

Vibrations and Waves

Dominik Szablonski

April 30, 2024

Contents

1	Simple Harmonic Motion	4
1.1	Energy in SHM	4
1.2	SHM In All Systems	5
1.2.1	Examples of Other Forms of SHM	5
1.3	General Form of SHM	6
1.3.1	Period Independent of Mass: Proof by Dimensional Analysis	6
2	Damped Harmonic Motion	7
2.1	Light Damping	7
2.2	Heavy Damping	7
2.3	Critical Dampinig	8
2.4	Quality Factor	8
2.5	LCR Circuit	8
2.6	General Damped Harmonic Oscillator	9
3	Driven, Damped Harmonic Oscillator	10
3.1	Transient Response	11
3.2	AC Circuit	11
3.2.1	Impedance	11
3.3	Absorbed power	12
4	Coupled Oscillators	14
4.1	Equations of motion in coupled pendula	14
4.1.1	Normal Modes for higher n	15
4.2	Transverse Oscillations	17
4.2.1	Single Oscillator	17
4.2.2	2 Oscillators	18
4.2.3	Higher n	18
4.3	Forced Coupled Oscillators	19
5	Waves	20
5.1	1D Wave Equation	20
5.1.1	Sinusoidal Waves	20
5.1.2	General solution to 1D wave equation	21
5.1.3	Velocity of travelling waves on a stretched string	21
5.1.4	Energy Stored in a Wave	22
5.2	Waves at Boundaries	23
5.3	Standing Waves	24
5.3.1	Discrete normal modes	25
5.4	Sound Waves	26
5.5	The wave equation in 2 dimensions	26
5.5.1	Solutions in exponential form	27
5.5.2	2D wave equation in polar coordinates	27

6	Dispersion and Interference	29
6.1	Phase and group velocities in a non-dispersive medium	29
6.2	Dispersion	30
6.2.1	Light dispersion	30
6.3	Wave interference	30
6.3.1	Conditions for destructive and constructive interference	30
6.3.2	Huygen's Principle	30
6.3.3	Young's Double Slit Experiment	31
6.4	Wave Diffraction	32
6.4.1	Circular aperture	33
6.4.2	Combining with the double slit experiment	33
7	The Schrodinger equation	35
7.1	Properties of the wavefunction	35
7.2	Infinite potential well	35
7.2.1	Time dependent solution	36
7.2.2	Stationary states	36
7.3	Tunelling	36
7.3.1	Tunelling probability	37
7.4	Finite potential well	38
7.4.1	Bound states	39
7.4.2	Unbound states	39
A	Common Waveforms and Examples	41
A.1	Gaussian Pulse	41
A.2	Waves on a square membrane	41
B	Trig Identities	43
C	Solving the time dependent Schrodiner equation	44

Chapter 1

Simple Harmonic Motion

Definition (Simple Harmonic Motion). *The oscillation of a quantity about some equilibrium point due to being trapped in a potential well.*

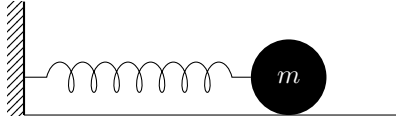


Figure 1.1: A mass resting on a friction-less surface attached to a spring.

Let us first consider a mass extending a spring while on a friction-less surface. We know by Hooke's law (for small displacements),

$$F = -kx.$$

We can apply Newton's second law to get,

$$m\ddot{x} = -kx, \tag{1.1}$$

which we can rewrite,

$$\ddot{x} = -\omega^2 x \tag{1.2}$$

where,

$$\omega^2 = \frac{k}{m}. \tag{1.3}$$

Equation (1.2) has a general solution,

$$x(t) = A \cos(\omega t + \phi) \tag{1.4}$$

where A is the amplitude of the motion and ϕ is the phase.

1.1 Energy in SHM

We know that generally,

$$E = K + U. \tag{1.5}$$

For a spring, this is,

$$E = \frac{1}{2}mv^2 + \frac{1}{2}kx^2, \tag{1.6}$$

which, when substituting equation (1.4), equation (1.6) reduces to,

$$E = \frac{1}{2}kA^2. \tag{1.7}$$

Thus, we can conclude that $E \propto A^2$. Therefore, the more we initially extend the spring, the more potential energy we can store in it.

1.2 SHM In All Systems

Provided that the amplitude of oscillation is small, a system will exhibit simple harmonic motion. This is because small displacements from a system's stable equilibrium resemble a parabola. We can prove this by considering the Maclaurin expansion of the potential energy,

$$U(x) = U_0 + x \left(\frac{dU}{dx} \right)_{x=0} + \frac{1}{2} x^2 \left(\frac{d^2U}{dx^2} \right)_{x=0}.$$

U_0 is an additive constant, so can be set to 0. Since the system is at equilibrium at $x = 0$, $\left(\frac{dU}{dx} \right)_{x=0}$ can also be set to 0. Thus, the potential energy is given by,

$$U(x) \approx \frac{1}{2} x^2 \left(\frac{d^2U}{dx^2} \right)_{x=0}. \quad (1.8)$$

Force is negative gradient of potential, so,

$$F = -x \left(\frac{d^2U}{dx^2} \right)_{x=0}, \quad (1.9)$$

which is the form of the equation of SHM.

1.2.1 Examples of Other Forms of SHM

Pendulum

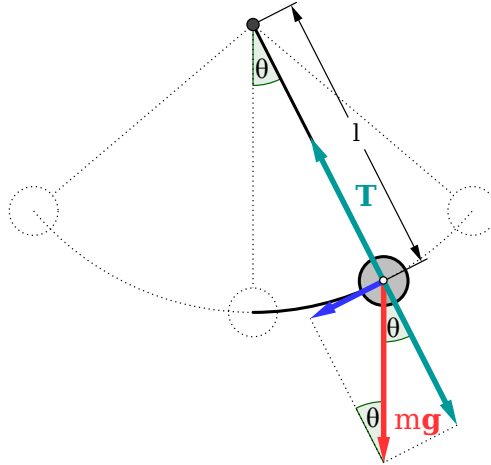


Figure 1.2: Diagram of a pendulum.

By using the small angle approximation, we are able to write the equation of motion of the pendulum,

$$\ddot{\theta} = -\frac{g}{l}\theta. \quad (1.10)$$

for which the energy will be equal to,

$$\begin{aligned} E &= \frac{1}{2}mv^2 + \frac{1}{2} \frac{mgx^2}{l} \\ &= \frac{1}{2} \frac{mgA^2}{l} \end{aligned} \quad (1.11)$$

where x is the horizontal distance from the equilibrium that the pendulum moves.

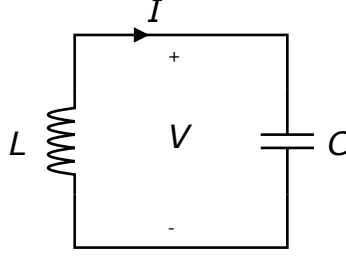


Figure 1.3: Diagram of the LC circuit we are considering.

LC Circuit

Figure 1.3 represents an inducer circuit charging a capacitor. We should expect the current to oscillate between the plates of the conductor. We can begin by understanding that the total charge in the system is equal to 0, which we can write as,

$$L \frac{dI}{dt} + \frac{q}{C} = 0. \quad (1.12)$$

Differentiating this, we get the expression,

$$\frac{d^2 I}{dt^2} = -\frac{1}{LC} I, \quad (1.13)$$

which is an equation of SHM.

1.3 General Form of SHM

From the examples, it is clear that any oscillating quantity can exhibit SHM. We can define the general differential equation as,

$$\alpha \frac{d^2 Z}{dt^2} = -\beta Z. \quad (1.14)$$

whose solution will be,

$$Z(t) = Z_0 \cos \left(\sqrt{\frac{\beta}{\alpha}} t + \phi \right), \quad (1.15)$$

and whose energy can be written as,

$$E = \frac{1}{2} \alpha \left(\frac{dZ}{dt} \right)^2 + \frac{1}{2} \beta Z^2. \quad (1.16)$$

The symbols are defined as follows,

Definition (Z). *This is the oscillating quantity.*

Definition (α). *This is a constant representing the inertia of a system (e.g., mass, moment of inertia, etc.).*

Definition (β). *This is a constant representing the restoring force per unit displacement through which the potential energy is stored (i.e., force due to gravity in a pendulum).*

Definition (Z_0). *The amplitude of the oscillating quantity.*

1.3.1 Period Independent of Mass: Proof by Dimensional Analysis

Chapter 2

Damped Harmonic Motion

Most oscillatory motion feels a damping force, such that,

$$F_D = -bv \quad (2.1)$$

where b is a positive constant. Including this in the equation of motion, we get

$$\ddot{x} + \gamma\dot{x} + \omega_0^2 x = 0. \quad (2.2)$$

where ω_0 is the **undamped natural frequency**, and

$$\gamma = \frac{b}{m}. \quad (2.3)$$

We should note that the the force being linearly proportional to velocity is only a very good approximation, and works when the velocity of the system is small.

2.1 Light Damping

For light damping, the oscillator retains its oscillatory behaviour but its amplitude decays with time. The solution to equation (2.2) is,

$$x(t) = Ae^{-\beta t} \cos(\omega t + \phi) \quad (2.4)$$

where,

$$\beta = \frac{\gamma}{2}, \quad (2.5)$$

and

$$\omega^2 = \omega_0^2 - \frac{\gamma^2}{4} \quad (2.6)$$

which is called the **damped angular frequency**. A key property of damped harmonic motion is that,

$$\omega_0^2 > \frac{\gamma^2}{4}. \quad (2.7)$$

2.2 Heavy Damping

In heavy damping, oscillations are suppressed and the system stops moving before it completes any oscillations. In this case, the solution will be in the form,

$$x(t) = e^{-\beta t} f(t) \quad (2.8)$$

where $f(t)$ is an unknown function. This gives way to a differential equation,

$$\ddot{f}(t) + \left(\omega_0^2 - \frac{\gamma^2}{4}\right) f(t) = 0. \quad (2.9)$$

Heavily damped motion will meet the condition,

$$\omega_0^2 < \frac{\gamma^2}{4}. \quad (2.10)$$

We will then define,

$$\alpha^2 = \frac{\gamma^2}{4} - \omega_0^2 \quad (2.11)$$

which allows us to find a solution,

$$f(t) = Ae^{\alpha t} + Be^{-\alpha t}, \quad (2.12)$$

and thus an expression for the displacement,

$$x(t) = e^{-\beta} (Ae^{\alpha t} + Be^{-\alpha t}). \quad (2.13)$$

2.3 Critical Dampinig

Critical damping occurs when the system returns to equilibrium in the shortest time without oscillating. It satisfies the condition, below, whose consequences then follow,

$$\frac{\gamma^2}{4} = \omega_0^2 \implies \alpha = 0 \implies \ddot{f}(t) = 0. \quad (2.14)$$

Thus,

$$f(t) = At + B \quad (2.15)$$

and

$$x(t) = e^{-\beta}(At + B). \quad (2.16)$$

2.4 Quality Factor

Definition. *The ratio of energy stored to energy lost in 1 radian cycle of oscillation.*

$$Q = \frac{\omega_0}{\gamma} \quad (2.17)$$

We can write a relation between the damped angular frequency and the natural frequency using Q ,

$$\omega = \omega_0 \left(1 - \frac{1}{4}Q^{-2}\right)^{\frac{1}{2}}. \quad (2.18)$$

