeneous boundary conditions. The second section will present the method of eigenfunctions expansion for the inhomogeneous heat equation in one space variable. The third section will give the solution of the wave equation in two dimensions. We close the chapter with the solution of Poisson’s equation.

5.1 Inhomogeneous Boundary Conditions

Consider the following inhomogeneous heat conduction problem:

\[ u_t = ku_{xx} + S(x,t), \quad 0 < x < L \]  

(5.1.1)

subject to the inhomogeneous boundary conditions

\[ u(0,t) = A(t), \]  

(5.1.2)

\[ u(L,t) = B(t), \]  

(5.1.3)

and an initial condition

\[ u(x,0) = f(x). \]  

(5.1.4)

Find a function \( w(x,t) \) satisfying the boundary conditions (5.1.2)-(5.1.3). It is easy to see that

\[ w(x,t) = A(t) + \frac{x}{L} (B(t) - A(t)) \]  

(5.1.5)

is one such function.

Let

\[ v(x,t) = u(x,t) - w(x,t) \]  

(5.1.6)

then clearly

\[ v(0,t) = u(0,t) - w(0,t) = A(t) - A(t) = 0 \]  

(5.1.7)

\[ v(L,t) = u(L,t) - w(L,t) = B(t) - B(t) = 0 \]  

(5.1.8)

i.e. the function \( v(x,t) \) satisfies homogeneous boundary conditions. The question is, what is the PDE satisfied by \( v(x,t) \)? To this end, we differentiate (5.1.6) twice with respect to \( x \) and once with respect to \( t \)

\[ v_x(x,t) = u_x - \frac{1}{L} (B(t) - A(t)) \]  

(5.1.9)

\[ v_{xx}(x,t) = u_{xx} - 0 = u_{xx} \]  

(5.1.10)

\[ v_t(x,t) = u_t - \frac{x}{L} \left( \dot{B}(t) - \dot{A}(t) \right) - \dot{A}(t) \]  

(5.1.11)
and substitute in (5.1.1)

\[ v_t + \dot{A}(t) + \frac{x}{L} \left( \dot{B}(t) - \dot{A}(t) \right) = kv_{xx} + S(x, t). \]  \hspace{1cm} (5.1.12)

Thus

\[ v_t = kv_{xx} + \dot{S}(x, t) \]  \hspace{1cm} (5.1.13)

where

\[ \dot{S}(x, t) = S(x, t) - \dot{A}(t) - \frac{x}{L} \left( \dot{B}(t) - \dot{A}(t) \right). \]  \hspace{1cm} (5.1.14)

The initial condition (5.1.4) becomes

\[ v(x, 0) = f(x) - A(0) - \frac{x}{L} (B(0) - A(0)) = \dot{f}(x). \]  \hspace{1cm} (5.1.15)

Therefore, we have to solve an inhomogeneous PDE (5.1.13) subject to homogeneous boundary conditions (5.1.7)-(5.1.8) and the initial condition (5.1.15).

If the boundary conditions were of a different type, the idea will still be the same. For example, if

\[ u(0, t) = A(t) \]  \hspace{1cm} (5.1.16)
\[ u_x(L, t) = B(t) \]  \hspace{1cm} (5.1.17)

then we try

\[ w(x, t) = \alpha(t)x + \beta(t). \]  \hspace{1cm} (5.1.18)

At \( x = 0 \),

\[ A(t) = w(0, t) = \beta(t) \]

and at \( x = L \),

\[ B(t) = w_x(L, t) = \alpha(t). \]

Thus

\[ w(x, t) = B(t)x + A(t) \]  \hspace{1cm} (5.1.19)

satisfies the boundary conditions (5.1.16)-(5.1.17).

Remark: If the boundary conditions are independent of time, we can take the steady state solution as \( w(x) \).
Problems

1. For each of the following problems obtain the function \( w(x, t) \) that satisfies the boundary conditions and obtain the PDE

a.
\[
    u_t(x, t) = k u_{xx}(x, t) + x, \quad 0 < x < L \\
    u_x(0, t) = 1, \\
    u(L, t) = t.
\]

b.
\[
    u_t(x, t) = k u_{xx}(x, t) + x, \quad 0 < x < L \\
    u(0, t) = 1, \\
    u_x(L, t) = 1.
\]

c.
\[
    u_t(x, t) = k u_{xx}(x, t) + x, \quad 0 < x < L \\
    u_x(0, t) = t, \\
    u_x(L, t) = t^2.
\]

2. Same as problem 1 for the wave equation
\[
    u_{tt} - c^2 u_{xx} = xt, \quad 0 < x < L
\]
subject to each of the boundary conditions

a.
\[
    u(0, t) = 1 \quad u(L, t) = t
\]

b.
\[
    u_x(0, t) = t \quad u_x(L, t) = t^2
\]

c.
\[
    u(0, t) = 0 \quad u_x(L, t) = t
\]

d.
\[
    u_x(0, t) = 0 \quad u_x(L, t) = 1
\]
5.2 Method of Eigenfunction Expansions

In this section, we consider the solution of the inhomogeneous heat equation

\[
  u_t = ku_{xx} + S(x,t), \quad 0 < x < L
\]  

\[
  u(0,t) = 0, \tag{5.2.2}
\]

\[
  u(L,t) = 0, \tag{5.2.3}
\]

\[
  u(x,0) = f(x). \tag{5.2.4}
\]

The solution of the homogeneous PDE leads to the eigenfunctions

\[
  \phi_n(x) = \sin \left( \frac{n\pi x}{L} \right), \quad n = 1, 2, \ldots
\]  

\[
  \lambda_n = \left( \frac{n\pi}{L} \right)^2, \quad n = 1, 2, \ldots
\]

Clearly the eigenfunctions depend on the boundary conditions and the PDE. Having the eigenfunctions, we now expand the source term

\[
  S(x,t) = \sum_{n=1}^{\infty} s_n(t) \phi_n(x), \tag{5.2.7}
\]

where

\[
  s_n(t) = \frac{\int_0^L S(x,t) \phi_n(x) dx}{\int_0^L \phi_n^2(x) dx}. \tag{5.2.8}
\]

Let

\[
  u(x,t) = \sum_{n=1}^{\infty} u_n(t) \phi_n(x), \tag{5.2.9}
\]

then

\[
  f(x) = u(x,0) = \sum_{n=1}^{\infty} u_n(0) \phi_n(x). \tag{5.2.10}
\]

Since \( f(x) \) is known, we have

\[
  u_n(0) = \frac{\int_0^L f(x) \phi_n(x) dx}{\int_0^L \phi_n^2(x) dx}. \tag{5.2.11}
\]

Substitute \( u(x,t) \) from (5.2.9) and its derivatives and \( S(x,t) \) from (5.2.7) into (5.2.1), we have

\[
  \sum_{n=1}^{\infty} \ddot{u}_n(t) \phi_n(x) = \sum_{n=1}^{\infty} (-k \lambda_n) u_n(t) \phi_n(x) + \sum_{n=1}^{\infty} s_n(t) \phi_n(x). \tag{5.2.12}
\]

Recall that \( u_{xx} \) gives a series with \( \phi_n''(x) \) which is \(-\lambda_n \phi_n\), since \( \lambda_n \) are the eigenvalues corresponding to \( \phi_n \). Combining all three sums in (5.2.12), one has

\[
  \sum_{n=1}^{\infty} \{ \ddot{u}_n(t) + k \lambda_n u_n(t) - s_n(t) \} \phi_n(x) = 0. \tag{5.2.13}
\]
Therefore
\[ \dot{u}_n(t) + k\lambda_n u_n(t) = s_n(t), \quad n = 1, 2, \ldots \] (5.2.14)

This inhomogeneous ODE should be combined with the initial condition (5.2.11).

The solution of (5.2.14), (5.2.11) is obtained by the method of variation of parameters (see e.g. Boyce and DiPrima)
\[ u_n(t) = u_n(0)e^{-\lambda_n kt} + \int_0^t s_n(\tau)e^{-\lambda_n k(t-\tau)}d\tau. \] (5.2.15)

It is easy to see that \( u_n(t) \) above satisfies (5.2.11) and (5.2.14). We summarize the solution by (5.2.9),(5.2.15),(5.2.11) and (5.2.8).

**Example**
\[ u_t = u_{xx} + 1, \quad 0 < x < 1 \] (5.2.16)
\[ u_x(0, t) = 2, \] (5.2.17)
\[ u(1, t) = 0, \] (5.2.18)
\[ u(x, 0) = x(1 - x). \] (5.2.19)

The function \( w(x, t) \) to satisfy the inhomogeneous boundary conditions is
\[ w(x, t) = 2x - 2. \] (5.2.20)

The function
\[ v(x, t) = u(x, t) - w(x, t) \] (5.2.21)
satisfies the following PDE
\[ v_t = v_{xx} + 1, \] (5.2.22)

since \( w_t = w_{xx} = 0 \). The initial condition is
\[ v(x, 0) = x(1 - x) - (2x - 2) = x(1 - x) + 2(1 - x) = (x + 2)(1 - x) \] (5.2.23)

and the homogeneous boundary conditions are
\[ v_x(0, t) = 0, \] (5.2.24)
\[ v(1, t) = 0. \] (5.2.25)

The eigenfunctions \( \phi_n(x) \) and eigenvalues \( \lambda_n \) satisfy
\[ \phi_n''(x) + \lambda_n \phi_n = 0, \] (5.2.26)
\[ \phi_n'(0) = 0, \] (5.2.27)
\[ \phi_n(1) = 0. \] (5.2.28)

Thus
\[ \phi_n(x) = \cos\left(n - \frac{1}{2}\right)\pi x, \quad n = 1, 2, \ldots \] (5.2.29)
\[ \lambda_n = \left[ (n - \frac{1}{2})\pi \right]^2. \]  
(5.2.30)

Expanding \( S(x, t) = 1 \) and \( v(x, t) \) in these eigenfunctions we have

\[ 1 = \sum_{n=1}^{\infty} s_n \phi_n(x) \]  
(5.2.31)

where

\[ s_n = \int_0^1 \frac{1}{\cos^2(n - \frac{1}{2})\pi x} \, dx = \frac{4(-1)^{n-1}}{(2n-1)\pi}, \]  
(5.2.32)

and

\[ v(x, t) = \sum_{n=1}^{\infty} v_n(t) \cos(n - \frac{1}{2})\pi x. \]  
(5.2.33)

The partial derivatives of \( v(x, t) \) required are

\[ v_t(x, t) = \sum_{n=1}^{\infty} \dot{v}_n(t) \cos(n - \frac{1}{2})\pi x, \]  
(5.2.34)

\[ v_{xx}(x, t) = -\sum_{n=1}^{\infty} \left( n - \frac{1}{2} \right)^2 \pi^2 v_n(t) \cos(n - \frac{1}{2})\pi x. \]  
(5.2.35)

Thus, upon substituting (5.2.34),(5.2.35) and (5.2.31) into (5.2.22), we get

\[ \dot{v}_n(t) + \left( n - \frac{1}{2} \right)^2 \pi^2 v_n(t) = s_n. \]  
(5.2.36)

The initial condition \( v_n(0) \) is given by the eigenfunction expansion of \( v(x, 0) \), i.e.

\[ (x + 2)(1 - x) = \sum_{n=1}^{\infty} v_n(0) \cos(n - \frac{1}{2})\pi x \]  
(5.2.37)

so

\[ v_n(0) = \frac{\int_0^1 (x + 2)(1 - x) \cos(n - \frac{1}{2})\pi x \, dx}{\int_0^1 \cos^2(n - \frac{1}{2})\pi x \, dx}. \]  
(5.2.38)

The solution of (5.2.36) is

\[ v_n(t) = v_n(0)e^{-\left( n - \frac{1}{2} \frac{\pi}{\pi} \right)^2 t} + s_n \int_0^t e^{-\left( n - \frac{1}{2} \frac{\pi}{\pi} \right)^2 (t-\tau)} \, d\tau \]

Performing the integration

\[ v_n(t) = v_n(0)e^{-\left( n - \frac{1}{2} \frac{\pi}{\pi} \right)^2 t} + s_n \frac{1 - e^{-\left( n - \frac{1}{2} \frac{1}{2} \pi \right)^2 t}}{\left( n - \frac{1}{2} \frac{1}{2} \pi \right)^2} \]  
(5.2.39)

where \( v_n(0), s_n \) are given by (5.2.38) and (5.2.32) respectively.
Problems

1. Solve the heat equation

\[ u_t = k u_{xx} + x, \quad 0 < x < L \]

subject to the initial condition

\[ u(x, 0) = x(L - x) \]

and each of the boundary conditions

a.

\[ u_x(0, t) = 1, \]
\[ u(L, t) = t. \]

b.

\[ u(0, t) = 1, \]
\[ u_x(L, t) = 1. \]

c.

\[ u_x(0, t) = t, \]
\[ u_x(L, t) = t^2. \]

2. Solve the heat equation

\[ u_t = u_{xx} + e^{-t}, \quad 0 < x < \pi, \quad t > 0, \]

subject to the initial condition

\[ u(x, 0) = \cos 2x, \quad 0 < x < \pi, \]

and the boundary condition

\[ u_x(0, t) = u_x(\pi, t) = 0. \]
5.3 Forced Vibrations

In this section we solve the inhomogeneous wave equation in two dimensions describing the forced vibrations of a membrane.

\[ u_{tt} = c^2 \nabla^2 u + S(x, y, t) \]  \hspace{1cm} (5.3.1)

subject to the boundary condition

\[ u(x, y, t) = 0, \text{ on the boundary}, \]  \hspace{1cm} (5.3.2)

and initial conditions

\[ u(x, y, 0) = \alpha(x, y), \]  \hspace{1cm} (5.3.3)
\[ u_t(x, y, 0) = \beta(x, y). \]  \hspace{1cm} (5.3.4)

Since the boundary condition is homogeneous, we can expand the solution \( u(x, y, t) \) and the forcing term \( S(x, y, t) \) in terms of the eigenfunctions \( \phi_n(x, y) \), i.e.

\[ u(x, y, t) = \sum_{i=1}^{\infty} u_i(t) \phi_i(x, y), \]  \hspace{1cm} (5.3.5)
\[ S(x, y, t) = \sum_{i=1}^{\infty} s_i(t) \phi_i(x, y), \]  \hspace{1cm} (5.3.6)

where

\[ \nabla^2 \phi_i = -\lambda_i \phi_i, \]  \hspace{1cm} (5.3.7)
\[ \phi_i(x, y) = 0, \text{ on the boundary}, \]  \hspace{1cm} (5.3.8)

and

\[ s_i(t) = \frac{\iint S(x, y, t) \phi_i(x, y) dx dy}{\iint \phi_i^2(x, y) dx dy}. \]  \hspace{1cm} (5.3.9)

Substituting (5.3.5),(5.3.6) into (5.3.1) we have

\[ \sum_{i=1}^{\infty} \ddot{u}_i(t) \phi_i(x, y) = c^2 \sum_{i=1}^{\infty} u_i(t) \nabla^2 \phi_i + \sum_{i=1}^{\infty} s_i(t) \phi_i(x, y). \]

Using (5.3.7) and combining all the sums, we get an ODE for the coefficients \( u_i(t) \),

\[ \ddot{u}_i(t) + c^2 \lambda_i u_i(t) = s_i(t). \]  \hspace{1cm} (5.3.10)

The solution can be found in any ODE book,

\[ u_i(t) = c_1 \cos c \sqrt{\lambda_i} t + c_2 \sin c \sqrt{\lambda_i} t + \int_0^t s_i(\tau) \sin c \sqrt{\lambda_i} (t - \tau) \frac{d\tau}{c \sqrt{\lambda_i}}. \]  \hspace{1cm} (5.3.11)

The initial conditions (5.3.3)-(5.3.4) imply

\[ u_i(0) = c_1 = \frac{\int \int \alpha(x, y) \phi_i(x, y) dx dy}{\int \int \phi_i^2(x, y) dx dy}, \]  \hspace{1cm} (5.3.12)
\[ \dot{u}_i(0) = c_2 c \sqrt{\lambda_i} = \frac{\int \int \beta(x, y) \phi_i(x, y) dx dy}{\int \int \phi_i^2(x, y) dx dy}. \]  \hspace{1cm} (5.3.13)

Equations (5.3.12)-(5.3.13) can be solved for \( c_1 \) and \( c_2 \). Thus the solution \( u(x, y, t) \) is given by (5.3.5) with \( u_i(t) \) given by (5.3.11)-(5.3.13) and \( s_i(t) \) are given by (5.3.9).
5.3.1 Periodic Forcing

If the forcing $S(x, y, t)$ is a periodic function in time, we have an interesting case. Suppose

$$S(x, y, t) = \sigma(x, y) \cos \omega t,$$  \hspace{1cm} (5.3.1.1)

then by (5.3.9) we have

$$s_i(t) = \sigma_i \cos \omega t,$$  \hspace{1cm} (5.3.1.2)

where

$$\sigma_i(t) = \frac{\iint \sigma(x, y) \phi_i(x, y) dx dy}{\iint \phi_i^2(x, y) dx dy}.$$

(5.3.1.3)

The ODE for the unknown $u_i(t)$ becomes

$$\ddot{u}_i(t) + c^2 \lambda_i u_i(t) = \sigma_i \cos \omega t.$$  \hspace{1cm} (5.3.1.4)

In this case the particular solution of the nonhomogeneous is

$$u_i(t) = \frac{\sigma_i}{c^2 \lambda_i - \omega^2} \cos \omega t.$$  \hspace{1cm} (5.3.1.5)

and thus

$$u_i(t) = c_1 \cos c \sqrt{\lambda_i} t + c_2 \sin c \sqrt{\lambda_i} t + \frac{\sigma_i}{c^2 \lambda_i - \omega^2} \cos \omega t.$$  \hspace{1cm} (5.3.1.6)

The amplitude $u_i(t)$ of the mode $\phi_i(x, y)$ is decomposed to a vibration at the natural frequency $c \sqrt{\lambda_i}$ and a vibration at the forcing frequency $\omega$. What happens if $\omega$ is one of the natural frequencies, i.e.

$$\omega = c \sqrt{\lambda_i} \quad \text{for some} \quad i.$$  \hspace{1cm} (5.3.1.7)

Then the denominator in (5.3.1.6) vanishes. The particular solution should not be (5.3.1.5) but rather

$$\frac{\sigma_i}{2 \omega} t \sin \omega t.$$  \hspace{1cm} (5.3.1.8)

The amplitude is growing linearly in $t$. This is called resonance.
Problems

1. Consider a vibrating string with time dependent forcing

\[ u_{tt} - c^2 u_{xx} = S(x, t), \quad 0 < x < L \]

subject to the initial conditions

\[ u(x, 0) = f(x), \]
\[ u_t(x, 0) = 0, \]

and the boundary conditions

\[ u(0, t) = u(L, t) = 0. \]

a. Solve the initial value problem.

b. Solve the initial value problem if \( S(x, t) = \cos \omega t \). For what values of \( \omega \) does resonance occur?

2. Consider the following damped wave equation

\[ u_{tt} - c^2 u_{xx} + \beta u_t = \cos \omega t, \quad 0 < x < \pi, \]

subject to the initial conditions

\[ u(x, 0) = f(x), \]
\[ u_t(x, 0) = 0, \]

and the boundary conditions

\[ u(0, t) = u(\pi, t) = 0. \]

Solve the problem if \( \beta \) is small \((0 < \beta < 2c)\).

3. Solve the following

\[ u_{tt} - c^2 u_{xx} = S(x, t), \quad 0 < x < L \]

subject to the initial conditions

\[ u(x, 0) = f(x), \]
\[ u_t(x, 0) = 0, \]

and each of the following boundary conditions

a. 
\[ u(0, t) = A(t) \quad u(L, t) = B(t) \]

b. 
\[ u(0, t) = 0 \quad u_x(L, t) = 0 \]

c. 
\[ u_x(0, t) = A(t) \quad u(L, t) = 0. \]
4. Solve the wave equation

\[ u_{tt} - c^2 u_{xx} = x t, \quad 0 < x < L, \]

subject to the initial conditions

\[ u(x, 0) = \sin x \]
\[ u_t(x, 0) = 0 \]

and each of the boundary conditions

a. \[ u(0, t) = 1, \]
\[ u(L, t) = t. \]

b. \[ u_x(0, t) = t, \]
\[ u_x(L, t) = t^2. \]

c. \[ u(0, t) = 0, \]
\[ u_x(L, t) = t. \]

d. \[ u_x(0, t) = 0, \]
\[ u_x(L, t) = 1. \]

5. Solve the wave equation

\[ u_{tt} - u_{xx} = 1, \quad 0 < x < L, \]

subject to the initial conditions

\[ u(x, 0) = f(x) \]
\[ u_t(x, 0) = g(x) \]

and the boundary conditions

\[ u(0, t) = 1, \]
\[ u_x(L, t) = B(t). \]
5.4 Poisson’s Equation

In this section we solve Poisson’s equation subject to homogeneous and nonhomogeneous boundary conditions. In the first case we can use the method of eigenfunction expansion in one dimension and two.

5.4.1 Homogeneous Boundary Conditions

Consider Poisson’s equation

$$\nabla^2 u = S, \quad (5.4.1.1)$$

subject to homogeneous boundary condition, e.g.

$$u = 0, \quad \text{on the boundary.} \quad (5.4.1.2)$$

The problem can be solved by the method of eigenfunction expansion. To be specific we suppose the domain is a rectangle of length \(L\) and height \(H\), see figure 27.

We first consider the one dimensional eigenfunction expansion, i.e.

$$\phi_n(x) = \sin \frac{n\pi}{L}x, \quad (5.4.1.3)$$

and

$$u(x, y) = \sum_{n=1}^{\infty} u_n(y) \sin \frac{n\pi}{L}x. \quad (5.4.1.4)$$

Substitution in Poisson’s equation, we get

$$\sum_{n=1}^{\infty} \left[ u_n''(y) - \left(\frac{n\pi}{L}\right)^2 u_n(y) \right] \sin \frac{n\pi}{L}x = \sum_{n=1}^{\infty} s_n(y) \sin \frac{n\pi}{L}x, \quad (5.4.1.5)$$

where

$$s_n(y) = \frac{2}{L} \int_0^L S(x, y) \sin \frac{n\pi}{L}x dx. \quad (5.4.1.6)$$

The other boundary conditions lead to

$$u_n(0) = 0, \quad (5.4.1.7)$$

$$u_n(H) = 0. \quad (5.4.1.8)$$

So we end up with a boundary value problem for \(u_n(y)\), i.e.

$$u_n''(y) - \left(\frac{n\pi}{L}\right)^2 u_n(y) = s_n(y), \quad (5.4.1.9)$$

subject to \((5.4.1.7)-(5.4.1.8)\).

It requires a lengthy algebraic manipulation to show that the solution is

$$u_n(y) = \frac{\sinh \frac{n\pi(H-y)}{L}}{\frac{n\pi}{L} \sinh \frac{n\pi H}{L}} \int_0^y s_n(\xi) \sinh \frac{n\pi}{L}\xi d\xi + \frac{\sinh \frac{n\pi y}{L}}{\frac{n\pi}{L} \sinh \frac{n\pi H}{L}} \int_y^H s_n(\xi) \sinh \frac{n\pi}{L}(H-\xi) d\xi. \quad (5.4.1.10)$$
Figure 27: Rectangular domain

So the solution is given by (5.4.1.4) with \( u_n(y) \) and \( s_n(y) \) given by (5.4.1.10) and (5.4.1.6) respectively.

Another approach, related to the first, is the use of two dimensional eigenfunctions. In the example,

\[
\phi_{nm} = \sin \frac{n\pi}{L} x \sin \frac{m\pi}{H} y, \quad (5.4.1.11)
\]

\[
\lambda_{nm} = \left( \frac{n\pi}{L} \right)^2 + \left( \frac{m\pi}{H} \right)^2. \quad (5.4.1.12)
\]

We then write the solution

\[
u(x, y) = \sum_{n=1}^{\infty} \sum_{m=1}^{\infty} u_{nm} \phi_{nm}(x, y). \quad (5.4.1.13)
\]

Substituting (5.4.1.13) into the equation, we get

\[
\sum_{n=1}^{\infty} \sum_{m=1}^{\infty} (-u_{nm})\lambda_{nm}\sin \frac{n\pi}{L} x \sin \frac{m\pi}{H} y = S(x, y). \quad (5.4.1.14)
\]

Therefore \(-u_{nm}\lambda_{nm}\) are the coefficients of the double Fourier series expansion of \( S(x, y) \), that is

\[
u_{nm} = \frac{\int_0^L \int_0^H S(x, y) \sin \frac{n\pi}{L} x \sin \frac{m\pi}{H} y dy dx}{-\lambda_{nm} \int_0^L \int_0^H \sin^2 \frac{n\pi}{L} x \sin^2 \frac{m\pi}{H} y dy dx}. \quad (5.4.1.15)
\]

This double series may converge slower than the previous solution.
5.4.2 Inhomogeneous Boundary Conditions

The problem is then
\[ \nabla^2 u = S, \]  
(5.4.2.1)
subject to inhomogeneous boundary condition, e.g.,
\[ u = \alpha, \quad \text{on the boundary.} \]  
(5.4.2.2)

The eigenvalues \( \lambda_i \) and the eigenfunctions \( \phi_i \) satisfy
\[ \nabla^2 \phi_i = -\lambda_i \phi_i, \]  
(5.4.2.3)
\[ \phi_i = 0, \quad \text{on the boundary.} \]  
(5.4.2.4)

Since the boundary condition (5.4.2.2) is not homogeneous, we cannot differentiate the infinite series term by term. But note that the coefficients \( u_n \) of the expansion are given by:
\[ u_n = \frac{\int \int u(x,y) \phi_n(x,y) dxdy}{\int \int \phi_n^2(x,y) dxdy} = -\frac{1}{\lambda_n} \frac{\int \int u \nabla^2 \phi_n dxdy}{\int \int \phi_n^2 dxdy}. \]  
(5.4.2.5)

Using Green’s formula, i.e.,
\[ \int \int u \nabla^2 \phi_n dxdy = \int \int \phi_n \nabla^2 u dxdy + \oint (u \nabla \phi_n - \phi_n \nabla u) \cdot \mathbf{n} ds, \]
substituting from (5.4.2.1), (5.4.2.2) and (5.4.2.4)
\[ = \int \int \phi_n S dxdy + \oint \alpha \nabla \phi_n \cdot \mathbf{n} ds \]  
(5.4.2.6)

Therefore the coefficients \( u_n \) become (combining (5.4.2.5)-(5.4.2.6))
\[ u_n = -\frac{1}{\lambda_n} \frac{\int \int S \phi_n dxdy + \oint \alpha \nabla \phi_n \cdot \mathbf{n} ds}{\int \int \phi_n^2 dxdy}. \]  
(5.4.2.7)

If \( \alpha = 0 \) we get (5.4.1.15). The case \( \lambda = 0 \) will not be discussed here.
Problems

1. Solve
\[ \nabla^2 u = S(x, y), \quad 0 < x < L, \quad 0 < y < H, \]
a.
\[ u(0, y) = u(L, y) = 0 \]
\[ u(x, 0) = u(x, H) = 0 \]
Use a Fourier sine series in \( y \).
b.
\[ u(0, y) = 0 \quad u(L, y) = 1 \]
\[ u(x, 0) = u(x, H) = 0 \]
Hint: Do NOT reduce to homogeneous boundary conditions.
c.
\[ u_x(0, y) = u_x(L, y) = 0 \]
\[ u_y(x, 0) = u_y(x, H) = 0 \]
In what situations are there solutions?

2. Solve the following Poisson's equation
\[ \nabla^2 u = e^{2y} \sin x, \quad 0 < x < \pi, \quad 0 < y < L, \]
\[ u(0, y) = u(\pi, y) = 0, \]
\[ u(x, 0) = 0, \]
\[ u(x, L) = f(x). \]
5.4.3 One Dimensional Boundary Value Problems

As a special case of Poisson's equation, we briefly discuss here the solution of boundary value problems, e.g.

\[ y''(x) = 1, \quad 0 < x < 1 \]  

subject to

\[ y(0) = y(1) = 0. \]  

Clearly one can solve this trivial problem by integrating twice

\[ y(x) = \frac{1}{2}x^2 - \frac{1}{2}x. \]  

One can also use the method of eigenfunctions expansion. In this case the eigenvalues and eigenfunctions are obtained by solving

\[ y''(x) = \lambda y(x), \quad 0 < x < 1 \]  

subject to the same boundary conditions. The eigenvalues are

\[ \lambda_n = -(n\pi)^2, \quad n = 1, 2, \ldots \]  

and the eigenfunctions are

\[ \phi_n = \sin n\pi x, \quad n = 1, 2, \ldots. \]  

Now we expand the solution \( y(x) \) and the right hand side in terms of these eigenfunctions

\[ y(x) = \sum_{n=1}^{\infty} y_n \sin n\pi x, \]  

\[ 1 = \sum_{n=1}^{\infty} r_n \sin n\pi x. \]  

The coefficients \( r_n \) can be found easily (see Chapter 3)

\[ r_n = \begin{cases} 
\frac{4}{n\pi} & n \text{ odd} \\
0 & n \text{ even} 
\end{cases} \]  

Substituting the expansions in the equation and comparing coefficients, we get

\[ y_n = \frac{r_n}{\lambda_n} \]  

that is

\[ y_n = \begin{cases} 
\frac{4}{(n\pi)^3} & n \text{ odd} \\
0 & n \text{ even} 
\end{cases} \]  

This is the Fourier sine series representation of the solution given earlier.

Can do lab 3
Problems

1. Find the eigenvalues and corresponding eigenfunctions in each of the following boundary value problems.
   
   (a) \( y'' - \lambda^2 y = 0 \quad 0 < x < a \quad y'(0) = y'(a) = 0 \)
   
   (b) \( y'' - \lambda^2 y = 0 \quad 0 < x < a \quad y(0) = 0 \quad y(a) = 1 \)
   
   (c) \( y'' + \lambda^2 y = 0 \quad 0 < x < a \quad y(0) = y'(a) = 0 \)
   
   (d) \( y'' + \lambda^2 y = 0 \quad 0 < x < a \quad y(0) = 1 \quad y'(a) = 0 \)

2. Find the eigenfunctions of the following boundary value problem.
   
   \( y'' + \lambda^2 y = 0 \quad 0 < x < 2\pi \quad y(0) = y(2\pi) \quad y'(0) = y'(2\pi) \)

3. Obtain the eigenvalues and eigenfunctions of the problem.
   
   \( y'' + y' + (\lambda + 1)y = 0 \quad 0 < x < \pi \quad y(0) = y(\pi) = 0 \)

4. Obtain the orthonormal set of eigenfunctions for the problem.
   
   (a) \( y'' + \lambda y = 0 \quad 0 < x < \pi \quad y'(0) = 0 \quad y(\pi) = 0 \)
   
   (b) \( y'' + (1 + \lambda)y = 0 \quad 0 < x < \pi \quad y(0) = 0 \quad y(\pi) = 0 \)
   
   (c) \( y'' + \lambda y = 0 \quad -\pi < x < \pi \quad y'(\pi) = 0 \quad y'(\pi) = 0 \)
SUMMARY
Nonhomogeneous problems
  1. Find a function $w$ that satisfies the inhomogeneous boundary conditions (except for
     Poisson’s equation).
  2. Let $v = u - w$, then $v$ satisfies an inhomogeneous PDE with homogeneous boundary
     conditions.
  3. Solve the homogeneous equation with homogeneous boundary conditions to obtain
     eigenvalues and eigenfunctions.
  4. Expand the solution $v$, the right hand side (source/sink) and initial condition(s) in
     eigenfunctions series.
  5. Solve the resulting inhomogeneous ODE.
6 Classification and Characteristics

In this chapter we classify the linear second order PDEs. This will require a discussion of transformations, characteristic curves and canonical forms. We will show that there are three types of PDEs and establish that these three cases are in a certain sense typical of what occurs in the general theory. The type of equation will turn out to be decisive in establishing the kind of initial and boundary conditions that serve in a natural way to determine a solution uniquely (see e.g. Garabedian (1964)).

6.1 Physical Classification

Partial differential equations can be classified as equilibrium problems and marching problems. The first class, equilibrium or steady state problems are also known as elliptic. For example, Laplace’s or Poisson’s equations are of this class. The marching problems include both the parabolic and hyperbolic problems, i.e. those whose solution depends on time.

6.2 Classification of Linear Second Order PDEs

Recall that a linear second order PDE in two variables is given by

\[ Au_{xx} + Bu_{xy} + Cu_{yy} + Du_x + Eu_y + Fu = G \]  \hspace{1cm} (6.2.1)

where all the coefficients \( A \) through \( F \) are real functions of the independent variables \( x, y \).

