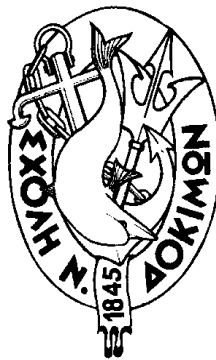


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One-Dimensional Newtonian Systems

- Conservative and Periodic Systems
 - Oscillatory Systems



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One-dimensional Newtonian systems

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The cases of conservative and oscillatory Newtonian systems in one dimension are studied. Certain unique properties of simple harmonic motion are noted.

A. One-dimensional conservative systems

1. The general solution to the problem

Consider a particle of mass m , moving along the x -axis under the action of a total force $F(x)$. The position $x(t)$ of the particle as a function of time is found by integrating the second-order differential equation (Newton's second law)

$$m d^2x / dt^2 = F(x) \quad (1)$$

for given initial conditions $x(t_0)=x_0$, $v(t_0)=v_0$, where $v=dx/dt$ is the velocity of the particle.

Define the auxiliary function $U(x)$ (potential energy of the particle) by

$$U(x) = -\int_0^x F(x') dx' \Leftrightarrow F(x) = -dU/dx \quad (2)$$

Then (1) is written

$$m d^2x / dt^2 + dU/dx = 0 .$$

We multiply by $v=dx/dt$, which plays the role of an integrating factor:

$$(dx/dt)(m d^2x / dt^2 + dU/dx) = 0 .$$

By noticing that

$$(dx/dt)(m d^2x / dt^2) = v(m dv/dt) = (d/dt)(m v^2/2)$$

and that $(dx/dt)(dU/dx) = dU/dt$, we have: $(d/dt)(m v^2/2 + U) = 0 \Rightarrow$

$$m v^2/2 + U(x) \equiv T + U = E = const. \quad (3)$$

(where T =kinetic energy) which expresses conservation of total mechanical energy.

From relation (3) we get

$$(dx/dt)^2 = (2/m)[E-U(x)] \Rightarrow dx/dt = \pm \{(2/m)[E-U(x)]\}^{1/2} .$$

Integrating this first-order differential equation and taking into account the initial condition $x=x_0$ for $t=t_0$, we have:

$$\int_{x_0}^x \frac{\pm dx}{\left\{ \frac{2}{m} [E - U(x)] \right\}^{1/2}} = t - t_0 \quad (4)$$

where the plus sign is chosen for motion in the *positive* direction ($v > 0$, $x > x_0$) while the minus sign applies to motion in the *negative* direction ($v < 0$, $x < x_0$).

The value of the constant E may be determined by applying the given initial conditions to (3):

$$E = m v_0^2 / 2 + U(x_0) \quad (5)$$

(although, as we will see, other physical considerations may also be used).

2. The case of periodic motion

Let us now assume that the potential energy $U(x)$ has the form of a U-shaped potential well (Fig. 1) such that $U(0)=0$ and $U(x)>0$ for $x \neq 0$ (this arrangement is always possible because of the arbitrariness in the definition of the zero-level of the potential energy). In general, the graph of $U(x)$ need not be symmetric with respect to the axis $x=0$.

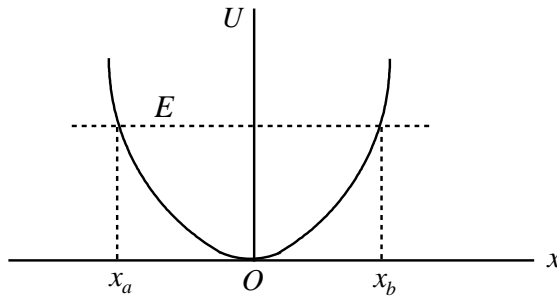


Fig. 1

Let E be the total mechanical energy of the particle. Since $E=T+U$ with $T \geq 0$, it follows that $E \geq U(x)$ for any physical motion. The motion is thus *bounded* between the points x_a and x_b of the x -axis, these points being *turning points* at which the particle stops momentarily ($E=U \Rightarrow T=0 \Rightarrow v_a=v_b=0$). The time it takes for a complete journey from x_a to x_b and back to x_a is found by using (4) with the appropriate sign for each direction of motion:

$$P = \int_{x_a}^{x_b} \frac{dx}{\left\{ \dots \right\}^{1/2}} + \int_{x_b}^{x_a} \frac{-dx}{\left\{ \dots \right\}^{1/2}} \Rightarrow$$

$$P = 2 \int_{x_a}^{x_b} \frac{dx}{\left\{ \frac{2}{m} [E - U(x)] \right\}^{1/2}} \quad (6)$$

Since P is fixed for given x_a and x_b , the motion is *periodic* with *period* P . Generally, the period depends on the limits of integration x_a and x_b and therefore it depends on the total energy E of the particle. An exception where P does *not* depend on E is *simple harmonic motion*, as we now show.

3. Simple harmonic motion (SHM)

In SHM the potential energy is of *parabolic* form: $U(x)=kx^2/2$, which is symmetric with respect to the axis $x=0$ (see Fig. 1). The total force is a *restoring force* given by

$$F(x) = -dU/dx = -kx \quad (7)$$

If frictional (damping) forces are present, the total force also contains a velocity-dependent term $-\lambda v = -\lambda dx/dt$ and the system is no longer conservative.