2.5 LCR Circuit

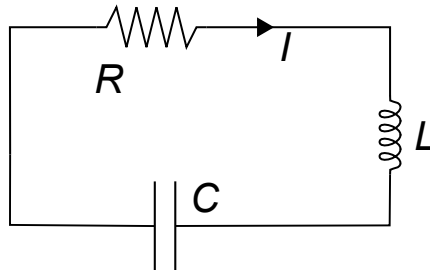


Figure 2.1: LCR Circuit

By adding a resistor to an LC circuit we can add a form of damping. The voltage across the resistor is given by,

$$V_R = RI = R\dot{q}. \quad (2.19)$$

So, the ODE becomes,

$$L\ddot{q} + R\dot{q} + \frac{1}{C}q = 0. \quad (2.20)$$

We can then conclude that the, for light damping,

$$q = q_0 \exp\left(\frac{R}{2L}t\right) \cos(\omega t) \quad (2.21)$$

where,

$$\omega^2 = \frac{1}{LC} - \frac{R^2}{4L^2}, \quad (2.22)$$

and

$$Q = \frac{1}{R} \sqrt{\frac{L}{C}}. \quad (2.23)$$

2.6 General Damped Harmonic Oscillator

Just as with SHM, we can describe a general equation of a damped harmonic oscillator,

$$\alpha\ddot{Z} + \kappa\dot{Z} + \beta Z = 0 \quad (2.24)$$

whose solution will be,

$$Z(t) = Z_0 \exp\left(\frac{\kappa}{2\alpha}t\right) \cos\left(\sqrt{\frac{\beta}{\alpha}} + \phi\right). \quad (2.25)$$

We introduced a new symbol,

Definition (κ). *The constant representing the magnitude of the damping force.*

Chapter 3

Driven, Damped Harmonic Oscillator

This is oscillatory motion where a driving force is provided. The driving force is given by,

$$F(t) = F_0 \cos(\omega t). \quad (3.1)$$

We can then modify the equation of motion,

$$\ddot{x} + \gamma \dot{x} + \omega_0^2 x = \frac{F_0}{m} \cos \omega t, \quad (3.2)$$

which we can rewrite in its exponential form,

$$\ddot{z} + \gamma \dot{z} + \omega_0^2 z = \frac{F_0}{m} e^{i\omega t}. \quad (3.3)$$

NOTE: When we write functions in exponential form, we represent physical values with **only** the real part. We can then intuit an expression for z ,

$$z(t) = A(\omega) e^{(i\omega t + \delta)} \quad (3.4)$$

where δ is the phase between the actual oscillations and the driving force. Differentiating, substituting back into equation (3.3), and then dividing by $e^{(i\omega t + \delta)}$. The equation then has real and imaginary parts,

$$(\omega_0^2 - \omega^2)A(\omega) = \frac{F_0}{m} \cos \delta \quad \text{Re}, \quad (3.5)$$

$$\gamma\omega A(\omega) = \frac{F_0}{m} \sin \delta \quad \text{Im}. \quad (3.6)$$

The phase is then given by,

$$\tan \delta = \frac{\gamma\omega}{\omega_0^2 - \omega^2}, \quad (3.7)$$

and the amplitude function is given by,

$$A(\omega) = \frac{F_0}{m\sqrt{(\omega_0^2 - \omega^2)^2 + (\omega\gamma)^2}}. \quad (3.8)$$

We must consider the limits of the function,

$$\omega \rightarrow 0 \quad A(\omega) \rightarrow \frac{F_0}{m\omega_0^2} \quad (3.9)$$

$$\omega \rightarrow \infty \quad A(\omega) \rightarrow 0 \quad (3.10)$$

$$\omega \rightarrow \omega_0 \quad A(\omega) \rightarrow \frac{F_0}{m\omega_0\gamma} \quad (3.11)$$

$$A_{\text{MAX}} = \frac{F_0}{m\omega_0\gamma} \frac{1}{\sqrt{1 - \left(\frac{\gamma}{2\omega_0}\right)^2}} \quad A_{\text{MAX}} \approx A(\omega_0) \rightarrow \text{Light damping} \quad (3.12)$$

3.1 Transient Response

The previous solution is what dominates. There is also the immediate transient response, which dies down pretty quickly. This will simply be the homogeneous solution to the equation, which is the solution to the damped harmonic oscillator.

3.2 AC Circuit

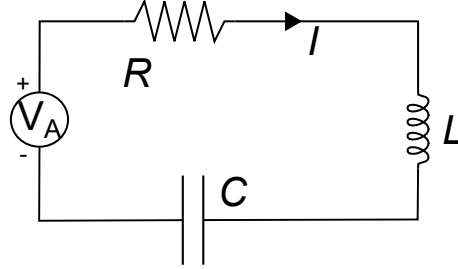


Figure 3.1: An LCR circuit with added AC.

We can add a driving AC, with a voltage,

$$V_A = V_0 e^{i\omega t}, \quad (3.13)$$

which modifies the equation as so,

$$L\ddot{q} + R\dot{q} + \frac{1}{C}q = V_0 e^{i\omega t} \quad (3.14)$$

whose solution is,

$$q = A(\omega) e^{i(\omega t - \delta)} \quad (3.15)$$

where

$$A(\omega) = \frac{V_0}{\sqrt{(\omega_0^2 - \omega^2)^2 + (\gamma\omega)^2}} \quad (3.16)$$

and

$$\omega_0^2 = \frac{1}{LC} \quad \gamma = \frac{R}{L}. \quad (3.17)$$

3.2.1 Impedance

The voltage and current in the resistor are in phase with the AC supply, however in the inductor and capacitor they are not. We can find the phase of the current by,

$$I(t) = I_0 e^{i\omega t}, \quad (3.18)$$

then the voltage is given by,

$$V_L = L \frac{dI}{dt} = i\omega L I_0 e^{i\omega t} \quad (3.19)$$

which implies that V_L leads I by a phase of $\frac{\pi}{2}$.

We are able to do the same for the capacitance,

$$\begin{aligned} V_C &= \frac{1}{C} q \\ &= \frac{1}{C} \int I dt \\ &= \frac{1}{i\omega C} I_0 e^{i\omega t} = -\frac{i}{\omega C} I. \end{aligned} \quad (3.20)$$

Thus, V_C lags I by $\frac{\pi}{2}$.

We can represent this lead and lag through Ohm's law,

$$V = IZ \quad (3.21)$$

where Z is the **impedance**. This is a complex function where the real part represents the **resistance** and the imaginary part represents the **reactance**.

$$Z = R + i \left(\omega L - \frac{1}{\omega C} \right). \quad (3.22)$$

3.3 Absorbed power

The power absorbed by a driven, damped oscillator is equal to the rate at which energy is dissipated. So,

$$P = \frac{W}{\Delta t} = F_D v. \quad (3.23)$$

We can then derive the instantaneous power,

$$P(t) = b(v_0(\omega))^2 \sin^2(\omega t - \delta). \quad (3.24)$$

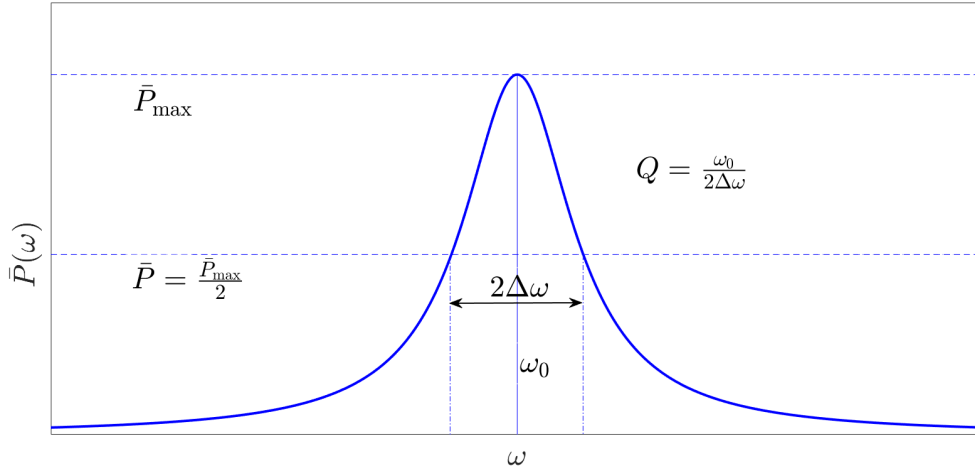
We can then get the average power absorbed per cycle by evaluating the integral,

$$\bar{P}(\omega) = \frac{1}{T} \int_{t_0}^{t_0+T} b(v_0(\omega))^2 \sin^2(\omega t - \delta), \quad (3.25)$$

which gives,

$$\bar{P}(\omega) = \frac{1}{2} b(v_0(\omega))^2 - \frac{\omega^2 F_0^2 \gamma^2}{2m \left[(\omega^2 - \omega_0^2)^2 + \omega^2 \gamma^2 \right]}. \quad (3.26)$$

From this function, we can obtain a power-resonance curve.



The width of the curve where $\omega = \omega_0$ defines the sharpness of the resonance. When close to the peak, we say $\omega \sim \omega_0$, so $\omega + \omega_0 \sim 2\omega_0$ and we can define $\Delta\omega = \omega - \omega_0$. Further,

$$\omega_0^2 - \omega^2 = (\omega_0 + \omega)(\omega_0 - \omega) = 2\omega_0 \Delta\omega. \quad (3.27)$$

So, we can rewrite our power curve equation as,

$$\bar{P}(\omega) = \frac{F_0^2}{2m\gamma \left(\frac{4\Delta\omega^2}{\gamma^2} + 1 \right)}. \quad (3.28)$$

The maximum value of power is then,

$$\bar{P}_{\max} = \frac{F_0^2}{2m\gamma}. \quad (3.29)$$

We can define the full width at half height (ω_{fwhh}) where $\bar{P}(\omega) = \frac{1}{2}\bar{P}_{\max}$. Thus,

$$\omega_{\text{fwhh}} = 2\Delta\omega = \gamma. \quad (3.30)$$

This allows us to redefine the quality factor,

$$Q = \frac{\omega_0}{\gamma} = \frac{\omega_0}{\omega_{\text{fwhh}}}. \quad (3.31)$$

Chapter 4

Coupled Oscillators

Coupled oscillators are joined together. We will first consider two identical pendula, A and B of mass m and length l coupled by a spring of spring constant k . We will consider two normal modes of oscillation (see figure 4.1).

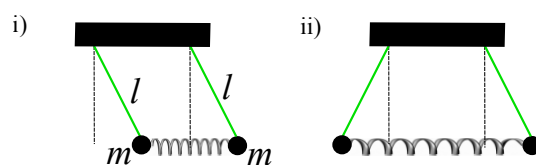


Figure 4.1

	Description	Implications
i)	Each pendulum is displaced by the same amount in the same direction	<ul style="list-style-type: none"> • They will oscillate with the same frequency $\omega_1^2 = \frac{g}{l}$. • They will be in phase.
ii)	Each pendulum is displaced by the same amount but in opposite directions.	<ul style="list-style-type: none"> • They oscillate with the same frequency, ω_2, which is different from ω_1. • They will be out of phase.