Define a discriminant \( \Delta(x, y) \) by

\[ \Delta(x_0, y_0) = B^2(x_0, y_0) - 4A(x_0, y_0)C(x_0, y_0). \]  \hspace{1cm} (6.2.2)

(Notice the similarity to the discriminant defined for conic sections.)

**Definition 7.** An equation is called hyperbolic at the point \( (x_0, y_0) \) if \( \Delta(x_0, y_0) > 0 \). It is parabolic at that point if \( \Delta(x_0, y_0) = 0 \) and elliptic if \( \Delta(x_0, y_0) < 0 \).

The classification for equations with more than two independent variables or with higher order derivatives are more complicated. See Courant and Hilbert [5].

**Example.**

\[ u_{tt} - c^2u_{xx} = 0 \]

\[ A = 1, \ B = 0, \ C = -c^2 \]

Therefore,

\[ \Delta = 0^2 - 4 \cdot 1(-c^2) = 4c^2 \geq 0 \]

Thus the problem is hyperbolic for \( c \neq 0 \) and parabolic for \( c = 0 \).

The transformation leads to the discovery of special loci known as characteristic curves along which the PDE provides only an incomplete expression for the second derivatives. Before we discuss transformation to canonical forms, we will motivate the name and explain why such transformation is useful. The name canonical form is used because this form
corresponds to particularly simple choices of the coefficients of the second partial derivatives. Such transformation will justify why we only discuss the method of solution of three basic equations (heat equation, wave equation and Laplace’s equation). Sometimes, we can obtain the solution of a PDE once it is in a canonical form (several examples will be given later in this chapter). Another reason is that characteristics are useful in solving first order quasilinear and second order linear hyperbolic PDEs, which will be discussed in the next chapter. (In fact nonlinear first order PDEs can be solved that way, see for example F. John (1982).)

To transform the equation into a canonical form, we first show how a general transformation affects equation (6.2.1). Let $\xi, \eta$ be twice continuously differentiable functions of $x, y$

$$\xi = \xi(x, y), \quad \eta = \eta(x, y).$$

Suppose also that the Jacobian $J$ of the transformation defined by

$$J = \begin{vmatrix} \xi_x & \xi_y \\ \eta_x & \eta_y \end{vmatrix}$$

is non zero. This assumption is necessary to ensure that one can make the transformation back to the original variables $x, y$.

Use the chain rule to obtain all the partial derivatives required in (6.2.1). It is easy to see that

$$u_x = u_\xi \xi_x + u_\eta \eta_x$$

$$u_y = u_\xi \xi_y + u_\eta \eta_y.$$  

The second partial derivatives can be obtained as follows:

$$u_{xy} = (u_x)_y = (u_\xi \xi_x + u_\eta \eta_x)_y$$

$$= (u_\xi \xi_x)_y + (u_\eta \eta_x)_y$$

$$= (u_\xi)_y \xi_x + u_\xi \xi_{xy} + (u_\eta)_y \eta_x + u_\eta \eta_{xy}$$

Now use (6.2.7)

$$u_{xy} = (u_\xi \xi_x + u_\xi \eta_y) \xi_x + u_\xi \xi_{xy} + (u_\eta \xi_x + u_\eta \eta_y) \eta_x + u_\eta \eta_{xy}.$$  

Reorganize the terms

$$u_{xy} = u_\xi \xi_x \xi_y + u_\xi \eta_y (\xi_x \eta_y + \xi_y \eta_x) + u_\eta \eta_x \eta_y + u_\xi \xi_{xy} + u_\eta \eta_{xy}.$$  

In a similar fashion we get $u_{xx}, u_{yy}$

$$u_{xx} = u_\xi \xi_x^2 + 2u_\xi \xi_x \eta_x + u_\eta \eta_x^2 + u_\xi \xi_{xx} + u_\eta \eta_{xx},$$

$$u_{yy} = u_\xi \xi_y^2 + 2u_\xi \xi_y \eta_y + u_\eta \eta_y^2 + u_\xi \xi_{yy} + u_\eta \eta_{yy}.$$  

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Introducing these into (6.2.1) one finds after collecting like terms

\[ A^* u_{\xi\xi} + B^* u_{\xi\eta} + C^* u_{\eta\eta} + D^* u_{\xi} + E^* u_{\eta} + F^* u = G^* \]  

(6.2.11)

where all the coefficients are now functions of \( \xi, \eta \) and

\[
\begin{align*}
A^* &= A_{xx}^2 + B_{xx} \xi_x \xi_y + C_{xx}^2, \\
B^* &= 2A_{xx} \xi_x \eta_x + B(\xi_x \eta_y + \xi_y \eta_x) + 2C_{xx} \eta_y, \\
C^* &= A_{xx}^2 + B_{xx} \eta_x \eta_y + C_{xx}^2, \\
D^* &= A_{xx}^2 + B_{xx} \xi_x + C_{xx} \xi_y + D_{xx} \xi_x + E_{xx} \xi_y, \\
E^* &= A_{xx} \eta_x + B_{xx} \eta_y + C_{xx} \eta_y + D_{xx} \eta_x + E_{xx} \eta_y, \\
F^* &= F, \\
G^* &= G.
\end{align*}
\]

(6.2.12) - (6.2.18)

The resulting equation (6.2.11) is in the same form as the original one. The type of the equation (hyperbolic, parabolic or elliptic) will not change under this transformation. The reason for this is that

\[ \Delta^* = (B^*)^2 - 4A^*C^* = J^2(B^2 - 4AC) = J^2 \Delta \]  

(6.2.19)

and since \( J \neq 0 \), the sign of \( \Delta^* \) is the same as that of \( \Delta \). Proving (6.2.19) is not complicated but definitely messy. It is left for the reader as an exercise using a symbolic manipulator such as MACSYMA or MATHEMATICA.

The classification depends only on the coefficients of the second derivative terms and thus we write (6.2.1) and (6.2.11) respectively as

\[ Au_{xx} + Bu_{xy} + C u_{yy} = H(x, y, u, u_x, u_y) \]  

(6.2.20)

and

\[ A^* u_{\xi\xi} + B^* u_{\xi\eta} + C^* u_{\eta\eta} = H^*(\xi, \eta, u, u_{\xi}, u_{\eta}). \]  

(6.2.21)
Problems

1. Classify each of the following as hyperbolic, parabolic or elliptic at every point \((x, y)\) of the domain

   a. \(x u_{xx} + u_{yy} = x^2\)
   b. \(x^2 u_{xx} - 2xy u_{xy} + y^2 u_{yy} = e^x\)
   c. \(e^x u_{xx} + e^y u_{yy} = u\)
   d. \(u_{xx} + u_{xy} - xu_{yy} = 0\) in the left half plane \((x \leq 0)\)
   e. \(x^2 u_{xx} + 2xy u_{xy} + y^2 u_{yy} + xy u_x + y^2 u_y = 0\)
   f. \(u_{xx} + xu_{yy} = 0\) (Tricomi equation)

2. Classify each of the following constant coefficient equations

   a. \(4u_{xx} + 5u_{xy} + u_{yy} + u_x + u_y = 2\)
   b. \(u_{xx} + u_{xy} + u_{yy} + u_x = 0\)
   c. \(3u_{xx} + 10u_{xy} + 3u_{yy} = 0\)
   d. \(u_{xx} + 2u_{xy} + 3u_{yy} + 4u_x + 5u_y + u = e^x\)
   e. \(2u_{xx} - 4u_{xy} + 2u_{yy} + 3u = 0\)
   f. \(u_{xx} + 5u_{xy} + 4u_{yy} + 7u_y = \sin x\)

3. Use any symbolic manipulator (e.g. MACSYMA or MATHEMATICA) to prove (6.2.19). This means that a transformation does NOT change the type of the PDE.
6.3 Canonical Forms

In this section we discuss canonical forms, which correspond to particularly simple choices of the coefficients of the second partial derivatives of the unknown. To obtain a canonical form, we have to transform the PDE which in turn will require the knowledge of characteristic curves. Three equivalent properties of characteristic curves, each can be used as a definition:

1. Initial data on a characteristic curve cannot be prescribed freely, but must satisfy a compatibility condition.
2. Discontinuities (of a certain nature) of a solution cannot occur except along characteristics.
3. Characteristics are the only possible “branch lines” of solutions, i.e. lines for which the same initial value problems may have several solutions.

We now consider specific choices for the functions $\xi, \eta$. This will be done in such a way that some of the coefficients $A^*, B^*, C^*$ in (6.2.21) become zero.

6.3.1 Hyperbolic

Note that $A^*, C^*$ are similar and can be written as

$$A\zeta_x^2 + B\zeta_x\zeta_y + C\zeta_y^2$$

in which $\zeta$ stands for either $\xi$ or $\eta$. Suppose we try to choose $\xi, \eta$ such that $A^* = C^* = 0$. This is of course possible only if the equation is hyperbolic. (Recall that $\Delta^* = (B^*)^2 - 4A^*C^*$ and for this choice $\Delta^* = (B^*)^2 > 0$. Since the type does not change under the transformation, we must have a hyperbolic PDE.) In order to annihilate $A^*$ and $C^*$ we have to find $\zeta$ such that

$$A\zeta_x^2 + B\zeta_x\zeta_y + C\zeta_y^2 = 0.$$  \hspace{1cm} (6.3.1.2)

Dividing by $\zeta_y^2$, the above equation becomes

$$A \left( \frac{\zeta_x}{\zeta_y} \right)^2 + B \left( \frac{\zeta_x}{\zeta_y} \right) + C = 0.$$ \hspace{1cm} (6.3.1.3)

Along the curve

$$\zeta(x, y) = \text{constant},$$ \hspace{1cm} (6.3.1.4)

we have

$$d\zeta = \zeta_x dx + \zeta_y dy = 0.$$ \hspace{1cm} (6.3.1.5)

Therefore,

$$\frac{\zeta_x}{\zeta_y} = -\frac{dy}{dx}.$$ \hspace{1cm} (6.3.1.6)

and equation (6.3.1.3) becomes

$$A \left( \frac{dy}{dx} \right)^2 - B \frac{dy}{dx} + C = 0.$$ \hspace{1cm} (6.3.1.7)
This is a quadratic equation for $\frac{dy}{dx}$ and its roots are

$$\frac{dy}{dx} = \frac{B \pm \sqrt{B^2 - 4AC}}{2A}.$$  

(6.3.1.8)

These equations are called characteristic equations and are ordinary differential equations for families of curves in $x, y$ plane along which $\zeta = constant$. The solutions are called characteristic curves. Notice that the discriminant is under the radical in (6.3.1.8) and since the problem is hyperbolic, $B^2 - 4AC > 0$, there are two distinct characteristic curves. We can choose one to be $\xi(x, y)$ and the other $\eta(x, y)$. Solving the ODEs (6.3.1.8), we get

$$\phi_1(x, y) = C_1,$$

(6.3.1.9)

$$\phi_2(x, y) = C_2.$$

(6.3.1.10)

Thus the transformation

$$\xi = \phi_1(x, y)$$

(6.3.1.11)

$$\eta = \phi_2(x, y)$$

(6.3.1.12)

will lead to $A^* = C^* = 0$ and the canonical form is

$$B^*u_{\xi\eta} = H^*$$

(6.3.1.13)

or after division by $B^*$

$$u_{\xi\eta} = \frac{H^*}{B^*}.$$

(6.3.1.14)

This is called the first canonical form of the hyperbolic equation.

Sometimes we find another canonical form for hyperbolic PDEs which is obtained by making a transformation

$$\alpha = \xi + \eta$$

(6.3.1.15)

$$\beta = \xi - \eta.$$  

(6.3.1.16)

Using (6.3.1.6)-(6.3.1.8) for this transformation one has

$$u_{\alpha\alpha} - u_{\beta\beta} = H^{**}(\alpha, \beta, u, u_{\alpha}, u_{\beta}).$$

(6.3.1.17)

This is called the second canonical form of the hyperbolic equation.

Example

$$y^2u_{xx} - x^2u_{yy} = 0 \quad \text{for} \quad x > 0, y > 0$$

(6.3.1.18)

$$A = y^2$$

$$B = 0$$

$$C = -x^2$$

$$\Delta = 0 - 4y^2(-x^2) = 4x^2y^2 > 0$$

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The equation is hyperbolic for all \( x, y \) of interest.

The characteristic equation

\[
\frac{dy}{dx} = 0 \pm \frac{\sqrt{4x^2y^2}}{2y^2} = \frac{\pm 2xy}{2y^2} = \pm \frac{x}{y}.
\]

These equations are separable ODEs and the solutions are

\[
\frac{1}{2}y^2 - \frac{1}{2}x^2 = c_1
\]

\[
\frac{1}{2}y^2 + \frac{1}{2}x^2 = c_2
\]

The first is a family of hyperbolas and the second is a family of circles (see figure 28).

![Figure 28: The families of characteristics for the hyperbolic example](image)

We take then the following transformation

\[
\xi = \frac{1}{2}y^2 - \frac{1}{2}x^2 \quad \text{(6.3.1.20)}
\]

\[
\eta = \frac{1}{2}y^2 + \frac{1}{2}x^2 \quad \text{(6.3.1.21)}
\]

Evaluate all derivatives of \( \xi, \eta \) necessary for (6.2.6) - (6.2.10)

\[
\xi_x = -x, \quad \xi_y = y, \quad \xi_{xx} = -1, \quad \xi_{xy} = 0, \quad \xi_{yy} = 1
\]

\[
\eta_x = x, \quad \eta_y = y, \quad \eta_{xx} = 1, \quad \eta_{xy} = 0, \quad \eta_{yy} = 1.
\]

Substituting all these in the expressions for \( B^*, D^* \) (you can check that \( A^* = C^* = 0 \))

\[
B^* = 2y^2(-x)x + 2(-x^2)y \cdot y = -2x^2y^2 - 2x^2y^2 = -4x^2y^2.
\]

\[
D^* = y^2(-1) + (-x^2) \cdot 1 = -x^2 - y^2.
\]

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Now solve (6.3.1.20) - (6.3.1.21) for \( x, y \)

\[
x^2 = \eta - \xi, \\
y^2 = \xi + \eta,
\]

and substitute in \( B^*, D^*, E^* \) we get

\[
-4(\eta - \xi)(\xi + \eta)u_{\xi\eta} + (-\eta + \xi - \xi - \eta)u_\xi + (\xi + \eta - \eta + \xi)u_\eta = 0 \\
4(\xi^2 - \eta^2)u_{\xi\eta} - 2\eta u_\xi + 2\xi u_\eta = 0 \\
u_{\xi\eta} = \frac{\eta}{2(\xi^2 - \eta^2)}u_\xi - \frac{\xi}{2(\xi^2 - \eta^2)}u_\eta
\]

This is the first canonical form of (6.3.1.18).

### 6.3.2 Parabolic

Since \( \Delta^* = 0, B^2 - 4AC = 0 \) and thus

\[
B = \pm 2\sqrt{A}\sqrt{C}.
\]

Clearly we cannot arrange for both \( A^* \) and \( C^* \) to be zero, since the characteristic equation (6.3.1.8) can have only one solution. That means that parabolic equations have only one characteristic curve. Suppose we choose the solution \( \phi_1(x, y) \) of (6.3.1.8)

\[
\frac{dy}{dx} = \frac{B}{2A}
\]

(6.3.2.2)

to define

\[
\xi = \phi_1(x, y).
\]

(6.3.2.3)

Therefore \( A^* = 0 \).

Using (6.3.2.1) we can show that

\[
0 = A^* = A\xi_x^2 + B\xi_x\xi_y + C\xi_y^2 \\
= A\xi_x^2 + 2\sqrt{A}\sqrt{C}\xi_x\xi_y + C\xi_y^2 \\
= \left(\sqrt{A}\xi_x + \sqrt{C}\xi_y\right)^2
\]

(6.3.2.4)

It is also easy to see that

\[
B^* = 2A\xi_x\eta_x + B(\xi_x\eta_y + \xi_y\eta_x) + 2C\xi_y\eta_y \\
= 2(\sqrt{A}\xi_x + \sqrt{C}\xi_y)(\sqrt{A}\eta_x + \sqrt{C}\eta_y) \\
= 0
\]
The last step is a result of (6.3.2.4). Therefore \( A^* = B^* = 0 \). To obtain the canonical form we must choose a function \( \eta(x, y) \). This can be taken judiciously as long as we ensure that the Jacobian is not zero.

The canonical form is then

\[ C^* u_{\eta\eta} = H^* \]

and after dividing by \( C^* \) (which cannot be zero) we have

\[ u_{\eta\eta} = \frac{H^*}{C^*}. \tag{6.3.2.5} \]

If we choose \( \eta = \phi_1(x, y) \) instead of (6.3.2.3), we will have \( C^* = 0 \). In this case \( B^* = 0 \) because the last factor \( \sqrt{A\eta_x} + \sqrt{C\eta_y} \) is zero. The canonical form in this case is

\[ u_{\xi\xi} = \frac{H^*}{A^*}. \tag{6.3.2.6} \]

**Example**

\[ x^2 u_{xx} - 2xy u_{xy} + y^2 u_{yy} = e^x \] \tag{6.3.2.7}

\[ A = x^2 \]
\[ B = -2xy \]
\[ C = y^2 \]
\[ \Delta = (-2xy)^2 - 4 \cdot x^2 \cdot y^2 = 4x^2y^2 - 4x^2y^2 = 0. \]

Thus the equation is parabolic for all \( x, y \). The characteristic equation (6.3.2.2) is

\[ \frac{dy}{dx} = \frac{-2xy}{2x^2} = -\frac{y}{x}. \tag{6.3.2.8} \]

Solve

\[ \frac{dy}{y} = -\frac{dx}{x} \]
\[ \ln y + \ln x = C \]

In figure 29 we sketch the family of characteristics for (6.3.2.7). Note that since the problem is parabolic, there is ONLY one family.

Therefore we can take \( \xi \) to be this family

\[ \xi = \ln y + \ln x \] \tag{6.3.2.9}

and \( \eta \) is arbitrary as long as \( J \neq 0 \). We take

\[ \eta = x. \] \tag{6.3.2.10}
Figure 29: The family of characteristics for the parabolic example

Computing the necessary derivatives of $\xi, \eta$ we have

$$\xi_x = \frac{1}{x}, \quad \xi_y = \frac{1}{y}, \quad \xi_{xx} = -\frac{1}{x^2}, \quad \xi_{xy} = 0, \quad \xi_{yy} = -\frac{1}{y^2}$$

$$\eta_x = 1, \quad \eta_y = \eta_{xx} = \eta_{xy} = \eta_{yy} = 0.$$  

Substituting these derivatives in the expressions for $C^*, D^*, E^*$ (recall that $A^* = B^* = 0$)

$$C^* = x^2 \cdot 1$$

$$D^* = x^2 \cdot \left( -\frac{1}{x^2} \right) - 2xy \cdot 0 + y^2 \left( -\frac{1}{y^2} \right) = -1 - 1 = -2$$

$$E^* = 0.$$  

The equation in the canonical form ($H^* = -D^*u_\xi + G^*$ in this case)

$$u_{\eta\eta} = \frac{2u_\xi + e^x}{x^2}$$

Now we must eliminate the old variables. Since $x = \eta$ we have

$$u_{\eta\eta} = \frac{2}{\eta^2}u_\xi + \frac{1}{\eta^2}e^\eta.$$  \hspace{1cm} (6.3.2.11)

Note that a different choice for $\eta$ will lead to a different right hand side in (6.3.2.11).

6.3.3 Elliptic

This is the case that $\Delta < 0$ and therefore there are NO real solutions to the characteristic equation (6.3.1.8). Suppose we solve for the complex valued functions $\xi$ and $\eta$. We now define

$$\alpha = \frac{\xi + \eta}{2}$$ \hspace{1cm} (6.3.3.1)

$$\beta = \frac{\xi - \eta}{2i}$$ \hspace{1cm} (6.3.3.2)
that is $\alpha$ and $\beta$ are the real and imaginary parts of $\xi$. Clearly $\eta$ is the complex conjugate of $\xi$ since the coefficients of the characteristic equation are real. If we use these functions $\alpha(x, y)$ and $\beta(x, y)$ we get an equation for which

$$B^{**} = 0, \quad A^{**} = C^{**}.$$ \hfill (6.3.3.3)

To show that (6.3.3.3) is correct, recall that our choice of $\xi, \eta$ led to $A^* = C^* = 0$. These are

$$A^* = (A\alpha^2_x + B\alpha_x \alpha_y + C\alpha^2_y) - (A\beta^2_x + B\beta_x \beta_y + C\beta^2_y) + i[2A\alpha_x \beta_x + B(\alpha_x \beta_y + \alpha_y \beta_x) + 2C\alpha_y \beta_y] = 0$$

$$C^* = (A\alpha^2_x + B\alpha_x \alpha_y + C\alpha^2_y) - (A\beta^2_x + B\beta_x \beta_y + C\beta^2_y) - i[2A\alpha_x \beta_x + B(\alpha_x \beta_y + \alpha_y \beta_x) + 2C\alpha_y \beta_y] = 0$$

Note the similarity of the terms in each bracket to those in (6.3.1.12)-(6.3.1.14)

$$A^* = (A^{**} - C^{**}) + iB^{**} = 0$$

$$C^* = (A^{**} - C^{**}) - iB^{**} = 0$$

where the double starred coefficients are given as in (6.3.1.12)-(6.3.1.14) except that $\alpha, \beta$ replace $\xi, \eta$ correspondingly. These last equations can be satisfied if and only if (6.3.3.3) is satisfied.

Therefore

$$A^{**}u_{\alpha \alpha} + A^{**}u_{\beta \beta} = H^{**}(\alpha, \beta, u, u_\alpha, u_\beta)$$

and the canonical form is

$$u_{\alpha \alpha} + u_{\beta \beta} = \frac{H^{**}}{A^{**}}.$$ \hfill (6.3.3.4)

**Example**

$$e^x u_{xx} + e^y u_{yy} = u \hfill (6.3.3.5)$$

$$A = e^x$$

$$B = 0$$

$$C = e^y$$

$$\Delta = 0^2 - 4e^x e^y < 0, \quad \text{for all } x, y$$

The characteristic equation

$$\frac{dy}{dx} = \frac{0 \pm \sqrt{-4e^x e^y}}{2e^x} = \frac{\pm 2i \sqrt{e^x e^y}}{2e^x} = \pm i \sqrt{\frac{e^y}{e^x}}$$

$$\frac{dy}{e^{y/2}} = \pm i \frac{dx}{e^{x/2}}. \hfill (6.3.3.6)$$

Therefore

$$\xi = -2e^{-y/2} - 2ie^{-x/2}$$

$$\eta = -2e^{-y/2} + 2ie^{-x/2}$$
The real and imaginary parts are:

\[ \alpha = -2e^{-y/2} \tag{6.3.3.6} \]
\[ \beta = -2e^{-x/2}. \tag{6.3.3.7} \]

Evaluate all necessary partial derivatives of \( \alpha, \beta \)

\[ \alpha_x = 0, \quad \alpha_y = e^{-y/2}, \quad \alpha_{xx} = 0, \quad \alpha_{xy} = 0, \quad \alpha_{yy} = -\frac{1}{2}e^{-y/2} \]
\[ \beta_x = e^{-x/2}, \quad \beta_y = 0, \quad \beta_{xx} = -\frac{1}{2}e^{-x/2}, \quad \beta_{xy} = 0, \quad \beta_{yy} = 0 \]

Now, instead of using both transformations, we recall that (6.3.1.12)-(6.3.1.18) are valid with \( \alpha, \beta \) instead of \( \xi, \eta \). Thus

\[ A^* = e^x \cdot 0 + 0 + e^y \left( e^{-y/2} \right)^2 = 1 \]
\[ B^* = 0 + 0 + 0 = 0 \quad \text{as can be expected} \]
\[ C^* = e^x \left( e^{-x/2} \right)^2 + 0 + 0 = 1 \quad \text{as can be expected} \]
\[ D^* = 0 + 0 + e^y \left( -\frac{1}{2}e^{-y/2} \right) = -\frac{1}{2}e^{y/2} \]
\[ E^* = e^x \left( -\frac{1}{2}e^{-x/2} \right) + 0 + 0 = -\frac{1}{2}e^{x/2} \]
\[ F^* = -1 \]
\[ H^* = -D^*u_\alpha - E^*u_\beta - F^*u = \frac{1}{2}e^{y/2}u_\alpha + \frac{1}{2}e^{x/2}u_\beta + u. \]

Thus

\[ u_{\alpha\alpha} + u_{\beta\beta} = \frac{1}{2}e^{y/2}u_\alpha + \frac{1}{2}e^{x/2}u_\beta + u. \]

Using (6.3.3.6)-(6.3.3.7) we have

\[ e^{x/2} = -\frac{2}{\beta} \]
\[ e^{y/2} = -\frac{2}{\alpha} \]

and therefore the canonical form is

\[ u_{\alpha\alpha} + u_{\beta\beta} = \frac{1}{\alpha}u_\alpha - \frac{1}{\beta}u_\beta + u. \tag{6.3.3.8} \]
Problems

1. Find the characteristic equation, characteristic curves and obtain a canonical form for each

   a. \( xu_{xx} + u_{yy} = x^2 \)
   b. \( u_{xx} + u_{xy} - xu_{yy} = 0 \quad (x \leq 0, \text{ all } y) \)
   c. \( x^2u_{xx} + 2xyu_{xy} + y^2u_{yy} + xyu_x + y^2u_y = 0 \)
   d. \( u_{xx} + xu_{yy} = 0 \)
   e. \( u_{xx} + y^2u_{yy} = y \)
   f. \( \sin^2 xu_{xx} + \sin 2xu_{xy} + \cos^2 xu_{yy} = x \)

2. Use Maple to plot the families of characteristic curves for each of the above.

3. Classify the following PDEs:
   (a) \( \frac{\partial^2 u}{\partial t^2} + \frac{\partial^2 u}{\partial x^2} + \frac{\partial u}{\partial x} = -e^{-kt} \)
   (b) \( \frac{\partial^2 u}{\partial x^2} - \frac{\partial^2 u}{\partial x \partial y} + \frac{\partial u}{\partial y} = 4 \)

4. Find the characteristics of each of the following PDEs:
   (a) \( \frac{\partial^2 u}{\partial x^2} + 3\frac{\partial^2 u}{\partial x \partial y} + 2\frac{\partial^2 u}{\partial y^2} = 0 \)
   (b) \( \frac{\partial^2 u}{\partial x^2} - 2\frac{\partial^2 u}{\partial x \partial y} + \frac{\partial^2 u}{\partial y^2} = 0 \)

5. Obtain the canonical form for the following elliptic PDEs:
   (a) \( \frac{\partial^2 u}{\partial x^2} + \frac{\partial^2 u}{\partial x \partial y} + \frac{\partial^2 u}{\partial y^2} = 0 \)
   (b) \( \frac{\partial^2 u}{\partial x^2} - 2\frac{\partial^2 u}{\partial x \partial y} + 5\frac{\partial^2 u}{\partial y^2} + \frac{\partial u}{\partial y} = 0 \)

6. Transform the following parabolic PDEs to canonical form:
   (a) \( \frac{\partial^2 u}{\partial x^2} - 6\frac{\partial^2 u}{\partial x \partial y} + 9\frac{\partial^2 u}{\partial y^2} + \frac{\partial u}{\partial x} - e^{xy} = 1 \)
   (b) \( \frac{\partial^2 u}{\partial x^2} + 2\frac{\partial^2 u}{\partial x \partial y} + \frac{\partial^2 u}{\partial y^2} + 7\frac{\partial u}{\partial x} - 8\frac{\partial u}{\partial y} = 0 \)
6.4 Equations with Constant Coefficients

In this case the discriminant is constant and thus the type of the equation is the same everywhere in the domain. The characteristic equation is easy to integrate.

6.4.1 Hyperbolic

The characteristic equation is

\[
\frac{dy}{dx} = \frac{B \pm \sqrt{\Delta}}{2A}.
\]

Thus

\[
\frac{dy}{dx} = \frac{B \pm \sqrt{\Delta}}{2A} dx
\]

and integration yields two families of straight lines

\[
\xi = y - \frac{B + \sqrt{\Delta}}{2A} x
\]

(6.4.1.2)

\[
\eta = y - \frac{B - \sqrt{\Delta}}{2A} x.
\]

(6.4.1.3)

Notice that if \( A = 0 \) then (6.4.1.1) is not valid. In this case we recall that (6.4.1.2) is

\[
B \zeta_x \zeta_y + C \zeta_y^2 = 0
\]

(6.4.1.4)

If we divide by \( \zeta_y^2 \) as before we get

\[
B \frac{\zeta_x}{\zeta_y} + C = 0
\]

(6.4.1.5)

which is only linear and thus we get only one characteristic family. To overcome this difficulty we divide (6.4.1.4) by \( \zeta_x^2 \) to get

\[
B \frac{\zeta_y}{\zeta_x} + C \left( \frac{\zeta_y}{\zeta_x} \right)^2 = 0
\]

(6.4.1.6)

which is quadratic. Now

\[
\frac{\zeta_y}{\zeta_x} = -\frac{dx}{dy}
\]

and so

\[
\frac{dx}{dy} = \frac{B \pm \sqrt{B^2 - 4 \cdot 0 \cdot C}}{2C} = \frac{B \pm B}{2C}
\]

or

\[
\frac{dx}{dy} = 0, \quad \frac{dx}{dy} = \frac{B}{C}.
\]

(6.4.1.7)

The transformation is then

\[
\xi = x,
\]

(6.4.1.8)

\[
\eta = x - \frac{B}{C} y.
\]

(6.4.1.9)

The canonical form is similar to (6.3.1.14).
6.4.2 Parabolic

The only solution of (6.4.1.1) is

$$\frac{dy}{dx} = \frac{B}{2A}.$$  

Thus

$$\xi = y - \frac{B}{2A}x. \quad (6.4.2.1)$$

Again $\eta$ is chosen judiciously but in such a way that the Jacobian of the transformation is not zero.

Can $A$ be zero in this case? In the parabolic case $A = 0$ implies $B = 0$ (since $\Delta = B^2 - 4 \cdot 0 \cdot C$ must be zero.) Therefore the original equation is

$$Cu_{yy} + Du_x + Eu_y + Fu = G$$

which is already in canonical form

$$u_{yy} = -\frac{D}{C}u_x - \frac{E}{C}u_y - \frac{F}{C}u + \frac{G}{C}. \quad (6.4.2.2)$$

6.4.3 Elliptic

Now we have complex conjugate functions $\xi, \eta$

$$\xi = y - \frac{B + i\sqrt{-\Delta}}{2A}x, \quad (6.4.3.1)$$

$$\eta = y - \frac{B - i\sqrt{-\Delta}}{2A}x. \quad (6.4.3.2)$$

Therefore

$$\alpha = y - \frac{B}{2A}x, \quad (6.4.3.3)$$

$$\beta = \frac{-\sqrt{-\Delta}}{2A}x. \quad (6.4.3.4)$$

(Note that $-\Delta > 0$ and the radical yields a real number.) The canonical form is similar to (6.3.3.4).

Example

$$u_{tt} - c^2u_{xx} = 0 \quad \text{(wave equation)} \quad (6.4.3.5)$$

$$A = 1$$

$$B = 0$$

$$C = -c^2$$

$$\Delta = 4c^2 > 0 \quad \text{(hyperbolic).}$$
The characteristic equation is
\[ \left(\frac{dx}{dt}\right)^2 - c^2 = 0 \]
and the transformation is
\[ \xi = x + ct, \quad \eta = x - ct. \]
\[ (6.4.3.6) \]
\[ (6.4.3.7) \]
The canonical form can be obtained as in the previous examples
\[ u_{\xi\eta} = 0. \]
\[ (6.4.3.8) \]
This is exactly the example from Chapter 1 for which we had
\[ u(\xi, \eta) = F(\xi) + G(\eta). \]
\[ (6.4.3.9) \]
The solution in terms of \(x, t\) is then (use (6.4.3.6)-(6.4.3.7))
\[ u(x, t) = F(x + ct) + G(x - ct). \]
\[ (6.4.3.10) \]
Problems

1. Find the characteristic equation, characteristic curves and obtain a canonical form for
   a. \( 4u_{xx} + 5u_{xy} + u_{yy} + u_x + u_y = 2 \)
   b. \( u_{xx} + u_{xy} + u_{yy} + u_x = 0 \)
   c. \( 3u_{xx} + 10u_{xy} + 3u_{yy} = x + 1 \)
   d. \( u_{xx} + 2u_{xy} + 3u_{yy} + 4u_x + 5u_y + u = e^x \)
   e. \( 2u_{xx} - 4u_{xy} + 2u_{yy} + 3u = 0 \)
   f. \( u_{xx} + 5u_{xy} + 4u_{yy} + 7u_y = \sin x \)

2. Use Maple to plot the families of characteristic curves for each of the above.
6.5 Linear Systems

In general, linear systems can be written in the form:

\[
\frac{\partial \mathbf{u}}{\partial t} + A \frac{\partial \mathbf{u}}{\partial x} + B \frac{\partial \mathbf{u}}{\partial y} + \mathbf{r} = 0
\]  

(6.5.1)

where \( \mathbf{u} \) is a vector valued function of \( t, x, y \).