According to Fig. 1 the motion takes place between $x_a = -A$ and $x_b = A$, where $A \geq 0$ is the *amplitude* of oscillation. At the two extreme points the kinetic energy T vanishes momentarily and the total energy, which is equal to $E = T + U$ and which retains a fixed value during the motion, is equal to the potential energy: $E = U(\pm A) = kA^2/2$. Since E is the same at all points x , we conclude that

$$E = mv^2/2 + kx^2/2 = kA^2/2 \quad (8)$$

The period of oscillation is found by using (6):

$$P = 2 \int_{-A}^A \left\{ \frac{2}{m} (E - kx^2/2) \right\}^{-1/2} dx.$$

Substituting for E from (8), we find:

$$P = \frac{2}{\omega} \int_{-A}^A (A^2 - x^2)^{-1/2} dx$$

where we have set $\omega = (k/m)^{1/2}$ (angular frequency). Putting $x/A = u$ and using the integral formula

$$\int \frac{du}{\sqrt{1-u^2}} = \arcsin u + C$$

we finally find (see Appendix):

$$P = 2\pi / \omega = 2\pi (m/k)^{1/2}.$$

We conclude that, if the potential energy is of parabolic form: $U(x)=kx^2/2$, the period P of motion is independent of the amplitude A , thus independent of the total energy $E=kA^2/2$.

But, what if $U(x)$ is like that in Fig. 1 but *not* parabolic? For example, let U be of the form $U(x)=\lambda x^4/4$, so that $F(x) = -dU/dx = -\lambda x^3$. Since $U(x)$ is symmetric with respect to the axis $x=0$, the periodic motion will take place between the points $x_a = -A$ and $x_b = A$ and the total energy will be equal to $E = U(\pm A) = \lambda A^4/4$. The period is

$$P = 2 \int_{-A}^A \left\{ \frac{2}{m} (E - \lambda x^4/4) \right\}^{-1/2} dx = \frac{2}{\mu} \int_{-A}^A (A^4 - x^4)^{-1/2} dx = \frac{2}{\mu A} \int_{-1}^1 \frac{du}{\sqrt{1-u^4}}$$

where we have set $u=x/A$ and $\mu=(\lambda/2m)^{1/2}$. Obviously, P depends on the amplitude A , thus on the total energy E . (A more general proof regarding non-parabolic potential energies, in general, is given in the Appendix.)

Returning to SHM, we may obtain the equation of motion $x=x(t)$ by using (4) with $U(x)=kx^2/2$ and $E=kA^2/2$. Let us assume first that the motion is in the positive direction, so that $x>x_0$. Setting $\omega=(k/m)^{1/2}$, we have:

$$\int_{x_0}^x (A^2 - x^2)^{-1/2} dx = \omega (t - t_0).$$

Using the integral formula

$$\int (A^2 - x^2)^{-1/2} dx = \arcsin(x/A) + C$$

and making appropriate substitutions for constants, we find an equation of the form¹

$$\arcsin(x/A) = \omega t + \alpha \Rightarrow x = A \sin(\omega t + \alpha).$$

For motion in the negative direction ($x<x_0$) we choose the minus sign in (4), so that

$$\int_{x_0}^x (A^2 - x^2)^{-1/2} dx = -\omega (t - t_0).$$

This yields a result of the form²

$$\arcsin(x/A) = -\omega t + \beta \Rightarrow x = -A \sin(\omega t - \beta).$$

Since the constant β is arbitrary (being dependent on the arbitrary constants x_0 and t_0) we may set $-\beta \equiv \pi + \alpha$, so that $x = A \sin(\omega t + \alpha)$, as before.

Thus, the general solution for SHM is $x(t) = A \sin(\omega t + \alpha)$. Physically, A is the *amplitude* of oscillation, ω is the *angular frequency* and α is the *initial phase* (i.e., the *phase* $\omega t + \alpha$ at $t=0$).

4. Motion under a constant force of gravity

A projectile of mass m is fired straight upward at time $t_0=0$ from the point $x=0$ of the vertical x -axis, with initial velocity $v_0>0$ (we choose the positive direction of the x -axis to be upward). The constant acceleration of gravity is directed downward, so that $a=dv/dt=-g$. The total force on the particle (assuming no air resistance) and the corresponding potential energy of the particle are given by

$$F(x) = ma = -mg \Leftrightarrow U(x) = mgx \quad [\text{we assume that } U(0)=0].$$

Relation (4) (with the plus sign for upward motion) is written

$$\int_0^x \frac{dx}{(E - mgx)^{1/2}} = (2/m)^{1/2} t.$$

¹ Explicitly: $\alpha = \arcsin(x_0/A) - \omega t_0$.

² Explicitly: $\beta = \arcsin(x_0/A) + \omega t_0$.

By (5) and by using the initial conditions we have that $E = mv_0^2/2 + U(0) = mv_0^2/2$ (since $U=0$ for $x_0=0$). Thus, the requirement $E - mgx \geq 0$ yields $x \leq v_0^2/2g$. Physically this means that the particle will reach a maximum height $h = v_0^2/2g$ where it will stop momentarily before it starts to move downward (i.e., in the negative direction).

With this restriction on the acceptable values of x , the integration may be performed to give

$$(E - mgx)^{1/2} = E^{1/2} - (m/2)^{1/2} gt.$$

Squaring this, we find:

$$x = (2E/m)^{1/2} t - gt^2/2.$$

But, $E = mv_0^2/2 \Rightarrow (2E/m)^{1/2} = v_0$ (since $v_0 > 0$). Thus, finally,

$$x = v_0 t - gt^2/2$$

which is, of course, a familiar result.