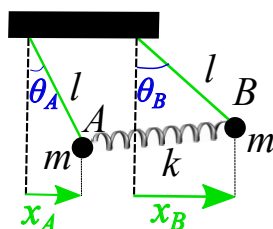


Figure 4.2

4.1 Equations of motion in coupled pendula

We will now consider the case where $x_B > x_A$ as in figure 4.2. The tension in the spring is given by,

$$T = k(x_B - x_A) \quad (4.1)$$

so the restoring force on B and A, using the small angle approximation, is given by,

$$F_B = -\frac{mgx_B}{l} - k(x_B - x_A) \quad F_A = -\frac{mgx_A}{l} + k(x_B + x_A). \quad (4.2)$$

This gives the equations of motion,

$$\ddot{x}_A + \frac{g}{l}x_A - \frac{k}{m}(x_B - x_A) = 0 \quad (4.3)$$

$$\ddot{x}_B + \frac{g}{l}x_B - \frac{k}{m}(x_B - x_A) = 0. \quad (4.4)$$

Now, let us define,

$$q_n = \begin{cases} n=1 & x_A + x_B \\ n=2 & x_A - x_B \end{cases} \quad (4.5)$$

which are the **normal coordinates** for our setup. We can then perform (4.3) + (4.4) and (4.3) - (4.4) respectively,

$$\ddot{q}_1 + \frac{g}{l}q_1 = 0 \quad (4.6)$$

$$\ddot{q}_2 + \left(\frac{g}{l} + \frac{2k}{m}\right)q_2 = 0, \quad (4.7)$$

where we can say,

$$\omega_n^2 = \begin{cases} n=1 & \frac{g}{l} \\ n=2 & \left(\frac{g}{l} + \frac{2k}{m}\right) \end{cases}. \quad (4.8)$$

As we can see, the normal coordinates are useful as they allow us to write decoupled versions of the equations of motion, representing the two normal modes of the motion of the pendulums. It can be viewed experimentally that when we excite a coupled system in a certain normal mode, it will stay in that normal mode, or a combination of normal modes as it cannot transfer energy between them. Thus, we can solve and model *any* motion of the coupled pendula by using a number of linear combinations of the normal mode equations.

For the example presented here, the most general solutions take the form,

$$x_a = C \cos(\omega_1 t - \delta_1) + D \cos(\omega_2 t - \delta_2) \quad (4.9)$$

$$x_b = C \cos(\omega_1 t - \delta_1) - D \cos(\omega_2 t - \delta_2) \quad (4.10)$$

The sign change occurs because we know for the normal mode ω_1 , $x_a = x_b$, so the signs are the same. We change sign for the term representing the second normal mode ω_2 as we have the condition $x_a = -x_b$.

4.1.1 Normal Modes for higher n

We previously considered only the normal mode of a coupled oscillator by a single mass, $n = 1$. Higher n require Eigen analysis.

2 Masses - $n = 2$

We will be considering the setup in figure 4.3. We understand the governing equations are,

$$m\ddot{x}_A = -kx_A + k'(x_B - x_A) \quad (4.11)$$

$$m\ddot{x}_B = -kx_B - k'(x_B - x_A). \quad (4.12)$$

We will denote,

$$\omega^2 = \frac{k}{m} \quad \omega'^2 = \frac{k'}{m},$$

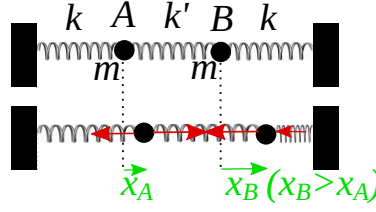


Figure 4.3: Two masses connected by springs.

We then rewrite (4.11) and (4.12),

$$\ddot{x}_A + \omega^2 x_A + \omega'^2 (x_A - x_B) = 0 \quad (4.13)$$

$$\ddot{x}_B + \omega^2 x_B + \omega'^2 (x_B - x_A) = 0. \quad (4.14)$$

We can assume that both of the masses will have the same frequency, Ω . We can then assume trial equations,

$$x_A = X_A e^{i\Omega t}; \quad x_B = X_B e^{i\Omega t}. \quad (4.15)$$

Substituting (4.15) into (4.13) and (4.14), and dividing by $e^{i\Omega t}$,

$$-\Omega^2 X_A + \omega^2 X_A + \omega'^2 (X_A - X_B) = 0 \quad (4.16)$$

$$-\Omega^2 X_B + \omega^2 X_B + \omega'^2 (X_B - X_A) = 0. \quad (4.17)$$

Then, collecting like terms,

$$(\omega^2 + \omega'^2) X_A - \omega'^2 X_B = \Omega^2 X_A \quad (4.18)$$

$$-\omega'^2 X_A + (\omega^2 + \omega'^2) X_B = \Omega^2 X_B. \quad (4.19)$$

We can rewrite this in vector form,

$$\begin{pmatrix} \omega^2 + \omega'^2 & -\omega'^2 \\ -\omega'^2 & \omega^2 + \omega'^2 \end{pmatrix} \begin{pmatrix} X_A \\ X_B \end{pmatrix} = \Omega^2 \begin{pmatrix} X_A \\ X_B \end{pmatrix}.$$

We will notice that this equation is in the form,

$$\mathbf{A}\mathbf{x} = \lambda\mathbf{x}$$

thus, we can perform Eigen analysis on this problem by solving the equation $|\mathbf{A} - \Omega^2 \mathbf{I}|$,

$$\begin{aligned} (\omega^2 + \omega'^2 - \Omega^2)^2 - \omega'^4 &= 0 \\ \omega^2 + \omega'^2 - \Omega^2 &= \pm \omega'^2. \end{aligned}$$

Thus, we have two Eigen values,

$$\Omega^2 = \omega^2 = \frac{k}{m} \quad (4.20)$$

$$\Omega^2 = \omega^2 + 2\omega'^2 = \frac{k}{m} + \frac{2k'}{m}. \quad (4.21)$$

We then have two normal modes for the system,

1. By substituting (4.20) into (4.18),

$$X_A = X_B \quad (4.22)$$

2. By substituting (4.21) into (4.18) or (4.19),

$$X_A = -X_B \quad (4.23)$$

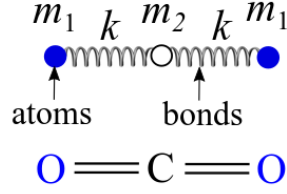


Figure 4.4

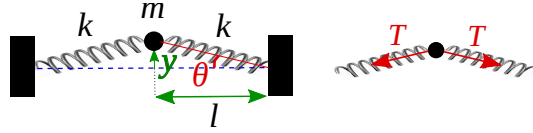


Figure 4.5

Linear Triatomic Molecule $n = 3$

Look at figure 4.4. We can model this as a 3-mass problem moving in 1 dimension. We will define the factors,

$$\omega^2 = \frac{k}{m_1} \qquad \omega'^2 = \frac{k}{m_2} \qquad x_n(t) = x_n e^{i\Omega t} \quad (4.24)$$

for $n = \{1, 2, 3\}$. When we write down the governing equations, we are able to once again perform eigen analysis on the system,

$$\begin{pmatrix} \omega^2 & -\omega^2 & 0 \\ -\omega'^2 & 2\omega'^2 & -\omega'^2 \\ 0 & -\omega^2 & \omega^2 \end{pmatrix} \begin{pmatrix} X_1 \\ X_2 \\ X_3 \end{pmatrix} = \Omega^2 \begin{pmatrix} X_1 \\ X_2 \\ X_3 \end{pmatrix}. \quad (4.25)$$

Performing Eigen analysis, we obtain 3 Eigen values and 3 normal modes.

1. $\Omega^2 = 0$, transnational motion.

$$\begin{aligned} \omega^2 X_1 - \omega^2 X_2 &= 0 \implies X_1 = X_2 \\ -\omega^2 X_2 + \omega^2 X_3 &= 0 \implies X_2 = X_3 \end{aligned}$$

2. $\Omega^2 = \omega^2$, symmetric stretching.

$$X_2 = 0, \qquad X_1 = -X_3. \quad (4.26)$$

3. $\Omega^2 = 2\omega'^2 + \omega^2$

$$X_3 = X_1 = -\frac{\omega^2}{2\omega'^2} X_2$$

4.2 Transverse Oscillations

These are oscillations which are displaced in the y direction.

4.2.1 Single Oscillator

We are modelling a setup like that in figure ???. We write this motion simply by asserting that T remains constant.

Proof. For a displacement y , each spring displaces by an amount Δl such that,

$$(l + \Delta l) \cos \theta = l$$

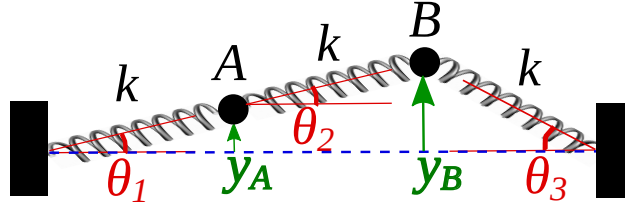


Figure 4.6

where

$$\theta = \arctan\left(\frac{y}{l}\right).$$

We can write,

$$\Delta l = l \left(\frac{1}{\cos \theta} - 1 \right).$$

For small θ , we can use the Taylor expansion of $\cos \theta$ to write Δl as,

$$\Delta l = l \frac{\theta^2}{2} \quad (4.27)$$

if θ is small. Thus, θ^2 can be neglected, and the tension T is constant. \square

Our equation of motion is then,

$$m\ddot{y} = -2T \sin \theta$$

and

$$\sin \theta \approx \tan \theta = \frac{y}{l} \quad (4.28)$$

so, we can rewrite,

$$\ddot{y} = -\frac{2T}{ml} y \quad (4.29)$$

which is an equation of SHM.

4.2.2 2 Oscillators

We are now modelling the setup in figure ???. Making the same assumptions as before, we can write the equations of motion,

$$m\ddot{y}_A = -T \sin \theta_1 + T \sin \theta_2 = \frac{T}{l}(y_B - 2y_A) \quad (4.30)$$

$$m\ddot{y}_B = -T \sin \theta_2 - T \sin \theta_3 = \frac{T}{l}(y_A - 2y_B). \quad (4.31)$$

We can once again apply Eigen analysis and find the Eigen values to find the normal modes,

$$\Omega^2 = \frac{T}{ml} \quad Y_A = Y_B \quad (4.32)$$

$$\Omega^2 = \frac{3T}{ml} \quad Y_A = -Y_B \quad (4.33)$$

4.2.3 Higher n

At higher n , the normal modes start being analogous to standing waves.

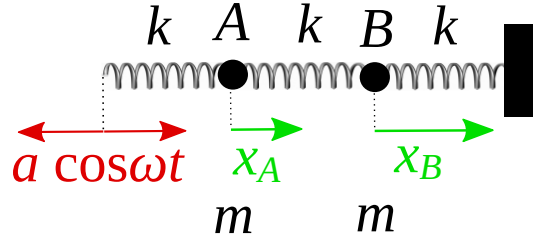


Figure 4.7

4.3 Forced Coupled Oscillators

Let us consider the setup as in figure 4.7. This is a 2 mass system, with A being connected to a sinusoidal drive with displacement $\epsilon = a \cos \omega t$. We have,

$$\ddot{x}_A + \frac{2k}{m}x_A - \frac{k}{m}x_B = \frac{F_0}{m} \cos \omega t \quad (4.34)$$

$$\ddot{x}_B - \frac{k}{m}x_A + \frac{2k}{m}x_B = 0 \quad (4.35)$$

where $F_0 = ka$. We can define the normal coordinates like we did in (4.5). Our equations of motion then become,

$$\ddot{q}_1 + \frac{k}{m}q_1 = \frac{F_0}{m} \cos \omega t \quad (4.36)$$

$$\ddot{q}_2 + \frac{3k}{m}q_2 = \frac{F_0}{m} \cos \omega t \quad (4.37)$$

These are both equations of undamped, forced harmonic motion. We know then that the solutions are,

$$q_1 = C_1(\omega) \cos(\omega t - \delta_1) \quad q_2 = C_2(\omega) \cos(\omega t - \delta_2) \quad (4.38)$$

where we define the following,

$$\omega_n^2 = \begin{cases} \frac{k}{m} & n = 1 \\ \frac{3k}{m} & n = 2 \end{cases} \quad (4.39)$$

$$\delta_n = \begin{cases} 0 & \omega < \omega_n \\ \pi & \omega > \omega_n \end{cases} \quad (4.40)$$

$$C_n = \frac{F_0/m}{\omega_n^2 - \omega^2} \quad (4.41)$$

Chapter 5

Waves

There are two types of waves, *standing* and *travelling* waves.

Travelling waves come in two forms,

1. **Longitudinal** - Change in physical quantity is in direction of travel.
2. **Transverse** - Change is orthogonal to direction of travel.