The system is called hyperbolic at a point \((t, x)\) if the eigenvalues of \( A \) are all real and distinct. Similarly at a point \((t, y)\) if the eigenvalues of \( B \) are real and distinct.

**Example** The system of equations

\[
v_t = cw_x
\]

(6.5.2)

\[
w_t = cv_x
\]

(6.5.3)

can be written in matrix form as

\[
\frac{\partial \mathbf{u}}{\partial t} + A \frac{\partial \mathbf{u}}{\partial x} = 0
\]

(6.5.4)

where

\[
\mathbf{u} = \begin{pmatrix} v \\ w \end{pmatrix}
\]

(6.5.5)

and

\[
A = \begin{pmatrix} 0 & -c \\ -c & 0 \end{pmatrix}.
\]

(6.5.6)

The eigenvalues of \( A \) are given by

\[
\lambda^2 - c^2 = 0
\]

(6.5.7)

or \( \lambda = c, -c \). Therefore the system is hyperbolic, which we knew in advance since the system is the familiar wave equation.

**Example** The system of equations

\[
xu_x = vy_y
\]

(6.5.8)

\[
v_y = -xv_x
\]

(6.5.9)

can be written in matrix form

\[
\frac{\partial \mathbf{w}}{\partial x} + A \frac{\partial \mathbf{w}}{\partial y} = 0
\]

(6.5.10)

where

\[
\mathbf{w} = \begin{pmatrix} u \\ v \end{pmatrix}
\]

(6.5.11)

and

\[
A = \begin{pmatrix} 0 & -1 \\ 1 & 0 \end{pmatrix}.
\]

(6.5.12)

The eigenvalues of \( A \) are given by

\[
\lambda^2 + 1 = 0
\]

(6.5.13)

or \( \lambda = i, -i \). Therefore the system is elliptic. In fact, this system is the same as Laplace's equation.
Problems

1. Classify the behavior of the following system of PDEs in \((t, x)\) and \((t, y)\) space:

\[
\frac{\partial u}{\partial t} + \frac{\partial v}{\partial x} - \frac{\partial u}{\partial y} = 0 \\
\frac{\partial v}{\partial t} - \frac{\partial u}{\partial x} + \frac{\partial v}{\partial y} = 0
\]
6.6 General Solution

As we mentioned earlier, sometimes we can get the general solution of an equation by transforming it to a canonical form. We have seen one example (namely the wave equation) in the last section.

Example

\[ x^2u_{xx} + 2xyu_{xy} + y^2u_{yy} = 0. \]  
(6.6.1)

Show that the canonical form is

\[ u_{\eta\eta} = 0 \quad \text{for } y \neq 0 \]  
(6.6.2)

\[ u_{xx} = 0 \quad \text{for } y = 0. \]  
(6.6.3)

To solve (6.6.2) we integrate with respect to \( \eta \) twice (\( \xi \) is fixed) to get

\[ u(\xi, \eta) = \eta F(\xi) + G(\xi). \]  
(6.6.4)

Since the transformation to canonical form is

\[ \xi = \frac{y}{x} \quad \eta = y \quad \text{(arbitrary choice for } \eta) \]  
(6.6.5)

then

\[ u(x, y) = y F\left(\frac{y}{x}\right) + G\left(\frac{y}{x}\right). \]  
(6.6.6)

Example

Obtain the general solution for

\[ 4u_{xx} + 5u_{xy} + u_{yy} + u_x + u_y = 2. \]  
(6.6.7)

(This example is taken from Myint-U and Debnath (19).) There is a mistake in their solution which we have corrected here. The transformation

\[ \xi = y - x, \]  

\[ \eta = y - \frac{x}{4}, \]  
(6.6.8)

leads to the canonical form

\[ u_{\xi\eta} = \frac{1}{3}u_\eta - \frac{8}{9}. \]  
(6.6.9)

Let \( v = u_\eta \) then (6.6.9) can be written as

\[ v_\xi = \frac{1}{3}v - \frac{8}{9}. \]  
(6.6.10)

which is a first order linear ODE (assuming \( \eta \) is fixed.) Therefore

\[ v = \frac{8}{3} + e^{\xi/3}\phi(\eta), \]  
(6.6.11)

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Now integrating with respect to \( \eta \) yields

\[
  u(\xi, \eta) = \frac{8}{3} \eta + G(\eta)e^{\xi/3} + F(\xi).
\]  

(6.6.12)

In terms of \( x, y \) the solution is

\[
  u(x, y) = \frac{8}{3} (y - \frac{x}{4}) + G \left( y - \frac{x}{4} \right) e^{(y-x)/3} + F(y - x).
\]  

(6.6.13)
Problems

1. Determine the general solution of
   a. \( u_{xx} - \frac{1}{3} u_{yy} = 0 \quad c = \text{constant} \)
   b. \( u_{xx} - 3u_{xy} + 2u_{yy} = 0 \)
   c. \( u_{xx} + u_{xy} = 0 \)
   d. \( u_{xx} + 10u_{xy} + 9u_{yy} = y \)

2. Transform the following equations to
   \( U_{\xi\eta} = cU \)
   by introducing the new variables
   \( U = ue^{-\alpha \xi - \beta \eta} \)
   where \( \alpha, \beta \) to be determined
   a. \( u_{xx} - u_{yy} + 3u_{x} - 2u_{y} + u = 0 \)
   b. \( 3u_{xx} + 7u_{xy} + 2u_{yy} + u_{yy} + u = 0 \)
   (Hint: First obtain a canonical form)

3. Show that
   \( u_{xx} = au_{t} + bu_{x} - \frac{b^2}{4}u + d \)
   is parabolic for \( a, b, d \) constants. Show that the substitution
   \( u(x,t) = v(x,t)e^{b \xi} \)
   transforms the equation to
   \( v_{xx} = av_{t} + de^{-\frac{b}{2}x} \)
Summary

Equation

\[ Au_{xx} + Bu_{xy} + Cu_{yy} = -Du_x - Eu_y - Fu + G = H(x, y, u, u_x, u_y) \]

Discriminant

\[ \Delta(x_0, y_0) = B^2(x_0, y_0) - 4A(x_0, y_0)C(x_0, y_0) \]

Class

\[ \Delta > 0 \quad \text{hyperbolic at the point } (x_0, y_0) \]
\[ \Delta = 0 \quad \text{parabolic at the point } (x_0, y_0) \]
\[ \Delta < 0 \quad \text{elliptic at the point } (x_0, y_0) \]

Transformed Equation

\[ A^*u_{\xi\xi} + B^*u_{\xi\eta} + C^*u_{\eta\eta} = -D^*u_\xi - E^*u_\eta - F^*u + G^* = H^*(\xi, \eta, u, u_\xi, u_\eta) \]

where

\[ A^* = A\xi_x^2 + B\xi_x\xi_y + C\xi_y^2 \]
\[ B^* = 2A\xi_x\eta_x + B(\xi_x\eta_y + \xi_y\eta_x) + 2C\xi_y\eta_y \]
\[ C^* = A\eta_x^2 + B\eta_x\eta_y + C\eta_y^2 \]
\[ D^* = A\xi_{xx} + B\xi_{xy} + C\xi_{yy} + D\xi_x + E\xi_y \]
\[ E^* = A\eta_{xx} + B\eta_{xy} + C\eta_{yy} + D\eta_x + E\eta_y \]
\[ F^* = F \]
\[ G^* = G \]
\[ H^* = -D^*u_\xi - E^*u_\eta - F^*u + G^* \]

\[ \frac{dy}{dx} = \frac{B \pm \sqrt{\Delta}}{2A} \quad \text{characteristic equation} \]
\[ u_{\xi\eta} = \frac{H^*}{B^*} \quad \text{first canonical form for hyperbolic} \]
\[ u_{\alpha\alpha} - u_{\beta\beta} = \frac{H^{**}}{B^{**}} \quad \alpha = \xi + \eta, \beta = \xi - \eta \quad \text{second canonical form for hyperbolic} \]
\[ u_{\xi\xi} = \frac{H^*}{A^*} \quad \text{a canonical form for parabolic} \]
\[ u_{\eta\eta} = \frac{H^*}{C^*} \quad \text{a canonical form for parabolic} \]
\[ u_{\alpha\alpha} + u_{\beta\beta} = \frac{H^{**}}{A^{**}} \quad \alpha = (\xi + \eta)/2, \beta = (\xi - \eta)/2i \quad \text{a canonical form for elliptic} \]
7 Method of Characteristics

In this chapter we will discuss a method to solve first order linear and quasilinear PDEs. This method is based on finding the characteristic curve of the PDE. We will also show how to generalize this method for a second order constant coefficients wave equation. The method of characteristics can be used only for hyperbolic problems which possess the right number of characteristic families. Recall that for second order parabolic problems we have only one family of characteristics and for elliptic PDEs no real characteristic curves exist.

7.1 Advection Equation (first order wave equation)

The one dimensional wave equation

\[
\frac{\partial^2 u}{\partial t^2} - c^2 \frac{\partial^2 u}{\partial x^2} = 0 \tag{7.1.1}
\]

can be rewritten as either of the following

\[
\left( \frac{\partial}{\partial t} + c \frac{\partial}{\partial x} \right) \left( \frac{\partial}{\partial t} - c \frac{\partial}{\partial x} \right) u = 0 \tag{7.1.2}
\]

\[
\left( \frac{\partial}{\partial t} - c \frac{\partial}{\partial x} \right) \left( \frac{\partial}{\partial t} + c \frac{\partial}{\partial x} \right) u = 0 \tag{7.1.3}
\]

since the mixed derivative terms cancel. If we let

\[
v = \frac{\partial u}{\partial t} - c \frac{\partial u}{\partial x} \tag{7.1.4}
\]

then (7.1.2) becomes

\[
\frac{\partial v}{\partial t} + c \frac{\partial v}{\partial x} = 0. \tag{7.1.5}
\]

Similarly (7.1.3) yields

\[
\frac{\partial w}{\partial t} - c \frac{\partial w}{\partial x} = 0 \tag{7.1.6}
\]

if

\[
w = \frac{\partial u}{\partial t} + c \frac{\partial u}{\partial x} \tag{7.1.7}
\]

The only difference between (7.1.5) and (7.1.6) is the sign of the second term. We now show how to solve (7.1.5) which is called the first order wave equation or advection equation (in Meteorology).

Remark: Although (7.1.4)-(7.1.5) or (7.1.6)-(7.1.7) can be used to solve the one dimensional second order wave equation (7.1.1), we will see in section 7.3 another way to solve (7.1.1) based on the results of Chapter 6.
To solve (7.1.5) we note that if we consider an observer moving on a curve \( x(t) \) then by the chain rule we get
\[
\frac{dv(x(t), t)}{dt} = \frac{\partial v}{\partial t} + \frac{\partial v}{\partial x} \frac{dx}{dt},
\]  
(7.1.8)

If the observer is moving at a rate \( \frac{dx}{dt} = c \), then by comparing (7.1.8) and (7.1.5) we find
\[
\frac{dv}{dt} = 0.
\]  
(7.1.9)

Therefore (7.1.5) can be replaced by a set of two ODEs
\[
\frac{dx}{dt} = c,
\]  
(7.1.10)

\[
\frac{dv}{dt} = 0.
\]  
(7.1.11)

These 2 ODEs are easy to solve. Integration of (7.1.10) yields
\[
x(t) = x(0) + ct
\]  
(7.1.12)

and the other one has a solution
\[v = \text{constant along the curve given in (7.1.12)}.\]

The curve (7.1.12) is a straight line. In fact, we have a family of parallel straight lines, called characteristics, see figure 30.

![Figure 30: Characteristics \( t = \frac{1}{c}x - \frac{1}{c}x(0) \)](image)

In order to obtain the general solution of the one dimensional equation (7.1.5) subject to the initial value
\[
v(x(0), 0) = f(x(0)),
\]  
(7.1.13)
we note that

\[ v = \text{constant along } x(t) = x(0) + ct \]

but that constant is \( f(x(0)) \) from (7.1.13). Since \( x(0) = x(t) - ct \), the general solution is then

\[ v(x,t) = f(x(t) - ct). \tag{7.1.14} \]

Let us show that (7.1.14) is the solution. First if we take \( t = 0 \), then (7.1.14) reduces to

\[ v(x,0) = f(x(0) - c \cdot 0) = f(x(0)). \]

To check the PDE we require the first partial derivatives of \( v \). Notice that \( f \) is a function of only one variable, i.e. of \( x - ct \). Therefore

\[
\frac{\partial v}{\partial t} = \frac{df(x - ct)}{dt} = \frac{df}{d(x - ct)} \frac{d(x - ct)}{dt} = -c \frac{df}{d(x - ct)}
\]

\[
\frac{\partial v}{\partial x} = \frac{df(x - ct)}{dx} = \frac{df}{d(x - ct)} \frac{d(x - ct)}{dx} = \frac{df}{d(x - ct)}.
\]

Substituting these two derivatives in (7.1.5) we see that the equation is satisfied.

Example 1

\[ \frac{\partial v}{\partial t} + 3 \frac{\partial v}{\partial x} = 0 \tag{7.1.15} \]

\[ v(x,0) = \begin{cases} \frac{1}{2} x & 0 < x < 1 \\ 0 & \text{otherwise.} \end{cases} \tag{7.1.16} \]

The two ODEs are

\[ \frac{dx}{dt} = 3, \tag{7.1.17} \]

\[ \frac{dv}{dt} = 0. \tag{7.1.18} \]

The solution of (7.1.17) is

\[ x(t) = x(0) + 3t \tag{7.1.19} \]

and the solution of (7.1.18) is

\[ v(x(t),t) = v(x(0),0) = \text{constant}. \tag{7.1.20} \]

Using (7.1.16) the solution is then

\[ v(x(t),t) = \begin{cases} \frac{1}{2} x(0) & 0 < x(0) < 1 \\ 0 & \text{otherwise.} \end{cases} \]

Substituting \( x(0) \) from (7.1.19) we have

\[ v(x,t) = \begin{cases} \frac{1}{2} (x - 3t) & 0 < x - 3t < 1 \\ 0 & \text{otherwise.} \end{cases} \tag{7.1.21} \]

The interpretation of (7.1.20) is as follows. Given a point \( x \) at time \( t \), find the characteristic through this point. Move on the characteristic to find the point \( x(0) \) and then use the initial value at that \( x(0) \) as the solution at \( (x,t) \). (Recall that \( v \) is constant along a characteristic.)
Let’s sketch the characteristics through the points $x = 0, 1$ (see (7.1.19) and Figure 31.)

![Figure 31: 2 characteristics for $x(0) = 0$ and $x(0) = 1$](image)

The initial solution is sketched in the next figure (32)

![Figure 32: Solution at time $t = 0$](image)

This shape is constant along a characteristic, and moving at the rate of 3 units. For example, the point $x = \frac{1}{2}$ at time $t = 0$ will be at $x = 3.5$ at time $t = 1$. The solution $v$ will be exactly the same at both points, namely $v = \frac{1}{3}$. The solution at several times is given in figure 33.

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The system of ODEs is
\[
\frac{du}{dt} = e^{2x} \quad (7.1.22)
\]
\[
u(x,0) = f(x). \quad (7.1.23)
\]

Solve (7.1.25) to get the characteristic curve
\[
x(t) = x(0) - 2t. \quad (7.1.26)
\]

Substituting the characteristic equation in (7.1.24) yields
\[
\frac{du}{dt} = e^{2(x(0)-2t)}. \quad (7.1.27)
\]

At \( t = 0 \)
\[
f(x(0)) = u(x(0), 0) = K - \frac{1}{4}e^{2x(0)}
\]
and therefore
\[
K = f(x(0)) + \frac{1}{4}e^{2x(0)}. \quad (7.1.28)
\]
Substitute $K$ in (7.1.27) we have

$$u(x, t) = f(x(0)) + \frac{1}{4}e^{2x(0)} - \frac{1}{4}e^{2x(0)-4t}.$$ 

Now substitute for $x(0)$ from (7.1.26) we get

$$u(x, t) = f(x + 2t) + \frac{1}{4}e^{2(x+2t)} - \frac{1}{4}e^{2x},$$

or

$$u(x, t) = f(x + 2t) + \frac{1}{4}e^{2x} \left( e^{4t} - 1 \right).$$

(7.1.29)
Problems

1. Solve
   \[ \frac{\partial w}{\partial t} - 3 \frac{\partial w}{\partial x} = 0 \]
   subject to
   \[ w(x, 0) = \sin x \]

2. Solve using the method of characteristics
   a. \[ \frac{\partial u}{\partial t} + e \frac{\partial u}{\partial x} = e^{2x} \]  subject to \( u(x, 0) = f(x) \)
   b. \[ \frac{\partial u}{\partial t} + x \frac{\partial u}{\partial x} = 1 \]  subject to \( u(x, 0) = f(x) \)
   c. \[ \frac{\partial u}{\partial t} + 3t \frac{\partial u}{\partial x} = u \]  subject to \( u(x, 0) = f(x) \)
   d. \[ \frac{\partial u}{\partial t} - 2 \frac{\partial u}{\partial x} = e^{2x} \]  subject to \( u(x, 0) = \cos x \)
   e. \[ \frac{\partial u}{\partial t} - t^2 \frac{\partial u}{\partial x} = -u \]  subject to \( u(x, 0) = 3e^x \)

3. Show that the characteristics of
   \[ \frac{\partial u}{\partial t} + 2u \frac{\partial u}{\partial x} = 0 \]
   \[ u(x, 0) = f(x) \]
   are straight lines.

4. Consider the problem
   \[ \frac{\partial u}{\partial t} + 2u \frac{\partial u}{\partial x} = 0 \]
   \[ u(x, 0) = f(x) = \begin{cases} 
   1 & x < 0 \\
   1 + \frac{x}{L} & 0 < x < L \\
   2 & L < x 
\end{cases} \]
   a. Determine equations for the characteristics
   b. Determine the solution \( u(x, t) \)
   c. Sketch the characteristic curves.
   d. Sketch the solution \( u(x, t) \) for fixed \( t \).

5. Solve the initial value problem for the damped unidirectional wave equation
   \[ v_t + cv_x + \lambda v = 0 \quad v(x, 0) = F(x) \]
where $\lambda > 0$ and $F(x)$ is given.

6. (a) Solve the initial value problem for the inhomogeneous equation

$$v_t + cv_x = f(x,t) \quad v(x,0) = F(x)$$

where $f(x,t)$ and $F(x)$ are specified functions.

(b) Solve this problem when $f(x,t) = xt$ and $F(x) = \sin x$.

7. Solve the “signaling” problem

$$v_t + cv_x = 0 \quad v(0,t) = G(t) \quad -\infty < t < \infty$$

in the region $x > 0$.

8. Solve the initial value problem

$$v_t + ev_x = 0 \quad v(x,0) = x$$

9. Show that the initial value problem

$$u_t + u_x = x \quad u(x,x) = 1$$

has no solution. Give a reason for the problem.
7.2 Quasilinear Equations

The method of characteristics is the only method applicable for quasilinear PDEs. All other methods such as separation of variables, Green’s functions, Fourier or Laplace transforms cannot be extended to quasilinear problems.

In this section, we describe the use of the method of characteristics for the solution of
\[
\frac{\partial u}{\partial t} + c(u, x, t) \frac{\partial u}{\partial x} = S(u, x, t) \tag{7.2.1}
\]
\[u(x, 0) = f(x). \tag{7.2.2}\]

Such problems have applications in gas dynamics or traffic flow. Equation (7.2.1) can be rewritten as a system of ODEs
\[
\frac{dx}{dt} = c(u, x, t) \tag{7.2.3}
\]
\[
\frac{du}{dt} = S(u, x, t). \tag{7.2.4}
\]

The first equation is the characteristic equation. The solution of this system can be very complicated since \(u\) appears nonlinearly in both. To find the characteristic curve one must know the solution. Geometrically, the characteristic curve has a slope depending on the solution \(u\) at that point, see figure 34.

\[\text{Figure 34: } u(x_0, 0) = f(x_0)\]

The slope of the characteristic curve at \(x_0\) is
\[
\frac{1}{c(u(x_0), x_0, 0)} = \frac{1}{c(f(x_0), x_0, 0)}. \tag{7.2.5}\]

Now we can compute the next point on the curve, by using this slope (assuming a slow change of rate and that the point is close to the previous one). Once we have the point, we can then solve for \(u\) at that point.
7.2.1 The Case $S = 0$, $c = c(u)$

The quasilinear equation

\[ u_t + c(u)u_x = 0 \]  \hspace{1cm} (7.2.1.1)

subject to the initial condition

\[ u(x, 0) = f(x) \]  \hspace{1cm} (7.2.1.2)

is equivalent to

\[ \frac{dx}{dt} = c(u) \]  \hspace{1cm} (7.2.1.3)

\[ x(0) = \xi \]  \hspace{1cm} (7.2.1.4)

\[ \frac{du}{dt} = 0 \]  \hspace{1cm} (7.2.1.5)

\[ u(\xi, 0) = f(\xi). \]  \hspace{1cm} (7.2.1.6)

Thus

\[ u(x, t) = u(\xi, 0) = f(\xi) \]  \hspace{1cm} (7.2.1.7)

\[ \frac{dx}{dt} = c(f(\xi)) \]

\[ x = tc(f(\xi)) + \xi. \]  \hspace{1cm} (7.2.1.8)

Solve (7.2.1.8) for $\xi$ and substitute in (7.2.1.7) to get the solution.

To check our solution, we compute the first partial derivatives of $u$

\[ \frac{\partial u}{\partial t} = \frac{du}{d\xi} \frac{d\xi}{dt} \]

\[ \frac{\partial u}{\partial x} = \frac{du}{d\xi} \frac{d\xi}{dx}. \]

Differentiating (7.2.1.8) with respect to $x$ and $t$ we have

\[ 1 = tc'(f(\xi))f'(\xi)\xi_x + \xi_x \]

\[ 0 = c(f(\xi)) + tc'(f(\xi))f'(\xi)\xi_t + \xi_t \]

correspondingly.

Thus when recalling that $\frac{du}{d\xi} = f'(\xi)$

\[ u_t = -\frac{c(f(\xi))}{1 + tc'(f(\xi))f'(\xi)} f'(\xi) \]  \hspace{1cm} (7.2.1.9)

\[ u_x = \frac{1}{1 + tc'(f(\xi))f'(\xi)} f'(\xi). \]  \hspace{1cm} (7.2.1.10)

Substituting these expressions in (7.2.1.1) results in an identity. The initial condition (7.2.1.2) is exactly (7.2.1.7).
Example 3

\[
\frac{\partial u}{\partial t} + u \frac{\partial u}{\partial x} = 0 \tag{7.2.1.11}
\]

\[
u(x, 0) = 3x. \tag{7.2.1.12}
\]

The equivalent system of ODEs is

\[
\frac{du}{dt} = 0 \tag{7.2.1.13}
\]

\[
\frac{dx}{dt} = u. \tag{7.2.1.14}
\]

Solving the first one yields

\[
u(x, t) = u(x(0), 0) = 3x(0). \tag{7.2.1.15}
\]

Substituting this solution in (7.2.1.14)

\[
\frac{dx}{dt} = 3x(0)
\]

which has a solution

\[
x = 3x(0)t + x(0). \tag{7.2.1.16}
\]

Solve (7.2.1.16) for \(x(0)\) and substitute in (7.2.1.15) gives

\[
u(x, t) = \frac{3x}{3t + 1}. \tag{7.2.1.17}
\]
Problems

1. Solve the following
   a. \( \frac{\partial u}{\partial t} = 0 \) subject to \( u(x, 0) = g(x) \)
   b. \( \frac{\partial u}{\partial t} = -3xu \) subject to \( u(x, 0) = g(x) \)

2. Solve
   \( \frac{\partial u}{\partial t} = u \)
   subject to
   \( u(x, t) = 1 + \cos x \) along \( x + 2t = 0 \)

3. Let
   \( \frac{\partial u}{\partial t} + c \frac{\partial u}{\partial x} = 0 \quad c = \text{constant} \)
   a. Solve the equation subject to \( u(x, 0) = \sin x \)
   b. If \( c > 0 \), determine \( u(x, t) \) for \( x > 0 \) and \( t > 0 \) where
      \[
      u(x, 0) = f(x) \quad \text{for} \ x > 0 \\
      u(0, t) = g(t) \quad \text{for} \ t > 0
      \]

4. Solve the following linear equations subject to \( u(x, 0) = f(x) \)
   a. \( \frac{\partial u}{\partial t} + c \frac{\partial u}{\partial x} = e^{-3x} \)
   b. \( \frac{\partial u}{\partial t} + t \frac{\partial u}{\partial x} = 5 \)
   c. \( \frac{\partial u}{\partial t} - t^2 \frac{\partial u}{\partial x} = -u \)
   d. \( \frac{\partial u}{\partial t} + x \frac{\partial u}{\partial x} = t \)
   e. \( \frac{\partial u}{\partial t} + x \frac{\partial u}{\partial x} = x \)

5. Determine the parametric representation of the solution satisfying \( u(x, 0) = f(x) \),
   a. \( \frac{\partial u}{\partial t} - u^2 \frac{\partial u}{\partial x} = 3u \)
b. \[ \frac{\partial u}{\partial t} + t^2 u \frac{\partial u}{\partial x} = -u \]

6. Solve \[ \frac{\partial u}{\partial t} + t^2 u \frac{\partial u}{\partial x} = 5 \]
subject to \[ u(x, 0) = x. \]

7. Using implicit differentiation, verify that \( u(x, t) = f(x - tu) \) is a solution of \[ u_t + uu_x = 0 \]

8. Consider the damped quasilinear wave equation
\[ u_t + uu_x + cu = 0 \]
where \( c \) is a positive constant.

(a) Using the method of characteristics, construct a solution of the initial value problem with \( u(x, 0) = f(x) \), in implicit form. Discuss the wave motion and the effect of the damping.

(b) Determine the breaking time of the solution by finding the envelope of the characteristic curves and by using implicit differentiation. With \( \tau \) as the parameter on the initial line, show that unless \( f'(\tau) < -c \), no breaking occurs.

9. Consider the one-dimensional form of Euler’s equations for isentropic flow and assume that the pressure \( p \) is a constant. The equations reduce to
\[ \rho_t + \rho u_x + u \rho_x = 0 \quad u_t + uu_x = 0 \]

Let \( u(x, 0) = f(x) \) and \( \rho(x, 0) = g(x) \). By first solving the equation for \( u \) and then the equation for \( \rho \), obtain the implicit solution
\[ u = f(x - ut) \quad \rho = \frac{g(x - ut)}{1 + tf'(x - ut)} \]
7.2.2 Graphical Solution

Graphically, one can obtain the solution as follows:

![Graphical Solution](image)

Figure 35: Graphical solution

Suppose the initial solution \( u(x, 0) \) is sketched as in figure 35. We know that each \( u(x_0) \) stays constant moving at its own constant speed \( c(u(x_0)) \). At time \( t \), it moved from \( x_0 \) to \( x_0 + tc(f(x_0)) \) (horizontal arrow). This process should be carried out to enough points on the initial curve to get the solution at time \( t \). Note that the lengths of the arrows are different and depend on \( c \).

7.2.3 Numerical Solution

Here we discuss a general linear first order hyperbolic

\[
a(x,t)u_x + b(x,t)u_t = c(x,t)u + d(x,t). \tag{7.2.3.1}
\]

Note that since \( b(x,t) \) may vanish, we cannot in general divide the equation by \( b(x,t) \) to get it in the same form as we had before. Thus we parametrize \( x \) and \( t \) in terms of a parameter \( s \), and instead of taking the curve \( x(t) \), we write it as \( x(s), t(s) \).

The characteristic equation is now a system

\[
\frac{dx}{ds} = a(x(s), t(s)) \tag{7.2.3.2}
\]

\[
x(0) = \xi \tag{7.2.3.3}
\]

\[
\frac{dt}{ds} = b(x(s), t(s)) \tag{7.2.3.4}
\]

\[
t(0) = 0 \tag{7.2.3.5}
\]

\[
\frac{du}{ds} = c(x(s), t(s))u(x(s), t(s)) + d(x(s), t(s)) \tag{7.2.3.6}
\]

\[
u(\xi, 0) = f(\xi) \tag{7.2.3.7}
\]
This system of ODEs need to be solved numerically. One possibility is the use of Runge-Kutta method, see Lab 4. This idea can also be used for quasilinear hyperbolic PDEs.

7.2.4 Fan-like Characteristics

Since the slope of the characteristic, $\frac{1}{c}$, depends in general on the solution, one may have characteristic curves intersecting or curves that fan-out. We demonstrate this by the following example.

**Example 4**

$$u_t + uu_x = 0$$  \hspace{1cm} (7.2.4.1)

$$u(x,0) = \begin{cases} 
1 & \text{for } x < 0 \\
2 & \text{for } x > 0.
\end{cases} \hspace{1cm} (7.2.4.2)$$

The system of ODEs is

$$\frac{dx}{dt} = u,$$  \hspace{1cm} (7.2.4.3)

$$\frac{du}{dt} = 0.$$  \hspace{1cm} (7.2.4.4)

The second ODE satisfies

$$u(x,t) = u(x(0),0)$$  \hspace{1cm} (7.2.4.5)

and thus the characteristics are

$$x = u(x(0),0)t + x(0)$$  \hspace{1cm} (7.2.4.6)

or

$$x(t) = \begin{cases} 
t + x(0) & \text{if } x(0) < 0 \\
2t + x(0) & \text{if } x(0) > 0.
\end{cases} \hspace{1cm} (7.2.4.7)$$

Let’s sketch those characteristics (Figure 36). If we start with a negative $x(0)$ we obtain a straight line with slope 1. If $x(0)$ is positive, the slope is $\frac{1}{2}$.

Since $u(x(0),0)$ is discontinuous at $x(0) = 0$, we find there are no characteristics through $t = 0$, $x(0) = 0$. In fact, we imagine that there are infinitely many characteristics with all possible slopes from $\frac{1}{2}$ to 1. Since the characteristics fan out from $x = t$ to $x = 2t$ we call these fan-like characteristics. The solution for $t < x < 2t$ will be given by (7.2.4.6) with $x(0) = 0$, i.e.

$$x = ut$$

or

$$u = \frac{x}{t} \hspace{1cm} \text{for } t < x < 2t.$$  \hspace{1cm} (7.2.4.8)
Figure 36: The characteristics for Example 4

To summarize the solution is then

\[
  u = \begin{cases} 
  1 & x(0) = x - t < 0 \\
  2 & x(0) = x - 2t > 0 \\
  \frac{x}{t} & t < x < 2t 
  \end{cases} \quad (7.2.4.9)
\]

The sketch of the solution is given in figure 37.

Figure 37: The solution of Example 4

7.2.5 Shock Waves

If the initial solution is discontinuous, but the value to the left is larger than that to the right, one will see intersecting characteristics.

Example 5

\[
  u_t + uu_x = 0 \quad (7.2.5.1)
\]
\[ u(x, 0) = \begin{cases} 2 & x < 1 \\ 1 & x > 1 \end{cases} \]  

(7.2.5.2)

The solution is as in the previous example, i.e.

\[ x(t) = u(x(0), 0)t + x(0) \]  

(7.2.5.3)

\[ x(t) = \begin{cases} 2t + x(0) & \text{if } x(0) < 1 \\ t + x(0) & \text{if } x(0) > 1 \end{cases} \]  

(7.2.5.4)

The sketch of the characteristics is given in figure 38.

Figure 38: Intersecting characteristics

Since there are two characteristics through a point, one cannot tell on which characteristic to move back to \( t = 0 \) to obtain the solution. In other words, at points of intersection the solution \( u \) is multi-valued. This situation happens whenever the speed along the characteristic on the left is larger than the one along the characteristic on the right, and thus catching up with it. We say in this case to have a shock wave. Let \( x_1(0) < x_2(0) \) be two points at \( t = 0 \), then

\[
x_1(t) = c\left(f(x_1(0))\right)t + x_1(0) \\
x_2(t) = c\left(f(x_2(0))\right)t + x_2(0).
\]  

(7.2.5.5)

If \( c(f(x_1(0))) > c(f(x_2(0))) \) then the characteristics emanating from \( x_1(0) \), \( x_2(0) \) will intersect. Suppose the points are close, i.e. \( x_2(0) = x_1(0) + \Delta x \), then to find the point of intersection we equate \( x_1(t) = x_2(t) \). Solving this for \( t \) yields

\[
t = -\frac{\Delta x}{-c(f(x_1(0))) + c(f(x_1(0) + \Delta x))}. 
\]  

(7.2.5.6)

If we let \( \Delta x \) tend to zero, the denominator (after dividing through by \( \Delta x \)) tends to the derivative of \( c \), i.e.

\[
t = -\frac{1}{dc(f(x_1(0)))}. 
\]  

(7.2.5.7)
Since \( t \) must be positive at intersection (we measure time from zero), this means that

\[
\frac{dc}{dx_1} < 0. \tag{7.2.5.8}
\]

So if the characteristic velocity \( c \) is locally decreasing then the characteristics will intersect. This is more general than the case in the last example where we have a discontinuity in the initial solution. One can have a continuous initial solution \( u(x, 0) \) and still get a shock wave. Note that (7.2.5.7) implies that

\[
1 + t \frac{dc(f)}{dx} = 0
\]

which is exactly the denominator in the first partial derivative of \( u \) (see (7.2.1.9)-(7.2.1.10)).