5. Phase curves of a one-dimensional conservative system

Newton's law for one-dimensional motion: $m d^2x/dt^2 = F(x)$, a second-order differential equation, may be rewritten as a system of first-order equations:

$$dx/dt = v, \quad m dv/dt = F(x) \quad (9)$$

Dividing these equations in order to eliminate dt , we have:

$$m v dv = F(x) dx = -dU$$

where

$$U(x) = -\int_0^x F(x') dx' \Leftrightarrow F(x) = -dU/dx.$$

Thus, $m v dv + dU = d(mv^2/2 + U) = 0 \Rightarrow$

$$mv^2/2 + U(x) = E \equiv \text{const.} \quad (10)$$

For each value of the constant E (total energy), Eq. (10) defines a curve in the 2-dimensional *phase space* with coordinates (x, v) . This curve is called a *phase curve*. The value of E is uniquely determined by the initial conditions of the system, according to (5). Since the solution of the system (9) is unique for given initial conditions, *no two phase curves may intersect* in phase space. Let us see two examples:

1. Simple harmonic motion (cf. Sec. 3)

Conservation of mechanical energy in SHM is expressed by $mv^2/2 + kx^2/2 = E \Rightarrow$

$$\frac{x^2}{2E/k} + \frac{v^2}{2E/m} = 1 \quad (\text{equation of an ellipse})$$

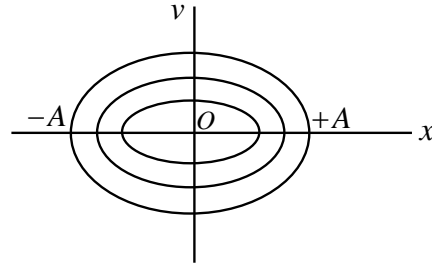


Fig. 2

Figure 2 shows a family of ellipses in phase space, corresponding to different values of E . Notice that, for $v=0 \Rightarrow x = \pm(2E/k)^{1/2} \equiv \pm A$, so that $E=kA^2/2$. Note also that the equations of motion, $\{dx/dt = v, dv/dt = -kx/m\}$, endow the phase curves with a sense of direction for increasing t (i.e., for $dt > 0$). Indeed, the velocity v is positive (negative) for increasing (decreasing) x , while v decreases (increases) *algebraically* for positive (negative) x . This indicates that the phase curves are described *clockwise*.

2. Vertical motion under the force of gravity (cf. Sec. 4)

Conservation of mechanical energy is expressed by $mv^2/2 + mgx = E \Rightarrow$

$$v^2 = (2/m)(E - mgx) \quad (\text{equation of a parabola})$$

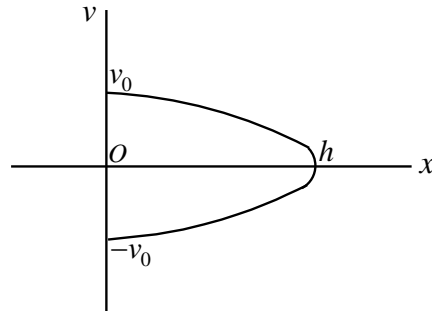


Fig. 3

Since $v^2 \geq 0$, we must have $E - mgx \geq 0 \Rightarrow x \leq E/mg$. Physically, this means that the particle will reach a maximum height $h=E/mg$ where it will stop momentarily and then its direction of motion will be reversed. On the other hand, at $x=0$ the velocity is $\pm v_0$ (see Fig. 3) where $v_0^2=2E/m \Rightarrow E=mv_0^2/2$. The maximum height is thus $h=v_0^2/2g$.

B. Oscillatory motion of (generally) non-conservative systems

1. Second-order linear differential equations

A second-order linear differential equation (DE) has the general form

$$y'' + a(x)y' + b(x)y = f(x) \quad (1)$$

where $y=y(x)$ and where $a(x)$, $b(x)$, $f(x)$ are given functions. If $f(x)\equiv 0$, the DE (1) is called *homogeneous linear*:

$$y'' + a(x)y' + b(x)y = 0 \quad (2)$$

As is easy to prove, if a function $y_1(x)$ is a solution of (2), then so is the function $y_2(x)=Cy_1(x)$ ($C=\text{const.}$). More generally, the following is true:

Theorem 1: If $y_1(x)$, $y_2(x)$,... are solutions of the homogeneous DE (2), then every linear combination of the form $y=C_1 y_1(x)+C_2 y_2(x)+\dots$ (where C_1 , C_2 , ... are constants) also is a solution of (2).

Proof: By substituting for y on the left-hand side of (2) and by taking into account that each of the $y_1(x)$, $y_2(x)$,... satisfies this DE, we have:

$$y'' + a(x)y' + b(x)y = C_1 (y_1'' + a y_1' + b y_1) + C_2 (y_2'' + a y_2' + b y_2) + \dots = 0.$$

Let $y_1(x)$ and $y_2(x)$ be two non-vanishing solutions of the homogeneous DE (2) [notice that the zero function $y(x)\equiv 0$ is a particular solution of (2)]. We say that the functions y_1 and y_2 are *linearly independent* if one is not a scalar multiple of the other. To put it in more formal terms, linear independence of y_1 and y_2 means that a relation of the form $C_1 y_1(x)+C_2 y_2(x)\equiv 0$ can only be true if $C_1=C_2=0$.

If we manage to find two linearly independent solutions $y_1(x)$ and $y_2(x)$ of the homogeneous DE (2) (I can assure you that no other solution linearly independent of the former two exists!) then the *general solution* of (2) is the linear combination

$$y = C_1 y_1(x) + C_2 y_2(x) \quad (3)$$

where C_1 , C_2 are arbitrary constants.