Further, these may either be,

1. **Mechanical** - requires a medium.
2. **Non-mechanical** - does not require a medium.

5.1 1D Wave Equation

All waves are generated by a wave equation. For a distribution $y(x, t)$, the wave equation is given by,

$$\frac{\partial^2 y}{\partial t^2} = v^2 \frac{\partial^2 y}{\partial x^2}, \quad (5.1)$$

where v is the velocity at which the wave propagates.

5.1.1 Sinusoidal Waves

Sinusoidal waves can be given by,

$$y(x, t) = A \sin \left[\frac{2\pi}{\lambda} (x - vt) \right]. \quad (5.2)$$

A travelling sinusoidal wave will have frequency,

$$\nu = \frac{v}{\lambda}. \quad (5.3)$$

We can then define other properties of the wave,

$$\nu = \frac{1}{T} \quad \omega = 2\pi\nu. \quad (5.4)$$

We can define the wave number, $k = \frac{2\pi}{\lambda}$. We can then re-write the equation as,

$$y(x, t) = A \sin(kx - \omega t) \quad (5.5)$$

and we can rewrite the velocity of propagation as,

$$v = \frac{\omega}{k}. \quad (5.6)$$

5.1.2 General solution to 1D wave equation

The general solution takes the form,

$$y(x, t) = f(x - vt) + g(x + vt). \quad (5.7)$$

Proof. Let us first define,

$$u = x - vt. \quad (5.8)$$

We then differentiate f with respect to x and t .

$$1. \quad \frac{\partial f}{\partial x} = \frac{df}{du} \frac{\partial u}{\partial x} \quad (5.9)$$

$$2. \quad \frac{\partial^2 f}{\partial x^2} = \frac{\partial}{\partial x} \left(\frac{\partial f}{\partial u} \frac{\partial u}{\partial x} \right) = \frac{d^2 f}{du^2} \left(\frac{\partial u}{\partial x} \right)^2 + \frac{df}{du} \frac{\partial^2 u}{\partial x^2} \quad (5.10)$$

$$\frac{\partial u}{\partial x} = 1 \quad \frac{\partial^2 u}{\partial x^2} = 0 \quad (5.11)$$

$$\frac{\partial^2 f}{\partial x^2} = \frac{d^2 f}{du^2} \quad (5.12)$$

$$3. \quad \begin{aligned} \frac{\partial f}{\partial t} &= \frac{df}{du} \frac{\partial u}{\partial t} = -v \frac{df}{du} \\ \Rightarrow \frac{\partial^2 f}{\partial t^2} &= v^2 \frac{d^2 f}{du^2} \end{aligned} \quad (5.13)$$

Combining (5.13) and (5.12),

$$\frac{\partial^2 f}{\partial t^2} = v^2 \frac{\partial^2 f}{\partial x^2}. \quad (5.14)$$

We similarly get,

$$\frac{\partial^2 g}{\partial t^2} = v^2 \frac{\partial^2 g}{\partial x^2}. \quad (5.15)$$

Adding (5.14) and (5.15), we get,

$$\frac{\partial^2 f + g}{\partial t^2} = v^2 \frac{\partial^2 f + g}{\partial x^2}, \quad (5.16)$$

which is the same form as the 1D wave equation. So, (5.7) is a general solution of the 1D wave equation. \square

5.1.3 Velocity of travelling waves on a stretched string

A string under tension T and density μ can support transverse, travelling waves. We will consider a length of element of the string between x and $x + \delta x$. We will consider that at x the string is at an angle θ to the horizontal, and at an angle $\theta + \delta\theta$ at a distance $x + \delta x$. (See figure 5.1).

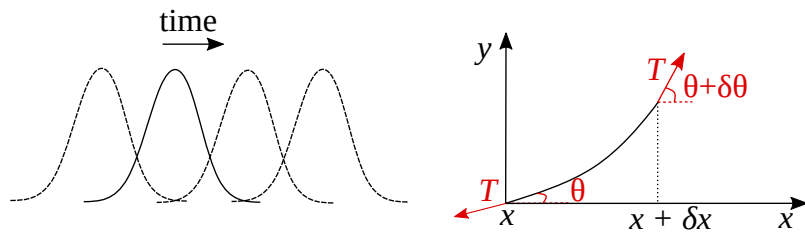


Figure 5.1: Stretched strings allow for waves to travel along their length.

The tension in the string at x ,

$$\begin{aligned} T \sin \theta &\sim \theta \sim \tan \theta \approx \frac{\delta y}{\delta x} \implies \sin \theta = \frac{\partial y}{\partial x} \\ T \sin \theta &= T \left(\frac{\partial y}{\partial x} \right)_x. \end{aligned} \quad (5.17)$$

Similarly, tension at $x + \delta x$ is given by,

$$T \sin(\theta + \delta\theta) \sim T \left(\frac{\partial y}{\partial x} \right)_{x+\delta x}. \quad (5.18)$$

We will now perform a first order Taylor expansion of (5.18) about x ,

$$\begin{aligned} T \left(\frac{\partial y}{\partial x} \right)_{x+\delta x} &= T \left\{ \left(\frac{\partial y}{\partial x} \right)_x + \frac{\partial}{\partial x} \left(\frac{\partial y}{\partial x} \right)_x \delta x \right\} \\ &= T \left\{ \left(\frac{\partial y}{\partial x} \right)_x + \left(\frac{\partial^2 y}{\partial x^2} \right)_x \delta x \right\}. \end{aligned} \quad (5.19)$$

The net transverse force F_y is the difference between the tensions at x and $x + \delta x$,

$$F_y = T \left\{ \left(\frac{\partial y}{\partial x} \right)_x + \left(\frac{\partial^2 y}{\partial x^2} \right)_x \delta x - \left(\frac{\partial y}{\partial x} \right)_x \right\} = T \frac{\partial^2 y}{\partial x^2} \delta x. \quad (5.20)$$

Similarly, for the x component,

$$F_x = T \cos(\theta + \delta\theta) - T \cos \theta \simeq T - T \simeq 0 \quad (5.21)$$

because,

$$\cos(\theta + \delta\theta) \simeq \cos \theta \simeq 1 \text{ for small } \theta. \quad (5.22)$$

We can apply a transverse disturbance to the string by accelerating it in the y direction. By Newton's second law,

$$\mu \delta x \frac{\partial^2 y}{\partial t^2} = T \frac{\partial^2 y}{\partial x^2} \delta x \implies \frac{\partial^2 y}{\partial t^2} = \frac{T}{\mu} \frac{\partial^2 y}{\partial x^2} \quad (5.23)$$

which is the 1D wave equation for a wave with velocity

$$v = \sqrt{\frac{T}{\mu}}. \quad (5.24)$$

5.1.4 Energy Stored in a Wave

Let us consider a string under tension. Let us consider a horizontal element δx which has mass $\mu \delta x$ which oscillates in the y direction. Its kinetic energy is given by,

$$\delta K = \frac{1}{2} (\mu \delta x) \left(\frac{\partial y}{\partial t} \right)^2. \quad (5.25)$$

Further, when not at equilibrium, the string is stretched so that it has an elemental, stretched length δs under a tension T . Its potential energy can be then given by,

$$\delta U = (\delta s - \delta x) T. \quad (5.26)$$

We understand,

$$\delta s = (\delta x^2 + \delta y^2)^{\frac{1}{2}} \approx \delta x \left(1 + \frac{1}{2} \left(\frac{\partial y}{\partial x} \right)^2 \right), \quad (5.27)$$

by binomial expansion, valid for small displacements. Then,

$$\delta s - \delta x = \frac{1}{2} \left(\frac{\partial y}{\partial x} \right)^2 \delta x. \quad (5.28)$$

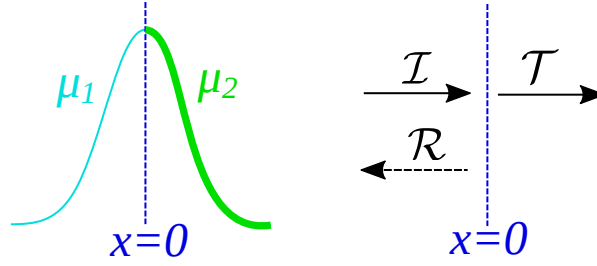


Figure 5.2

Thus, we can write the total energy element,

$$\delta E = \delta U + \delta K = \frac{1}{2}\mu\delta x \left(\frac{\partial y}{\partial t} \right)^2 + \frac{1}{2} \left(\frac{\partial y}{\partial x} \right)^2 T. \quad (5.29)$$

By (5.24),

$$\delta E = \frac{1}{2}\mu\delta x \left[\left(\frac{\partial y}{\partial t} \right)^2 + v^2 \left(\frac{\partial y}{\partial x} \right)^2 \right], \quad (5.30)$$

which we can integrate,

$$E = \frac{1}{2} \int_a^b \mu \left[\left(\frac{\partial y}{\partial t} \right)^2 + v^2 \left(\frac{\partial y}{\partial x} \right)^2 \right] dx \quad (5.31)$$

to find the energy stored in the string between $x = a \rightarrow x = b$.

Power of a sinusoidal wave

Let us consider a solution to the 1D wave equation $y = A \sin(kx - \omega t)$. By substituting this into (5.26) and (5.25) and integrating from $0 \rightarrow \lambda$, we can obtain the kinetic and potential energy over 1 wavelength,

$$K_\lambda = \frac{1}{4}\mu\omega^2 A^2 \lambda \quad U_\lambda = \frac{1}{4}\mu\omega^2 A^2 \lambda, \quad (5.32)$$

and adding these yields,

$$E_\lambda = \frac{1}{2}\mu\omega^2 A^2 \lambda. \quad (5.33)$$

The power is then $E_\lambda \times \frac{v}{\lambda}$, and thus,

$$P = \frac{1}{2}\mu\omega^2 A^2 v. \quad (5.34)$$

5.2 Waves at Boundaries

When waves encounter discontinuity, some fraction of the wave will be reflected back to us. At this discontinuity, there will be a *incident*, *reflected*, and *transmitted* wave at the discontinuity.

To illustrate this, consider figure 5.2, where the density of a string changes from μ_1 to μ_2 at the boundary $x = 0$. The incident wave is described by,

$$y_I = \mathcal{I} \cos(k_1 x - \omega t) \quad (5.35)$$

the reflected wave is described by,

$$y_R = \mathcal{R} \cos(k_2 x + \omega t) \quad (5.36)$$

and the transmitted wave is given by,

$$y_T = \mathcal{T} \cos(k_2 x - \omega' t). \quad (5.37)$$

At the boundary, μ changes, and thus v must also change. **However**, the string is continuous so the frequency of the exciting pulse must be equal in both regions $\omega = \omega'$, thus, k , and, equivalently, λ must change. We understand $\omega = vk$, so,

$$\omega = v_1 k_1 = \omega' = v_2 k_2 \quad (5.38)$$

hence,

$$\frac{k_2}{k_1} = \frac{v_1}{v_2} = \sqrt{\frac{\mu_2}{\mu_1}} = \frac{\lambda_1}{\lambda_2} = n \quad (5.39)$$

where n is known as the *refractive index*.