Example 6

\[
\begin{align*}
  &u_t + uu_x = 0 \quad \tag{7.2.5.9} \\
  &u(x, 0) = -x. \quad \tag{7.2.5.10}
\end{align*}
\]

The solution of the ODEs

\[
\begin{align*}
  &\frac{du}{dt} = 0, \quad \tag{7.2.5.11} \\
  &\frac{dx}{dt} = u,
\end{align*}
\]

is

\[
\begin{align*}
  &u(x, t) = u(x(0), 0) = -x(0), \quad \tag{7.2.5.12} \\
  &x(t) = -x(0)t + x(0) = x(0)(1 - t). \quad \tag{7.2.5.13}
\end{align*}
\]

Solving for \( x(0) \) and substituting in (7.2.5.12) yields

\[
u(x, t) = -\frac{x(t)}{1 - t}. \tag{7.2.5.14}
\]

This solution is undefined at \( t = 1 \). If we use (7.2.5.7) we get exactly the same value for \( t \), since

\[
\begin{align*}
  &f(x_0) = -x_0 \quad \text{(from (7.2.5.10)} \\
  &c(f(x_0)) = u(x_0) = -x_0 \quad \text{(from (7.2.5.9)} \\
  &\frac{dc}{dx_0} = -1 \\
  &t = -\frac{1}{-1} = 1.
\end{align*}
\]

In the next figure we sketch the characteristics given by (7.2.5.13). It is clear that all characteristics intersect at \( t = 1 \). The shock wave starts at \( t = 1 \). If the initial solution is discontinuous then the shock wave is formed immediately.
How do we find the shock position $x_s(t)$ and its speed? To this end, we rewrite the original equation in conservation law form, i.e.

$$u_t + \frac{\partial}{\partial x} q(u) = 0$$

or

$$\int_{\alpha}^{\beta} u_t dx = \frac{d}{dt} \int_{\alpha}^{\beta} u dx = -q|_{\beta}.$$  

This is equivalent to the quasilinear equation (7.2.5.9) if $q(u) = \frac{1}{2} u^2$.

The terms “conservative form”, “conservation-law form”, “weak form” or “divergence form” are all equivalent. PDEs having this form have the property that the coefficients of the derivative term are either constant or, if variable, their derivatives appear nowhere in the equation. Normally, for PDEs to represent a physical conservation statement, this means that the divergence of a physical quantity can be identified in the equation. For example, the conservation form of the one-dimensional heat equation for a substance whose density, $\rho$, specific heat, $c$, and thermal conductivity $K$, all vary with position is

$$\rho c \frac{\partial u}{\partial t} = \frac{\partial}{\partial x} \left( K \frac{\partial u}{\partial x} \right)$$

whereas a nonconservative form would be

$$\rho c \frac{\partial u}{\partial t} = \frac{\partial K}{\partial x} \frac{\partial u}{\partial x} + K \frac{\partial^2 u}{\partial x^2}.$$  

In the conservative form, the right hand side can be identified as the negative of the divergence of the heat flux (see Chapter 1).
Consider a discontinuous initial condition, then the equation must be taken in the integral form \((7.2.5.15)\). We seek a solution \(u\) and a curve \(x = x_s(t)\) across which \(u\) may have a jump. Suppose that the left and right limits are
\[
\lim_{x \to x_s(t)^-} u(x, t) = u_L
\]
\[
\lim_{x \to x_s(t)^+} u(x, t) = u_R
\]
and define the jump across \(x_s(t)\) by
\[
[u] = u_R - u_L.
\]
Let \([\alpha, \beta]\) be any interval containing \(x_s(t)\) at time \(t\). Then
\[
\frac{d}{dt} \int_{\alpha}^{\beta} u(x, t) dx = -[q(u(\beta, t)) - q(u(\alpha, t))].
\]
However the left hand side is
\[
\frac{d}{dt} \int_{x_s(t)^-}^{x_s(t)^+} u dx + \frac{d}{dt} \int_{x_s(t)^+}^{x_s(t)^-} u dx = \int_{x_s(t)^-}^{x_s(t)^+} \frac{dx_s}{dt} + \int_{x_s(t)^+}^{x_s(t)^-} \frac{dx_s}{dt} = u_L \frac{dx_s}{dt} - u_R \frac{dx_s}{dt}.
\]
Recall the rule to differentiate a definite integral when one of the endpoints depends on the variable of differentiation, i.e.
\[
\frac{d}{dt} \int_{a}^{\phi(t)} u(x, t) dx = \int_{a}^{\phi(t)} u_t(x, t) dx + u(\phi(t), t) \frac{d\phi}{dt}.
\]
Since \(u_t\) is bounded in each of the intervals separately, the integrals on the right hand side of \((7.2.5.19)\) tend to zero as \(\alpha \to x_s^-\) and \(\beta \to x_s^+\). Thus
\[
[u] \frac{dx_s}{dt} = [q].
\]
This gives the characteristic equation for shocks
\[
\frac{dx_s}{dt} = \frac{[q]}{[u]}.
\]
Going back to the example \((7.2.5.1)-(7.2.5.2)\) we find from \((7.2.5.1)\) that
\[
q = \frac{1}{2} u^2
\]
and from \((7.2.5.2)\)
\[
\begin{align*}
\frac{dx_s}{dt} &= \frac{1}{2} \cdot 1^2 - \frac{1}{2} \cdot 2^2 = -2 + \frac{1}{2} = \frac{3}{2} \\
\end{align*}
\]
Thus
\[
\frac{dx_s}{dt} = \frac{1}{1 - 2} = -1 = \frac{3}{2}
\]
\[
x_s(0) = 1 \quad \text{(where discontinuity starts)}.
\]
The solution is then
\[
x_s = \frac{3}{2} t + 1.
\]
We can now sketch this along with the other characteristics in figure 40. Any characteristic reaching the one given by \((7.2.5.21)\) will stop there. The solution is given in figure 41.
Figure 40: Shock characteristic for Example 5

Figure 41: Solution of Example 5
Problems

1. Consider Burgers’ equation

\[ \frac{\partial \rho}{\partial t} + u_{\text{max}} \left[ 1 - \frac{2\rho}{\rho_{\text{max}}} \right] \frac{\partial \rho}{\partial x} = \nu \frac{\partial^2 \rho}{\partial x^2} \]

Suppose that a solution exists as a density wave moving without change of shape at a velocity \( V, \rho(x, t) = f(x - Vt) \).

a. What ordinary differential equation is satisfied by \( f \)

b. Show that the velocity of wave propagation, \( V \), is the same as the shock velocity separating \( \rho = \rho_1 \) from \( \rho = \rho_2 \) (occurring if \( \nu = 0 \)).

2. Solve

\[ \frac{\partial \rho}{\partial t} + \rho^2 \frac{\partial \rho}{\partial x} = 0 \]

subject to

\[ \rho(x, 0) = \begin{cases} 4 & x < 0 \\ 3 & x > 0 \end{cases} \]

3. Solve

\[ \frac{\partial u}{\partial t} + 4u \frac{\partial u}{\partial x} = 0 \]

subject to

\[ u(x, 0) = \begin{cases} 3 & x < 1 \\ 2 & x > 1 \end{cases} \]

4. Solve the above equation subject to

\[ u(x, 0) = \begin{cases} 2 & x < -1 \\ 3 & x > -1 \end{cases} \]

5. Solve the quasilinear equation

\[ \frac{\partial u}{\partial t} + u \frac{\partial u}{\partial x} = 0 \]

subject to

\[ u(x, 0) = \begin{cases} 2 & x < 2 \\ 3 & x > 2 \end{cases} \]

6. Solve the quasilinear equation

\[ \frac{\partial u}{\partial t} + u \frac{\partial u}{\partial x} = 0 \]
subject to

\[ u(x, 0) = \begin{cases} 
0 & x < 0 \\
0 & 0 \leq x < 1 \\
1 & 1 \leq x
\end{cases} \]

7. Solve the inviscid Burgers’ equation

\[ u_t + uu_x = 0 \]

\[ u(x, 0) = \begin{cases} 
2 & x < 0 \\
1 & 0 < x < 1 \\
0 & x > 1
\end{cases} \]

Note that two shocks start at \( t = 0 \), and eventually intersect to create a third shock. Find the solution for all time (analytically), and graphically display your solution, labeling all appropriate bounding curves.
7.3 Second Order Wave Equation

In this section we show how the method of characteristics is applied to solve the second order wave equation describing a vibrating string. The equation is

\[ u_{tt} - c^2 u_{xx} = 0, \quad c = \text{constant}. \]  

(7.3.1)

For the rest of this chapter the unknown \( u(x, t) \) describes the displacement from rest of every point \( x \) on the string at time \( t \). We have shown in section 6.4 that the general solution is

\[ u(x, t) = F(x - ct) + G(x + ct). \]  

(7.3.2)

7.3.1 Infinite Domain

The problem is to find the solution of (7.3.1) subject to the initial conditions

\[ u(x, 0) = f(x) \quad -\infty < x < \infty \]  

(7.3.1.1)

\[ u_t(x, 0) = g(x) \quad -\infty < x < \infty. \]  

(7.3.1.2)

These conditions will specify the arbitrary functions \( F, G \). Combining the conditions with (7.3.2), we have

\[ F(x) + G(x) = f(x) \]  

(7.3.1.3)

\[ -c \frac{dF}{dx} + c \frac{dG}{dx} = g(x). \]  

(7.3.1.4)

These are two equations for the two arbitrary functions \( F \) and \( G \). In order to solve the system, we first integrate (7.3.1.4), thus

\[ -F(x) + G(x) = \frac{1}{c} \int_0^x g(\xi) d\xi. \]  

(7.3.1.5)

Therefore, the solution of (7.3.1.3) and (7.3.1.5) is

\[ F(x) = \frac{1}{2} f(x) - \frac{1}{2c} \int_0^x g(\xi) d\xi, \]  

(7.3.1.6)

\[ G(x) = \frac{1}{2} f(x) + \frac{1}{2c} \int_0^x g(\xi) d\xi. \]  

(7.3.1.7)

Combining these expressions with (7.3.2), we have

\[ u(x, t) = \frac{f(x + ct) + f(x - ct)}{2} + \frac{1}{2c} \int_{x-ct}^{x+ct} g(\xi) d\xi. \]  

(7.3.1.8)

This is d’Alembert’s solution to (7.3.1) subject to (7.3.1.1)-(7.3.1.2).

Note that the solution \( u \) at a point \((x_0, t_0)\) depends on \( f \) at the points \((x_0 + ct_0, 0)\) and \((x_0 - ct_0, 0)\), and on the values of \( g \) on the interval \((x_0 - ct_0, x_0 + ct_0)\). This interval is called
domain of dependence. In figure 42, we see that the domain of dependence is obtained by
drawing the two characteristics

\[
\begin{align*}
x - ct &= x_0 - ct_0 \\
x + ct &= x_0 + ct_0
\end{align*}
\]

through the point \((x_0, t_0)\). This behavior is to be expected because the effects of the initial
data propagate at the finite speed \(c\). Thus the only part of the initial data that can influence
the solution at \(x_0\) at time \(t_0\) must be within \(ct_0\) units of \(x_0\). This is precisely the data given
in the interval \((x_0 - ct_0, x_0 + ct_0)\).

![Figure 42: Domain of dependence](image)

The functions \(f(x)\), \(g(x)\) describing the initial position and speed of the string are defined
for all \(x\). The initial disturbance \(f(x)\) at a point \(x_1\) will propagate at speed \(c\) whereas the
effect of the initial velocity \(g(x)\) propagates at all speeds up to \(c\). This infinite sector (figure
43) is called the domain of influence of \(x_1\).

The solution (7.3.2) represents a sum of two waves, one is travelling at a speed \(c\) to the
right \((F(x - ct))\) and the other is travelling to the left at the same speed.
Figure 43: Domain of influence
Problems

1. Suppose that

\[ u(x, t) = F(x - ct). \]

Evaluate

a. \( \frac{\partial u}{\partial t}(x, 0) \)

b. \( \frac{\partial u}{\partial x}(0, t) \)

2. The general solution of the one dimensional wave equation

\[ u_{tt} - 4u_{xx} = 0 \]

is given by

\[ u(x, t) = F(x - 2t) + G(x + 2t). \]

Find the solution subject to the initial conditions

\[ u(x, 0) = \cos x \quad -\infty < x < \infty, \]

\[ u_t(x, 0) = 0 \quad -\infty < x < \infty. \]

3. In section 3.1, we suggest that the wave equation can be written as a system of two first order PDEs. Show how to solve

\[ u_{tt} - c^2 u_{xx} = 0 \]

using that idea.
7.3.2 Semi-infinite String

The problem is to solve the one-dimensional wave equation

$$u_{tt} - c^2 u_{xx} = 0, \quad 0 \leq x < \infty,$$

subject to the initial conditions

$$u(x, 0) = f(x), \quad 0 \leq x < \infty,$$

$$u_t(x, 0) = g(x), \quad 0 \leq x < \infty,$$

and the boundary condition

$$u(0, t) = h(t), \quad 0 \leq t.$$

Note that $f(x)$ and $g(x)$ are defined only for nonnegative $x$. Therefore, the solution (7.3.1.8) holds only if the arguments of $f(x)$ are nonnegative, i.e.

$$x - ct \geq 0$$

$$x + ct \geq 0$$

As can be seen in figure 44, the first quadrant must be divided to two sectors by the characteristic $x - ct = 0$. In the lower sector I, the solution (7.3.1.8) holds. In the other sector, one should note that a characteristic $x - ct = K$ will cross the negative $x$ axis and the positive $t$ axis.

![Figure 44: The characteristic $x - ct = 0$ divides the first quadrant](image)

The solution at point $(x_1, t_1)$ must depend on the boundary condition $h(t)$. We will show how the dependence presents itself.

For $x - ct < 0$, we proceed as follows:
- Combine (7.3.2.4) with the general solution (7.3.2) at \( x = 0 \)

\[
h(t) = F(-ct) + G(ct) \quad (7.3.2.6)
\]

- Since \( x - ct < 0 \) and since \( F \) is evaluated at this negative value, we use (7.3.2.6)

\[
F(-ct) = h(t) - G(ct) \quad (7.3.2.7)
\]

- Now let

\[
z = -ct < 0
\]

then

\[
F(z) = h\left(-\frac{z}{c}\right) - G(-z). \quad (7.3.2.8)
\]

So \( F \) for negative values is computed by (7.3.2.8) which requires \( G \) at positive values. In particular, we can take \( x - ct \) as \( z \), to get

\[
F(x - ct) = h\left(-\frac{x - ct}{c}\right) - G( ct - x). \quad (7.3.2.9)
\]

- Now combine (7.3.2.9) with the formula (7.3.1.7) for \( G \)

\[
F(x - ct) = h\left(t - \frac{x}{c}\right) - \left(\frac{1}{2} f(ct - x) + \frac{1}{2c} \int_{0}^{ct-x} g(\xi) d\xi\right)
\]

- The solution in sector II is then

\[
\begin{align*}
  u(x,t) = h \left(t - \frac{x}{c}\right) - \frac{1}{2} & f( ct - x ) - \frac{1}{2c} \int_{0}^{ct-x} g(\xi) d\xi + \frac{1}{2} f( x + ct) + \frac{1}{2c} \int_{0}^{x+ct} g(\xi) d\xi \\
  u(x,t) = \begin{cases} 
  \frac{f(x + ct) + f(x - ct)}{2} + \frac{1}{2c} \int_{x - ct}^{x+ct} g(\xi) d\xi & x - ct \geq 0 \\
  h \left(t - \frac{x}{c}\right) + \frac{f(x + ct) - f(x - ct)}{2} + \frac{1}{2c} \int_{ct-x}^{x+ct} g(\xi) d\xi & x - ct < 0
  \end{cases}
\end{align*}
\]

(7.3.2.10)

Note that the solution in sector II requires the knowledge of \( f(x) \) at point B (see Figure 45) which is the image of A about the \( t \) axis. The line BD is a characteristic (parallel to PC)

\[
x + ct = K.
\]

Therefore the solution at \((x_1, t_1)\) is a combination of a wave moving on the characteristic CP and one moving on BD and reflected by the wall at \( x = 0 \) to arrive at P along a characteristic

\[
x - ct = x_1 - ct_1.
\]
We now introduce several definitions to help us show that d’Alembert’s solution (7.3.1.8) holds in other cases.

**Definition 8.** A function \( f(x) \) is called an even function if
\[
f(-x) = f(x).
\]

**Definition 9.** A function \( f(x) \) is called an odd function if
\[
f(-x) = -f(x).
\]

Note that some functions are neither.

**Examples**
1. \( f(x) = x^2 \) is an even function.
2. \( f(x) = x^3 \) is an odd function.
3. \( f(x) = x - x^2 \) is neither odd nor even.

**Definition 10.** A function \( f(x) \) is called a periodic function of period \( p \) if
\[
f(x + p) = f(x) \quad \text{for all } x.
\]

The smallest such real number \( p \) is called the fundamental period.

**Remark:** If the boundary condition (7.3.2.4) is
\[
u(0, t) = 0,
\]
then the solution for the semi-infinite interval is the same as that for the infinite interval with \( f(x) \) and \( g(x) \) being extended as odd functions for \( x < 0 \). Since if \( f \) and \( g \) are odd functions then
\[
f(-z) = -f(z), \quad g(-z) = -g(z),
\]
(7.3.2.11)
The solution for $x - ct$ is now
\[ u(x, t) = \frac{f(x + ct) - f(-(x - ct))}{2} + \frac{1}{2c} \left( \int_{ct-x}^{0} g(\xi) d\xi + \int_{0}^{x+ct} g(\xi) d\xi \right). \] (7.3.2.12)

But if we let $\zeta = -\xi$ then
\[ \int_{ct-x}^{0} g(\xi) d\xi = \int_{x-ct}^{0} g(-\zeta)(-d\zeta) = \int_{x-ct}^{0} g(\zeta) d\zeta. \]

Now combine this integral with the last term in (7.3.2.12) to have
\[ u(x, t) = \frac{f(x + ct) + f(x - ct)}{2} + \frac{1}{2c} \int_{x-ct}^{x+ct} g(\xi) d\xi \]
which is exactly the same formula as for $x - ct \geq 0$. Therefore we have shown that for a semi-infinite string with fixed ends, one can use d’Alembert’s solution (7.3.1.8) after extending $f(x)$ and $g(x)$ as odd functions for $x < 0$.

What happens if the boundary condition is
\[ u_x(0, t) = 0? \]

We claim that one has to extend $f(x)$, $g(x)$ as even functions and then use (7.3.1.8). The details will be given in the next section.

### 7.3.3 Semi Infinite String with a Free End

In this section we show how to solve the wave equation
\[ u_{tt} - c^2 u_{xx} = 0, \quad 0 \leq x < \infty, \] (7.3.3.1)
subject to
\[ u(x, 0) = f(x), \] (7.3.3.2)
\[ u_t(x, 0) = g(x), \] (7.3.3.3)
\[ u_x(0, t) = 0. \] (7.3.3.4)

Clearly, the general solution for $x - ct \geq 0$ is the same as before, i.e. given by (7.3.1.8). For $x - ct < 0$, we proceed in a similar fashion as last section. Using the boundary condition (7.3.3.4)
\[ 0 = u_x(0, t) = \frac{dF}{dx} \bigg|_{x=0} + \frac{dG}{dx} \bigg|_{x=0} = F'(-ct) + G'(ct). \]

Therefore
\[ F'(-ct) = -G'(ct). \] (7.3.3.5)
Let \( z = -ct < 0 \) and integrate over \([0, z]\)

\[
F(z) - F(0) = G(-z) - G(0).
\]  
(7.3.3.6)

From (7.3.1.6)-(7.3.1.7) we have

\[
F(0) = G(0) = \frac{1}{2} f(0).
\]  
(7.3.3.7)

Replacing \( z \) by \( x - ct < 0 \), we have

\[
F(x - ct) = G(-(x - ct)),
\]
or

\[
F(x - ct) = \frac{1}{2} f(ct - x) + \frac{1}{2c} \int_{x}^{x-ct} g(\xi) d\xi.
\]  
(7.3.3.8)

To summarize, the solution is

\[
u(x, t) = \begin{cases} 
\frac{f(x + ct) + f(x - ct)}{2} + \frac{1}{2c} \int_{x-ct}^{x+ct} g(\xi) d\xi, & x \geq ct \\
\frac{f(x + ct) + f(x - ct)}{2} + \frac{1}{2c} \int_{0}^{x} g(\xi) d\xi + \frac{1}{2c} \int_{0}^{ct-x} g(\xi) d\xi, & x < ct.
\end{cases}
\]  
(7.3.3.9)

Remark: If \( f(x) \) and \( g(x) \) are extended for \( x < 0 \) as even functions then

\[
f(ct - x) = f(-(x - ct)) = f(x - ct)
\]
and

\[
\int_{0}^{ct-x} g(\xi) d\xi = \int_{0}^{x-ct} g(\zeta) d\zeta = \int_{x-ct}^{0} g(\zeta) d\zeta
\]
where \( \zeta = -\xi \).

Thus the integrals can be combined to one to give

\[
\frac{1}{2c} \int_{x-ct}^{x+ct} g(\xi) d\xi.
\]

Therefore with this extension of \( f(x) \) and \( g(x) \) we can write the solution in the form (7.3.1.8).
Problems

1. Solve by the method of characteristics

\[ \frac{\partial^2 u}{\partial t^2} - c^2 \frac{\partial^2 u}{\partial x^2} = 0, \quad x > 0 \]

subject to

\[ u(x, 0) = 0, \]
\[ \frac{\partial u}{\partial t}(x, 0) = 0, \]
\[ u(0, t) = h(t). \]

2. Solve

\[ \frac{\partial^2 u}{\partial t^2} - c^2 \frac{\partial^2 u}{\partial x^2} = 0, \quad x < 0 \]

subject to

\[ u(x, 0) = \sin x, \quad x < 0 \]
\[ \frac{\partial u}{\partial t}(x, 0) = 0, \quad x < 0 \]
\[ u(0, t) = e^{-t}, \quad t > 0. \]

3. a. Solve

\[ \frac{\partial^2 u}{\partial t^2} - c^2 \frac{\partial^2 u}{\partial x^2} = 0, \quad 0 < x < \infty \]

subject to

\[ u(x, 0) = \begin{cases} 
0 & 0 < x < 2 \\
1 & 2 < x < 3 \\
0 & 3 < x 
\end{cases} \]
\[ \frac{\partial u}{\partial t}(x, 0) = 0, \]
\[ \frac{\partial u}{\partial x}(0, t) = 0. \]

b. Suppose \( u \) is continuous at \( x = t = 0 \), sketch the solution at various times.

4. Solve

\[ \frac{\partial^2 u}{\partial t^2} - c^2 \frac{\partial^2 u}{\partial x^2} = 0, \quad x > 0, \quad t > 0 \]

subject to

\[ u(x, 0) = 0, \]
\[ \frac{\partial u}{\partial t}(x, 0) = 0, \]
\[ \frac{\partial u}{\partial x}(0, t) = h(t). \]

5. Give the domain of influence in the case of semi-infinite string.
7.3.4 Finite String

This problem is more complicated because of multiple reflections. Consider the vibrations of a string of length $L$,

$$u_{tt} - c^2 u_{xx} = 0, \quad 0 \leq x \leq L,$$

subject to

$$u(x,0) = f(x), \quad (7.3.4.2)$$

$$u_t(x,0) = g(x), \quad (7.3.4.3)$$

$$u(0,t) = 0, \quad (7.3.4.4)$$

$$u(L,t) = 0. \quad (7.3.4.5)$$

From the previous section, we can write the solution in regions 1 and 2 (see figure 46), i.e.

![Figure 46: Reflected waves reaching a point in region 5](image)

$u(x,t)$ is given by (7.3.1.8) in region 1 and by (7.3.2.10) with $h \equiv 0$ in region 2. The solution in region 3 can be obtained in a similar fashion as (7.3.2.10), but now use the boundary condition (7.3.4.5).

In region 3, the boundary condition (7.3.4.5) becomes

$$u(L,t) = F(L - ct) + G(L + ct) = 0. \quad (7.3.4.6)$$

Since $L + ct \geq L$, we solve for $G$

$$G(L + ct) = -F(L - ct).$$
Let
\[ z = L + ct \geq L, \]
then
\[ L - ct = 2L - z \leq L. \]
Thus
\[ G(z) = -F(2L - z) \]
or
\[ G(x + ct) = -F(2L - x - ct) = -\frac{1}{2} f(2L - x - ct) + \frac{1}{2c} \int_{0}^{2L-x-ct} g(\xi)d\xi \]
and so adding \( F(x - ct) \) given by (7.3.1.6) to the above we get the solution in region 3,
\[ u(x, t) = \frac{f(x - ct) - f(2L - x - ct)}{2} + \frac{1}{2c} \int_{0}^{x-c(t)} g(\xi)d\xi + \frac{1}{2c} \int_{0}^{2L-x-c(t)} g(\xi)d\xi. \]
In other regions multiply reflected waves give the solution. (See figure 46, showing doubly reflected waves reaching points in region 5.)

As we remarked earlier, the boundary condition (7.3.4.4) essentially say that the initial conditions were extended as odd functions for \( x < 0 \) (in this case for \([-L \leq x \leq 0]\). The other boundary condition means that the initial conditions are extended again as odd functions to the interval \([L, 2L]\), which is the same as saying that the initial conditions on the interval \([-L, L]\) are now extended periodically everywhere. Once the functions are extended to the real line, one can use (7.3.1.8) as a solution. A word of caution, this is true only when the boundary conditions are given by (7.3.4.4)-(7.3.4.5).

Figure 47: Parallelogram rule
Parallelogram Rule

If the four points $A$, $B$, $C$, and $D$ form the vertices of a parallelogram whose sides are all segments of characteristic curves, (see figure 47) then the sums of the values of $u$ at opposite vertices are equal, i.e.

$$u(A) + u(C) = u(B) + u(D).$$

This rule is useful in solving a problem with both initial and boundary conditions.

In region $R_1$ (see figure 47) the solution is defined by d’Alembert’s formula. For $A = (x, t)$ in region $R_2$, let us form the parallelogram $ABCD$ with $B$ on the $t$-axis and $C$ and $D$ on the characteristic curve from $(0, 0)$. Thus

$$u(A) = -u(C) + u(B) + u(D)$$

Figure 48: Use of parallelogram rule to solve the finite string case

$u(B)$ is a known boundary value and the others are known from $R_1$. We can do this for any point $A$ in $R_2$. Similarly for $R_3$. One can use the solutions in $R_2$, $R_3$ to get the solution in $R_4$ and so on. The limitation is that $u$ must be given on the boundary. If the boundary conditions are not of Dirichlet type, this rule is not helpful.
SUMMARY
Linear:
\[ u_t + c(x,t)u_x = S(u,x,t) \]
\[ u(x(0), 0) = f(x(0)) \]
Solve the characteristic equation
\[ \frac{dx}{dt} = c(x,t) \]
\[ x(0) = x_0 \]
then solve
\[ \frac{du}{dt} = S(u,x,t) \]
\[ u(x(0), 0) = f(x(0)) \] on the characteristic curve
Quasilinear:
\[ u_t + c(u,x,t)u_x = S(u,x,t) \]
\[ u(x(0), 0) = f(x(0)) \]
Solve the characteristic equation
\[ \frac{dx}{dt} = c(u,x,t) \]
\[ x(0) = x_0 \]
then solve
\[ \frac{du}{dt} = S(u,x,t) \]
\[ u(x(0), 0) = f(x(0)) \] on the characteristic curve
fan-like characteristics
shock waves
Second order hyperbolic equations:
Infinite string
\[ u_{tt} - c^2 u_{xx} = 0 \quad c = \text{constant,} \quad -\infty < x < \infty \]
\[ u(x,0) = f(x), \]
\[ u_t(x,0) = g(x), \]
\[ u(x,t) = \frac{f(x+ct) + f(x-ct)}{2} + \frac{1}{2c} \int_{x-ct}^{x+ct} g(\xi)d\xi. \]
Semi infinite string

\[
\begin{align*}
    u_{tt} - c^2 u_{xx} &= 0 \quad c = \text{constant}, \quad 0 \leq x < \infty \\
    u(x,0) &= f(x), \\
    u_t(x,0) &= g(x), \\
    u(0,t) &= h(t), \quad 0 \leq t.
\end{align*}
\]

\[
u(x,t) = \begin{cases} 
    \frac{f(x+ct) + f(x-ct)}{2} + \frac{1}{2c} \int_{x-ct}^{x+ct} g(\xi) \, d\xi, & x - ct \geq 0, \\
    h \left( t - \frac{x}{c} \right) + \frac{f(x+ct) - f(ct-x)}{2} + \frac{1}{2c} \int_{ct-x}^{x+ct} g(\xi) \, d\xi, & x - ct < 0.
\end{cases}
\]

Semi infinite string - free end

\[
\begin{align*}
    u_{tt} - c^2 u_{xx} &= 0 \quad c = \text{constant}, \quad 0 \leq x < \infty, \\
    u(x,0) &= f(x), \\
    u_t(x,0) &= g(x), \\
    u_x(0,t) &= h(t).
\end{align*}
\]

\[
u(x,t) = \begin{cases} 
    \frac{f(x+ct) + f(x-ct)}{2} + \frac{1}{2c} \int_{x-ct}^{x+ct} g(\xi) \, d\xi, & x \geq ct, \\
    \int_0^{x-ct} h(-z/c) \, dz + \frac{f(x+ct) + f(ct-x)}{2} + \frac{1}{2c} \int_0^{x+ct} g(\xi) \, d\xi + \frac{1}{2c} \int_{ct-x}^{x+ct} g(\xi) \, d\xi, & x < ct.
\end{cases}
\]
8 Finite Differences

8.1 Taylor Series

In this chapter we discuss finite difference approximations to partial derivatives. The approximations are based on Taylor series expansions of a function of one or more variables.

Recall that the Taylor series expansion for a function of one variable is given by

\[ f(x + h) = f(x) + \frac{h}{1!} f'(x) + \frac{h^2}{2!} f''(x) + \cdots \]  

(8.1.1)

The remainder is given by

\[ f^{(n)}(\xi) \frac{h^n}{n!}, \quad \xi \in (x, x + h). \]  

(8.1.2)

For a function of more than one independent variable we have the derivatives replaced by partial derivatives. We give here the case of 2 independent variables

\[ f(x + h, y + k) = f(x, y) + \frac{h}{1!} f_x(x, y) + \frac{k}{1!} f_y(x, y) + \frac{h^2}{2!} f_{xx}(x, y) \]

\[ + \frac{2hk}{2!} f_{xy}(x, y) + \frac{k^2}{2!} f_{yy}(x, y) + \frac{h^3}{3!} f_{xxx}(x, y) + \frac{3h^2k}{3!} f_{xxy}(x, y) \]

\[ + \frac{3h^2k^2}{3!} f_{xyy}(x, y) + \frac{k^3}{3!} f_{yyy}(x, y) + \cdots \]  

(8.1.3)

The remainder can be written in the form

\[ \frac{1}{n!} \left( h \frac{\partial}{\partial x} + k \frac{\partial}{\partial y} \right)^n f(x + \theta h, y + \theta k), \quad 0 \leq \theta \leq 1. \]  

(8.1.4)

Here we used a subscript to denote partial differentiation. We will be interested in obtaining approximation about the point \((x_i, y_j)\) and we use a subscript to denote the function values at the point, i.e. \(f_{ij} = f(x_i, y_j)\).

The Taylor series expansion for \(f_{i+1}\) about the point \(x_i\) is given by

\[ f_{i+1} = f_i + h f_i' + \frac{h^2}{2!} f_i'' + \frac{h^3}{3!} f_i''' + \cdots \]  

(8.1.5)

The Taylor series expansion for \(f_{i+1,j+1}\) about the point \((x_i, y_j)\) is given by

\[ f_{i+1,j+1} = f_{ij} + (h_x f_{x} + h_y f_{y})_{ij} + \left( \frac{h_x^2}{2} f_{xx} + h_x h_y f_{xy} + \frac{h_y^2}{2} f_{yy} \right)_{ij} + \cdots \]  

(8.1.6)

Remark: The expansion for \(f_{i+1,j}\) about \((x_i, y_j)\) proceeds as in the case of a function of one variable.