Theorem 2: The general solution of the non-homogeneous DE (1) is the sum of the general solution (3) of the corresponding homogeneous equation (2) and *any particular solution* of (1).

Analytically: Let $y_1(x)$, $y_2(x)$ be two linearly independent solutions of the homogeneous DE (2), and let $y_0(x)$ be any particular solution of (1). Then, the general solution of (1) is

$$y = C_1 y_1(x) + C_2 y_2(x) + y_0(x) \quad (4)$$

This practically means that, for any chosen y_0 , any other particular solution of (1) can be derived from (4) by properly choosing the constants C_1 and C_2 . Since (4) contains the totality of particular solutions of (1), it must be the general solution of (1).

2. Homogeneous linear equation with constant coefficients

This DE has the form

$$y'' + ay' + by = 0 \quad (5)$$

with constant a and b . It will be assumed that a and b are real numbers.

Theorem 3: If the complex function $y=u(x)+iv(x)$ satisfies the DE (5), then the same is true for each of the real functions $y_1=u(x)$ and $y_2=v(x)$ (real and imaginary part of y , respectively).

Proof: Putting $y=u+iv$ into (5), we find:

$$(u'' + au' + bu) + i(v'' + av' + bv) = 0,$$

which is true iff $u'' + au' + bu = 0$ and $v'' + av' + bv = 0$.

The standard method for solving (5) is the following: We try an exponential solution of the form $y=e^{kx}$. Then, $y'=ke^{kx}$, $y''=k^2e^{kx}$, and (5) yields (after eliminating e^{kx}):

$$k^2 + ak + b = 0 \quad (\text{characteristic equation}) \quad (6)$$

We distinguish the following cases:

1. Eq. (6) has real and distinct roots k_1, k_2 . Then, the functions e^{k_1x} and e^{k_2x} are linearly independent and, according to (3), the general solution of (5) is of the form

$$y = C_1 e^{k_1x} + C_2 e^{k_2x} \quad (7)$$

2. Eq. (6) has real and equal roots, $k_1 = k_2 \equiv k$. The general solution of (5) is, in this case (check!),

$$y = (C_1 + C_2 x) e^{kx} \quad (8)$$

3. Eq. (6) has complex conjugate roots $k_1 = \alpha + i\beta$, $k_2 = \alpha - i\beta$ (where α, β are real). The general solution of (5) is

$$y = C_1 e^{k_1x} + C_2 e^{k_2x} = e^{\alpha x} (C_1 e^{i\beta x} + C_2 e^{-i\beta x}).$$

By Euler's formula, $e^{\pm i\beta x} = \cos \beta x \pm i \sin \beta x$. We thus have:

$$y = e^{\alpha x} [(C_1 + C_2) \cos \beta x + i(C_1 - C_2) \sin \beta x].$$

Since the (generally complex) constants C_1 and C_2 are arbitrary, we may put C_1 in place of $C_1 + C_2$ and C_2 in place of $i(C_1 - C_2)$, so that, finally,

$$y = e^{\alpha x} (C_1 \cos \beta x + C_2 \sin \beta x) \quad (9)$$

In any case, the general solution of (5) contains two arbitrary constants C_1 and C_2 . Upon assigning specific values to C_1 and C_2 we get a *particular solution* of (5). The values of C_1 and C_2 (and thus the particular solution itself) are determined from the

general solution if we are given two *initial conditions* that the sought-for particular solution must obey. There are two kinds of initial conditions:

(a) We are given the values of $y(x)$ and $y'(x)$ for some value $x=x_0$ of x .

(b) We are given the values of $y(x)$ for $x=x_1$ and $x=x_2$.

Examples:

1. $y'' - y' - 2y = 0 \Rightarrow a = -1, b = -2$. The characteristic equation (6) is written: $k^2 - k - 2 = 0$, with real roots $k_1=2, k_2=-1$. The general solution (7) is $y = C_1 e^{2x} + C_2 e^{-x}$. Assume the initial conditions $y=2$ and $y' = -5$ when $x=0$. Then, $C_1 = -1, C_2=3$ (show it!) and we get the *particular* solution $y = -e^{2x} + 3e^{-x}$.

2. $y'' - 6y' + 9y = 0 \Rightarrow a = -6, b=9$. The characteristic equation (6) is written: $k^2 - 6k + 9 = 0$, with real and equal roots $k_1=k_2=3$. The general solution (8) is $y = (C_1 + C_2 x)e^{3x}$.

3. $y'' - 4y' + 13y = 0 \Rightarrow a = -4, b=13$. The characteristic equation (6) is written: $k^2 - 4k + 13 = 0$, with complex conjugate roots $k_1=2+3i, k_2=2-3i$. The general solution (9) is (with $\alpha=2, \beta=3$): $y = e^{2x} (C_1 \cos 3x + C_2 \sin 3x)$. (Show that essentially the same result is found by making the alternative choice $\alpha=2, \beta=-3$.)

3. Harmonic oscillation

In a harmonic oscillation along the x -axis the total force on the oscillating body (of mass m) is $F = -kx$ ($k>0$), where x is the momentary displacement of the body from the position of equilibrium ($x=0$). By Newton's second law we have that $F=ma$, where a is the acceleration of the body: $a=d^2x/dt^2$. Therefore,

$$m d^2x / dt^2 = -kx$$

or, setting $k/m \equiv \omega^2$ (where we assume that $\omega>0$),

$$x'' + \omega^2 x = 0 \tag{10}$$

Eq. (10) is a homogeneous linear DE of the form (5) with x in place of y and t in place of x (notice that the first-derivative term is missing in this case). The characteristic equation (6) is written: $k^2 + \omega^2 = 0$ (or, analytically, $k^2 + 0k + \omega^2 = 0$), with complex roots $k = \pm i\omega$ (analytically, $k_1=0+i\omega, k_2=0-i\omega$). The general solution of (10) is given by (9), with $\alpha=0$ and $\beta=\omega$:

$$x = C_1 \cos \omega t + C_2 \sin \omega t \tag{11}$$

where we assume that the constant coefficients C_1 and C_2 are real in order for the solution (11) to have physical meaning.