We can further state that $\frac{\partial y}{\partial x}$ is continuous across the boundary, because the restoring force is continuous across the boundary and the tension T is constant. We can apply the continuity of this gradient as well as the continuity of y in order to better describe the waves,

$$y_I + y_R = y_T \quad (5.40)$$

which, when substituting and evaluating at $x = 0$, yields the first condition,

$$\mathcal{I} + \mathcal{R} = \mathcal{T}. \quad (5.41)$$

We can obtain the second condition by evaluating $\frac{\partial y}{\partial x}$ at $x = 0$,

$$-k_1 \mathcal{I} + k_1 \mathcal{R} = -k_2 \mathcal{T}. \quad (5.42)$$

Using these we are able to express the reflected and transmitted amplitudes as a function of incident amplitude,

$$\frac{\mathcal{T}}{\mathcal{I}} = \frac{2k_1}{k_1 + k_2} = \frac{2}{n + 1} = \frac{2\sqrt{\mu_1}}{\sqrt{\mu_1} + \sqrt{\mu_2}} = \mathcal{T}_c \quad \text{Transmission Coefficient} \quad (5.43)$$

$$\frac{\mathcal{R}}{\mathcal{I}} = \frac{k_1 - k_2}{k_1 + k_2} = \frac{1 - n}{1 + n} = \mathcal{R}_c \quad \text{Reflection Coefficient} \quad (5.44)$$

Below are listed some properties of the coefficients,

1. $0 < \mathcal{T}_c < 2 \implies$ Transmitted wave is always in phase with the incident wave
2. $-1 < \mathcal{R}_c < 1$ and $\mathcal{R}_c < 0 \iff k_2 > k_1 \implies$ The reflected wave is only in phase with the incident wave if $n < 1$.
3. If $k_1 \sim k_2$, \mathcal{R}_c is small, \implies there is hardly any reflection.
4. When $k_2 \rightarrow \infty$, $\mathcal{T}_c \rightarrow 0$, and $\mathcal{R}_c \rightarrow -1 \implies$ no transmission occurs, we get total reflection, which results in *standing waves*.
5. $\mathcal{T}_c = 1 + \mathcal{R}_c$.

5.3 Standing Waves

Before we talk about standing waves, we must understand the **principle of superposition**. It states that if $y_1(x, t)$ and $y_2(x, t)$ are both solutions to the wave equation, the linear combination,

$$y(x, y) = A_1 y_1(x, t) + A_2 y_2(x, t) \quad (5.45)$$

is also a solution. This follows from the fact that the wave equation is linear.

Proof. Consider the wave equation for y_1 and y_2 . Multiply each by A_1 and A_2 respectively and add them together,

$$\begin{aligned} A_1 \frac{\partial^2 y_1}{\partial t^2} + A_2 \frac{\partial^2 y_2}{\partial t^2} &= v^2 A_1 \frac{\partial^2 A_1}{\partial x^2} = v^2 A_2 \frac{\partial^2 y_2}{\partial x^2} \\ \implies \frac{\partial^2}{\partial t^2} (A_1 y_1 + A_2 y_2) &= v^2 \frac{\partial^2}{\partial x^2} (A_1 y_1 + A_2 y_2) \end{aligned} \quad (5.46)$$

thus $A_1 y_1 + A_2 y_2$ is also a solution to the wave equation. \square

Considering standing waves, let us imagine holding a string at a fixed length from $x = 0$ to $x = L$. We have boundary conditions such that $\mathcal{T}_c = 0$ and $\mathcal{R}_c = -1$. The wave can be described the superposition of two waves of the same frequency and amplitude travelling in opposite directions. Then,

$$\begin{aligned} y &= y_{\mathcal{I}} + y_{\mathcal{R}} = Ie^{i(\omega t - kx)} - Ie^{i(\omega t + kx)} \\ &= Ie^{i\omega t} (e^{-ikx} - e^{ikx}) \\ &= f(t)g(x). \end{aligned} \quad (5.47)$$

Thus, we have separated the wave into its time and spatial components. Given a fixed x , all points vibrate at the same frequency, but the spatial patterns vary along the string.

EXAMPLE: Consider the incident wave as, $y_{\mathcal{I}}(x, t) = \frac{A}{2} \sin(kx - \omega t)$. The superposition of the incident and reflected waves gives,

$$y = y_{\mathcal{I}} + y_{\mathcal{R}} = \frac{A}{2} \sin(kx - \omega t) + \frac{A}{2} \sin(kx + \omega t). \quad (5.48)$$

Using the trigonometric identity,

$$\sin(\alpha + \beta) + \sin(\alpha - \beta) = 2 \sin \alpha \cos \beta \quad (5.49)$$

we get,

$$y = A \sin kx \cos \omega t. \quad (5.50)$$

5.3.1 Discrete normal modes

Given that the string is fixed, y must satisfy the boundary conditions,

$$y(0, t) = 0 \quad \implies \quad \sin(k \cdot 0) = 0 \quad (5.51)$$

$$y(L, t) = 0 \quad \implies \quad \sin(kL) = 0 \quad (5.52)$$

In order for this to be possible, the values of k must be restricted to,

$$k_n = \frac{n\pi}{L}, n \in \mathbb{Z}. \quad (5.53)$$

So, we have a *discrete set of normal modes of oscillation*, such that,

$$y_n = A_n \sin\left(\frac{n\pi}{L}\right) \cos(\omega_n t) \quad (5.54)$$

where

$$\omega_n = vk_n = \frac{vn\pi}{L} \quad (5.55)$$

which is the angular frequency of the n^{th} mode and A_n is the amplitude of the n^{th} mode.

Wavelength of standing waves

We have,

$$k_n = \frac{2\pi}{\lambda_n} = \frac{n\pi}{L} \quad (5.56)$$

so,

$$\lambda_n = \frac{2L}{n} \quad \implies \quad L = n \frac{\lambda_n}{2} \quad (5.57)$$

thus, a standing wave is obtained if an integer number of half-wavelengths fit onto a string. The frequency of the standing wave is given by,

$$\nu = \frac{n}{2L} \sqrt{\frac{T}{\mu}}. \quad (5.58)$$

Superposition of normal modes

Most of the time, motion on a string will be given by the superposition of normal modes given by,

$$y(x, t) = \sum_n A_n \sin\left(\frac{n\pi x}{L}\right) \cos(\omega_n t). \quad (5.59)$$

5.4 Sound Waves

Sound waves are longitudinal, which means that the direction of motion is the same as the direction of propagation. This means that their wave equation is given by,

$$\frac{\partial^2 f}{\partial t^2} = \frac{\gamma P_0}{\rho} \frac{\partial^2 f}{\partial x^2} \quad (5.60)$$

where P_0 is the equilibrium pressure, ρ is the density, and γ is the adiabatic index. The velocity of sound waves is then,

$$v = \sqrt{\frac{\gamma P_0}{\rho}}. \quad (5.61)$$

5.5 The wave equation in 2 dimensions

The wave equation on a membrane can be given by,

$$\nabla^2 f = \frac{\sigma}{T} \frac{\partial^2 f}{\partial t^2}. \quad (5.62)$$

where $f(x, y, t)$ is the displacement and T is the tension. Expanding this, the equation we wish to solve in cartesian coordinates is,

$$\left(\frac{\partial^2}{\partial x^2} + \frac{\partial^2}{\partial y^2} \right) f(x, y, t) = v^{-2} \frac{\partial^2 f(x, y, t)}{\partial t^2}. \quad (5.63)$$

We will now assume a solution,

$$f(x, y, t) = f_x(x) f_y(y) f_t(t). \quad (5.64)$$

Substituting this in, we get,

$$\begin{aligned} \left(\frac{\partial^2}{\partial x^2} + \frac{\partial^2}{\partial y^2} \right) (f_x(x) f_y(y) f_t(t)) &= v^{-2} \frac{\partial^2 f_x(x) f_y(y) f_t(t)}{\partial t^2} \\ \implies f_y(y) f_t(t) \frac{d^2 f_x(x)}{dx^2} + f_x(x) f_t(t) \frac{d^2 f_y(y)}{dy^2} &= v^{-2} f_x(x) f_y(y) \frac{d^2 f_t(t)}{dt^2}. \end{aligned} \quad (5.65)$$

We can then divide through by $f_x(x) f_y(y) f_t(t)$, where we will find that each side only depends on its own variables, so we can equate the expression to a constant,

$$\frac{1}{f_x(x)} \frac{d^2 f_x(x)}{dx^2} + \frac{1}{f_y(y)} \frac{d^2 f_y(y)}{dy^2} = v^{-2} \frac{1}{f_t(t)} \frac{d^2 f_t(t)}{dt^2} = -k^2. \quad (5.66)$$

We can further separate the x and y components,

$$\frac{1}{f_x(x)} \frac{d^2 f_x(x)}{dx^2} = -k^2 - \frac{1}{f_y(y)} \frac{d^2 f_y(y)}{dy^2} = k_x^2, \quad (5.67)$$

where equation for f_y becomes,

$$\frac{1}{f_y(y)} \frac{d^2 f_y(y)}{dy^2} = -k_y^2 \quad (5.68)$$

so we have the relation,

$$k^2 = k_x^2 + k_y^2. \quad (5.69)$$

Thus, our problem turns into 3 ordinary, second order differential equations,

$$\frac{d^2 f_x(x)}{dx^2} = -k_x^2 f_x(x) \quad (5.70)$$

$$\frac{d^2 f_y(y)}{dy^2} = -k_y^2 f_y(y) \quad (5.71)$$

$$\frac{1}{v^2} \frac{\partial^2 f_t(t)}{\partial t^2} = -k^2 f_t(t). \quad (5.72)$$

For which the solutions are,

$$f_x(x) = A_x \cos k_x x + B_x \sin k_x x \quad (5.73)$$

$$f_y(y) = A_y \cos k_y y + B_y \sin k_y y \quad (5.74)$$

$$f_t(t) = A_t \cos \omega t + B_t \sin \omega t \quad (5.75)$$

where $\omega^2 = k^2 v^2$. The complete solution then takes the form,

$$f(x, y, t) = f_x(x) f_y(y) f_t(t). \quad (5.76)$$

5.5.1 Solutions in exponential form

$$f_x(x) = c_x e^{\pm i k_x x} \quad (5.77)$$

$$f_y(y) = c_y e^{\pm i k_y y} \quad (5.78)$$

$$f_t(t) = c_t e^{\pm i \omega t} \quad (5.79)$$

where the full solution becomes,

$$f(x, y, t) = C e^{i(k_x x + k_y y - \omega t)} = C e^{i(\mathbf{k} \cdot \mathbf{x} - \omega t)} \quad (5.80)$$

where,

$$\mathbf{k} = (k_x, k_y) \quad \mathbf{x} = (x, y). \quad (5.81)$$

5.5.2 2D wave equation in polar coordinates

In polar coordinates, the wave equation takes the form,

$$\frac{\partial^2 f}{\partial r^2} + \frac{1}{r} + \frac{1}{r^2} \frac{\partial^2 f}{\partial \theta^2} = \frac{\sigma}{T} \frac{\partial^2 f}{\partial t^2}. \quad (5.82)$$

In our case, we would like an axisymmetric solution without any dependence on θ . So, we set $\frac{\partial^2 f}{\partial \theta^2} = 0$, and (5.82) becomes,

$$\frac{\partial^2 f}{\partial r^2} + \frac{1}{r} = \frac{\sigma}{T} \frac{\partial^2 f}{\partial t^2}. \quad (5.83)$$

We will impose boundary conditions $f(R, t) = 0$ and assume $f = f_r(r) f_t(t)$. Substituting this into (5.83),

$$\frac{1}{f_r} \frac{d^2 f_r}{dr^2} + \frac{1}{r f_r} \frac{df_r}{dr} = \frac{\sigma}{r} = \frac{\sigma}{T} \frac{1}{f_t} \frac{d^2 f_t}{dt^2} = -k^2. \quad (5.84)$$

We can rewrite (5.84),

$$\frac{d^2 f_r}{dr^2} + \frac{1}{r} \frac{df_r}{dr} = -k^2 f_r. \quad (5.85)$$

(5.85) is known as *Bessel's equation* and the solutions are known as *Bessel functions*. The axisymmetric solution takes the form,

$$f(r, t) = f_r f_t = J_0(kr)(C \cos \omega t + D \sin \omega t) \quad (5.86)$$

where $J_0(kr)$ is a zero'th order Bessel function of the first kind, plotted in figure 5.3. We can quantise the wavenumber by imposing previously stated boundary conditions.