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8.2 Finite Differences

An infinite number of difference representations can be found for the partial derivatives of \( f(x,y) \). Let us use the following operators:

- **forward difference operator**
  \[ \Delta_x f_{ij} = f_{i+1j} - f_{ij} \] (8.2.1)

- **backward difference operator**
  \[ \nabla_x f_{ij} = f_{ij} - f_{i-1j} \] (8.2.2)

- **centered difference**
  \[ \delta_x f_{ij} = f_{i+1j} - f_{i-1j} \] (8.2.3)
  \[ \delta_x f_{ij} = f_{i+1/2j} - f_{i-1/2j} \] (8.2.4)

- **averaging operator**
  \[ \mu_x f_{ij} = (f_{i+1/2j} + f_{i-1/2j})/2 \] (8.2.5)

Note that \( \delta_x = \mu_x \delta_x \). (8.2.6)

In a similar fashion we can define the corresponding operators in \( y \).

In the following table we collected some of the common approximations for the first derivative.

<table>
<thead>
<tr>
<th>Finite Difference</th>
<th>Order (see next chapter)</th>
</tr>
</thead>
<tbody>
<tr>
<td>( \frac{1}{h_x} \Delta_x f_{ij} )</td>
<td>( O(h_x) )</td>
</tr>
<tr>
<td>( \frac{1}{h_x} \nabla_x f_{ij} )</td>
<td>( O(h_x) )</td>
</tr>
<tr>
<td>( \frac{1}{2h_x} \delta_x f_{ij} )</td>
<td>( O(h_x^2) )</td>
</tr>
<tr>
<td>( \frac{1}{2h_x} (-3f_{ij} + 4f_{i+1j} - f_{i+2j}) = \frac{1}{h_x} (\Delta_x - \frac{1}{2}\Delta_x^2) f_{ij} )</td>
<td>( O(h_x^2) )</td>
</tr>
<tr>
<td>( \frac{1}{2h_x} (3f_{ij} - 4f_{i-1j} + f_{i-2j}) = \frac{1}{h_x} (\nabla_x + \frac{1}{2}\nabla_x^2) f_{ij} )</td>
<td>( O(h_x^2) )</td>
</tr>
<tr>
<td>( \frac{1}{h_x} (\mu_x \delta_x - \frac{1}{3} \mu_x \delta_x^3) f_{ij} )</td>
<td>( O(h_x^3) )</td>
</tr>
<tr>
<td>( \frac{1}{2h_x} \delta_x f_{ij} )</td>
<td>( O(h_x^4) )</td>
</tr>
</tbody>
</table>

Table 1: Order of approximations to \( f_x \)

The compact fourth order three point scheme deserves some explanation. Let \( f_x \) be \( v \), then the method is to be interpreted as

\[ (1 + \frac{1}{6} \delta_x^2) v_{ij} = \frac{1}{2h_x} \delta_x f_{ij} \] (8.2.7)

or

\[ \frac{1}{6} (v_{i+1j} + 4v_{ij} + v_{i-1j}) = \frac{1}{2h_x} \delta_x f_{ij}. \] (8.2.8)
This is an implicit formula for the derivative \( \frac{\partial f}{\partial x} \) at \((x_i, y_j)\). The \( v_{ij} \) can be computed from the \( f_{ij} \) by solving a tridiagonal system of algebraic equations.

The most common second derivative approximations are

\[
f_{xx|i,j} = \frac{1}{h_x^2} \left( f_{ij} - 2f_{i+1,j} + f_{i+2,j} \right) + O(h_x) \tag{8.2.9}
\]

\[
f_{xx|i,j} = \frac{1}{h_x^2} \left( f_{ij} - 2f_{i-1,j} + f_{i-2,j} \right) + O(h_x) \tag{8.2.10}
\]

\[
f_{xx|i,j} = \frac{1}{h_x^2} \delta_x^2 f_{ij} + O(h_x^2) \tag{8.2.11}
\]

\[
f_{xx|i,j} = \frac{1}{h_x^2} \frac{\delta_x^2 f_{ij}}{1 + \frac{\delta_x^2}{12}} + O(h_x^4) \tag{8.2.12}
\]

Remarks:
1. The order of a scheme is given for a uniform mesh.
2. Tables for difference approximations using more than three points and approximations of mixed derivatives are given in Anderson, Tannehill and Pletcher (1984, p.45).
3. We will use the notation
\[
\delta_x^2 = \frac{\delta_x^2}{h_x^2} \tag{8.2.13}
\]

The centered difference operator can be written as a product of the forward and backward operator, i.e.

\[
\delta_x^2 f_{ij} = \nabla_x \Delta_x f_{ij} \tag{8.2.14}
\]

This is true since on the right we have

\[
\nabla_x (f_{i+1,j} - f_{ij}) = f_{i+1,j} - f_{ij} - (f_{ij} - f_{i-1,j})
\]

which agrees with the right hand side of (8.2.14). This idea is important when one wants to approximate \((p(x)y'(x))'\) at the point \(x_i\) to a second order. In this case one takes the forward difference inside and the backward difference outside (or vice versa)

\[
\nabla_x \left( \frac{p_i y_{i+1} - y_i}{\Delta x} \right) \tag{8.2.15}
\]

and after expanding again

\[
\frac{p_i y_{i+1} - y_i}{\Delta x} - \frac{p_{i-1} y_i - y_{i-1}}{\Delta x} \tag{8.2.16}
\]

or

\[
\frac{p_i y_{i+1} - (p_i + p_{i-1}) y_i + p_{i-1} y_{i-1}}{(\Delta x)^2} \tag{8.2.17}
\]

Note that if \(p(x) \equiv 1\) then we get the well known centered difference.
Problems

1. Verify that
\[ \frac{\partial^3 u}{\partial x^3} u_{i,j} = \frac{\Delta_x^3 u_{i,j}}{(\Delta x)^3} + O(\Delta x). \]

2. Consider the function \( f(x) = e^x \). Using a mesh increment \( \Delta x = 0.1 \), determine \( f'(x) \) at \( x = 2 \) with forward-difference formula, the central-difference formula, and the second order three-point formula. Compare the results with the exact value. Repeat the comparison for \( \Delta x = 0.2 \). Have the order estimates for truncation errors been a reliable guide? Discuss this point.

3. Develop a finite difference approximation with T.E. of \( O(\Delta y) \) for \( \partial^2 u / \partial y^2 \) at point \((i, j)\) using \( u_{i,j}, u_{i,j+1}, u_{i,j-1} \) when the grid spacing is not uniform. Use the Taylor series method. Can you devise a three point scheme with second-order accuracy with unequal spacing? Before you draw your final conclusions, consider the use of compact implicit representations.

4. Establish the T.E. for the following finite difference approximation to \( \partial u / \partial y \) at the point \((i, j)\) for a uniform mesh:
\[ \frac{\partial u}{\partial y} \approx -3u_{i,j} + 4u_{i,j+1} - u_{i,j+2} \]
\[ 2\Delta y. \]
What is the order?
Clearly it is more convenient to use a uniform mesh and it is more accurate in some cases. However, in many cases this is not possible due to boundaries which do not coincide with the mesh or due to the need to refine the mesh in part of the domain to maintain the accuracy. In the latter case one is advised to use a coordinate transformation.

In the former case several possible cures are given in, e.g. Anderson et al (1984). The most accurate of these is a development of a finite difference approximation which is valid even when the mesh is nonuniform. It can be shown that

\[ u_{xx} \bigg|_O \approx \frac{2}{(1 + \alpha)h_x} \left( \frac{u_c - u_O}{\alpha h_x} - \frac{u_O - u_A}{h_x} \right) \]  

(8.3.1)

Similar formula for \( u_{yy} \). Note that for \( \alpha = 1 \) one obtains the centered difference approximation.

We now develop a three point second order approximation for \( \frac{\partial f}{\partial x} \) on a nonuniform mesh. \( \frac{\partial f}{\partial x} \) at point \( O \) can be written as a linear combination of values of \( f \) at \( A, O, \) and \( B \),

\[ \frac{\partial f}{\partial x} \bigg|_O = C_1 f(A) + C_2 f(O) + C_3 f(B). \]  

(8.3.2)

We use Taylor series to expand \( f(A) \) and \( f(B) \) about the point \( O \),

\[ f(A) = f(O - \Delta x) = f(O) - \Delta x f'(O) + \frac{\Delta x^2}{2} f''(O) - \frac{\Delta x^3}{6} f'''(O) \pm \cdots \]  

(8.3.3)
\[ f(B) = f(O + \alpha \Delta x) = f(O) + \alpha \Delta x f'(O) + \frac{\alpha^2 \Delta x^2}{2} f''(O) + \frac{\alpha^3 \Delta x^3}{6} f'''(O) + \cdots \]  

(8.3.4)

Thus

\[
\left. \frac{\partial f}{\partial x} \right|_O = (C_1 + C_2 + C_3) f(O) + (\alpha C_3 - C_1) \Delta x \left. \frac{\partial f}{\partial x} \right|_O + (C_1 + \alpha^2 C_3) \frac{\Delta x^2}{2} \left. \frac{\partial^2 f}{\partial x^2} \right|_O \\
+ (\alpha^3 C_3 - C_1) \frac{\Delta x^3}{6} \left. \frac{\partial^3 f}{\partial x^3} \right|_O + \cdots
\]

(8.3.5)

This yields the following system of equations

\[
C_1 + C_2 + C_3 = 0
\]

(8.3.6)

\[
-C_1 + \alpha C_3 = \frac{1}{\Delta x}
\]

(8.3.7)

\[
C_1 + \alpha^2 C_3 = 0
\]

(8.3.8)

The solution is

\[
C_1 = \frac{-\alpha}{(\alpha + 1) \Delta x}, \quad C_2 = \frac{\alpha - 1}{\alpha \Delta x}, \quad C_3 = \frac{1}{\alpha (\alpha + 1) \Delta x}
\]

(8.3.9)

and thus

\[
\left. \frac{\partial f}{\partial x} \right|_O = -\frac{\alpha^2 f(A) + (\alpha^2 - 1)f(O) + f(B)}{\alpha (\alpha + 1) \Delta x} + \frac{\alpha^3 \Delta x^3}{6} \left. \frac{\partial^3 f}{\partial x^3} \right|_O + \cdots
\]

(8.3.10)

Note that if the grid is uniform then \( \alpha = 1 \) and this becomes the familiar centered difference.
Problems

1. Develop a finite difference approximation with T.E. of \( O(\Delta y)^2 \) for \( \partial T / \partial y \) at point \((i, j)\) using \( T_{i,j}, T_{i,j+1}, T_{i,j+2} \) when the grid spacing is \textbf{not} uniform.

2. Determine the T.E. of the following finite difference approximation for \( \partial u / \partial x \) at point \((i, j)\) when the grid space is \textbf{not} uniform:

\[
\frac{\partial u}{\partial x} \bigg|_{i,j} \approx \frac{u_{i+1,j} - (\Delta x_+/\Delta x_-)^2 u_{i-1,j} - [1 - (\Delta x_+/\Delta x_-)^2] u_{i,j}}{\Delta x_- (\Delta x_+/\Delta x_-)^2 + \Delta x_+}
\]
8.4 Thomas Algorithm

This is an algorithm to solve a tridiagonal system of equations

\[
\begin{pmatrix}
  d_1 & a_1 \\
  b_2 & d_2 & a_2 \\
  & b_3 & d_3 & a_3 \\
  & & & & \ddots
\end{pmatrix}
\begin{pmatrix}
  u_1 \\
  u_2 \\
  u_3 \\
  \vdots
\end{pmatrix}
= \begin{pmatrix}
  c_1 \\
  c_2 \\
  c_3 \\
  \vdots
\end{pmatrix}
\]  

(8.4.1)

The first step of Thomas algorithm is to bring the tridiagonal $M$ by $M$ matrix to an upper triangular form

\[
d_i \leftarrow d_i - \frac{b_i}{d_{i-1}}a_{i-1}, \quad i = 2, 3, \ldots, M
\]

(8.4.2)

\[
c_i \leftarrow c_i - \frac{b_i}{d_{i-1}}c_{i-1}, \quad i = 2, 3, \ldots, M.
\]

(8.4.3)

The second step is to backsolve

\[
u_M = \frac{c_M}{d_M}
\]

(8.4.4)

\[
u_j = \frac{c_j - a_j u_{j+1}}{d_j}, \quad j = M - 1, \ldots, 1.
\]

(8.4.5)

The following subroutine solves a tridiagonal system of equations:

```
subroutine tridg(il,iu,rl,d,ru,r)
  integer il, iu
  real rl(1), d(1), ru(1), r(1)
  c the equations are
  C rl(i)*u(i-1)+d(i)*u(i)+ru(i)*u(i+1)=r(i)
  C il subscript of first equation
  C iu subscript of last equation
  C
  ilp=il+1
  do 1 i=ilp,iu
    g=rl(i)/d(i-1)
    d(i)=d(i)-g*ru(i-1)
    r(i)=r(i)-g*r(i-1)
  1 continue
```

C

8.5 Methods for Approximating PDEs

In this section we discuss several methods to approximate PDEs. These are certainly not all the possibilities.

8.5.1 Undetermined coefficients

In this case, we approximate the required partial derivative by a linear combination of function values. The weights are chosen so that the approximation is of the appropriate order. For example, we can approximate \( u_{xx} \) at \( x_i, y_j \) by taking the three neighboring points,

\[
 u_{xx}|_{i,j} = Au_{i+1,j} + Bu_{i,j} + Cu_{i-1,j} \tag{8.5.1.1}
\]

Now expand each of the terms on the right in Taylor series and compare coefficients (all terms are evaluated at \( i,j \))

\[
 u_{xx} = A \left( u + hu_x + \frac{h^2}{2}u_{xx} + \frac{h^3}{6}u_{xxx} + \frac{h^4}{24}u_{xxxx} + \cdots \right) + B u + C \left( u - hu_x + \frac{h^2}{2}u_{xx} - \frac{h^3}{6}u_{xxx} + \frac{h^4}{24}u_{xxxx} + \cdots \right) \tag{8.5.1.2}
\]

Upon collecting coefficients, we have

\[
 A + B + C = 0 \tag{8.5.1.3}
\]

\[
 A - C = 0 \tag{8.5.1.4}
\]

\[
 (A + C) \frac{h^2}{2} = 1 \tag{8.5.1.5}
\]

This yields

\[
 A = C = \frac{1}{h^2} \tag{8.5.1.6}
\]

\[
 B = \frac{-2}{h^2} \tag{8.5.1.7}
\]
The error term, is the next nonzero term, which is

\[(A + C)\frac{h^4}{24}u_{xxxx} = \frac{h^2}{12}u_{xxxx}.\]  

(8.5.1.8)

We call the method second order, because of the \(h^2\) factor in the error term. This is the centered difference approximation given by (8.2.11).

### 8.5.2 Polynomial Fitting

We demonstrate the use of polynomial fitting on Laplace’s equation.

\[ u_{xx} + u_{yy} = 0 \]  

(8.5.2.1)

The solution can be approximated **locally** by a polynomial, say

\[ u(x, y_0) = a + bx + cx^2. \]  

(8.5.2.2)

Suppose we take \( x = 0 \) at \((x_i, y_j)\), then

\[ \frac{\partial u}{\partial x} = b \]  

(8.5.2.3)

\[ \frac{\partial^2 u}{\partial x^2} = 2c. \]  

(8.5.2.4)

To find \( a, b, c \) in terms of grid values we have to assume which points to use. Clearly

\[ u_{i,j} = a. \]  

(8.5.2.5)

Suppose we use the points \( i+1\ j \) and \( i-1\ j \) (i.e. centered differencing, then

\[ u_{i+1,j} = a + b\Delta x + c(\Delta x)^2 \]  

(8.5.2.6)

\[ u_{i-1,j} = a - b\Delta x + c(\Delta x)^2 \]  

(8.5.2.7)

Subtracting these two equations, we get

\[ b = \frac{u_{i+1,j} - u_{i-1,j}}{2\Delta x} \]  

(8.5.2.8)

and substituting \( a \) and \( b \) in the equation for \( u_{i+1,j} \), we get

\[ 2c = \frac{u_{i+1,j} - 2u_{i,j} + u_{i-1,j}}{(\Delta x)^2} \]  

(8.5.2.9)

but we found earlier that \(2c\) is \(u_{xx}\), this gives the centered difference approximation for \(u_{xx}\). Similarly for \(u_{yy}\), now taking a quadratic polynomial in \(y\).
8.5.3 Integral Method

The strategy here is to develop an algebraic relationship among the values of the unknowns at neighboring grid points, by integrating the PDE. We demonstrate this on the heat equation integrated around the point \((x_j, t_n)\). The solution at this point can be related to neighboring values by integration, e.g.

\[ \int_{x_j - \Delta x/2}^{x_j + \Delta x/2} \left( \int_{t_n}^{t_n + \Delta t} u_t \, dt \right) \, dx = \alpha \int_{t_n}^{t_n + \Delta t} \left( \int_{x_j - \Delta x/2}^{x_j + \Delta x/2} u_{xx} \, dx \right) \, dt. \tag{8.5.3.1} \]

Note the order of integration on both sides.

\[ \int_{x_j - \Delta x/2}^{x_j + \Delta x/2} \left( u(x, t_n + \Delta t) - u(x, t_n) \right) \, dx = \alpha \int_{t_n}^{t_n + \Delta t} \left( u_x(x_j + \Delta x/2, t) - u_x(x_j - \Delta x/2, t) \right) \, dt. \tag{8.5.3.2} \]

Now use the mean value theorem, choosing \(x_j\) as the intermediate point on the left and \(t_n + \Delta t\) as the intermediate point on the right,

\[ (u(x_j, t_n + \Delta t) - u(x_j, t_n)) \, \Delta x = \alpha \left( u_x(x_j + \Delta x/2, t_n + \Delta t) - u_x(x_j - \Delta x/2, t_n + \Delta t) \right) \, \Delta t. \tag{8.5.3.3} \]

Now use a centered difference approximation for the \(u_x\) terms and we get the fully implicit scheme, i.e.

\[ \frac{u_j^{n+1} - u_j^n}{\Delta t} = \alpha \frac{u_{j+1}^{n+1} - 2u_j^{n+1} + u_{j-1}^{n+1}}{(\Delta x)^2}. \tag{8.5.3.4} \]

8.6 Eigenpairs of a Certain Tridiagonal Matrix

Let \(A\) be an \(M\) by \(M\) tridiagonal matrix whose elements on the diagonal are all \(a\), on the superdiagonal are all \(b\) and on the subdiagonal are all \(c\),

\[ A = \begin{pmatrix} a & b \\ c & a & b \\ & c & a & b \\ & & c & a \end{pmatrix}. \tag{8.6.1} \]

Let \(\lambda\) be an eigenvalue of \(A\) with an eigenvector \(v\), whose components are \(v_i\). Then the eigenvalue equation

\[ Av = \lambda v \tag{8.6.2} \]

can be written as follows

\[ (a - \lambda)v_1 + bv_2 = 0 \]
\[ cv_1 + (a - \lambda)v_2 + bv_3 = 0 \]

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\[
\ldots
\]
\[
cv_{j-1} + (a - \lambda)v_j + bv_{j+1} = 0
\]
\[
\ldots
\]
\[
cv_{M-1} + (a - \lambda)v_M = 0.
\]
If we let \(v_0 = 0\) and \(v_{M+1} = 0\), then all the equations can be written as
\[
cv_{j-1} + (a - \lambda)v_j + bv_{j+1} = 0, \quad j = 1, 2, \ldots, M. \tag{8.6.3}
\]
The solution of such second order difference equation is
\[
v_j = Bm_1^j + Cm_2^j \tag{8.6.4}
\]
where \(m_1\) and \(m_2\) are the solutions of the characteristic equation
\[
c + (a - \lambda)m + bm^2 = 0. \tag{8.6.5}
\]
It can be shown that the roots are distinct (otherwise \(v_j = (B + Cj)m_1^j\) and the boundary conditions forces \(B = C = 0\)). Using the boundary conditions, we have
\[
B + C = 0 \tag{8.6.6}
\]
and
\[
Bm_1^{M+1} + Cm_2^{M+1} = 0. \tag{8.6.7}
\]
Hence
\[
\left(\frac{m_1}{m_2}\right)^{M+1} = 1 = e^{2\pi si}, \quad s = 1, 2, \ldots, M. \tag{8.6.8}
\]
Therefore
\[
\frac{m_1}{m_2} = e^{2\pi i/(M+1)}. \tag{8.6.9}
\]
From the characteristic equation, we have
\[
m_1m_2 = \frac{c}{b}, \tag{8.6.10}
\]
eliminating \(m_2\) leads to
\[
m_1 = \sqrt{\frac{c}{b}}e^{s\pi i/(M+1)}. \tag{8.6.11}
\]
Similarly for \(m_2\),
\[
m_2 = \sqrt{\frac{c}{b}}e^{-s\pi i/(M+1)}. \tag{8.6.12}
\]
Again from the characteristic equation
\[
m_1 + m_2 = (\lambda - a)/b, \tag{8.6.13}
\]
giving
\[
\lambda = a + b\sqrt{\frac{c}{b}}\left(e^{s\pi i/(M+1)} + e^{-s\pi i/(M+1)}\right). \tag{8.6.14}
\]
Hence the $M$ eigenvalues are

$$\lambda_s = a + 2b \sqrt{e b} \cos \frac{s\pi}{M+1}, \quad s = 1, 2, \ldots, M. \quad (8.6.15)$$

The $j^{th}$ component of the eigenvector is

$$v_j = Bm_1^j + Cm_2^j = B \left( \frac{c}{b} \right)^{j/2} \left( e^{j\pi i/(M+1)} - e^{-j\pi i/(M+1)} \right), \quad (8.6.16)$$

that is

$$v_j = 2iB \left( \frac{c}{b} \right)^{j/2} \sin \frac{j\pi}{M+1}. \quad (8.6.17)$$

Use centered difference to approximate the second derivative in $X'' + \lambda X = 0$ to estimate the eigenvalues assuming $X(0) = X(1) = 0$. 
9 Finite Differences

9.1 Introduction

In previous chapters we introduced several methods to solve linear first and second order PDEs and quasilinear first order hyperbolic equations. There are many problems we cannot solve by those analytic methods. Such problems include quasilinear or nonlinear PDEs which are not hyperbolic. We should remark here that the method of characteristics can be applied to nonlinear hyperbolic PDEs. Even some linear PDEs, we cannot solve analytically. For example, Laplace’s equation

$$u_{xx} + u_{yy} = 0$$

(9.1.1)

inside a rectangular domain with a hole (see figure 51)

![Figure 51: Rectangular domain with a hole](image)

or a rectangular domain with one of the corners clipped off.

For such problems, we must use numerical methods. There are several possibilities, but here we only discuss finite difference schemes.

One of the first steps in using finite difference methods is to replace the continuous problem domain by a difference mesh or a grid. Let $f(x)$ be a function of the single independent variable $x$ for $a \leq x \leq b$. The interval $[a, b]$ is discretized by considering the nodes $a = x_0 < x_1 < \cdots < x_N < x_{N+1} = b$, and we denote $f(x_i)$ by $f_i$. The mesh size is $x_{i+1} - x_i$. 

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and we shall assume for simplicity that the mesh size is a constant

\[ h = \frac{b - a}{N + 1} \]  

and

\[ x_i = a + ih \quad i = 0, 1, \ldots, N + 1 \]  

In the two dimensional case, the function \( f(x, y) \) may be specified at nodal point \((x_i, y_j)\) by \( f_{ij} \). The spacing in the \( x \) direction is \( h_x \) and in the \( y \) direction is \( h_y \).

### 9.2 Difference Representations of PDEs

#### I. Truncation error

The difference approximations for the derivatives can be expanded in Taylor series. The truncation error is the difference between the partial derivative and its finite difference representation. For example

\[
\left. f_x \right|_{ij} = f_x \left( x_i, y_j \right) \approx f_i + \frac{1}{h_x} \Delta_x f_{ij} = f_x \left( x_i, y_j \right) - \frac{f_{i+1,j} - f_{i,j}}{h_x}
\]

\[
= -f_{xx} \left|_{ij} \right. \frac{h_x}{2!} - \cdots
\]

We use \( O(h_x) \) which means that the truncation error satisfies \( |T.E.| \leq K|h_x| \) for \( h_x \to 0 \), sufficiently small, where \( K \) is a positive real constant. Note that \( O(h_x) \) does not tell us the exact size of the truncation error. If another approximation has a truncation error of \( O(h_x^2) \), we might expect that this would be smaller **only** if the mesh is sufficiently fine.

We define the order of a method as the lowest power of the mesh size in the truncation error. Thus Table 1 (Chapter 8) gives first through fourth order approximations of the first derivative of \( f \).

The truncation error for a finite difference approximation of a given PDE is defined as the difference between the two. For example, if we approximate the advection equation

\[
\frac{\partial F}{\partial t} + c \frac{\partial F}{\partial x} = 0, \quad c > 0
\]

by centered differences

\[
\frac{F_{i,j+1} - F_{i,j-1}}{2\Delta t} + c \frac{F_{i+1,j} - F_{i-1,j}}{2\Delta x} = 0
\]

then the truncation error is

\[
T.E. = \left( \frac{\partial F}{\partial t} + c \frac{\partial F}{\partial x} \right)_{ij} - \frac{F_{i,j+1} - F_{i,j-1}}{2\Delta t} - c \frac{F_{i+1,j} - F_{i-1,j}}{2\Delta x}
\]

\[
= -\frac{1}{6} \Delta t^2 \frac{\partial^3 F}{\partial t^3} - c \frac{1}{6} \Delta x^2 \frac{\partial^3 F}{\partial x^3} - \text{higher powers of } \Delta t \text{ and } \Delta x.
\]
We will write
\[ T.E. = O(\Delta t^2, \Delta x^2) \]  
(9.2.6)

In the case of the simple explicit method
\[ \frac{u_j^{n+1} - u_j^n}{\Delta t} = k \frac{u_{j+1}^{n+1} - 2u_j^n + u_{j-1}^n}{(\Delta x)^2} \]  
(9.2.7)
for the heat equation
\[ u_t = ku_{xx} \]  
(9.2.8)
one can show that the truncation error is
\[ T.E. = O(\Delta t, \Delta x^2) \]  
(9.2.9)
since the terms in the finite difference approximation (9.2.7) can be expanded in Taylor series to get
\[ u_t - ku_{xx} + u_{tt} \frac{\Delta t}{2} - ku_{xxx} \frac{(\Delta x)^2}{12} + \cdots \]
All the terms are evaluated at \( x_j, t_n \). Note that the first two terms are the PDE and all other terms are the truncation error. Of those, the ones with the lowest order in \( \Delta t \) and \( \Delta x \) are called the leading terms of the truncation error.

Remark: See lab3 (3243taylor.ms) for the use of Maple to get the truncation error.

II. Consistency
A difference equation is said to be consistent or compatible with the partial differential equation when it approaches the latter as the mesh sizes approaches zero. This is equivalent to
\[ T.E. \to 0 \quad \text{as mesh sizes} \quad \to 0. \]
This seems obviously true. One can mention an example of an inconsistent method (see e.g. Smith (1985)). The DuFort-Frankel scheme for the heat equation (9.2.8) is given by
\[ \frac{u_j^{n+1} - u_j^{n-1}}{2\Delta t} = k \frac{u_{j+1}^{n+1} - u_j^{n+1} - u_j^{n-1} + u_{j-1}^n}{\Delta x^2}. \]  
(9.2.10)
The truncation error is
\[ \frac{k}{12} \frac{\partial^4 u^n}{\partial x^4} \Delta x^2 - \frac{\partial^2 u^n}{\partial t^2} \frac{\Delta t}{\Delta x} \left( \frac{\Delta t}{\Delta x} \right)^2 - \frac{1}{6} \frac{\partial^3 u^n}{\partial t^3} \left( \frac{\Delta t}{\Delta x} \right)^2 + \cdots \]  
(9.2.11)
If \( \Delta t, \Delta x \) approach zero at the same rate such that \( \frac{\Delta t}{\Delta x} = \text{constant} = \beta \), then the method is inconsistent (we get the PDE
\[ u_t + \beta^2 u_{tt} = ku_{xx} \]
instead of (9.2.8).)

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III. Stability

A numerical scheme is called stable if errors from any source (e.g., truncation, round-off, errors in measurements) are not permitted to grow as the calculation proceeds. One can show that DuFort-Frankel scheme is unconditionally stable. Richtmeyer and Morton give a less stringent definition of stability. A scheme is stable if its solution remains a uniformly bounded function of the initial state for all sufficiently small $\Delta t$.

The problem of stability is very important in numerical analysis. There are two methods for checking the stability of linear difference equations. The first one is referred to as Fourier or von Neumann assumes the boundary conditions are periodic. The second one is called the matrix method and takes care of contributions to the error from the boundary.

von Neumann analysis

Suppose we solve the heat equation (9.2.8) by the simple explicit method (9.2.7). If a term (a single term of Fourier and thus the linearity assumption)

$$e_j^n = e^{at_n} e^{ik_m x_j}$$

is substituted into the difference equation, one obtains after dividing through by $e^{at_n} e^{ik_m x_j}$

$$e^{\Delta t} = 1 + 2r (\cos \beta - 1) = 1 - 4r \sin^2 \frac{\beta}{2}$$

where

$$r = k \frac{\Delta t}{(\Delta x)^2}$$

$$\beta = k_m \Delta x, \quad k_m = \frac{2\pi m}{2L}, m = 0, \ldots, M,$$

where $M$ is the number of $\Delta x$ units contained in $L$. The stability requirement is

$$|e^{\Delta t}| \leq 1$$

implies

$$r \leq \frac{1}{2}.$$  

The term $|e^{\Delta t}|$ also denoted $G$ is called the amplification factor. The simple explicit method is called conditionally stable, since we had to satisfy the condition (9.2.17) for stability.

One can show that the simple implicit method for the same equation is unconditionally stable. Of course the price in this case is the need to solve a system of equations at every time step. The following method is an example of an unconditionally unstable method:

$$u_j^{n+1} - u_j^{n-1} \frac{2 \Delta t}{2 \Delta t} = k \frac{u_j^{n+1} - 2u_j^n + u_j^{n-1}}{(\Delta x)^2}.$$  

This method is second order in time and space but useless. The DuFort Frankel is a way to stabilize this second order in time scheme.

IV. Convergence
A scheme is called convergent if the solution to the finite difference equation approaches the exact solution to the PDE with the same initial and boundary conditions as the mesh sizes approach zero. Lax has proved that under appropriate conditions a consistent scheme is convergent if and only if it is stable.

Lax equivalence theorem
Given a properly posed linear initial value problem and a finite difference approximation to it that satisfies the consistency condition, stability (a-la Richtmeyer and Morton (1967)) is the necessary and sufficient condition for convergence.

V. Modified Equation
The importance of the modified equation is in helping to analyze the numerical effects of the discretization. The way to obtain the modified equation is by starting with the truncation error and replacing the time derivatives by spatial differentiation using the equation obtained from truncation error. It is easier to discuss the details on an example. For the heat equation

\[
\frac{u_{j}^{n+1} - u_{j}^{n}}{\Delta t} - k \frac{u_{j+1}^{n} - 2u_{j}^{n} + u_{j-1}^{n}}{(\Delta x)^{2}} = 0.
\] (9.2.19)

The truncation error is (all terms are given at \(t_n, x_j\))

\[
u_t - k u_{xx} = -\frac{\Delta t}{2} u_{tt} + \frac{(\Delta x)^2}{12} k u_{xxxx} \pm \ldots
\] (9.2.20)

This is the equation we have to use to eliminate the time derivatives. After several differentiations and substitutions, we get

\[
u_t - k u_{xx} = \left[ -\frac{1}{2} k^2 \Delta t + k \frac{(\Delta x)^2}{12} \right] u_{xxxx} + \left[ \frac{1}{3} k^3 (\Delta t)^2 - \frac{1}{12} k^2 (\Delta t) (\Delta x)^2 + \frac{1}{360} k (\Delta x)^4 \right] u_{xxxxx} + \ldots
\]

It is easier to organize the work in a tabular form. We will show that later when discussing first order hyperbolic.

Note that for \(r = \frac{1}{6}\), the truncation error is \(O(\Delta t^2, \Delta x^4)\). The problem is that one has to do 3 times the number of steps required by the limit of stability, \(r = \frac{1}{2}\).

Note also there are NO odd derivative terms, that is no dispersive error (dispersion means that phase relation between various waves are distorted, or the same as saying that the amplification factor has no imaginary part.)