The general solution (11) can be put in different but equivalent form by setting

$$C_1 = A \sin \varphi, \quad C_2 = A \cos \varphi \quad (A > 0) \Leftrightarrow A = (C_1^2 + C_2^2)^{1/2}, \quad \tan \varphi = C_1 / C_2.$$

Then,

$$x = A \sin(\omega t + \varphi) \tag{12}$$

The positive constant A is called the *amplitude* of the oscillation, while the angle φ is called the *initial phase* (the value of the *phase* $\omega t + \varphi$ at time $t=0$). The positive constant ω is the *angular frequency* of oscillation, to be called just “*frequency*” in the sequel.

Notice that, if we set $C_1 = A \cos \varphi$, $C_2 = -A \sin \varphi$ in (11), we will get the general solution of (10) in the form

$$x = A \cos(\omega t + \varphi) \tag{13}$$

which is equivalent to (12). Indeed, equation (13) follows directly from (12) by putting $\varphi + (\pi/2)$ in place of φ (which is arbitrary anyway) in the latter equation.

4. Damped oscillation

In a damped oscillation, in addition to the restoring force $-kx$, opposite to the displacement x from the equilibrium position, there is a frictional force $-\lambda v = -\lambda dx/dt$ ($\lambda > 0$) opposite to the velocity v . The total force on the body is³ $F = -kx - \lambda dx/dt$. By Newton’s law, $F = m d^2x/dt^2$. Hence,

$$m d^2x/dt^2 = -kx - \lambda dx/dt.$$

We set

$$k/m \equiv \omega_0^2 \quad (\omega_0 = \text{natural frequency of oscillation without damping}), \quad \lambda/m \equiv 2\gamma,$$

so that

$$x'' + 2\gamma x' + \omega_0^2 x = 0 \tag{14}$$

Eq. (14) is a homogeneous linear DE. The characteristic equation (6) is

$$k^2 + 2\gamma k + \omega_0^2 = 0 \quad \Rightarrow \quad k = -\gamma \pm (\gamma^2 - \omega_0^2)^{1/2}.$$

We distinguish the following cases:

1. *Large damping* $\Leftrightarrow \gamma > \omega_0$. We have two real solutions:

$$k_1 = -\gamma + (\gamma^2 - \omega_0^2)^{1/2}, \quad k_2 = -\gamma - (\gamma^2 - \omega_0^2)^{1/2}.$$

The general solution of (14) is of the form (7):

$$x = C_1 e^{k_1 t} + C_2 e^{k_2 t} \tag{15}$$

³ Note that a velocity-dependent force is *not* conservative. Thus, conservation of energy methods do not apply in this case.

Let us assume that $C_1 > 0$ and $C_2 > 0$. Given that $k_1 < 0$ and $k_2 < 0$ (why?) we see that $x > 0$ at all times t and, moreover, $x \rightarrow 0$ as $t \rightarrow \infty$. That is, as the time t increases, the moving object approaches the equilibrium position $x=0$ without ever crossing it. The motion is therefore *non-oscillatory*.

2. *Critical damping* $\Leftrightarrow \gamma = \omega_0$. Then, $k_1 = k_2 = -\gamma$, and the general solution of (14) is of the form (8):

$$x = (C_1 + C_2 t) e^{-\gamma t} = (C_1 + C_2 t) e^{-\gamma t} \quad (16)$$

If we assume that $C_1 > 0$ and $C_2 > 0$, we see again that $x > 0$ at all t and that $x \rightarrow 0$ as $t \rightarrow \infty$. (For the term $t e^{-\gamma t} = t / e^{\gamma t}$ we may use L'Hospital's rule for the indeterminate form ∞/∞ ; show this!) Thus, there is no oscillation in this case either.

3. *Small damping* $\Leftrightarrow \gamma < \omega_0$. We have two complex conjugate solutions:

$$k = -\gamma \pm i \omega_1 \quad \text{where} \quad \omega_1 = (\omega_0^2 - \gamma^2)^{1/2}.$$

The general solution will be of the form (9), with $\alpha = -\gamma$ and $\beta = \omega_1$:

$$x = e^{-\gamma t} (C_1 \cos \omega_1 t + C_2 \sin \omega_1 t),$$

or, by setting $C_1 = A \sin \varphi$, $C_2 = A \cos \varphi$ ($A > 0$),

$$x = A e^{-\gamma t} \sin(\omega_1 t + \varphi) \quad (17)$$

We notice that the amplitude $A e^{-\gamma t}$ decreases exponentially with time (Fig. 1). Thus, strictly speaking, damped oscillatory motion is *not* periodic.