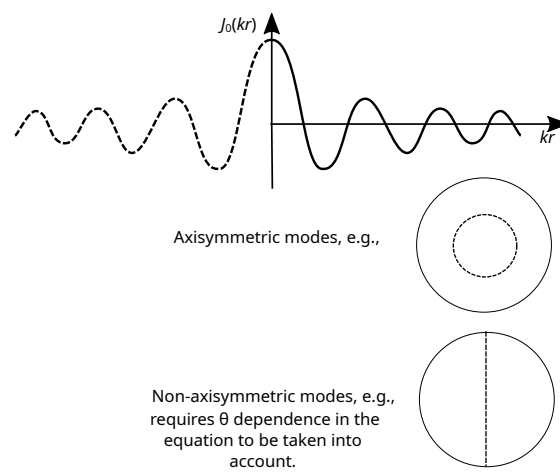


Figure 5.3: A zero'th order Bessel function of the first kind.

Chapter 6

Dispersion and Interference

In general, the velocity of a wave is not independent of frequency. This is called dispersion. Let us consider two waves propagating in the same medium with the same amplitude, K , at slightly different frequencies. We will assume they have different velocities, so,

$$v_1 = \frac{\omega_1}{k_1} \qquad v_2 = \frac{\omega_2}{k_2} \qquad v_1 \neq v_2. \quad (6.1)$$

The sum of the two waves is given by,

$$y = K [\sin(k_1x - \omega_1t) + \sin(k_2x - \omega_2t)] \quad (6.2)$$

Using identity (B.1), and defining,

$$\alpha + \beta = k_2x\omega_2t \qquad \alpha - \beta = k_1x - \omega_1t \quad (6.3)$$

so that,

$$\alpha = \frac{1}{2}(k_2 + k_1)x - \frac{1}{2}(\omega_2 + \omega_1)t \quad (6.4)$$

$$\beta = \frac{1}{2}(k_2 - k_1)x - \frac{1}{2}(\omega_2 - \omega_1)t. \quad (6.5)$$

We can then define the mean of the frequencies and wavevectors,

$$k_0 = \frac{k_1 + k_2}{2} \qquad \omega_0 = \frac{\omega_2 + \omega_1}{2} \quad (6.6)$$

and then write the difference of wavevectors and frequencies as,

$$\Delta k = \frac{k_2 - k_1}{2} \qquad \Delta\omega = \frac{\omega_2 - \omega_1}{2}. \quad (6.7)$$

Thus, we can rewrite (6.2) as,

$$y = 2K \sin(k_0x - \omega_0t) \cos(\Delta kx - \Delta\omega t). \quad (6.8)$$

we assume that Δk and $\Delta\omega$ is small. Thus, (6.8) describes a wave moving with *phase velocity* $v_p = \frac{\omega_0}{k_0}$ modulated by a slowly varying envelope wave, which has a *group velocity* $v_g = \frac{\Delta\omega}{\Delta k}$. So, we can rewrite (6.8) as,

$$y = 2KA(x, t) \sin(k_0x - \omega_0t), \quad (6.9)$$

where A is the envelope wave moving at speed v_g . Taking limits, we also have the relation,

$$v_g = \frac{d\omega}{dk}. \quad (6.10)$$

6.1 Phase and group velocities in a non-dispersive medium

In the case of a non-dispersive medium, both waves move at the same velocity v , and we find that,

$$v = v_p = v_g \quad (6.11)$$

6.2 Dispersion

In a dispersive medium, the phase velocity varies with k . We can call $v_p = v_p(k)$ the dispersion relation. We will use the fact that,

$$\lambda = \frac{2\pi}{k} \implies \frac{d\lambda}{dk} = -\frac{2\pi}{k^2} = -\frac{\lambda}{k}, \quad (6.12)$$

then the group velocity is found by,

$$v_g = \frac{d\omega}{dk} = d k(v_p k) = v_p - \lambda \frac{dv_p}{d\lambda}. \quad (6.13)$$

In the case where $\frac{dv_p}{d\lambda} = 0$, $v_g = v_p$. In the case where $\frac{dv_p}{d\lambda} > 0$ we get *normal dispersion*. In the case where $\frac{dv_p}{d\lambda} < 0$ we get *anomalous dispersion*.

6.2.1 Light dispersion

If we consider light moving through a prism, we can define the *refractive index* as the ratio of the speed of light in air and in glass. We have $n = \frac{c}{v_p}$, which controls how the velocity depends on λ ,

$$\frac{dv_p}{d\lambda} = \frac{dv_p}{dn} \frac{dn}{d\lambda} = -\frac{c}{n^2} \frac{dn}{d\lambda} = -\frac{v_p}{n} \frac{dn}{d\lambda}, \quad (6.14)$$

so the group velocity is given by,

$$v_g = v_p \left(1 + \frac{\lambda}{n} \frac{dn}{d\lambda} \right). \quad (6.15)$$

So, for $\frac{dn}{d\lambda} > 0$ we get anomalous dispersion. For $\frac{dn}{d\lambda} < 0$ we get normal dispersion. This then implies **shorter wavelengths have a greater refractive index, thus bend more in a prism.**

6.3 Wave interference

6.3.1 Conditions for destructive and constructive interference

Let us consider two monochromatic waves originating from the same source. They follow different paths and meet at a particular point in space. If their difference in path lengths from the common source is given by an integer number of wavelengths,

$$s = n\lambda, n \in \mathbb{Z} \quad (6.16)$$

the crests and troughs of each wave line up exactly with each other, as their phase difference is,

$$\phi = \frac{2\pi s}{\lambda} = 2n\pi \quad (6.17)$$

and so, constructive interference occurs.

Destructive interference occurs when,

$$s = \left(n + \frac{1}{2} \right) \lambda \quad \phi = (2n + 1)\pi \quad (6.18)$$

is satisfied, i.e., the waves are a phase π apart.

6.3.2 Huygen's Principle

Every point on a primary wavefront acts as a source of secondary wavefronts. To construct wavelets at a time Δt , we draw circles at the source of the wavelets at with a radius $v\Delta t$. If a wavefront encounters an aperture on an opaque barrier, the points on the wavefront across the aperture act like sources of secondary wavelets.

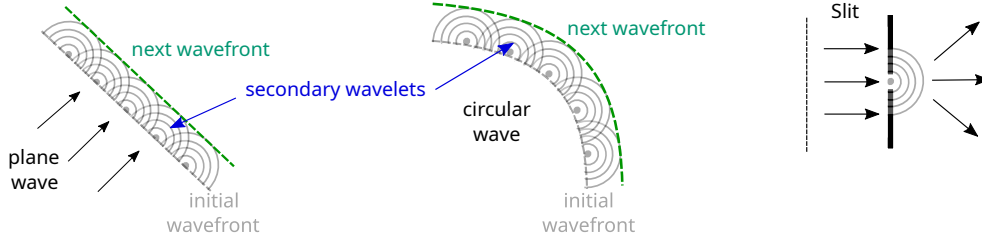


Figure 6.1: Wavefronts and wavelets illustrated by Huygen's principle.

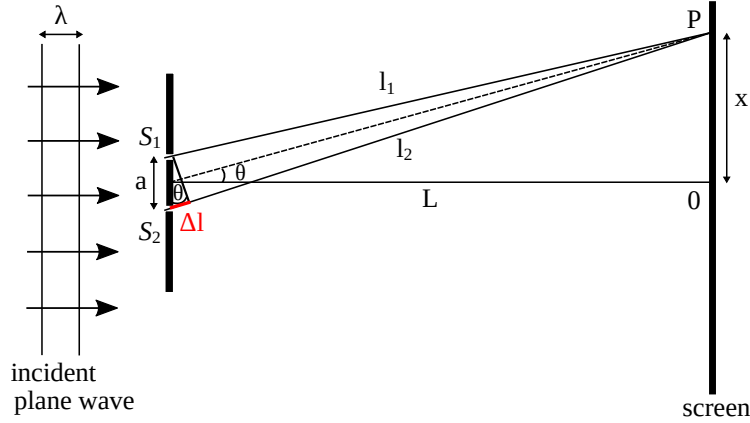


Figure 6.2: A diagram of Young's double slit.

6.3.3 Young's Double Slit Experiment

Let us consider a setup like in figure 6.2. A monochromatic wave is incident on an opaque barrier with 2 slits, S_1 and S_2 , separated by a distance a . The two slits act as a source of secondary wavelets. The disturbance beyond the barrier is the superposition of all wavelets spreading out from the two slits. Since both wavelets originate from the same source, there is a clear phase relation between them. This phase difference gives rise to an interference pattern.

The wavefronts will fall on a screen a distance L away. If we consider a point P , let l_1 and l_2 be the distance from S_1 and S_2 to this point. The extra path length to P from S_2 is given by $\Delta l = l_1 - l_2 = a \sin \theta$. Thus, constructive interference occurs when,

$$a \sin \theta = n\lambda \quad (6.19)$$

where there are bright fringes. Dark fringes occurs when,

$$a \sin \theta = \left(n + \frac{1}{2}\right) \lambda. \quad (6.20)$$

We are now interested in the intensity map produced on the screen. We know $I \propto A^2$. Let us define the amplitude at any point P on the screen,

$$\begin{aligned} R &= A [\cos(\omega t - k l_1) + \cos(\omega t - k l_2)] \\ &= 2A \cos\left(\omega t - k \frac{l_1 + l_2}{2}\right) \cos\left(k \frac{l_1 - l_2}{2}\right). \end{aligned} \quad (6.21)$$

Given $L \gg a$, $l_1 \sim l_2 \sim L$, so we can reasonable estimate,

$$\frac{l_1 + l_2}{2} \simeq L. \quad (6.22)$$

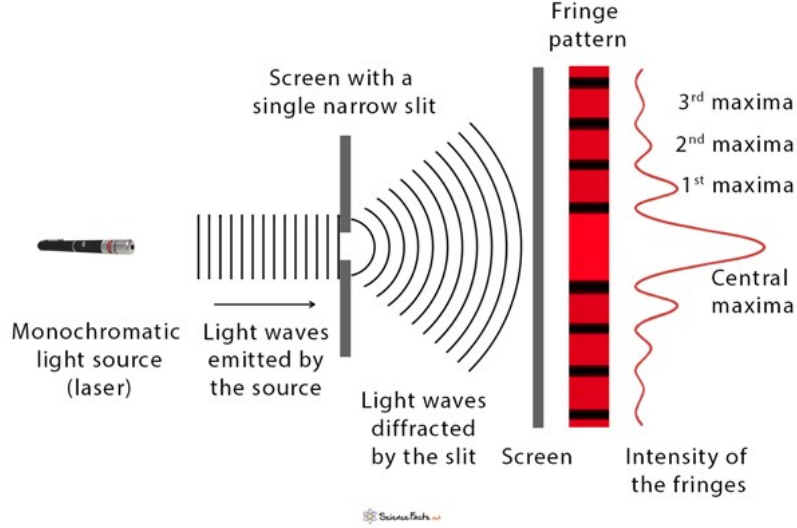


Figure 6.3: A slit has finite width, so each infinitesimal point on the slit acts as a source of wavelets, so there is interference. This then leads to a diffraction pattern.