Note that the exact amplification can be obtained as the quotient

\[
G_{\text{exact}} = \frac{u(t + \Delta t, x)}{u(t, x)} = e^{-r \beta^2}
\] (9.2.21)

See figure 53 for a plot of the amplification factor \(G\) versus \(\beta\).
Figure 53: Amplification factor for simple explicit method

Problems

1. Utilize Taylor series expansions about the point \((n + \frac{1}{2}, j)\) to determine the T.E. of the Crank Nicolson representation of the heat equation. Compare these results with the T.E. obtained from Taylor series expansion about the point \((n, j)\).

2. The DuFort Frankel method for solving the heat equation requires solution of the difference equation

\[
\frac{u_{j}^{n+1} - u_{j}^{n-1}}{2 \Delta t} = \frac{\alpha}{(\Delta x)^2} \left( u_{j+1}^{n} - u_{j-1}^{n} - u_{j}^{n+1} + u_{j}^{n-1} \right)
\]

Develop the stability requirements necessary for the solution of this equation.
9.3 Heat Equation in One Dimension

In this section we apply finite differences to obtain an approximate solution of the heat equation in one dimension,

\[ u_t = ku_{xx}, \quad 0 < x < 1, \quad t > 0, \quad (9.3.1) \]

subject to the initial and boundary conditions

\[ u(x,0) = f(x), \quad (9.3.2) \]
\[ u(0,t) = u(1,t) = 0. \quad (9.3.3) \]

Using forward approximation for \( u_t \) and centered differences for \( u_{xx} \) we have

\[ u^n_j + 1 - u^n_j = \frac{k}{\Delta x^2} (u^n_{j-1} - 2u^n_j + u^n_{j+1}), \quad j = 1, 2, \ldots, N - 1, \quad n = 0, 1, \ldots \quad (9.3.4) \]

where \( u^n_j \) is the approximation to \( u(x_j, t_n) \), the nodes \( x_j, t_n \) are given by

\[ x_j = j\Delta x, \quad j = 0, 1, \ldots, N \quad (9.3.5) \]
\[ t_n = n\Delta t, \quad n = 0, 1, \ldots \quad (9.3.6) \]

and the mesh spacing

\[ \Delta x = \frac{1}{N}, \quad (9.3.7) \]

see figure 54.

![Figure 54: Uniform mesh for the heat equation](image)

The solution at the points marked by * is given by the initial condition

\[ u^0_j = u(x_j, 0) = f(x_j), \quad j = 0, 1, \ldots, N \quad (9.3.8) \]
and the solution at the points marked by $\oplus$ is given by the boundary conditions

$$u(0, t_n) = u(x_N, t_n) = 0,$$

or

$$u^n_0 = u^n_N = 0. \quad (9.3.9)$$

The solution at other grid points can be obtained from (9.3.4)

$$u_{j+1}^{n+1} = r u_j^n + (1 - 2r) u_j^n + r u_{j+1}^n,$$  \( (9.3.10) \)

where $r$ is given by (9.2.14). The implementation of (9.3.10) is easy. The value at any grid point requires the knowledge of the solution at the three points below. We describe this by the following computational molecule (figure 55).

![Figure 55: Computational molecule for explicit solver](image)

We can compute the solution at the leftmost grid point on the horizontal line representing $t_1$ and continue to the right. Then we can advance to the next horizontal line representing $t_2$ and so on. Such a scheme is called explicit.

The time step $\Delta t$ must be chosen in such a way that stability is satisfied, that is

$$\Delta t \leq \frac{k}{2} (\Delta x)^2. \quad (9.3.11)$$

We will see in the next sections how to overcome the stability restriction and how to obtain higher order method.

[Can do Lab 5]
Problems

1. Use the simple explicit method to solve the 1-D heat equation on the computational grid (figure 56) with the boundary conditions

\[ u^n_1 = 2 = u^n_3 \]

and initial conditions

\[ u^1_1 = 2 = u^1_3, \quad u^1_2 = 1. \]

Show that if \( r = \frac{1}{4} \), the steady state value of \( u \) along \( j = 2 \) becomes

\[ u_{2\text{steady state}} = \lim_{n \to \infty} \sum_{k=1}^{n} \frac{1}{2^{k-1}} \]

Note that this infinite series is geometric that has a known sum.

![Figure 56: domain for problem 1 section 9.3](image-url)
9.3.1 Implicit method

One of the ways to overcome this restriction is to use an implicit method

\[ u_j^{n+1} - u_j^n = k \frac{\Delta t}{(\Delta x)^2} (u_{j-1}^{n+1} - 2u_j^{n+1} + u_{j+1}^{n+1}), \quad j = 1, 2, \ldots, N-1, \quad n = 0, 1, \ldots \]  

(9.3.1.1)

The computational molecule is given in figure 57. The method is unconditionally stable, since the amplification factor is given by

\[ G = \frac{1}{1 + 2r(1 - \cos \beta)} \]  

(9.3.1.2)

which is \( \leq 1 \) for any \( r \). The price for this is having to solve a tridiagonal system for each time step. The method is still first order in time. See figure 58 for a plot of \( G \) for explicit and implicit methods.

Figure 57: Computational molecule for implicit solver

9.3.2 DuFort Frankel method

If one tries to use centered difference in time and space, one gets an unconditionally unstable method as we mentioned earlier. Thus to get a stable method of second order in time, DuFort Frankel came up with:

\[ \frac{u_j^{n+1} - u_j^{n-1}}{2\Delta t} = k \frac{u_{j+1}^n - u_j^{n+1} - u_j^{n-1} + u_{j-1}^n}{\Delta x^2} \]  

(9.3.2.1)

We have seen earlier that the method is explicit with a truncation error

\[ T.E. = O \left( \Delta t^2, \Delta x^2, \left( \frac{\Delta t}{\Delta x} \right)^2 \right). \]  

(9.3.2.2)

The modified equation is
The amplification factor is given by

\[ G = \frac{2r \cos \beta \pm \sqrt{1 - 4r^2 \sin^2 \beta}}{1 + 2r} \]  

and thus the method is unconditionally stable.

The only drawback is the requirement of an additional starting line.

### 9.3.3 Crank-Nicolson method

Another way to overcome this stability restriction, we can use Crank-Nicolson implicit scheme

\[ -ru_{j-1}^{n+1} + 2(1 + r)u_j^{n+1} - ru_{j+1}^{n+1} = ru_{j-1}^n + 2(1 - r)u_j^n + ru_{j+1}^n. \]  

This is obtained by centered differencing in time about the point \( x_j, t_{n+1/2} \). On the right we average the centered differences in space at time \( t_n \) and \( t_{n+1} \). The computational molecule is now given in the next figure (59).

The method is unconditionally stable, since the denominator is always larger than numerator in

\[ G = \frac{1 - r(1 - \cos \beta)}{1 + r(1 - \cos \beta)}. \]
It is second order in time (centered difference about \( x_j, t_{n+1/2} \)) and space. The modified equation is

\[
    u_t - ku_{xx} = \frac{k\Delta x^2}{12} u_{xxxx} + \left[ \frac{1}{12} k^3 \Delta t^2 + \frac{1}{360} k \Delta x^4 \right] u_{xxxxx} + \ldots
\]  

(9.3.3.3)

The disadvantage of the implicit scheme (or the price we pay to overcome the stability barrier) is that we require a solution of system of equations at each time step. The number of equations is \( N - 1 \).

We include in the appendix a Fortran code for the solution of (9.3.1)-(9.3.3) using the explicit and implicit solvers. We must say that one can construct many other explicit or implicit solvers. We allow for the more general boundary conditions

\[
    A_L u_x + B_L u = C_L, \quad \text{on the left boundary} \tag{9.3.3.4}
\]

\[
    A_R u_x + B_R u = C_R, \quad \text{on the right boundary}. \tag{9.3.3.5}
\]

Remark: For a more general boundary conditions, see for example Smith (1985), we need to finite difference the derivative in the boundary conditions.

### 9.3.4 Theta (\( \theta \)) method

All the method discussed above (except DuFort Frankel) can be written as

\[
    \frac{u_j^{n+1} - u_j^n}{\Delta t} = k \left[ \theta \left( u_{j+1}^{n+1} - 2u_j^{n+1} + u_{j-1}^{n+1} \right) + (1 - \theta) \left( u_{j+1}^n - 2u_j^n + u_{j-1}^n \right) \right] \quad \Delta x^2
\]

(9.3.4.1)

For \( \theta = 0 \) we get the explicit method (9.3.10), for \( \theta = 1 \), we get the implicit method (9.3.1.1) and for \( \theta = \frac{1}{2} \) we have Crank Nicolson (9.3.3.1).

The truncation error is

\[
    O \left( \Delta t, \Delta x^2 \right)
\]
except for Crank Nicolson as we have seen earlier (see also the modified equation below.) If one chooses \( \theta = \frac{1}{2} - \frac{\Delta x^2}{12k\Delta t} \) (the coefficient of \( u_{xxxx} \) vanishes), then we get \( O(\Delta t^2, \Delta x^4) \), and if we choose the same \( \theta \) with \( \frac{\Delta x^2}{k\Delta t} = \sqrt{20} \) (the coefficient of \( u_{xxxxx} \) vanishes), then \( O(\Delta t^2, \Delta x^6) \).

The method is conditionally stable for \( 0 \leq \theta < \frac{1}{2} \) with the condition

\[
r \leq \frac{1}{2 - 4\theta}
\]

and unconditionally stable for \( \frac{1}{2} \leq \theta \leq 1 \).

The modified equation is

\[
\begin{align*}
    u_t - ku_{xx} &= \left( \frac{1}{12} k \Delta x^2 + (\theta - \frac{1}{2})k^2 \Delta t \right) u_{xxxx} \\
    &+ \left[ (\theta^2 - \theta + \frac{1}{3})k^3 \Delta t^2 + \frac{1}{6}(\theta - \frac{1}{2})k^2 \Delta t \Delta x^2 + \frac{1}{360} k \Delta x^4 \right] u_{xxxxx} + \ldots
\end{align*}
\]

(9.3.4.3)

**9.3.5 An example**

We have used the explicit solver program to approximate the solution of

\[
u_t = u_{xx}, \quad 0 < x < 1, \quad t > 0
\]

(9.3.5.1)

\[
u(x, 0) = \begin{cases} 
2x & 0 < x < \frac{1}{2} \\
2(1 - x) & \frac{1}{2} < x < 1 
\end{cases}
\]

(9.3.5.2)

\[
u(0, t) = u(1, t) = 0,
\]

(9.3.5.3)

using a variety of values of \( r \). The results are summarized in the following figures.

The analytic solution (using separation of variables) is given by

\[
u(x, t) = \sum_{n=0}^{\infty} a_n e^{-(n\pi)^2 t} \sin n\pi x,
\]

(9.3.5.4)

where \( a_n \) are the Fourier coefficients for the expansion of the initial condition (9.3.5.2),

\[
a_n = \frac{8}{(n\pi)^2} \sin \frac{n\pi}{2}, \quad n = 1, 2, \ldots
\]

(9.3.5.5)

The analytic solution (9.3.5.4) and the numerical solution (using \( \Delta x = .1, r = .5 \)) at times \( t = .025 \) and \( t = .5 \) are given in the two figures 60, 61. It is clear that the error increases in time but still smaller than \( .5 \times 10^{-4} \).
Figure 60: Numerical and analytic solution with $r = .5$ at $t = .025$

Figure 61: Numerical and analytic solution with $r = .5$ at $t = .5$

On the other hand, if $r = .51$, we see oscillations at time $t = .0255$ (figure 62) which become very large at time $t = .255$ (figure 63) and the temperature becomes negative at $t = .459$ (figure 64).

Clearly the solution does not converge when $r > .5$.

The implicit solver program was used to approximate the solution of (9.3.5.1) subject to

$$u(x, 0) = 100 - 10|x - 10|$$  \hspace{1cm} (9.3.5.6)

and

$$u_x(0, t) = .2(u(0, t) - 15),$$  \hspace{1cm} (9.3.5.7)

$$u(1, t) = 100.$$  \hspace{1cm} (9.3.5.8)

Notice that the boundary and initial conditions do not agree at the right boundary. Because of the type of boundary condition at $x = 0$, we cannot give the eigenvalues explicitly. Notice
that the problem is also having inhomogeneous boundary conditions. To be able to compare the implicit and explicit solvers, we have used Crank-Nicolson to solve (9.3.5.1)-(9.3.5.3). We plot the analytic and numerical solution with \( r = 1 \) at time \( t = .5 \) to show that the method is stable (compare the following figure 65 to the previous one with \( r = .51 \)).
Figure 63: Numerical and analytic solution with $r = .51$ at $t = .255$

Figure 64: Numerical and analytic solution with $r = .51$ at $t = .459$
Figure 65: Numerical (implicit) and analytic solution with $r = 1. \text{ at } t = .5$
9.3.6 Unbounded Region - Coordinate Transformation

Suppose we have to solve a problem on unbounded domain, e.g.

\[ u_t = u_{xx}, \quad 0 \leq x < \infty \]

subject to

\[ u(x, 0) = g(x) \]
\[ u(0, t) = f(t). \]

There is no difficulty with the unbounded domain if we use one sided approximation for \( u_{xx} \), i.e.

\[ u_{xx} = \frac{u_i^n - 2u_{i-1}^n + u_{i-2}^n}{h^2} \]

which is first order accurate. If one decides to use second order centered differences then an unclosed set of equations are obtained (always need a point to the right). The most obvious way to overcome this is to impose a boundary condition at an artificial boundary \( x = L \), such as

\[ u(L, t) = 0. \]

Another way is to transform the domain to a finite interval, say \([0, 1]\) by using one of these transformations:

\[ z = 1 - e^{-x/L}, \]

or

\[ z = \frac{x}{x + L}, \]

for some scale factor \( L \). This, of course, will affect the equation.

9.4 Two Dimensional Heat Equation

In this section, we generalize the solution of the heat equation obtained in section 9.3 to two dimensions. The problem of heat conduction in a rectangular membrane is described by

\[ u_t = \alpha(u_{xx} + u_{yy}), \quad 0 < x < L, \quad 0 < y < H, \quad t > 0 \]

subject to

\[ u(x, y, t) = g(x, y, t), \quad \text{on the boundary} \]
\[ u(x, y, 0) = f(x, y), \quad 0 < x < L, \quad 0 < y < H. \]
9.4.1 Explicit

To obtain an explicit scheme, we use forward difference in time and centered differences in space. Thus

\[
\frac{u_{ij}^{n+1} - u_{ij}^n}{\Delta t} = \alpha \left( \frac{u_{i-1,j}^n - 2u_{ij}^n + u_{i+1,j}^n}{(\Delta x)^2} + \frac{u_{ij-1}^n - 2u_{ij}^n + u_{ij+1}^n}{(\Delta y)^2} \right)
\]

(9.4.1.1)

or

\[
u_{ij}^{n+1} = r_x u_{i-1,j}^n + (1 - 2r_x - 2r_y) u_{ij}^n + r_x u_{i+1,j}^n + r_y u_{ij-1}^n + r_y u_{ij+1}^n,
\]

(9.4.1.2)

where \( u_{ij}^n \) is the approximation to \( u(x_i, y_j, t_n) \) and

\[ r_x = \alpha \frac{\Delta t}{(\Delta x)^2}, \]

(9.4.1.3)

\[ r_y = \alpha \frac{\Delta t}{(\Delta y)^2}. \]

(9.4.1.4)

The stability condition imposes a limit on the time step

\[
\alpha \Delta t \left( \frac{1}{(\Delta x)^2} + \frac{1}{(\Delta y)^2} \right) \leq \frac{1}{2}
\]

(9.4.1.5)

For the case \( \Delta x = \Delta y = d \), we have

\[
\Delta t \leq \frac{1}{4\alpha} d^2
\]

(9.4.1.6)

which is more restrictive than in the one dimensional case. The solution at any point \( (x_i, y_j, t_n) \) requires the knowledge of the solution at all 5 points at the previous time step (see next figure 66).

![Diagram](image)  

Figure 66: Computational molecule for the explicit solver for 2D heat equation

Since the solution is known at \( t = 0 \), we can compute the solution at \( t = \Delta t \) one point at a time.

To overcome the stability restriction, we can use Crank-Nicolson implicit scheme. The matrix in this case will be banded of higher dimension and wider band. There are other implicit schemes requiring solution of smaller size systems, such as alternating direction. In the next section we will discuss Crank-Nicolson and ADI (Alternating Direction Implicit).
9.4.2 Crank Nicolson

One way to overcome this stability restriction is to use Crank-Nicolson implicit scheme

\[
\frac{u_{ij}^{n+1} - u_{ij}^n}{\Delta t} = \alpha \frac{\delta_x^2 u_{ij}^n + \delta_y^2 u_{ij}^{n+1}}{2(\Delta x)^2} + \alpha \frac{\delta_x^2 u_{ij}^n + \delta_y^2 u_{ij}^{n+1}}{2(\Delta y)^2}
\]  

(9.4.2.1)

The method is unconditionally stable. It is second order in time (centered difference about \(x_i, y_j, t_{n+1/2}\)) and space.

It is important to order the two subscript in one dimensional index in the right direction (if the number of grid point in \(x\) and \(y\) is not identical), otherwise the bandwidth will increase.

Note that the coefficients of the banded matrix are independent of time (if \(\alpha\) is not a function of \(t\)), and thus one have to factor the matrix only once.

9.4.3 Alternating Direction Implicit

The idea here is to alternate direction and thus solve two one-dimensional problem at each time step. The first step to keep \(y\) fixed

\[
\frac{u_{ij}^{n+1/2} - u_{ij}^n}{\Delta t/2} = \alpha \left( \delta_x^2 u_{ij}^{n+1/2} + \delta_y^2 u_{ij}^n \right)
\]

(9.4.3.1)

In the second step we keep \(x\) fixed

\[
\frac{u_{ij}^{n+1/2} - u_{ij}^{n+1/2}}{\Delta t/2} = \alpha \left( \delta_x^2 u_{ij}^{n+1/2} + \delta_y^2 u_{ij}^{n+1/2} \right)
\]

(9.4.3.2)

So we have a tridiagonal system at every step. We have to order the unknown differently at every step.

The method is second order in time and space and it is unconditionally stable, since the denominator is always larger than numerator in

\[
G = \frac{1 - r_x (1 - \cos \beta_x)}{1 + r_x (1 - \cos \beta_x)} \frac{1 - r_y (1 - \cos \beta_y)}{1 + r_y (1 - \cos \beta_y)}
\]

(9.4.3.3)

The obvious extension to three dimensions is only first order in time and conditionally stable. Douglas & Gunn developed a general scheme called approximate factorization to ensure second order and unconditional stability.

Let

\[
\Delta u_{ij} = u_{ij}^{n+1} - u_{ij}^n
\]

(9.4.3.4)

Substitute this into the two dimensional Crank Nicolson

\[
\Delta u_{ij} = \frac{\alpha \Delta t}{2} \left\{ \delta_x^2 \Delta u_{ij} + \delta_y^2 \Delta u_{ij} + 2\delta_x^2 u_{ij}^n + 2\delta_y^2 u_{ij}^n \right\}
\]

(9.4.3.5)

Now rearrange,

\[
\left( 1 - \frac{r_x \delta^2}{2 \delta_x} - \frac{r_y \delta^2}{2 \delta_y} \right) \Delta u_{ij} = \left( r_x \delta^2 + r_y \delta^2 \right) u_{ij}^n
\]

(9.4.3.6)
The left hand side operator can be factored

\[ 1 - \frac{r_x}{2} \frac{\delta_x^2}{\delta_x^2} - \frac{r_y}{2} \frac{\delta_y^2}{\delta_y^2} = \left( 1 - \frac{r_x}{2} \frac{\delta_x^2}{\delta_x^2} \right) \left( 1 - \frac{r_y}{2} \frac{\delta_y^2}{\delta_y^2} \right) - \frac{r_x r_y}{4} \frac{\delta_x^2 \delta_y^2}{\delta_x^2 \delta_y^2} \]  

(9.4.3.7)

The last term can be neglected because it is of higher order. Thus the method for two dimensions becomes

\[ \left( 1 - \frac{r_x}{2} \frac{\delta_x^2}{\delta_x^2} \right) \Delta u_{ij}^* = \left( r_x \delta_x^2 + r_y \delta_y^2 \right) u_{ij}^n \]  

(9.4.3.8)

\[ \left( 1 - \frac{r_y}{2} \frac{\delta_y^2}{\delta_y^2} \right) \Delta u_{ij} = \Delta u_{ij}^* \]  

(9.4.3.9)

\[ u_{ij}^{n+1} = u_{ij}^n + \Delta u_{ij} \]  

(9.4.3.10)
Problems

1. Apply the ADI scheme to the 2-D heat equation and find $u^{n+1}$ at the internal grid points in the mesh shown in figure 67 for $r_x = r_y = 2$. The initial conditions are

$$u^n = 1 - \frac{x}{3\Delta x} \quad \text{along } y = 0$$

$$u^n = 1 - \frac{y}{2\Delta y} \quad \text{along } x = 0$$

$$u^n = 0 \quad \text{everywhere else}$$

and the boundary conditions remain fixed at their initial values.

Figure 67: domain for problem 1 section 9.4.2
9.4.4 Alternating Direction Implicit for Three Dimensional Problems

Here we extend Douglas & Gunn method to three dimensions

\[
\left(1 - \frac{r_x^2}{2\delta_x^2}\right) \Delta u_{ijk} = \left( r_x \delta_x^2 + r_y \delta_y^2 + r_z \delta_z^2 \right) u_{ijk}^n \tag{9.4.4.1}
\]

\[
\left(1 - \frac{r_y^2}{2\delta_y^2}\right) \Delta u_{ijk}^* = \Delta u_{ijk}^* \tag{9.4.4.2}
\]

\[
\left(1 - \frac{r_z^2}{2\delta_z^2}\right) \Delta u_{ijk} = \Delta u_{ijk}^{**} \tag{9.4.4.3}
\]

\[
u_{ijk}^{n+1} = u_{ijk}^n + \Delta u_{ijk}. \tag{9.4.4.4}
\]

9.5 Laplace’s Equation

In this section, we discuss the approximation of the steady state solution inside a rectangle

\[
u_{xx} + \nu_{yy} = 0, \quad 0 < x < L, \quad 0 < y < H, \tag{9.5.1}
\]

subject to Dirichlet boundary conditions

\[
u(x, y) = f(x, y), \quad \text{on the boundary.} \tag{9.5.2}
\]

![Uniform grid on a rectangle](image)

Figure 68: Uniform grid on a rectangle

We impose a uniform grid on the rectangle with mesh spacing \(\Delta x, \Delta y\) in the \(x, y\) directions, respectively. The finite difference approximation is given by

\[
\frac{u_{i-1,j} - 2u_{ij} + u_{i+1,j}}{(\Delta x)^2} + \frac{u_{ij-1} - 2u_{ij} + u_{ij+1}}{(\Delta y)^2} = 0, \tag{9.5.3}
\]
or
\[
\left[ \frac{2}{(\Delta x)^2} + \frac{2}{(\Delta y)^2} \right] u_{ij} = \frac{u_{i-1j} + u_{i+1j}}{(\Delta x)^2} + \frac{u_{ij-1} + u_{ij+1}}{(\Delta y)^2}.
\]
(9.5.4)

For $\Delta x = \Delta y$ we have
\[
4u_{ij} = u_{i-1j} + u_{i+1j} + u_{ij-1} + u_{ij+1}.
\]
(9.5.5)

The computational molecule is given in the next figure (69). This scheme is called five point star because of the shape of the molecule.

![Computational molecule for Laplace's equation](image)

Figure 69: Computational molecule for Laplace’s equation

The truncation error is
\[
T.E. = O \left( \Delta x^2, \Delta y^2 \right)
\]
and the modified equations is
\[
u_{xx} + u_{yy} = -\frac{1}{12} \left( \Delta x^2 u_{xxxx} + \Delta y^2 u_{yyyy} \right) + \cdots
\]
(9.5.7)

Remark: To obtain a higher order method, one can use the nine point star, which is of sixth order if $\Delta x = \Delta y = d$, but otherwise it is only second order. The nine point star is given by

\[
u_{i+1,j+1} + u_{i-1,j+1} + u_{i+1,j-1} + u_{i-1,j-1} - 2 \frac{\Delta x^2 - 5\Delta y^2}{\Delta x^2 + \Delta y^2} (u_{i+1,j} + u_{i-1,j})
\]
\[
+ 25 \frac{\Delta x^2 - \Delta y^2}{\Delta x^2 + \Delta y^2} (u_{i,j+1} + u_{i,j-1}) - 20u_{ij} = 0
\]
(9.5.8)

For three dimensional problem the equivalent to five point star is seven point star. It is given by

\[
u_{i-1,j,k} - 2u_{ij,k} + u_{i+1,j,k} + u_{i,j-1,k} - 2u_{ij,k} + u_{i,j+1,k} + u_{i,j,k-1} - 2u_{ij,k} + u_{ij,k+1} = 0.
\]
(9.5.9)

The solution is obtained by solving the linear system of equations
\[
Au = b
\]
(9.5.10)
where the block banded matrix $A$ is given by
\[
A = \begin{bmatrix}
  T & B & 0 & \cdots & 0 \\
  B & T & B \\
  0 & B & T & B \\
  \vdots \\
  0 & \cdots & 0 & B & T
\end{bmatrix}
\]
and the matrices $B$ and $T$ are given by
\[
B = -I
\]
\[
T = \begin{bmatrix}
  4 & -1 & 0 & \cdots & 0 \\
  -1 & 4 & -1 \\
  0 & -1 & 4 & -1 & 0 \\
  \vdots \\
  0 & \cdots & 0 & -1 & 4
\end{bmatrix}
\]
and the right hand side $b$ contains boundary values. If we have Poisson’s equation then $b$ will also contain the values of the right hand side of the equation evaluated at the center point of the molecule.

One can use Thomas algorithm for block tridiagonal matrices. The system could also be solved by an iterative method such as Jacobi, Gauss-Seidel or successive over relaxation (SOR). Such solvers can be found in many numerical analysis texts. In the next section, we give a little information on each.

Remarks:
1. The solution is obtained in one step since there is no time dependence.
2. One can use ELLPACK (ELLiptic PACKage, a research tool for the study of numerical methods for solving elliptic problems, see Rice and Boisvert (1984)) to solve any elliptic PDEs.

9.5.1 Iterative solution

The idea is to start with an initial guess for the solution and iterate using an easy system to solve. The sequence of iterates $x^{(i)}$ will converge to the answer under certain conditions on the iteration matrix. Here we discuss three iterative scheme. Let’s write the coefficient matrix $A$ as
\[
A = D - L - U
\]
then one can iterate as follows
\[
Dx^{(i+1)} = (L + U)x^{(i)} + b, \quad i = 0, 1, 2, \ldots
\]
This scheme is called Jacobi’s method. At each time step one has to solve a diagonal system. The convergence of the iterative procedure depends on the spectral radius of the iteration matrix
\[
J = D^{-1}(L + U).
\]
If \( \rho(J) < 1 \) then the iterative method converges (the speed depends on how small the spectral radius is. (spectral radius of a matrix is defined later and it relates to the modulus of the dominant eigenvalue.) If \( \rho(J) \geq 1 \) then the iterative method diverges.

Assuming that the new iterate is a better approximation to the answer, one comes up with Gauss-Seidel method. Here we suggest the use of the component of the new iterate as soon as they become available. Thus

\[
(D - L)x^{(i+1)} = Lx^{(i)} + b, \quad i = 0, 1, 2, \ldots
\]  

(9.5.1.4)

and the iteration matrix \( G \) is

\[
G = (D - L)^{-1}U
\]  

(9.5.1.5)

We can write Gauss Seidel iterative procedure also in componentwise

\[
x_k^{(i+1)} = \frac{1}{a_{kk}} \left( b_k - \sum_{j=1}^{k-1} a_{kj}x_j^{(i+1)} - \sum_{j=k+1}^{n} a_{kj}x_j^{(i)} \right)
\]  

(9.5.1.6)

It can be shown that if

\[
|a_{ii}| \geq \sum_{j \neq i} |a_{ij}|, \quad \text{for all } i
\]

and if for at least one \( i \) we have a strict inequality and the system is irreducible (i.e. can’t break to subsystems to be solved independently) then Gauss Seidel method converges. In the case of Laplace’s equation, these conditions are met.

The third method we mention here is called successive over relaxation or SOR for short. The method is based on Gauss-Seidel, but at each iteration we add a step

\[
u_{ij}^{(k+1)} = u_{ij}^{(k+1)} + \omega \left( u_{ij}^{(k+1)} - u_{ij}^{(k)} \right)
\]  

(9.5.1.7)

For \( 0 < \omega < 1 \) the method is really under relaxation. For \( \omega = 1 \) we have Gauss Seidel and for \( 1 < \omega < 2 \) we have over relaxation. There is no point in taking \( \omega \geq 2 \), because the method will diverge. It can be shown that for Laplace’s equation the best choice for \( \omega \) is

\[
\omega_{opt} = \frac{2}{1 + \sqrt{1 - \sigma^2}}
\]  

(9.5.1.8)

where

\[
\sigma = \frac{1}{1 + \beta^2} \left( \cos \frac{\pi}{p} + \beta^2 \cos \frac{\pi}{q} \right),
\]  

(9.5.1.9)

\[
\beta = \frac{\Delta x}{\Delta y}, \quad \text{grid aspect ratio}
\]  

(9.5.1.10)

and \( p, q \) are the number of \( \Delta x, \Delta y \) respectively.

Can do Lab 6
9.6 Vector and Matrix Norms

Norms have the following properties

Let

$$\vec{x}, \vec{y} \in \mathbb{R}^n \quad \vec{x} \neq \vec{0} \quad \alpha \in \mathbb{R}$$

1) \(\| \vec{x} \| > 0\)

2) \(\| \alpha \vec{x} \| = |\alpha| \| \vec{x} \|\)

3) \(\| \vec{x} + \vec{y} \| \leq \| \vec{x} \| + \| \vec{y} \|\)

Let \(\vec{x} = \begin{pmatrix} x_1 \\ x_2 \\ \vdots \\ x_n \end{pmatrix}\)

then the “integral” norms are:

\(\| \vec{x} \|_1 = \sum_{i=1}^{n} |x_i|\) one norm

\(\| \vec{x} \|_2 = \sqrt{\sum_{i=1}^{n} x_i^2}\) two norm (Euclidean norm)

\(\| \vec{x} \|_k = \left[ \sum_{i=1}^{n} |x_i|^k \right]^{1/k}\) k norm

\(\| \vec{x} \|_\infty = \max_{1 \leq i \leq n} |x_i|\) infinity norm

Example

\(\vec{x} = \begin{pmatrix} -3 \\ 4 \\ 5 \end{pmatrix}\)

\(\| \vec{x} \|_1 = 12\)

\(\| \vec{x} \|_2 = 5 \sqrt{2} \approx 7.071\)
Matrix Norms

Let $A$ be an $m \times n$ non-zero matrix (i.e. $A \in \mathbb{R}^m \times \mathbb{R}^n$). Matrix norms have the properties

1) $\| A \| \geq 0$
2) $\| \alpha A \| = |\alpha| \| A \|$
3) $\| A + B \| \leq \| A \| + \| B \|$

Definition

A matrix norm is consistent with vector norms $\| \cdot \|_a$ on $\mathbb{R}^n$ and $\| \cdot \|_b$ on $\mathbb{R}^m$ with $A \in \mathbb{R}^m \times \mathbb{R}^n$ if

$$\| A \vec{x} \|_b \leq \| A \| \| \vec{x} \|_a$$

and for the special case that $A$ is a square matrix

$$\| A \vec{x} \| \leq \| A \| \| \vec{x} \|$$

Definition

Given a vector norm, a corresponding matrix norm for square matrices, called the subordinate matrix norm is defined as

$$l. u. b. (A) = \max_{\vec{x} \neq 0} \left\{ \frac{\| A \vec{x} \|}{\| \vec{x} \|} \right\}$$

least upper bound

Note that this matrix norm is consistent with the vector norm because

$$\| A \vec{x} \| \leq l. u. b. (A) \cdot \| \vec{x} \|$$

by definition. Said another way, the $l. u. b. (A)$ is a measure of the greatest magnification a vector $\vec{x}$ can obtain, by the linear transformation $A$, using the vector norm $\| \cdot \|$.

Examples

For $\| \cdot \|_\infty$ the subordinate matrix norm is
\[ l.u.b_{\infty} (A) = \max_{x \neq 0} \frac{\| A \vec{x} \|_{\infty}}{\| \vec{x} \|_{\infty}} \]

\[ = \max_{x \neq 0} \left\{ \frac{\max_i \{ | \sum_{k=1}^{n} a_{ik} x_k | \} }{\max_k \{ | x_k | \} } \right\} \]

\[ = \max_i \{ \sum_{k=1}^{n} | a_{ik} | \} \]

where in the last equality, we’ve chosen \( x_k = \text{sign}(a_{ik}) \). The “inf”-norm is sometimes written

\[ \| A \|_{\infty} = \max_{1 \leq i \leq n} \sum_{j=1}^{n} | a_{ij} | \]

where it is readily seen to be the maximum row sum.