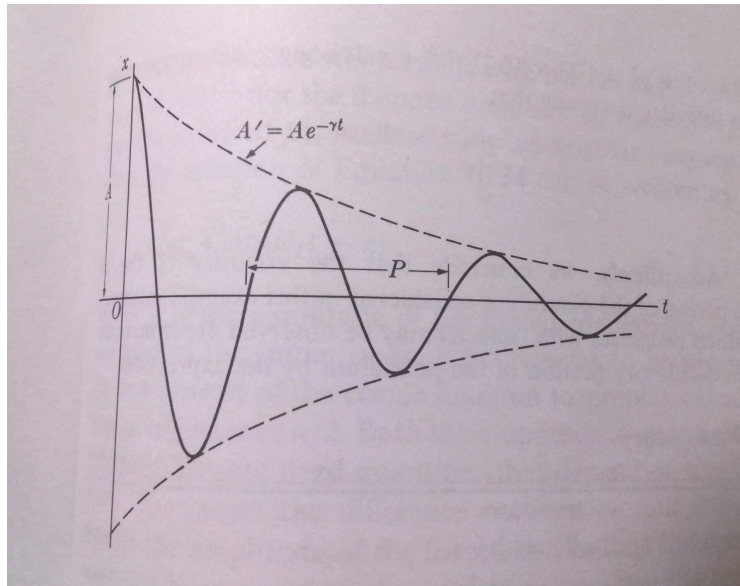


Fig. 1

5. Forced oscillation

In a forced oscillation, in addition to the restoring force $-kx$ and the frictional force $-\lambda v = -\lambda dx/dt$ the body is subject to an external force of the form

$$F(t) = F_0 \sin \omega_f t \quad (F_0 > 0).$$

The total force on the body is $F = -kx - \lambda dx/dt + F_0 \sin \omega_f t$. By Newton's law we have that

$$m d^2x/dt^2 = -kx - \lambda dx/dt + F_0 \sin \omega_f t.$$

We set

$$k/m \equiv \omega_0^2 \quad (\omega_0 = \text{natural frequency}), \quad \lambda/m \equiv 2\gamma, \quad F_0/m \equiv f_0,$$

so that

$$x'' + 2\gamma x' + \omega_0^2 x = f_0 \sin \omega_f t \quad (18)$$

Eq. (18) is a non-homogeneous linear DE. According to Theorem 2 of Sec. 1, its general solution is the sum of the general solution of the corresponding homogeneous equation,

$$x'' + 2\gamma x' + \omega_0^2 x = 0,$$

and *any particular solution* of (18). For small damping ($\gamma < \omega_0$) the general solution of the homogeneous equation is given by (17):

$$x = A_1 e^{-\gamma t} \sin(\omega_1 t + \varphi_1) \quad \text{where} \quad \omega_1 = (\omega_0^2 - \gamma^2)^{1/2}.$$

As can be verified, a particular solution of (18) is the following:

$$x = A \sin(\omega_f t + \varphi) \quad (19)$$

where

$$A = \frac{f_0}{\left[(\omega_f^2 - \omega_0^2)^2 + 4\gamma^2 \omega_f^2 \right]^{1/2}} \quad \text{and} \quad \tan \varphi = \frac{2\gamma \omega_f}{\omega_f^2 - \omega_0^2} \quad (20)$$

The general solution of (18) is, therefore,

$$x = A_1 e^{-\gamma t} \sin(\omega_1 t + \varphi_1) + A \sin(\omega_f t + \varphi) \quad (21)$$

with *arbitrary* A_1, φ_1 . The first term on the right in (21) decreases exponentially with time and dies out quickly. In a steady-state situation, therefore, what remains is the particular solution (19):

$$x = A \sin(\omega_f t + \varphi).$$

The amplitude A of oscillation is a function of the applied frequency ω_f , according to (20). This amplitude attains a maximum value when the denominator in the first relation (20) becomes minimum. This occurs when

$$\omega_f = (\omega_0^2 - 2\gamma^2)^{1/2} \equiv \omega_A \quad (22)$$

Proof: We set $\omega_f \equiv \omega$, for simplicity, and we consider the function

$$\Psi(\omega) = (\omega^2 - \omega_0^2)^2 + 4\gamma^2\omega^2,$$

so that $A = f_0 / [\Psi(\omega)]^{1/2}$. We can show that

$$\Psi'(\omega) = 0 \text{ for } \omega = (\omega_0^2 - 2\gamma^2)^{1/2} = \omega_A \text{ and } \Psi''(\omega_A) = 8\omega_A^2 > 0.$$

Thus, for small damping ($2\gamma^2 < \omega_0^2$) the function $\Psi(\omega)$ is *minimum*, hence the amplitude A is *maximum*, when $\omega_f = \omega_A$. This situation is called *amplitude resonance*.

In Fig. 2 it is assumed that $\lambda_1 < \lambda_2 \Leftrightarrow \gamma_1 < \gamma_2$. This means that, in accordance with (22), $\omega_{A,1} > \omega_{A,2}$. In the case of no damping ($\lambda = 0 \Leftrightarrow \gamma = 0$) Eq. (22) yields $\omega_A = \omega_0$. In other words, in an *undamped* forced oscillation the amplitude becomes maximum (in fact, infinite) when the applied frequency ω_f is equal to the natural frequency ω_0 of oscillation.