So, we can write the intensity as,

$$I = R^2 = 4A^2 \cos^2(\omega t - kL) \cos^2\left(\frac{k\Delta L}{2}\right). \quad (6.23)$$

We are able to average this expression over time by integrating it from $0 \rightarrow T$ and dividing by T to get,

$$I = I_0 \cos^2 \frac{ka \sin \theta}{2} \quad (6.24)$$

where

$$I_0 = 2A^2. \quad (6.25)$$

We can determine that the first bright fringe occurs at $a \sin \theta = \pi$. If θ is small, we can state,

$$a\theta \approx n\lambda \quad (6.26)$$

where $\theta \approx \frac{x}{L}$. This implies that we can write the distance to the n th peak as,

$$x_n = \frac{n\lambda L}{a} \quad (6.27)$$

and thus, the distance between peaks are given by,

$$x_{n+1} - x_n = \frac{\lambda L}{a}. \quad (6.28)$$

6.4 Wave Diffraction

In practice, a slit will have a finite length, hence, path lengths from different points on the slit to the point P on the screen will be different. So, the phase of the wavelets arriving at P will be different, giving rise to a diffraction pattern.

Consider a setup like in figure 6.4. As waves spread, each point on the slit acts as a new source of waves according to Huygen's principle. If we consider a small part of the slit dx , then the resultant amplitude at P can be considered as the superposition of each wavelet over all dx . We can consider $L \gg d$, and assume that all wavelets arriving at P have the same amplitude and the lines joining P to all points on the slit make the same angle θ with the horizontal. If their path difference is given by,

$$s = x \sin \theta \quad (6.29)$$

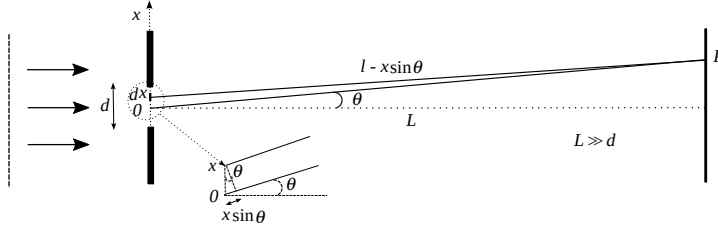


Figure 6.4: Diagram of single slit diffraction.

then the distance to the point P from a source of wavelets is given by,

$$d = \ell - x \sin \theta \quad (6.30)$$

where ℓ is the distance to P from the centre of the fringe. We can then consider the amplitude of a single wavelet dR ,

$$dR = \alpha dx \cos [\omega t - k(\ell - x \sin \theta)], \quad (6.31)$$

where α is a constant. We then integrate over the slit width,

$$R = \int_{-\frac{d}{2}}^{\frac{d}{2}} \frac{d}{2} \alpha dx \cos [\omega t - k(\ell - x \sin \theta)] \quad (6.32)$$

which can be evaluated to,

$$R = \frac{\alpha d}{\frac{kd}{2} \sin \theta} \sin \left[\frac{kd}{2} \sin \theta \right] \cos(\omega t - k\ell). \quad (6.33)$$

The diffraction pattern is then achieved by time averaging R^2 ,

$$I(\theta) = I_0 \frac{\sin^2 \left[\frac{kd}{2} \sin \theta \right]}{\left[\frac{kd}{2} \sin \theta \right]^2} \quad (6.34)$$

for,

$$I_0 = \frac{\alpha^2 d^2}{2} \quad \beta = \frac{kd}{2} \sin \theta. \quad (6.35)$$

Dark fringes are observed at $\sin \theta = \pm \frac{n\lambda}{d}$.

6.4.1 Circular aperture

The first minimum of a circular aperture is given by,

$$\sin \theta = 1.22 \frac{\lambda}{d}. \quad (6.36)$$

If we observe an image through a circular aperture, we observe considerable diffraction at close separation instead of a clear image. However, we can resolve the image of 2 closely spaced objects if the maximum diffraction pattern of the first object falls on the first minimum of the of the diffraction pattern of the second object. We can then resolve the two objects if they are observed with with a lens of diameter d and have an angular separation θ which satisfies (6.36).

6.4.2 Combining with the double slit experiment

We are able to combine the equation of diffraction with the young double slit interference equation. This equation is given by,

$$I(\theta) = I_0 \frac{\sin^2 \left[\frac{kd}{2} \sin \theta \right]}{\left[\frac{kd}{2} \sin \theta \right]^2} \cos^2 \frac{ka \sin \theta}{2} \quad (6.37)$$

which results in an intensity map like in figure 6.5, where the double slit pattern is modulated by the intensity pattern due to diffraction.

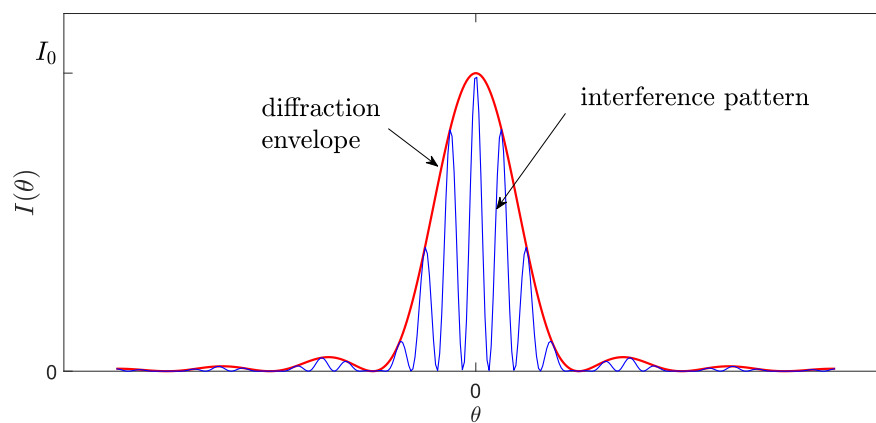


Figure 6.5: When combining the derivation of dispersion of a single slit with the double slit experiment, we find that the diffraction pattern acts as an envelope wave for the interference pattern.

Chapter 7

The Schrodinger equation

The 1D Schrodinger equation for a particle of mass m moving in a potential $V(x, t)$ is given by,

$$-\frac{\hbar^2}{2m} \frac{\partial^2 \Psi}{\partial x^2} + V(x, t) \Psi = i\hbar \frac{\partial \Psi}{\partial t}. \quad (7.1)$$

If the potential is only dependent on x we can solve the equation by separation of variables by assuming a solution in the form $\Psi(x, t) = \psi(x)T(t)$, and solving as we did for standing waves. However, in our case, the shared variable should be E , for energy. See appendix C for details.

7.1 Properties of the wavefunction

An interpretation of the wavefunction is that it can give the probability of finding a particle in a space between x and dx at time t (Born's interpretation),

$$|\Psi(x, t)|^2 dx = P(x, t) dx. \quad (7.2)$$

This would imply that,

$$\int_{-\infty}^{\infty} |\Psi(x, t)|^2 dx = 1 \quad (7.3)$$

and thus this places certain restrictions on the wave equation,

1. Ψ must be finite everywhere,
2. Ψ must be single valued,
3. Ψ and $\frac{\partial \Psi}{\partial x}$ must be continuous.

These ensure that only certain solutions are allowed.

7.2 Infinite potential well

Let us consider a potential well such that,

$$V(x) = \begin{cases} \infty & x < 0 \\ 0 & 0 \leq x \leq a \\ \infty & x > a \end{cases}. \quad (7.4)$$

This resembles an idealised potential seen by electrons in the low energy levels near the nucleus of an atom. We want to find the wavefunction and energy of the particle in this potential. We write the time independent Schrodinger equation,

$$\begin{aligned} & -\frac{\hbar^2}{2m} \frac{d^2 \Psi}{dx^2} + V(x) \Psi = E \Psi \\ \Rightarrow & \frac{d^2 \Psi}{dx^2} + \frac{2m}{\hbar^2} (E - V(x)) \Psi = 0. \end{aligned} \quad (7.5)$$

We require $\Psi(x) = 0$ in the regions $x < 0$ and $x > a$, as otherwise the LHS would be infinite. For $0 \leq x \leq a$, we have $V(x) = 0$, thus,

$$\frac{d^2\Psi}{dx^2} = -\frac{2mE}{\hbar^2}\Psi. \quad (7.6)$$

Let us define the wave number k ,

$$k^2 = \frac{2mE}{\hbar^2}. \quad (7.7)$$

We then have the general solution to (7.6),

$$\Psi(x) = A \cos kx + B \sin kx. \quad (7.8)$$

Applying boundary conditions,

$$\Psi(0) = 0 \implies A = 0 \quad (7.9)$$

$$\Psi(a) = 0 \implies B \sin ka = 0 \quad (7.10)$$

We require non-trivial solutions, thus $B \neq 0$, and,

$$\sin k_n a = 0 \implies k_n = \frac{n\pi}{a}, n \in \mathbb{N}. \quad (7.11)$$

This implies that only discrete values of k are allowed, which corresponds to discrete energy levels,

$$E_n = \frac{\hbar^2}{2m} k_n^2 = \frac{n^2 \hbar^2 \pi^2}{2ma^2} \quad (7.12)$$

and the wave functions,

$$\Psi(x) = A_n \sin n\pi x \quad (7.13)$$

are the energy eigen functions (stationary states) for a particle in a 1 dimensionall, infinite potential well.

7.2.1 Time dependent solution

By separation of variables, one finds the time dependent solution,

$$\Psi(x, t) = \begin{cases} 0 & x < 0 \\ A_n \sin\left(\frac{n\pi x}{a}\right) \exp\left(\frac{-ik_n t}{\hbar}\right) & 0 \leq x \leq a \\ 0 & x > a \end{cases} \quad (7.14)$$

We can then normalise the equation by applying (7.3) and finding that,

$$A_n = \sqrt{\frac{2}{a}} \quad (7.15)$$

7.2.2 Stationary states

A stationary state refers to an observable which is independent of time. The energy eigenfunction refers to the eigen vector of the energy operator,

$$\hat{E} = i\hbar \frac{\partial}{\partial t} \quad (7.16)$$

which acts on a given wavefunction.

7.3 Tunelling

For a finite barrier, there is a probability of finding a particle outside of its classically forbidden region. This phenomenon is referred to as tunnelling.

The finite barrier has the form,

$$V(x) = \begin{cases} 0 & x < 0 \\ V_B & 0 \leq x \leq a \\ 0 & x > a \end{cases} \quad (7.17)$$

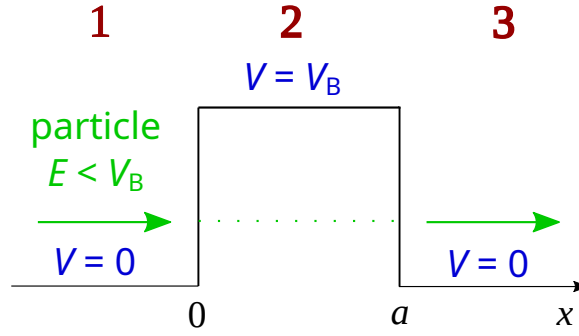


Figure 7.1

illustrated in figure 7.1.

Classically, a particle can only be transmitted if $E > V_B$, and is reflected otherwise. However, quantum mechanically it is possible for a particle to be transmitted if $E < V_B$.

Let us consider region 1. We have $V(x) = 0$,

$$\begin{aligned} \frac{d^2\Psi}{dx^2} &= -k^2\Psi \\ \implies \Psi(x) &= A_I e^{ikx} + A_R e^{-ikx} \end{aligned} \quad (7.18)$$

which can be thought of as the sum of incident and reflected waves.

Considering region 2, $V(x) = V_B$. Inside the barrier, we have,

$$\frac{d^2\Psi}{dx^2} = \beta^2\Psi \text{ for } E = -\frac{\hbar^2\beta^2}{2m} + V_B \quad (7.19)$$

which has solutions,

$$\Psi(x) = C e^{\beta x} + D e^{-\beta x}. \quad (7.20)$$

The solutions of region 3 have the same form as those of region 1, but only the transmitted particle needs to be represented,

$$\Psi(x) = A_T e^{ikx}. \quad (7.21)$$

The stationary function in a finite potential well then takes the form,

$$\Psi(x) = \begin{cases} A_I e^{ikx} + A_R e^{-ikx} & x < 0 \\ C e^{\beta x} + D e^{-\beta x} & 0 \leq x \leq a \\ A_T e^{ikx} & x > a \end{cases} \quad (7.22)$$

7.3.1 Tunelling probability

The continuity of $\frac{d\Psi}{dx}$ and Ψ at $x = 0$ and $x = a$ are our boundary conditions.