In a similar fashion, the “one”-norm of a matrix can be found, and is sometimes referred to as the column norm, since for a given \( m \times n \) matrix \( A \) it is

\[ \| A \|_1 = \max_{1 \leq j \leq n} \{ |a_{1j}| + |a_{2j}| + \cdots + |a_{mj}| \} \]

For \( \| \cdot \|_2 \) we have

\[ l.u.b_{2} (A) = \max_{x \neq 0} \frac{\| A \vec{x} \|_2}{\| \vec{x} \|_2} \]

\[ = \max_{x \neq 0} \sqrt{\frac{\vec{x}^T A^T A \vec{x}}{\vec{x}^T \vec{x}}} = \sqrt{\lambda_{\max} (A^T A)} \]

\[ = \sqrt{\rho(A^T A)} \]

where \( \lambda_{\max} \) is the magnitude of the largest eigenvalue of the symmetric matrix \( A^T A \), and where the notation \( \rho(A^T A) \) is referred to as the “spectral radius” of \( A^T A \). Note that if \( A = A^T \) then

\[ l.u.b_{2}(A) = \| A \|_2 = \sqrt{\rho^2(A)} = \rho(A) \]

The spectral radius of a matrix is smaller than any consistent matrix norm of that matrix. Therefore, the largest (in magnitude) eigenvalue of a matrix is the least upper bound of all consistent matrix norms. In mathematical terms,

\[ l.u.b_{\cdot} (\| A \|) = | \lambda_{\max} | = \rho(A) \]

where \( \| \cdot \| \) is any consistent matrix norm.
To see this, let \((\lambda_i, \vec{x}_i)\) be an eigenvalue/eigenvector pair of the matrix \(A\). Then we have
\[
A \vec{x}_i = \lambda_i \vec{x}_i
\]
Taking consistent matrix norms,
\[
\|A \vec{x}_i\| = \|\lambda_i \vec{x}_i\| = |\lambda_i| \|\vec{x}_i\|
\]
Because \(\|\cdot\|\) is a consistent matrix norm
\[
\|A\| \|\vec{x}_i\| \geq \|A \vec{x}_i\| = |\lambda_i| \|\vec{x}_i\|
\]
and dividing out the magnitude of the eigenvector (which must be other than zero), we have
\[
\|A\| \geq |\lambda_i| \quad \text{for all } \lambda_i
\]

**Example** Given the matrix

\[
A = \begin{pmatrix}
-12 & 4 & 3 & 2 & 1 \\
2 & 10 & 1 & 5 & 1 \\
3 & 3 & 21 & -5 & -4 \\
1 & -1 & 2 & 12 & -3 \\
5 & 5 & -3 & -2 & 20
\end{pmatrix}
\]

we can determine the various norms of the matrix \(A\).

The 1 norm of \(A\) is given by:
\[
\|A\|_1 = \max_j \{ |a_{1,j}| + |a_{2,j}| + \ldots + |a_{5,j}| \}
\]

The matrix \(A\) can be seen to have a 1-norm of 30 from the 3\(^{rd}\) column.

The \(\infty\) norm of \(A\) is given by:
\[
\|A\|_\infty = \max_i \{ |a_{i,1}| + |a_{i,2}| + \ldots + |a_{i,5}| \}
\]

and therefore has the \(\infty\) norm of 36 which comes from its 3\(^{rd}\) row.

To find the “two”-norm of \(A\), we need to find the eigenvalues of \(A^T A\) which are:

52.3239, 157.9076, 211.3953, 407.6951, and 597.6781

Taking the square root of the largest eigenvalue gives us the 2 norm: \(\|A\|_2 = 24.4475\).

To determine the spectral radius of \(A\), we find that \(A\) has the eigenvalues:

\(-12.8462, 9.0428, 12.9628, 23.0237, \text{ and } 18.8170\)

Therefore the spectral radius of \(A\), (or \(\rho(A)\)) is 23.0237, which is in fact less than all other norms of \(A\) (\(\|A\|_1 = 30, \|A\|_2 = 24.4475, \|A\|_\infty = 36\)).
Problems

1. Find the one-, two-, and infinity norms of the following vectors and matrices:

   \[
   \begin{align*}
   & (a) \begin{pmatrix} 1 & 2 & 3 \\ 2 & 5 & 6 \\ 3 & 6 & 9 \end{pmatrix} & (b) \begin{pmatrix} 3 \\ 4 \\ 5 \end{pmatrix} & (c) \begin{pmatrix} 1 & 6 \\ 7 & 3 \end{pmatrix}
   \end{align*}
   \]
9.7 Matrix Method for Stability

We demonstrate the matrix method for stability on two methods for solving the one-dimensional heat equation. Recall that the explicit method can be written in matrix form as

\[ u^{n+1} = Au^n + b \]  \hspace{1cm} (9.7.1)

where the tridiagonal matrix \( A \) have \( 1-2r \) on diagonal and \( r \) on the super- and sub-diagonal. The norm of the matrix dictates how fast errors are growing (the vector \( b \) doesn’t come into play). If we check the infinity or 1 norm we get

\[ ||A||_1 = ||A||_\infty = |1-2r| + |r| + |r| \]  \hspace{1cm} (9.7.2)

For \( 0 < r \leq 1/2 \), all numbers inside the absolute values are non negative and we get a norm of 1. For \( r > 1/2 \), the norms are \( 4r - 1 \) which is greater than 1. Thus we have conditional stability with the condition \( 0 < r \leq 1/2 \).

The Crank Nicolson scheme can be written in matrix form as follows

\[ (2I - rT)u^{n+1} = (2I + rT)u^n + b \]  \hspace{1cm} (9.7.3)

where the tridiagonal matrix \( T \) has -2 on diagonal and 1 on super- and sub-diagonals. The eigenvalues of \( T \) can be expressed analytically, based on results of section 8.6,

\[ \lambda_s(T) = -4 \sin^2 \frac{s\pi}{2N}, \quad s = 1, 2, \ldots, N - 1 \]  \hspace{1cm} (9.7.4)

Thus the iteration matrix is

\[ A = (2I - rT)^{-1}(2I + rT) \]  \hspace{1cm} (9.7.5)

for which we can express the eigenvalues as

\[ \lambda_s(A) = \frac{2 - 4r \sin^2 \frac{s\pi}{2N}}{2 + 4r \sin^2 \frac{s\pi}{2N}} \]  \hspace{1cm} (9.7.6)

All the eigenvalues are bounded by 1 since the denominator is larger than numerator. Thus we have unconditional stability.

9.8 Derivative Boundary Conditions

Derivative boundary conditions appear when a boundary is insulated

\[ \frac{\partial u}{\partial n} = 0 \]  \hspace{1cm} (9.8.1)

or when heat is transferred by radiation into the surrounding medium (whose temperature is \( v \))

\[ -k \frac{\partial u}{\partial n} = H(u - v) \]  \hspace{1cm} (9.8.2)
where $H$ is the coefficient of surface heat transfer and $k$ is the thermal conductivity of the material.

Here we show how to approximate these two types of boundary conditions in connection with the one dimensional heat equation

$$u_t = ku_{xx}, \quad 0 < x < 1$$

(9.8.3)

$$u(0, t) = g(t)$$

(9.8.4)

$$\frac{\partial u(1, t)}{\partial n} = -h(u(1, t) - v)$$

(9.8.5)

$$u(x, 0) = f(x)$$

(9.8.6)

Clearly one can use backward differences to approximate the derivative boundary condition on the right end ($x = 1$), but this is of first order which will degrade the accuracy in $x$ everywhere (since the error will propagate to the interior in time). If we decide to use a second order approximation, then we have

$$\frac{u^n_{N+1} - u^n_{N-1}}{2\Delta x} = -h(u^n_N - v)$$

(9.8.7)

where $x_{N+1}$ is a fictitious point outside the interval, i.e. $x_{N+1} = 1 + \Delta x$. This will require another equation to match the number of unknowns. We then apply the finite difference equation at the boundary. For example, if we are using explicit scheme then we apply the equation

$$u_{j+1}^{n+1} = ru_j^n + (1 - 2r)u_j^n + ru_{j+1}^n,$$

(9.8.8)

for $j = 1, 2, \ldots, N$. At $j = N$, we then have

$$u_{N+1}^{n+1} = ru_{N-1}^n + (1 - 2r)u_N^n + ru_{N+1}^n.$$  

(9.8.9)

Substitute the value of $u_{N+1}^n$ from (9.8.7) into (9.8.9) and we get

$$u_{N}^{n+1} = ru_{N-1}^n + (1 - 2r)u_N^n + r \left[ u_{N-1}^n - 2h\Delta x (u_N^n - v) \right].$$

(9.8.10)

This idea can be implemented with any finite difference scheme.

Suggested Problem: Solve Laplace’s equation on a unit square subject to given temperature on right, left and bottom and insulated top boundary. Assume $\Delta x = \Delta y = h = \frac{1}{4}$.

### 9.9 Hyperbolic Equations

An important property of hyperbolic PDEs can be deduced from the solution of the wave equation. As the reader may recall the definitions of domain of dependence and domain of influence, the solution at any point $(x_0, t_0)$ depends only upon the initial data contained in the interval

$$x_0 - ct_0 \leq x \leq x_0 + ct_0.$$  

As we will see, this will relate to the so called CFL condition for stability.
9.9.1 Stability

Consider the first order hyperbolic

\[ u_t + cu_x = 0 \]  \hspace{1cm} (9.9.1.1)

\[ u(x,0) = F(x). \]  \hspace{1cm} (9.9.1.2)

As we have seen earlier, the characteristic curves are given by

\[ x - ct = \text{constant} \]  \hspace{1cm} (9.9.1.3)

and the general solution is

\[ u(x,t) = F(x - ct). \]  \hspace{1cm} (9.9.1.4)

Now consider Lax method for the approximation of the PDE

\[ u_j^{n+1} - \frac{u_{j+1}^{n+1} + u_{j-1}^{n+1}}{2} + c \frac{\Delta t}{\Delta x} \left( \frac{u_{j+1}^n - u_{j-1}^n}{2} \right) = 0. \]  \hspace{1cm} (9.9.1.5)

To check stability, we can use either Fourier method or the matrix method. In the first case, we substitute a Fourier mode and find that

\[ G = e^{a \Delta t} = \cos \beta - i \nu \sin \beta \]  \hspace{1cm} (9.9.1.6)

where the Courant number \( \nu \) is given by

\[ \nu = c \frac{\Delta t}{\Delta x}. \]  \hspace{1cm} (9.9.1.7)

Thus, for the method to be stable, the amplification factor \( G \) must satisfy

\[ |G| \leq 1 \]

i.e.

\[ \sqrt{\cos^2 \beta + \nu^2 \sin^2 \beta} \leq 1 \]  \hspace{1cm} (9.9.1.8)

This holds if

\[ |\nu| \leq 1, \]  \hspace{1cm} (9.9.1.9)

or

\[ c \frac{\Delta t}{\Delta x} \leq 1. \]  \hspace{1cm} (9.9.1.10)

Compare this CFL condition to the domain of dependence discussion previously. Note that here we have a complex number for the amplification. Writing it in polar form,

\[ G = \cos \beta - i \nu \sin \beta = |G|e^{i \phi} \]  \hspace{1cm} (9.9.1.11)

where the phase angle \( \phi \) is given by

\[ \phi = \arctan(-\nu \tan \beta). \]  \hspace{1cm} (9.9.1.12)
A good understanding of the amplification factor comes from a polar plot of amplitude versus relative phase $r = -\frac{\phi}{\pi}$ for various $\nu$ (see figure 70).

Note that the amplitude for all these values of Courant number never exceeds 1. For $\nu = 1$, there is no attenuation. For $\nu < 1$, the low $(\phi = 0)$ and high $(\phi = -\pi)$ frequency components are mildly attenuated, while the mid range frequencies are severely attenuated.

Suppose now we solve the same equation using Lax method but we assume periodic boundary conditions, i.e.,

$$u^n_{m+1} = u^n_1$$

(9.9.13)

The system of equations obtained is

$$u^{n+1} = A u^n$$

(9.9.14)

where

$$u^n = \begin{bmatrix} u^n_1 \\ \vdots \\ u^n_m \end{bmatrix}$$

(9.9.15)

$$A = \begin{bmatrix} 0 & \frac{1 - \nu}{2} & 0 & \cdots & \frac{1 + \nu}{2} \\ \frac{1 + \nu}{2} & 0 & \frac{1 - \nu}{2} & \cdots & 0 \\ 0 & \frac{1 + \nu}{2} & 0 & \cdots & \frac{1 - \nu}{2} \\ \frac{1 - \nu}{2} & \cdots & 0 & \frac{1 + \nu}{2} & 0 \end{bmatrix}.$$  

(9.9.16)

It is clear that the eigenvalues of $A$ are

$$\lambda_j = \cos \frac{2\pi}{m} (j - 1) + i \nu \sin \frac{2\pi}{m} (j - 1), \quad j = 1, \ldots, m.$$  

(9.9.17)
Since the stability of the method depends on
\[ |\rho(A)| \leq 1, \]  
(9.9.1.18)

one obtains the same condition in this case. The two methods yield identical results for periodic boundary condition. It can be shown that this is not the case in general.

If we change the boundary conditions to
\[ u_1^{n+1} = u_1^n \]  
(9.9.1.19)

with
\[ u_4^{n+1} = u_3^n \]  
(9.9.1.20)

to match the wave equation, then the matrix becomes
\[
A = \begin{bmatrix}
1 & 0 & 0 & 0 \\
1 + \nu & 0 & 1 - \nu & 0 \\
2 & 0 & 2 & 1 - \nu \\
0 & 2 & 1 & 0 \\
\end{bmatrix},
\]  
(9.9.1.21)

The eigenvalues are
\[ \lambda_1 = 1, \quad \lambda_2 = 0, \quad \lambda_{3,4} = \pm \frac{1}{2}\sqrt{(1 - \nu)(3 + \nu)}. \]  
(9.9.1.22)

Thus the condition for stability becomes
\[ -\sqrt{8} - 1 \leq \nu \leq \sqrt{8} - 1. \]  
(9.9.1.23)

See work by Hirt (1968), Warning and Hyett (1974) and Richtmeyer and Morton (1967).
Problems

1. Use a von Neumann stability analysis to show for the wave equation that a simple explicit Euler predictor using central differencing in space is unstable. The difference equation is

\[ u_j^{n+1} = u_j^n - c \frac{\Delta t}{\Delta x} \left( \frac{u_{j+1}^n - u_{j-1}^n}{2} \right) \]

Now show that the same difference method is stable when written as the implicit formula

\[ u_j^{n+1} = u_j^n - c \frac{\Delta t}{\Delta x} \left( \frac{u_{j+1}^{n+1} - u_{j-1}^{n+1}}{2} \right) \]

2. Prove that the CFL condition is the stability requirement when the Lax Wendroff method is applied to solve the simple 1-D wave equation. The difference equation is of the form:

\[ u_j^{n+1} = u_j^n - \frac{c \Delta t}{2 \Delta x} \left( u_{j+1}^n - u_{j-1}^n \right) + \frac{c^2 (\Delta t)^2}{2 (\Delta x)^2} \left( u_{j+1}^n - 2u_j^n + u_{j-1}^n \right) \]

3. Determine the stability requirement to solve the 1-D heat equation with a source term

\[ \frac{\partial u}{\partial t} = \alpha \frac{\partial^2 u}{\partial x^2} + ku \]

Use the central-space, forward-time difference method. Does the von Neumann necessary condition make physical sense for this type of computational problem?

4. In attempting to solve a simple PDE, a system of finite-difference equations of the form

\[ u_j^{n+1} = \left[ \begin{array}{ccc} 1+\nu & 1+\nu & 0 \\ 0 & 1+\nu & \nu \\ -\nu & 0 & 1+\nu \end{array} \right] u_j^n. \]

Investigate the stability of the scheme.
9.9.2 Euler Explicit Method

Euler explicit method for the first order hyperbolic is given by (for $c > 0$)

$$\frac{u_{j}^{n+1} - u_{j}^{n}}{\Delta t} + c \frac{u_{j+1}^{n} - u_{j}^{n}}{\Delta x} = 0$$  \hspace{1cm} (9.9.2.1)

or

$$\frac{u_{j}^{n+1} - u_{j}^{n}}{\Delta t} + c \frac{u_{j+1}^{n} - u_{j-1}^{n}}{2\Delta x} = 0$$  \hspace{1cm} (9.9.2.2)

Both methods are explicit and first order in time, but also unconditionally unstable.

$$G = 1 - \frac{\nu}{2} \left(2i \sin \beta \right) \quad \text{for centred difference in space,} \hspace{1cm} (9.9.2.3)$$

$$G = 1 - \nu \left(2i \sin \frac{\beta}{2} \right) e^{i\beta/2} \quad \text{for forward difference in space.} \hspace{1cm} (9.9.2.4)$$

In both cases the amplification factor is always above 1. The only difference between the two is the spatial order.

9.9.3 Upstream Differencing

Euler’s method can be made stable if one takes backward differences in space in case $c > 0$ and forward differences in case $c < 0$. The method is called upstream differencing or upwind differencing. It is written as

$$\frac{u_{j}^{n+1} - u_{j}^{n}}{\Delta t} + c \frac{u_{j}^{n} - u_{j-1}^{n}}{\Delta x} = 0, \quad c > 0.$$  \hspace{1cm} (9.9.3.1)

The method is of first order in both space and time, it is conditionally stable for $0 \leq \nu \leq 1$. The truncation error can be obtained by substituting Taylor series expansions for $u_{j-1}^{n}$ and $u_{j+1}^{n}$ in (9.9.3.1).

$$\frac{1}{\Delta t} \left\{ \Delta t u_t + \frac{1}{2} \Delta t^2 u_{tt} + \frac{1}{6} \Delta t^3 u_{ttt} + \cdots \right\}$$

$$+ \frac{c}{\Delta x} \left\{ u - \left[ u - \Delta x u_x + \frac{1}{2} \Delta x^2 u_{xx} - \frac{1}{6} \Delta x^3 u_{xxx} + \cdots \right] \right\}$$

where all the terms are evaluated at $x_j, t_n$. Thus the truncation error is

$$u_t + cu_x = -\frac{\Delta t}{2} u_{tt} + \frac{\Delta x}{2} u_{xx} \hspace{1cm} (9.9.3.2)$$

$$-\frac{\Delta t^2}{6} u_{ttt} - \frac{\Delta x^2}{6} u_{xxx} \pm \cdots$$
The modified equation is
\[
\frac{\partial u}{\partial t} + c \frac{\partial u}{\partial x} = \frac{\Delta x}{2} (1 - \nu) u_{xx} - c \frac{\Delta x^2}{6} (2\nu^2 - 3\nu + 1) u_{xxx} + O \left[ \Delta x^3, \Delta t \Delta x^2, \Delta x \Delta t^2, \Delta t^3 \right]
\]
(9.9.3.3)

In the next table we organized the calculations. We start with the coefficients of truncation error, (9.9.3.2), after moving all terms to the left. These coefficients are given in the second row of the table. The first row give the partials of \(u\) corresponding to the coefficients. Now in order to eliminate the coefficient of \(u_{tt}\), we have to differentiate the first row and multiply by \(-\Delta t/2\). This will modify the coefficients of other terms. Next we eliminate the new coefficient of \(u_{tx}\), and so on. The last row shows the sum of coefficients in each column, which are the coefficients of the modified equation.

<table>
<thead>
<tr>
<th>(u_t)</th>
<th>(u_x)</th>
<th>(u_{tt})</th>
<th>(u_{tx})</th>
<th>(u_{xx})</th>
</tr>
</thead>
<tbody>
<tr>
<td>1</td>
<td>(c)</td>
<td>(\Delta t/2)</td>
<td>0</td>
<td>(-c\Delta t/2)</td>
</tr>
<tr>
<td>(-\frac{\Delta t}{2} \frac{\partial}{\partial t}) (9.9.3.2)</td>
<td>(-\Delta t/2)</td>
<td>(-c\Delta t/2)</td>
<td>0</td>
<td></td>
</tr>
<tr>
<td>(\frac{c}{2} \Delta t \frac{\partial}{\partial x}) (9.9.3.2)</td>
<td>(c\Delta t/2)</td>
<td>(\frac{c^2}{2} \Delta t)</td>
<td></td>
<td></td>
</tr>
<tr>
<td>(\frac{1}{12} \Delta t^2 \frac{\partial^2}{\partial t^2}) (9.9.3.2)</td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td>(-\frac{1}{3} c \Delta t^2 \frac{\partial^2}{\partial t \partial x}) (9.9.3.2)</td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td>(\left( \frac{1}{3} c^2 \Delta t^2 - c \frac{\Delta t \Delta x}{4} \right) \frac{\partial^2}{\partial x^2}) (9.9.3.2)</td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
</tbody>
</table>

**Sum of coefficients**

| 1 | \(c\) | 0 | 0 | \(c\Delta t/2(\nu - 1)\) |

Table 2: Organizing the calculation of the coefficients of the modified equation for upstream differencing.

The right hand side of (9.9.3.3) is the truncation error. The method is of first order. If \(\nu = 1\), the right hand side becomes zero and the equation is solved exactly. In this case the upstream method becomes

\[u_{j}^{n+1} = u_{j}^{n-1}\]

which is equivalent to the exact solution using the method of characteristics.

The lowest order term of the truncation error contains \(u_{xx}\), which makes this term similar to the viscous term in one dimensional fluid flow. Thus when \(\nu \neq 1\), the upstream differencing introduces an **artificial viscosity** into the solution. Artificial viscosity tends to reduce all gradients in the solution whether physically correct or numerically induced. This effect, which is the direct result of even order derivative terms in the truncation error is called **dissipation**.
A dispersion is a result of the odd order derivative terms. As a result of dispersion, phase relations between waves are distorted. The combined effect of dissipation and dispersion is called diffusion. Diffusion tends to spread out sharp dividing lines that may appear in the computational region.

The amplification factor for the upstream differencing is

$$ e^{\alpha \Delta t} - 1 + \nu \left(1 - e^{-i\beta}\right) = 0 $$

or

$$ G = (1 - \nu + \nu \cos \beta) - i\nu \sin \beta \quad (9.9.3.4) $$

The amplitude and phase are then

$$ |G| = \sqrt{(1 - \nu + \nu \cos \beta)^2 + (-\nu \sin \beta)^2} \quad (9.9.3.5) $$

$$ \phi = \arctan \frac{Im(G)}{Re(G)} = \arctan \frac{-\nu \sin \beta}{1 - \nu + \nu \cos \beta}. \quad (9.9.3.6) $$

See figure 71 for polar plot of the amplification factor modulus as a function of $\beta$ for various values of $\nu$. For $\nu = 1.25$, we get values outside the unit circle and thus we have instability ($|G| > 1$).

The amplification factor for the exact solution is

$$ G_e = \frac{u(t + \Delta t)}{u(t)} = \frac{e^{ik_m[x-c(t+\Delta t)]}}{e^{ik_m[x-c(t)]}} = e^{-ik_m c \Delta t} = e^{i \phi_e} \quad (9.9.3.7) $$
Note that the magnitude is 1, and
\[ \phi_e = -k_m c \Delta t = -\nu \beta. \]  

(9.9.3.8)

The total dissipation error in \( N \) steps is
\[ (1 - G^N) A_0 \]

and the total dispersion error in \( N \) steps is
\[ N(\phi_e - \phi). \]  

(9.9.3.10)

The relative phase shift in one step is
\[ \frac{\phi}{\phi_e} = \frac{\text{arctan} \left( -\nu \sin \beta \right)}{-\nu \beta}. \]  

(9.9.3.11)

See figure 72 for relative phase error of upstream differencing. For small \( \beta \) (wave number) the relative phase error is
\[ \frac{\phi}{\phi_e} \approx 1 - \frac{1}{6}(2\nu^2 - 3\nu + 1)\beta^2 \]  

(9.9.3.12)

If \( \frac{\phi}{\phi_e} > 1 \) for a given \( \beta \), the corresponding Fourier component of the numerical solution has a wave speed greater than the exact solution and this is a leading phase error, otherwise lagging phase error.

The upstream has a leading phase error for \( .5 < \nu < 1 \) (outside unit circle) and lagging phase error for \( \nu < .5 \) (inside unit circle).
To derive Lax Wendroff method, we use Taylor series

\[ u_{j}^{n+1} = u_{j}^{n} + \Delta t u_{t} + \frac{1}{2} (\Delta t)^2 u_{tt} + O \left( (\Delta t)^3 \right) \]  

Substitute for \( u_{t} \) from the PDE

\[ u_{t} = -c u_{x} \]  

and for \( u_{tt} \) from its derivative

\[ u_{tt} = -c u_{xt} = -c (-c u_{xx}) = c^2 u_{xx} \]

to get

\[ u_{j}^{n+1} = u_{j}^{n} - c \frac{\Delta t}{2 \Delta x} (u_{j+1}^{n} - u_{j-1}^{n}) + \frac{1}{2} c^2 \frac{(\Delta t)^2}{(\Delta x)^2} (u_{j+1}^{n} - 2u_{j}^{n} + u_{j-1}^{n}) \cdot \]

The method is explicit, one step, second order with truncation error

\[ T.E. = O \left( (\Delta x)^2, (\Delta t)^2 \right) . \]

The modified equation is

\[ u_{t} + cu_{x} = -c \frac{(\Delta x)^2}{6} (1 - \nu^2) u_{xxx} - c \frac{(\Delta x)^3}{8} \nu (1 - \nu^2) u_{xxxx} + \cdots \]

The amplification factor

\[ G = 1 - \nu^2 (1 - \cos \beta) - i \nu \sin \beta , \]
and the method is stable for

$$|\nu| \leq 1.$$  \hfill (9.9.4.8)

The relative phase error is

$$\frac{\phi}{\phi_e} = \frac{-\nu \sin \beta}{1 - \nu^2 (1 - \cos \beta)}.$$  \hfill (9.9.4.9)

See figure 73 for the amplification factor modulus and the relative phase error. The method is predominantly lagging phase except for $\sqrt{5} < \nu < 1$.

Figure 73: Amplification factor modulus (left) and relative phase error (right) of Lax Wendroff scheme
Problems

1. Derive the modified equation for the Lax Wendroff method.
For nonlinear equations such as the inviscid Burgers’ equation, a two step variation of this method can be used. For the first order wave equation (9.9.1.1) this explicit two-step three time level method becomes

\[
\frac{u_{j+1}^{n+1/2} - (u_j^{n+1} + u_j^n)/2}{\Delta t/2} + c \frac{u_j^{n+1} - u_j^n}{\Delta x} = 0 \quad (9.9.4.10)
\]

\[
\frac{u_j^{n+1} - u_j^n}{\Delta t} + c \frac{u_j^{n+1/2} - u_{j-1}^{n+1/2}}{\Delta x} = 0. \quad (9.9.4.11)
\]

This scheme is second order accurate with a truncation error

\[
T.E. = O \left( \Delta x^2, (\Delta t)^2 \right), \quad (9.9.4.12)
\]

and is stable for \(|\nu| \leq 1\). For the linear first order hyperbolic this scheme is equivalent to the Lax Wendroff method.

### 9.9.5 MacCormack Method

MacCormack method is a predictor-corrector type. The method consists of two steps, the first is called predictor (predicting the value at time \(t_{n+1}\) and the second is called corrector.

\[
\text{Predictor} : u_j^{n+1} = u_j^n - c \frac{\Delta t}{\Delta x} (u_{j+1}^n - u_j^n) \quad (9.9.5.1)
\]

\[
\text{Corrector} : u_j^{n+1} = \frac{1}{2} \left[ u_j^n + u_j^{n+1} - c \frac{\Delta t}{\Delta x} (u_{j+1}^{n+1} - u_{j-1}^{n+1}) \right]. \quad (9.9.5.2)
\]

In the predictor, a forward difference for \(u_x\) while in the corrector, a backward difference for \(u_x\). This differencing can be reversed and sometimes (moving discontinuities) it is advantageous.

For linear problems, this is equivalent to Lax Wendroff scheme and thus the truncation error, stability criterion, modified equation, and amplification factor are all identical. We can now turn to nonlinear wave equation. The problem we discuss is Burgers’ equation.

[Can do Lab 7]

### 9.10 Inviscid Burgers’ Equation

Fluid mechanics problems are highly nonlinear. The governing PDEs form a nonlinear system that must be solved for the unknown pressures, densities, temperatures and velocities. A single equation that could serve as a nonlinear analog must have terms that closely duplicate the physical properties of the fluid equations, i.e. the equation should have a convective terms \((uu_x)\), a diffusive or dissipative term \((\mu u_{xx})\) and a time dependent term \((u_t)\). Thus the equation

\[
u_t + uu_x = \mu u_{xx} \quad (9.10.1)
\]
is parabolic. If the viscous term is neglected, the equation becomes hyperbolic,

\[ u_t + uu_x = 0. \]  

(9.10.2)

This can be viewed as a simple analog of the Euler equations for the flow of an inviscid fluid. The vector form of Euler equations is

\[ \frac{\partial U}{\partial t} + \frac{\partial E}{\partial x} + \frac{\partial F}{\partial y} + \frac{\partial G}{\partial z} = 0 \]  

(9.10.3)

where the vectors \( U, E, F, \) and \( G \) are nonlinear functions of the density \( (\rho) \), the velocity components \((u, v, w)\), the pressure \((p)\) and the total energy per unit volume \((E_t)\).

\[
U = \begin{bmatrix}
\rho \\
\rho u \\
\rho v \\
\rho w \\
E_t
\end{bmatrix},
\]

(9.10.4)

\[
E = \begin{bmatrix}
\rho u \\
\rho u^2 + p \\
\rho uv \\
\rho uw \\
(E_t + p)u
\end{bmatrix},
\]

(9.10.5)

\[
F = \begin{bmatrix}
\rho v \\
\rho uv \\
\rho v^2 + p \\
\rho vw \\
(E_t + p)v
\end{bmatrix},
\]

(9.10.6)

\[
G = \begin{bmatrix}
\rho w \\
\rho uw \\
\rho vw \\
\rho w^2 + p \\
(E_t + p)w
\end{bmatrix}.
\]

(9.10.7)

In this section, we discuss the inviscid Burgers’ equation (9.10.2). As we have seen in a previous chapter, the characteristics may coalesce and discontinuous solution may form. We consider the scalar equation

\[ u_t + F(u)_x = 0 \]  

(9.10.8)

and if \( u \) and \( F \) are vectors

\[ u_t + Au_x = 0 \]  

(9.10.9)

where \( A(u) \) is the Jacobian matrix \( \frac{\partial F_i}{\partial u_j} \). Since the equation is hyperbolic, the eigenvalues of the Matrix \( A \) are all real. We now discuss various methods for the numerical solution of (9.10.2).
9.10.1 Lax Method

Lax method is first order, as in the previous section, we have

\[ u_{j+1}^{n+1} = \frac{u_j^n + u_{j-1}^n}{2} - \frac{\Delta t}{\Delta x} \frac{F_{j+1}^n - F_{j-1}^n}{2}. \]  

(9.10.1.1)

Figure 74: Solution of Burgers’ equation using Lax method

In Burgers’ equation

\[ F(u) = \frac{1}{2} u^2. \]  

(9.10.1.2)

The amplification factor is given by

\[ G = \cos \beta - i \frac{\Delta t}{\Delta x} A \sin \beta \]  

(9.10.1.3)

where \( A \) is the Jacobian \( \frac{dF}{du} \), which is just \( u \) for Burgers’ equation. The stability requirement is

\[ \left| \frac{\Delta t}{\Delta x} u_{\text{max}} \right| \leq 1, \]  

(9.10.1.4)

because \( u_{\text{max}} \) is the maximum eigenvalue of the matrix \( A \). See Figure 74 for the exact versus numerical solution with various ratios \( \frac{\Delta t}{\Delta x} \). The location of the moving discontinuity is correctly predicted, but the dissipative nature of the method is evident in the smearing of the discontinuity over several mesh intervals. This smearing becomes worse as the Courant number decreases. Compare the solutions in Figure 74.
9.10.2 Lax Wendroff Method

This is a second order method which one can develop using Taylor series expansion

\[ u(x, t + \Delta t) = u(x, t) + \Delta t \frac{\partial u}{\partial t} + \frac{1}{2}(\Delta t)^2 \frac{\partial^2 u}{\partial t^2} + \cdots \]  \hspace{1cm} (9.10.2.1)

Using Burgers’ equation and the chain rule, we have

\[ u_t = -F_x = -F_u u_x = -Au_x \]  \hspace{1cm} (9.10.2.2)

\[ u_{tt} = -F_{tx} = -F_{xt} = -(F_t)_x. \]

Now

\[ F_t = F_u u_t = Au_t = -AF_x. \]  \hspace{1cm} (9.10.2.3)

Therefore

\[ u_{tt} = -(-AF_x)_x = (AF_x)_x. \]  \hspace{1cm} (9.10.2.4)

Substituting in (9.10.2.1) we get

\[ u(x, t + \Delta t) = u(x, t) - \Delta t \frac{\partial F}{\partial x} + \frac{1}{2}(\Delta t)^2 \frac{\partial}{\partial x} \left( A \frac{\partial F}{\partial x} \right) + \cdots \]  \hspace{1cm} (9.10.2.5)

Now use centered differences for the spatial derivatives

\[ u_{j+1}^{n+1} = u_j^n - \frac{\Delta t}{\Delta x} \frac{F_{j+1}^n - F_{j-1}^n}{2} \]  \hspace{1cm} (9.10.2.6)

\[ + \frac{1}{2} \left( \frac{\Delta t}{\Delta x} \right)^2 \left\{ A_{j+1/2}^n \left( F_{j+1}^n - F_j^n \right) - A_{j-1/2}^n \left( F_j^n - F_{j-1}^n \right) \right\} \]

where

\[ A_{j+1/2}^n = A \left( \frac{u_j^n + u_{j+1}^n}{2} \right). \]  \hspace{1cm} (9.10.2.7)

For Burgers’ equation, \( F = \frac{1}{2} u^2 \), thus \( A = u \) and

\[ A_{j+1/2}^n = \frac{u_j^n + u_{j+1}^n}{2}, \]  \hspace{1cm} (9.10.2.8)

\[ A_{j-1/2}^n = \frac{u_j^n + u_{j-1}^n}{2}. \]  \hspace{1cm} (9.10.2.9)

The amplification factor is given by

\[ G = 1 - 2 \left( \frac{\Delta t}{\Delta x} A \right)^2 (1 - \cos \beta) - 2i \frac{\Delta t}{\Delta x} A \sin \beta. \]  \hspace{1cm} (9.10.2.10)

Thus the condition for stability is

\[ \left| \frac{\Delta t}{\Delta x} u_{\text{max}} \right| \leq 1. \]  \hspace{1cm} (9.10.2.11)
The numerical solution is given in Figure 75. The right moving discontinuity is correctly positioned and sharply defined. The dispersive nature is evidenced in the oscillation near the discontinuity.