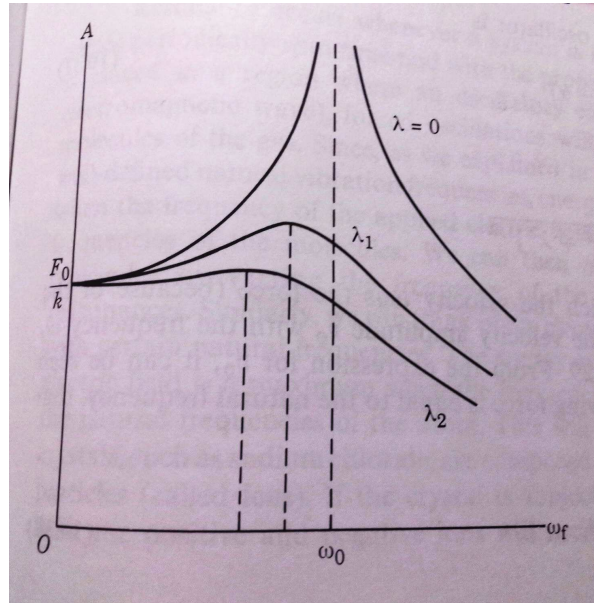


Fig. 2

By differentiating (19) we find the velocity of the oscillating body:

$$v = dx/dt = \omega_f A \cos(\omega_f t + \varphi) \equiv v_0 \cos(\omega_f t + \varphi)$$

where, by (20),

$$v_0 = \omega_f A = \frac{f_0}{\left[\left(1 - \frac{\omega_0^2}{\omega_f^2} \right)^2 + 4\gamma^2 \right]^{1/2}} .$$

The velocity amplitude v_0 becomes maximum when the denominator on the right is minimum, which occurs for $\omega_f = \omega_0$ (Fig. 3). The kinetic energy $mv_0^2/2$ then reaches its maximum value and there is *energy resonance*.

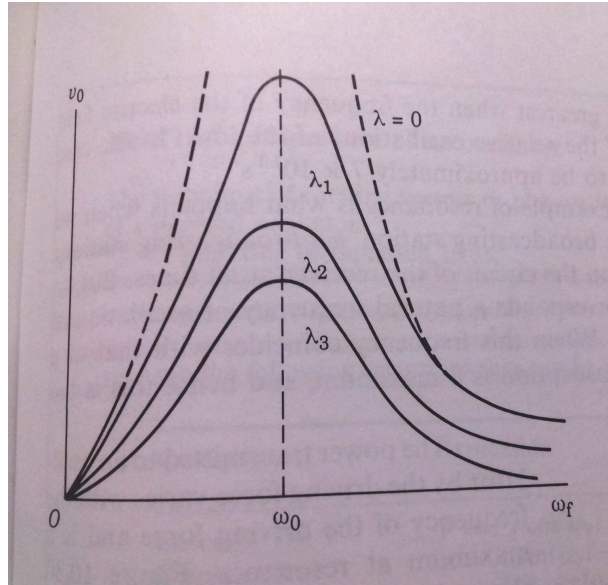


Fig. 3

Note that, in contrast to amplitude resonance, the frequency ω_f for energy resonance is independent of the damping factor λ and is always equal to the *natural frequency* ω_0 of the oscillator. At this frequency the work supplied by the external force $F(t)$ to the oscillator per unit time is maximum. That is, the oscillator absorbs the largest possible power from the external agent that exerts the force F .

Notice also that, in the case of zero damping ($\lambda=0 \Leftrightarrow \gamma=0$) the velocity amplitude v_0 becomes *infinite* at energy resonance, i.e., for $\omega_f = \omega_0$. This rather unphysical situation is, of course, purely theoretical since a mechanical motion with no friction whatsoever is practically impossible!

Appendix: Amplitude dependence of period

As we have shown, the general solution to the one-dimensional conservative Newtonian problem is

$$\int_{x_0}^x \frac{\pm dx}{\left\{ \frac{2}{m} [E - U(x)] \right\}^{1/2}} = t - t_0 \quad (1)$$

where the plus sign is chosen for motion in the positive direction ($v > 0$, $x > x_0$) while the minus sign applies to motion in the negative direction ($v < 0$, $x < x_0$).

Let us assume that the potential energy $U(x)$ has the form of a U-shaped potential well (Fig. 1) such that $U(0) = 0$ and $U(x) > 0$ for $x \neq 0$. The graph of $U(x)$ is assumed to be symmetric with respect to the axis $x = 0$, which means that $U(x)$ is an *even* function: $U(-x) = U(x)$.

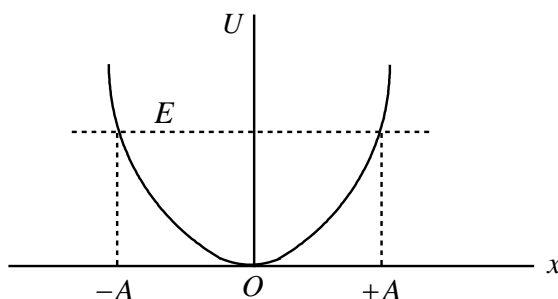


Fig. 1

If E is the total mechanical energy of the particle, then, according to Fig. 1, the motion is bounded between the points $-A$ and $+A$ of the x -axis, which are turning points at which the particle stops momentarily. Since E is constant, its value at all points equals its value at the turning points; i.e.,

$$E = U(\pm A) \quad (2)$$

The time it takes for a complete journey from $-A$ to $+A$ and back to $-A$ is found by using (1) with the appropriate sign for each direction of motion:

$$P = \int_{-A}^A \frac{dx}{\left\{ \dots \right\}^{1/2}} + \int_A^{-A} \frac{-dx}{\left\{ \dots \right\}^{1/2}} \Rightarrow$$

$$P = 2 \int_{-A}^A \frac{dx}{\left\{ \frac{2}{m} [E - U(x)] \right\}^{1/2}} = (2m)^{1/2} \int_{-A}^A [E - U(x)]^{-1/2} dx \quad (3)$$

Since P is fixed for a given A , the motion is periodic about the point $x = 0$, with amplitude equal to A and with period P . It follows from (2) and (3) that the period P depends on A and thus on the total energy E of the particle. We will now show that an exception where P does *not* depend on A (thus on E also) is simple harmonic motion.