1. Continuity of Ψ at $x = 0$,

$$\implies A_I + A_R = C + D \quad (7.23)$$

2. Continuity of $\frac{d\Psi}{dx}$ at $x = 0$,

$$\implies ikA_I - ikA_R = \beta C - \beta D \quad (7.24)$$

- We can write (7.23) as,

$$ikA_I + ikA_R = ikC + ikD \quad (7.25)$$

- (7.25) + (7.23),

$$\implies 2kA_I = (\beta + ik)C - (\beta - ik)D. \quad (7.26)$$

3. Continuity of Ψ and $\frac{d\Psi}{dx}$ at $x = a$ gives,

$$Ce^{\beta a} + De^{-\beta a} = A_T e^{ika} \quad (7.27)$$

$$C\beta e^{\beta a} - \beta D e^{-\beta a} = ik A_T e^{ika}. \quad (7.28)$$

We seek A_T and C in terms of D ,

$$A_T e^{ika} = \frac{a\beta}{\beta - ik} D e^{-\beta a} \quad (7.29)$$

$$C = D e^{-2\beta a} \left(\frac{\beta + ik}{\beta - ik} \right) \quad (7.30)$$

If the barrier is wide, then,

$$e^{-2\beta a} \ll 1 \quad C \ll D. \quad (7.31)$$

From (7.26),

$$2ik A_I \approx -D(\beta - ik) \quad (7.32)$$

and (7.29),

$$A_T e^{ika} \approx \frac{2\beta}{\beta - ik} \frac{-2ik A_I}{\beta - ik} e^{-\beta a}, \quad (7.33)$$

hence,

$$\frac{A_T}{A_I} \approx \frac{-4ik\beta e^{-(\beta+ik)a}}{(\beta - ik)^2}. \quad (7.34)$$

we can then define a transmission coefficient by taking the ratio of transmitted particle flux, $v|A_T|^2$ and the incident particle flux $v|A_I|^2$ which gives the *tunnelling probability*,

$$T = \frac{|A_T|^2}{|A_I|^2} = \left[\frac{16k^2\beta^2}{(\beta^2 + k^2)^2} \right] e^{-2\beta a}. \quad (7.35)$$

Using definitions of k and β ,

$$k = \frac{\sqrt{2mE}}{\hbar} \quad \beta = \frac{\sqrt{2m(V_B - E)}}{\hbar} \quad (7.36)$$

to define T in terms of E and V_B ,

$$T = \left[\frac{16E(V_B - E)}{V_B^2} \right] e^{-2\beta a} \quad (7.37)$$

which is valid for a wide barrier where the wavefunction is dominated by the $D^{-\beta x}$ term.

The tunnelling probability then depends on,

1. The width of the barrier,
2. The mass of the particle,
3. The height of the barrier compared to the energy of the particle.

7.4 Finite potential well

There is a chance a particle may be found outside of the potential well if its walls are finite. A finite potential well is illustrated in figure 7.2.

There are two types of solutions for this potential,

1. Bound solutions $\rightarrow |E| < |V_0|$,
2. Unbound solutions $\rightarrow |E| > |V_0|$.

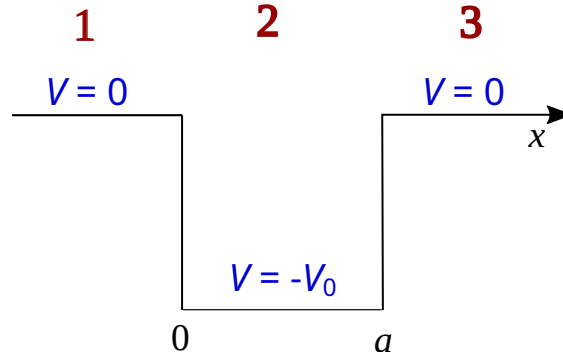


Figure 7.2

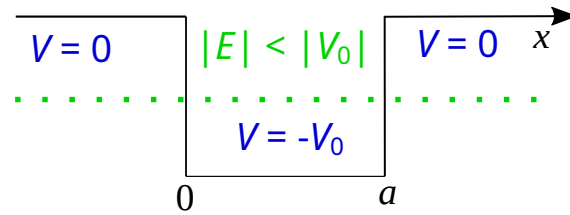


Figure 7.3

7.4.1 Bound states

This type of solution is illustrated in figure 7.3. Within the well, the wavefunctions which satisfy the independent Schrodinger equation are oscillatory standing wave functions,

$$\frac{d^2\Psi}{dx^2} = -k_0^2\Psi \quad k_0 = \frac{\sqrt{2m(V_0 + E)}}{\hbar}. \quad (7.38)$$

We find that the wavefunction decays exponentially inside the walls of the well and the distance which that wavefunction goes into the classically forbidden region depends on the height of the barrier.

7.4.2 Unbound states

This setup is illustrated in figure (7.4). In this case, the solutions are travelling waves where the particle isn't confined to the well, but is affected by the presence of the potential.

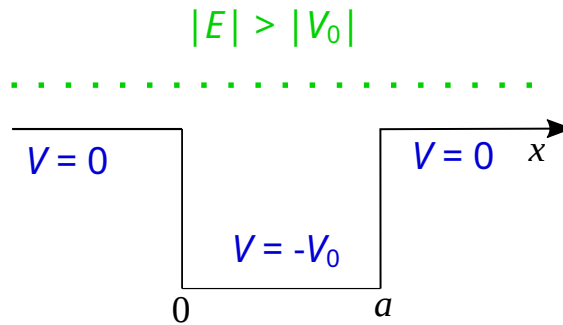


Figure 7.4

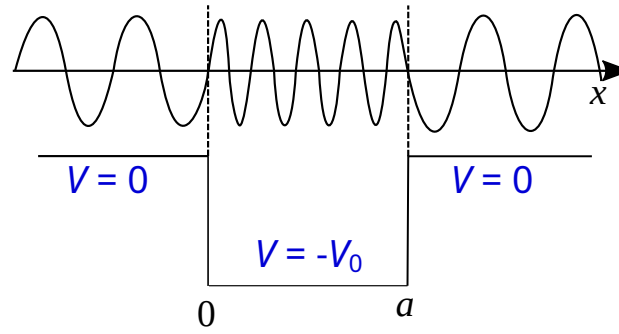


Figure 7.5

If k describes the wave outside the wall, and k_0 inside, then,

$$E = \frac{\hbar^2 k^2}{2m} = \frac{\hbar^2 k_0^2}{2m} - V_0 \quad (7.39)$$

hence, $k_0 > k$, and λ is greater outside the wavefunction than inside, illustrated by figure 7.5.

Appendix A

Common Waveforms and Examples

A.1 Gaussian Pulse

$$y(x, t) = A \exp\left(-\frac{1}{a}(x - vt)^2\right) \quad (\text{A.1})$$

A.2 Waves on a square membrane

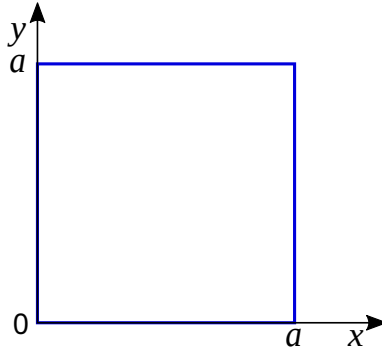


Figure A.1: Model of a square membrane.

In order to solve this problem, we must specify boundary conditions for the membrane. For this problem, we often clamp the edges of the membrane so that its edges are fixed. So,

$$f(0, y, t) = 0 \quad f(a, y, t) = 0 \quad (\text{A.2})$$

$$f(x, 0, t) = 0 \quad f(x, a, t) = 0 \quad (\text{A.3})$$

By imposing the initial conditions on (5.73), we find,

$$A = 0 \quad B \sin k_x a = 0. \quad (\text{A.4})$$

For a non trivial solution to B , we can state,

$$k_x^{(n)} = \frac{n\pi}{a}, n \in \mathbb{Z}^+, \quad (\text{A.5})$$

so that the solution for x is,

$$f_x(x) = B \sin\left(\frac{n\pi x}{a}\right). \quad (\text{A.6})$$

Similarly for y ,

$$f_y(y) = B \sin\left(\frac{m\pi y}{a}\right), m \in \mathbb{Z}^+. \quad (\text{A.7})$$

Solving for t ,

$$f_t(t) = A_t \cos(\omega_n^{(m)} t) + B_t \sin(\omega_n^{(m)} t) \quad (\text{A.8})$$

Where $\omega_n^{(n)}$ depends on n and m . Let us recall the dispersion relation $\omega^2 = k^2 v^2$, and $k^2 = k_x^2 + k_y^2$. So,

$$k^2 = \left(\frac{n\pi}{a}\right)^2 + \left(\frac{m\pi}{a}\right)^2. \quad (\text{A.9})$$

Thus we can obtain characteristic frequencies of the normal modes n and m ,

$$\omega_n^{(m)} = v\pi \sqrt{\frac{n^2}{a^2} + \frac{m^2}{a^2}}. \quad (\text{A.10})$$

We then obtain the full solution for the normal modes of the wave,

$$f_n^{(m)}(x, y, t) = \sin\left(\frac{n\pi x}{a}\right) \sin\left(\frac{m\pi y}{a}\right) \left(a \cos \omega_n^{(m)} t + a \sin \omega_n^{(m)} t\right) \quad (\text{A.11})$$

The general solution is the sum over all normal modes.

$$F(x, y, t) = \sum_{n=1}^{\infty} \sum_{m=1}^{\infty} f_n^{(m)}(x, y, t). \quad (\text{A.12})$$

Appendix B

Trig Identities

$$\sin(\alpha + \beta) + \sin(\alpha - \beta) = 2 \sin \alpha \cos \beta \quad (\text{B.1})$$

$$\sin \alpha \sin \beta = \frac{1}{2}(\cos(\alpha - \beta) - \cos(\alpha + \beta)) \quad (\text{B.2})$$

$$\cos \alpha \cos \beta = \frac{1}{2}(\cos(\alpha - \beta) + \cos(\alpha + \beta)) \quad (\text{B.3})$$

$$\cos \alpha \sin \beta = \frac{1}{2}(\sin(\alpha + \beta) - \sin(\alpha - \beta)) \quad (\text{B.4})$$

Appendix C

Solving the time dependent Schrodiner equation

If we have that $\Psi(x, t) = \Psi(x)T(t)$, and substituing this into the time dependent schrodinger equation,

$$\left(\frac{-\hbar^2}{2m} \frac{d^2\Psi}{dx^2} + V(x)\Psi(x) \right) T(t) = i\hbar\Psi(x) \frac{dT(t)}{dt}, \quad (\text{C.1})$$

which rearranges,

$$\frac{1}{\Psi(x)} \frac{-\hbar^2}{2m} \frac{d^2\Psi}{dx^2} + V(x) = \frac{i\hbar}{T(t)} \frac{T(t)}{t}. \quad (\text{C.2})$$

We can equate both sides to a constant, which we can call E . This results in two differential equations, with corresponding solutions,

$$\frac{i\hbar}{T(t)} \frac{dT(t)}{dt} = E \quad \implies \quad T = T_o \exp\left(-\frac{iEt}{\hbar}\right) = T_o \exp(-i\omega t) \quad (\text{C.3})$$

$$-\frac{\hbar^2}{2m} \frac{d^2\Psi}{dx^2} + V(x)\Psi = E\Psi \quad \implies \quad \text{Stationary state solutions for given potential} \quad (\text{C.4})$$