The solution shows more oscillations when \( \nu = .6 \) than when \( \nu = 1 \). When \( \nu \) is reduced the quality of the solution is degraded.

The flux \( F(u) \) at \( x_j \) and the numerical flux \( f_{j+1/2} \), to be defined later, must be consistent with each other. The numerical flux is defined, depending on the scheme, by matching the method to

\[
    u_j^{n+1} = u_j^n - \frac{\Delta t}{\Delta x} \left[ f_{j+1/2}^{n+1} - f_{j+1/2}^{n} \right].
\]

In order to obtain the numerical flux for Lax Wendroff method for solving Burgers’ equation, let’s add and subtract \( F_j^n \) in the numerator of the first fraction on the right, and substitute \( u \) for \( A \)

\[
    u_j^{n+1} = u_j^n - \frac{\Delta t}{\Delta x} \left\{ \frac{F_{j+1}^n + F_j^n - F_{j-1}^n}{2} - \frac{1}{2} \frac{\Delta t}{\Delta x} \left[ \frac{u_j^n + u_{j+1}^n}{2} \left( F_{j+1}^n - F_j^n \right) - \frac{u_j^n + u_{j-1}^n}{2} \left( F_j^n - F_{j-1}^n \right) \right] \right\}.
\]

Recall that \( F(u) = \frac{1}{2}u^2 \), and factor the difference of squares to get

\[
    f_{j+1/2}^n = \frac{1}{2} (F_j^n + F_{j+1}^n) - \frac{1}{2} \frac{\Delta t}{\Delta x} (u_{j+1/2}^n) (u_{j+1}^n - u_j^n).
\]

The numerical flux for Lax method is given by

\[
    f_{j+1/2}^n = \frac{1}{2} \left[ F_j^n + F_{j+1}^n - \frac{\Delta x}{\Delta t} (u_{j+1}^n - u_j^n) \right].
\]
Lax method is monotone, and Gudonov showed that one cannot get higher order than first and keep monotonicity.

9.10.3 MacCormack Method

This method is different than other, it is a two step predictor corrector method. One predicts the value at time $n+1$ is the first step and then corrects it in the second step.

$$\text{Predictor}: \ u_{j}^{n+1} = u_{j}^{n} - \frac{\Delta t}{\Delta x} \left( F_{j+1}^{n} - F_{j}^{n} \right) \quad (9.10.3.1)$$

$$\text{Corrector}: \ u_{j}^{n+1} = \frac{1}{2} \left[ u_{j}^{n} + u_{j}^{n+1} - \frac{\Delta t}{\Delta x} \left( F_{j+1}^{n+1} - F_{j}^{n+1} \right) \right] . \quad (9.10.3.2)$$

Compare this to MacCormack method for the linear case where $F = cu$. The amplification factor and stability requirements are as in Lax Wendroff scheme. See figure 76 for the numerical solution of Burgers’ equation. Notice the oscillations only ahead of the jump. The difference is because of the switched differencing in the predictor-corrector.

![Figure 76: Solution of Burgers’ equation using MacCormack method](image)

Note: The best resolution of discontinuities occurs when the difference in the predictor is in the same direction of the propagation of discontinuity.
Problems

1. Determine the errors in amplitude and phase for $\beta = 90^\circ$ if the MacCormack scheme is applied to the wave equation for 10 time steps with $\nu = .5$. 
9.10.4 Implicit Method

A second order accurate implicit scheme results from

\[ u_{j}^{n+1} = u_{j}^{n} + \frac{\Delta t}{2} \left[ (u_{t})^{n} + (u_{t})^{n+1} \right]_{j} + O \left( (\Delta t)^{3} \right) \]  

(9.10.4.1)

which is based on the trapezoidal rule. Since

\[ u_{t} = -F_{x} \]

we can write

\[ u_{j}^{n+1} = u_{j}^{n} - \frac{\Delta t}{2} \left[ (F_{x})^{n} + (F_{x})^{n+1} \right]_{j} \].  

(9.10.4.2)

This is nonlinear in \( u_{j}^{n+1} \) and thus requires a linearization or an iterative process. Beam and Warming suggest to linearize in the following manner

\[ F^{n+1} = F^{n} + F^{n}_{u} \left( u^{n+1} - u^{n} \right) = F^{n} + A^{n} \left( u^{n+1} - u^{n} \right) \].  

(9.10.4.3)

Thus

\[ u_{j}^{n+1} = u_{j}^{n} - \frac{\Delta t}{2} \left\{ 2F_{x}^{n} + \frac{\partial}{\partial x} \left[ A \left( u^{n+1}_{j} - u^{n}_{j} \right) \right] \right\} . \]  

(9.10.4.4)

Now replace the spatial derivatives by centered differences and collect terms

\[ -\frac{\Delta t}{4\Delta x} A_{j-1}^{n} u_{j-1}^{n+1} + u_{j}^{n+1} + \frac{\Delta t}{4\Delta x} A_{j+1}^{n} u_{j+1}^{n+1} = \]  

\[ -\frac{\Delta t}{2} \frac{F^{n}_{j+1} - F^{n}_{j-1}}{\Delta x} - \frac{\Delta t}{4\Delta x} A_{j-1}^{n} u_{j-1}^{n} + u_{j}^{n} + \frac{\Delta t}{4\Delta x} A_{j+1}^{n} u_{j+1}^{n} \].  

(9.10.4.5)

This is a linear tridiagonal system for each time level. The entries of the matrix depend on time and thus we have to reconstruct it at each time level.

The modified equation contains no even order derivative terms, i.e., no dissipation. Figure 77 shows the exact solution of Burgers’ equation subject to the same initial condition as in previous figures along with the numerical solution. Notice how large is the amplitude of the oscillations. Artificial smoothing is added to right hand side

\[ -\frac{\omega}{8} \left( u_{j+2}^{n} - 4u_{j+1}^{n} + 6u_{j}^{n} - 4u_{j-1}^{n} + u_{j-2}^{n} \right) \]

where \( 0 < \omega \leq 1 \). This makes the amplitude of the oscillations smaller. In Figure 77, we have the solution without damping and with \( \omega = .5 \) after 20 time steps using \( \nu = .5 \)

Another implicit method due to Beam and Warming is based on Euler implicit:

\[ u^{n+1} = u^{n} + \Delta t \left( u_{t} \right)^{n+1} \]  

(9.10.4.6)

\[ u^{n+1} = u^{n} - \Delta t \left( F_{x} \right)^{n+1} \]  

(9.10.4.7)

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Figure 77: Solution of Burgers’ equation using implicit (trapezoidal) method

with the same linearization

\[-\frac{\Delta t}{2\Delta x} A^n_{j-1} u_{j-1}^{n+1} + u_j^{n+1} + \frac{\Delta t}{2\Delta x} A^n_{j+1} u_{j+1}^{n+1} =\]

\[-\frac{\Delta t}{\Delta x} \frac{F^n_{j+1} - F^n_{j-1}}{2} - \frac{\Delta t}{2\Delta x} A^n_{j-1} u_{j-1}^{n} + u_j^{n} + \frac{\Delta t}{2\Delta x} A^n_{j+1} u_{j+1}^{n}.\]  

(9.10.4.8)

Again we get a tridiagonal system and same smoothing must be added.
Problems

1. Apply the two-step Lax Wendorff method to the PDE

\[ \frac{\partial u}{\partial t} + \frac{\partial F}{\partial x} + u \frac{\partial^3 u}{\partial x^3} = 0 \]

where \( F = F(u) \). Develop the final finite difference equations.

2. Apply the Beam-Warming scheme with Euler implicit time differencing to the linearized Burgers’ equation on the computational grid given in Figure 78 and determine the steady state values of \( u \) at \( j = 2 \) and \( j = 3 \). The boundary conditions are

\[ u_1^n = 1, \quad u_4^n = 4 \]

and the initial conditions are

\[ u_2^1 = 0, \quad u_3^1 = 0 \]

Do not use a computer to solve this problem.

Figure 78: Computational Grid for Problem 2
9.11 Viscous Burgers’ Equation

Adding viscosity to Burgers’ equation we get

\[ u_t + uu_x = \mu u_{xx}. \] (9.11.1)

The equation is now parabolic. In this section we mention analytic solutions for several cases. We assume Dirichlet boundary conditions:

\[ u(0,t) = u_0, \] (9.11.2)
\[ u(L,t) = 0. \] (9.11.3)

The steady state solution (of course will not require an initial condition) is given by

\[ u = u_0 \hat{u} \left\{ \frac{1 - e^{\hat{u} Re_L(x/L-1)}}{1 + e^{\hat{u} Re_L(x/L-1)}} \right\} \] (9.11.4)

where

\[ Re_L = \frac{u_0 L}{\mu} \] (9.11.5)

and \( \hat{u} \) is the solution of the nonlinear equation

\[ \frac{\hat{u} - 1}{\hat{u} + 1} = e^{-\hat{u} Re_L}. \] (9.11.6)

The linearized equation (9.10.1) is

\[ u_t + cu_x = \mu u_{xx} \] (9.11.7)

and the steady state solution is now

\[ u = u_0 \left\{ \frac{1 - e^{R_L(x/L-1)}}{1 - e^{-R_L}} \right\} \] (9.11.8)

where

\[ R_L = \frac{cL}{\mu}. \] (9.11.9)

The exact unsteady solution with initial condition

\[ u(x,0) = \sin kx \] (9.11.10)

and periodic boundary conditions is

\[ u(x,t) = e^{-k^2 \mu t} \sin k(x - ct). \] (9.11.11)

The equations (9.10.1) and (9.11.7) can be combined into a generalized equation

\[ u_t + (c + bu)u_x = \mu u_{xx}. \] (9.11.12)
For $b = 0$ we get the linearized Burgers’ equation and for $c = 0$, $b = 1$, we get the nonlinear equation. For $c = \frac{1}{2}$, $b = -1$ the generalized equation (9.11.12) has a steady state solution

$$u = -\frac{c}{b} \left( 1 + \tanh \frac{c(x - x_0)}{2\mu} \right). \quad (9.11.13)$$

Hence if the initial $u$ is given by (9.11.13), then the exact solution does not vary with time.

For more exact solutions, see Benton and Platzman (1972).

The generalized equation (9.11.12) can be written as

$$u_t + \hat{F}_x = 0 \quad (9.11.14)$$

where

$$\hat{F} = cu + \frac{1}{2}bu^2 - \mu u_x, \quad (9.11.15)$$

or as

$$u_t + F_x = \mu u_{xx}, \quad (9.11.16)$$

where

$$F = cu + \frac{1}{2}bu^2, \quad (9.11.17)$$

or

$$u_t + A(u)u_x = \mu u_{xx}. \quad (9.11.18)$$

The various schemes described earlier for the inviscid Burgers’ equation can also be applied here, by simply adding an approximation to $u_{xx}$.

### 9.11.1 FTCS method

This is a Forward in Time Centered in Space (hence the name),

$$\frac{u_j^{n+1} - u_j^n}{\Delta t} + c \frac{u_{j+1}^n - u_{j-1}^n}{2\Delta x} = \mu \frac{u_{j+1}^n - 2u_j^n + u_{j-1}^n}{(\Delta x)^2}. \quad (9.11.1.1)$$

Clearly the method is one step explicit and the truncation error

$$T.E. = O \left( \Delta t, (\Delta x)^2 \right). \quad (9.11.1.2)$$

Thus it is first order in time and second order in space. The modified equation is given by

$$u_t + cu_x = \left( \mu - \frac{c^2 \Delta t}{2} \right) u_{xx} + c \frac{(\Delta x)^2}{3} \left( 3r - \nu^2 - \frac{1}{2} \right) u_{xxx} \quad (9.11.1.3)$$

$$+ c \frac{(\Delta x)^3}{12} \left( \frac{r}{\nu} - 3 \frac{r^2}{\nu} - 2\nu + 10
\nu r - 3\nu^3 \right) u_{xxxx} + \cdots$$
where as usual

\[ r = \mu \frac{\Delta t}{(\Delta x)^2}, \]

\[ \nu = r \frac{\Delta t}{\Delta x}. \]  

(9.11.1.4)  

(9.11.1.5)

If \( r = \frac{1}{2} \) and \( \nu = 1 \), the first two terms on the right hand side of the modified equation vanish. This is NOT a good choice because it eliminated the viscous term that was originally in the PDE.

![Figure 79: Stability of FTCS method](image)

We now discuss the stability condition. Using Fourier method, we find that the amplification factor is

\[ G = 1 + 2r(\cos \beta - 1) - i\nu \sin \beta. \]  

(9.11.1.6)

In figure 79 we see a polar plot of \( G \) as a function of \( \nu \) and \( \beta \) for \( \nu < 1 \) and \( r < \frac{1}{2} \) and \( \nu^2 > 2r \) (left) and \( \nu^2 < 2r \) (right). Notice that if we allow \( \nu^2 \) to exceed \( 2r \), the ellipse describing \( G \) will have parts outside the unit circle and thus we have instability. This means that taking the combination of the conditions from the hyperbolic part (\( \nu < 1 \)) and the parabolic part (\( r < \frac{1}{2} \)) is not enough. This extra condition is required to ensure that the coefficient of \( u_{xx} \) is positive, i.e.

\[ c^2 \frac{\Delta t}{2} \leq \mu. \]  

(9.11.1.7)

Let’s define the mesh Reynolds number

\[ Re_{\Delta x} = \frac{c\Delta x}{\mu} = \frac{\nu}{r}, \]  

(9.11.1.8)

then the above condition becomes

\[ Re_{\Delta x} \leq \frac{2}{\nu}. \]  

(9.11.1.9)

It turns out that the method is stable if

\[ \nu^2 \leq 2r, \quad \text{and} \quad r \leq \frac{1}{2}. \]  

(9.11.1.10)
This combination implies that \( \nu \leq 1 \). Therefore we have

\[
2\nu \leq \text{Re}_{\Delta x} \leq \frac{2}{\nu}.
\] (9.11.1.11)

For \( \text{Re}_{\Delta x} > 2 \), FTCS will produce undesirable oscillations. To explain the origin of these oscillations consider the following example. Find the steady state solution of (9.10.1) subject to the boundary conditions

\[
u(0, t) = 0, \quad \nu(1, t) = 1
\] (9.11.1.12)

and the initial condition

\[
u(x, 0) = 0,
\] (9.11.1.13)

using an 11 point mesh. Note that we can write FTCS in terms of mesh Reynolds number as

\[
u_{j+1}^{n+1} = \frac{r}{2} (2 - \text{Re}_{\Delta x}) \nu_{j+1}^{n} + (1 - 2r) \nu_{j}^{n} + \frac{r}{2} (2 + \text{Re}_{\Delta x}) \nu_{j-1}^{n}.
\] (9.11.1.14)

\begin{figure}[h]
\centering
\includegraphics[width=\textwidth]{solution.png}
\caption{Solution of example using FTCS method}
\end{figure}

For the first time step

\[
u_{j}^{1} = 0, \quad j < 10
\]

and

\[
u_{10}^{1} = \frac{r}{2} (2 - \text{Re}_{\Delta x}) < 0, \quad \nu_{11}^{1} = 1,
\]

and this will initiate the oscillation. During the next time step the oscillation will propagate to the left. Note that \( \text{Re}_{\Delta x} > 2 \) means that \( \nu_{j+1}^{n} \) will have a negative weight which is physically wrong.

To eliminate the oscillations we can replace the centered difference for \( cu_x \) term by a first order upwind which adds more dissipation. This is too much. Leonard (1979) suggested a third order upstream for the convective term (for \( c > 0 \))

\[
\frac{u_{j+1}^{n} - u_{j-1}^{n}}{2\Delta x} - \frac{u_{j+1}^{n} - 3u_{j}^{n} + 3u_{j-1}^{n} - u_{j-2}^{n}}{6\Delta x}.
\]
9.11.2 Lax Wendroff method

This is a two step method:

\[ u_j^{n+1/2} = \frac{1}{2} \left( u_{j+1/2}^n + u_{j-1/2}^n \right) - \frac{\Delta t}{\Delta x} \left( F_{j+1/2}^n - F_{j-1/2}^n \right) \]

\[ + r \left[ \left( u_{j-3/2}^n - 2u_{j-1/2}^n + u_{j+1/2}^n \right) + \left( u_{j+3/2}^n - 2u_{j+1/2}^n + u_{j-1/2}^n \right) \right] \]

(9.11.2.1)

The second step is

\[ u_j^{n+1} = u_j^n - \frac{\Delta t}{\Delta x} \left( F_{j+1/2}^{n+1/2} - F_{j-1/2}^{n+1/2} \right) + r \left( u_{j+1}^n - 2u_j^n + u_{j-1}^n \right) . \]

(9.11.2.2)

The method is first order in time and second order in space. The linear stability condition is

\[ \frac{\Delta t}{(\Delta x)^2} \left( A^2 \Delta t + 2 \mu \right) \leq 1. \]

(9.11.2.3)

Can do problem 43

9.11.3 MacCormack method

This method is similar to the inviscid case. The viscous term is approximated by centered differences.

**Predictor**: \[ u_j^{n+1} = u_j^n - \frac{\Delta t}{\Delta x} \left( F_{j+1}^n - F_j^n \right) + r \left( u_{j+1}^n - 2u_j^n + u_{j-1}^n \right) \]

(9.11.3.1)

**Corrector**: \[ u_j^{n+1} = \frac{1}{2} \left[ u_j^n + u_j^{n+1} \right] - \frac{\Delta t}{\Delta x} \left( F_{j+1}^{n+1} - F_j^{n+1} \right) + r \left( u_{j+1}^{n+1} - 2u_j^{n+1} + u_{j-1}^{n+1} \right) . \]

(9.11.3.2)

The method is second order in space and time. It is not possible to get a simple stability criterion. Tannehill et al (1975) suggest an empirical value

\[ \Delta t \leq \frac{(\Delta x)^2}{|A| \Delta x + 2 \mu} . \]

(9.11.3.3)

This method is widely used for Euler’s equations and Navier Stokes for laminar flow. In multidimensional problems there is a time-split MacCormack method. An interesting variation when using relaxation is as follows

**Predictor**: \[ \overline{u}_j^{n+1} = u_j^n - \frac{\Delta t}{\Delta x} \left( F_{j+1}^n - F_j^n \right) + r \left( u_{j+1}^n - 2u_j^n + u_{j-1}^n \right) \]

(9.11.3.4)

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\[
\overline{u}_j^{n+1} = \overline{u}_j^n + \omega^P \left( v_j^{n+1} - \overline{u}_j^n \right), \quad (9.11.3.5)
\]

Corrector: \[
v_j^{n+1} = \frac{\Delta t}{\Delta x} \left( \frac{F_j^{n+1} - F_{j-1}^{n+1}}{2} \right) + r \left( u_{j+1}^{n+1} - 2u_j^{n+1} + u_{j-1}^{n+1} \right), \quad (9.11.3.6)
\]
\[
u_j^{n+1} = u_j^n + \omega^C \left( v_j^{n+1} - u_j^n \right). \quad (9.11.3.7)
\]

Note that for \( \omega^P = 1, \omega^C = \frac{1}{2} \) one gets the original scheme. In order to preserve the order of the method we must have \[
\omega^C \omega^P = \left| \omega^P - \omega^C \right|. \quad (9.11.3.8)
\]

A necessary condition for stability
\[
\left| (\omega^P - 1)(\omega^C - 1) \right| \leq 1. \quad (9.11.3.9)
\]

This scheme accelerates the convergence of the original scheme by a factor of
\[
\frac{2\omega^P \omega^C}{1 - (\omega^P - 1)(\omega^C - 1)}. \quad (9.11.3.10)
\]

9.11.4 Time-Split MacCormack method

The time-split MacCormack method is specifically designed for multidimensional problems. Let’s demonstrate it on the two dimensional Burgers’ equation
\[
u_t + F(u)_x + G(u)_y = \mu \left( u_{xx} + u_{yy} \right) \quad (9.11.4.1)
\]
or
\[
u_t + Au_x + Bu_y = \mu \left( u_{xx} + u_{yy} \right). \quad (9.11.4.2)
\]
The exact steady state solution of the two dimensional linearized Burgers’ equation on a unit square subject to the boundary conditions
\[
u(x, 0, t) = \frac{1 - e^{(x-1)c/\mu}}{1 - e^{-c/\mu}}, \quad u(x, 1, t) = 0, \quad (9.11.4.3)
\]
\[
u(0, y, t) = \frac{1 - e^{(y-1)d/\mu}}{1 - e^{-d/\mu}}, \quad u(1, y, t) = 0, \quad (9.11.4.4)
\]
is
\[
u(x, y) = \nu(x, 0, t)\nu(0, y, t). \quad (9.11.4.5)
\]

All the methods mentioned for the one-dimensional case can be extended to higher dimensions but the stability condition is more restrictive for explicit schemes and the systems are no
longer tridiagonal for implicit methods. The time split MacCormack method splits the original MacCormack scheme into a sequence of one dimensional equations:

$$u_{ij}^{n+1} = L_y \left( \frac{\Delta t}{2} \right) L_x (\Delta t) L_y \left( \frac{\Delta t}{2} \right) u_{ij}^n,$$

(9.11.4.6)

where the operators $L_x (\Delta t)$ and $L_y (\Delta t)$ are each equivalent to the two step formula as follows

$$u_{ij}^* = L_x (\Delta t) u_{ij}^n,$$

(9.11.4.7)

means

$$\overline{u_{ij}} = u_{ij}^n - \frac{\Delta t}{\Delta x} \left( F_{i+1,j}^n - F_{i,j}^n \right) + \mu \Delta t \delta_x^2 u_{ij}^n,$$

(9.11.4.8)

$$u_{ij}^* = \frac{1}{2} \left[ u_{ij}^n + \overline{u_{ij}} - \frac{\Delta t}{\Delta x} \left( \overline{F_{i,j}} - \overline{F_{i-1,j}} \right) \right] + \mu \Delta t \delta_x^2 \overline{u_{ij}},$$

(9.11.4.9)

and

$$u_{ij}^* = L_y (\Delta t) u_{ij}^n,$$

(9.11.4.10)

means

$$\overline{u_{ij}} = u_{ij}^n - \frac{\Delta t}{\Delta y} \left( G_{i,j+1}^n - G_{i,j}^n \right) + \mu \Delta t \delta_y^2 u_{ij}^n,$$

(9.11.4.11)

$$u_{ij}^* = \frac{1}{2} \left[ u_{ij}^n + \overline{u_{ij}} - \frac{\Delta t}{\Delta y} \left( \overline{G_{i,j}} - \overline{G_{i,j-1}} \right) \right] + \mu \Delta t \delta_y^2 \overline{u_{ij}},$$

(9.11.4.12)

The truncation error is

$$T.E. = O \left( (\Delta t)^2, (\Delta x)^2, (\Delta y)^2 \right).$$

(9.11.4.13)

In general such a scheme is stable if the time step of each operator doesn’t exceed the allowable size for that operator, it is consistent if the sum of the time steps for each operator is the same and it is second order if the sequence is symmetric.
9.12 Appendix - Fortran Codes

C************************************************************************************************************
C* PROGRAM FOR THE EXPLICIT SOLVER FOR THE HEAT EQUATION *
C* IN ONE DIMENSION *
C* DIRICHLET BOUNDARY CONDITIONS *
C* LIST OF VARIABLES *
C* I LOCATION OF X GRID POINTS *
C* J LOCATION OF T GRID POINTS *
C* U(I,J) TEMPERATURE OF BAR AT GRID POINT I,J *
C* K TIME SPACING *
C* H X SPACING *
C* IH NUMBER OF X DIVISIONS *
C* R K/H**2 *
C* NT NUMBER OF TIME STEPS *
C* IFREQ HOW MANY TIME STEPS BETWEEN PRINTOUTS *
C************************************************************************************************************
DIMENSION U(501,2),X(501)
REAL K
C* SPACING
NT=100
IFREQ=5
R=.1
IH = 10
PRINT 999
READ (5,*) TF
999 FORMAT(1X,'PLEASE TYPE IN THE FINAL TIME OF INTEGRATION')
PRINT 998
READ(5,*) IH
998 FORMAT(1X,'PLEASE TYPE IN THE NUMBER OF INTERIOR GRID POINTS')
PRINT 997
READ (5,*) R
997 FORMAT(1X,'PLEASE TYPE IN THE RATIO R')
H = 1.0/IH
IH1 = IH+ 1
K = R*H**2
NT=TF/K+1
WRITE(6,*) ' K=',K,' H=',H,' R=',R,' TF=',TF,' NT=',NT
C CALCULATIONS
DO 25 I = 1,IH1
C INITIAL CONDITIONS
X(I)=(I-1.)*H
U(I,1) = 2.0*(1.0-X(I))
IF (X(I) .LE. .5 ) U(I,1) = 2.0*X(I)
25 CONTINUE
25 CONTINUE
   TIME = 0.
   WRITE(6,202)TIME,(U(L,1),L=1,1H)
202   FORMAT(//2X,'AT T =',F7.3/(1X,5E13.6))
   DO 10 J = 1,NT
   C BOUNDARY CONDITIONS
   U(1,2) = 0.0
   U(IH1,2) = 0.0
   DO 15 I = 2,IH
      U(I,2) = R*U(I-1,1) + (1.0-2.0*R)*U(I,1) + R*U(I+1,1)
15 CONTINUE
   DO 20 L=1,1H
   U(L,1)=U(L,2)
   TIME=TIME+K
   IF(J/IFREQ*IFREQ.EQ.J) WRITE(6,202)TIME,(U(L,1),L=1,1H)
20 CONTINUE
   RETURN
END
C THIS PROGRAM SOLVES THE HEAT EQUATION IN ONE DIMENSION
C USING CRANK-NICHOLSON IMPLICIT METHOD. THE TEMPERATURE AT
C EACH END IS DETERMINED BY A RELATION OF THE FORM AU+BU'=C
C PARAMETERS ARE -
C U VALUES OF TEMPERATURE AT NODES
C T TIME
C TF FINAL TIME VALUE FOR WHICH SOLUTION IS DESIRED
C DT DELTA T
C DX DELTA X
C N NUMBER OF X INTERVALS
C RATIO RATIO OF DT//DX**2
C COEF COEFFICIENT MATRIX FOR IMPLICIT EQUATIONS
REAL U(500),COEF(500,3),RHS(500),X(500)
DATA T///0/.,TF///1/0/0/0/.,N///2/0/.,RATIO///1/.
C THE FOLLOWING STATEMENT GIVE THE BOUNDARY CONDITION AT X=0.
C A/,B/,C ON THE LEFT
C U'=/.2*/(U/-/1/5/)
DATA AL///-.2///,BL///1./0///,CL///-.3/0 //
C THE FOLLOWING STATEMENT GIVE THE BOUNDARY CONDITION AT X=1.
C A/,B/,C ON THE RIGHT
C U=100
DATA AR///1/.,BR///0/.,CR///100/.
PRINT 999
READ (5,* TF
999 FORMAT(1X,'PLEASE TYPE IN THE FINAL TIME OF INTEGRATION')
PRINT 998
READ(5,* N
998 FORMAT(1X,'PLEASE TYPE IN THE NUMBER OF INTERIOR GRID POINTS')
PRINT 997
READ (5,* ) RATIO
997 FORMAT(1X,'PLEASE TYPE IN THE RATIO R'
JJ=0
DX=1./N
DT=RATIO*DX*DX
NP1=N+1
C EVALUATE THE MESH POINTS
DO 1 I=1,NP1
1   X(I)=(I-1)*DX
C WRITE OUT HEADING AND INITIAL VALUES
WRITE(6,201) DX
201   FORMAT(‘/2X,’F6.3)
C COMPUTES INITIAL VALUES
DO 2 I=1,NP1
C ESTABLISH COEFFICIENT MATRIX
C LET ALPHA = -A/B
C LET BETA = C/B
C AT LEFT (4-2*ALPHA*DX)*U(1,J+1)-2*U(2,J+1) =
   2*ALPHA*DX*U(1,J)+2*U(2,J)-4*BETA*DX
C AT INTERIOR -U(I-1,J+1)+4*U(I,J+1)-U(I+1,J+1) =
   U(I-1,I)+U(I+1,J)
C AT RIGHT (-2*U(N,J+1)+(4+2*ALPHA*DX)*U(N+1,J+1) =
   2*U(N,J)-2*ALPHA*DX*U(N+1,J)+4*BETA*DX

IF(BL.EQ.0.) GO TO 10
COEF(1,2) = 2./RATIO+2.-2.*AL*DX/BL
COEF(1,3) = -2.
GO TO 20
10 COEF(1,2) = 1.
COEF(1,3) = 0.
20 DO 25 I = 2, N
   COEF(I,1) = -1.
   COEF(I,2) = 2./RATIO+2.
   COEF(I,3) = -1.
25 CONTINUE
IF(BR.EQ.0.) GO TO 30
COEF(N+1,1) = -2.
COEF(N+1,2) = 2./RATIO+2.+2.*AR*DX/BR
GO TO 40
30 COEF(N+1,1) = 0.
COEF(N+1,2) = 1.
C GET THE LU DECOMPOSITION
40 DO 50 I = 2, NP1
   COEF(I-1,3) = COEF(I-1,3)/COEF(I-1,2)
   COEF(I,2) = COEF(I,2) - COEF(I,1)*COEF(I-1,3)
50 CONTINUE
C CALCULATE THE R.H.S. VECTOR - FIRST THE TOP AND BOTTOM ROWS
55 IF(BL.EQ.0.) GO TO 60
   RHS(1) = (2./RATIO-2.+2.*AL*DX/BL)*U(1)+2.*U(2)-
      4.*CL*DX/BL
   GO TO 70
60 RHS(1) = CL/AL
70 IF(BR.EQ.0.) GO TO 80
   RHS(N+1) = 2.*U(N)+(2./RATIO-2.*AR*DX/BR)*U(N+1)+
      4.*CR*DX/BR
GO TO 90
80 RHS(N+1)=CR/AR
C NOW FOR THE OTHER ROWS OF THE RHS VECTOR
90 DO 100 I=2,N
100 RHS(I)=U(I-1)+(2./RATIO-2.)*U(I)+U(I+1)
C GET THE SOLUTION FOR THE CURRENT TIME
   U(1)=RHS(1)/COEF(1,2)
   DO 110 I=2,NP1
110 U(I)=(RHS(I)-COEF(I,1)*U(I-1))/COEF(I,2)
   DO 120 I=1,N
   JROW=N-I+1
120 U(JROW)=U(JROW)-COEF(JROW,3)*U(JROW+1)
C WRITE OUT THE SOLUTION
   T=T+DT
   JJ=JJ+1
   WRITE(6,202) T,(U(I),I=1,NP1)
   IF(T.LT.TF) GO TO 55
   STOP
END
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