Since $U(x)$ is an even function with $U(0)=0$, it can be expanded into a Maclaurin series of the form

$$U(x) = \sum_{l=1}^{\infty} a_l x^{2l} \quad (4)$$

where the coefficients a_l are not necessarily all different from zero. From (2) we have

$$E = U(\pm A) = \sum_{l=1}^{\infty} a_l A^{2l}$$

so that

$$E - U(x) = \sum_{l=1}^{\infty} a_l (A^{2l} - x^{2l}).$$

Equation (3) then yields

$$P = (2m)^{1/2} \int_{-A}^A \left[\sum_{l=1}^{\infty} a_l (A^{2l} - x^{2l}) \right]^{-1/2} dx.$$

By setting $x/A=u \Leftrightarrow x=Au$, we get:

$$P = (2m)^{1/2} A \int_{-1}^1 \left[\sum_{l=1}^{\infty} a_l A^{2l} (1-u^{2l}) \right]^{-1/2} du \quad (5)$$

It is obvious that, in general, P depends on A . The only exception where P is *not* dependent on A is the case where the following condition is satisfied: $a_l=0$ for $l \neq 1$. That is, the only nonvanishing coefficient a_l in the series (4) is a_1 . By setting $a_1 = k/2$ the potential energy (4) reduces to $U(x) = kx^2/2$, which corresponds to a restoring force of the form

$$F(x) = -dU/dx = -kx \quad (6)$$

The periodic motion is then *simple harmonic motion* (SHM) and the period (5) reduces to

$$\begin{aligned} P &= 2(m/k)^{1/2} \int_{-1}^1 (1-u^2)^{-1/2} du = 2(m/k)^{1/2} [\arcsin u]_{-1}^1 \\ &= 2(m/k)^{1/2} \left[\frac{\pi}{2} - \left(-\frac{\pi}{2} \right) \right] \Rightarrow \\ P &= 2\pi \left(\frac{m}{k} \right)^{1/2} \equiv \frac{2\pi}{\omega} \quad \text{where} \quad \omega = \frac{2\pi}{P} = \left(\frac{k}{m} \right)^{1/2}. \end{aligned}$$

We notice that the period of SHM is amplitude-independent, hence also energy-independent.

It is of interest to examine a one-dimensional periodic motion that follows a *curved* path (where by “one-dimensional” we now mean that a single generalized coordinate – such as, e.g., an angle or a distance along the curve – is needed in order to specify the location of the particle). A nice example is that of an oscillating pendulum (Fig. 2; see also Sec. 5.5 and Problem 25 of [1]). The position of the mass m is specified by the arc length $OA=s=l\theta$ or, equivalently, by the angle θ (in rad). The algebraic value of the velocity of m is $v=ds/dt=l d\theta/dt$; it may be positive or negative, depending on the direction of motion relative to the unit tangent vector \hat{u}_T .

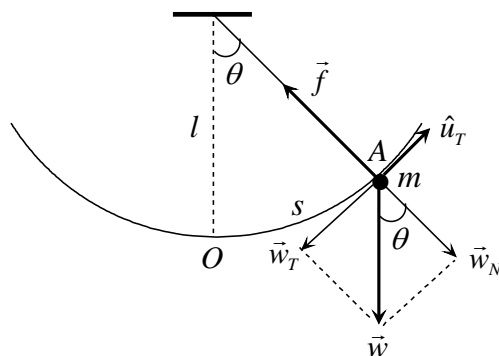


Fig. 2

The motion is governed by the tangential component $w_T = -mg \sin\theta$ (algebraic value) of the weight w . The tangential equation of motion of m is

$$m dv/dt = -mg \sin\theta \Rightarrow dv/dt = -g \sin\theta \quad (7)$$

We seek a conserved quantity that associates the velocity v with the position θ . We could, of course, work with (7) directly, but there is an easier way; namely, conservation of mechanical energy. This principle may be applied in view of the fact that the mass m is subject to the conservative force of gravity and the tension f of the string which, being normal to the velocity, produces no work (cf. Sec. 4.5 of [1]). The potential energy of m at point A (Fig. 2) is

$$U(\theta) = mg(l - l \cos\theta) = mgl(1 - \cos\theta),$$

where we have assumed that $U(0)=0$ (i.e., U is zero at the lowest point O). If α is the angular amplitude of oscillation (i.e., the maximum angle of deflection of the string from the vertical) then at $\theta = \pm\alpha$ the kinetic energy T vanishes and the total mechanical energy E is equal to $U(\pm\alpha)$. Applying conservation of mechanical energy between an arbitrary angle θ and the maximum angle $\theta = \alpha$, we have:

$$\begin{aligned} m v^2/2 + mgl(1 - \cos\theta) &= 0 + mgl(1 - \cos\alpha) \Rightarrow \text{(after eliminating } m) \\ v^2 &= 2gl(\cos\theta - \cos\alpha) \end{aligned} \quad (8)$$

Exercise: By differentiating (8) with respect to t and by using the fact that $v=l d\theta/dt$, recover the equation of motion (7). Conversely, show that (8) is a direct consequence of (7). [*Hint:* Multiply (7) by v .]

Setting $v=l d\theta/dt$ in (8), we get a first-order differential equation:

$$d\theta/dt = \pm [(2g/l)(\cos\theta - \cos\alpha)]^{1/2},$$

which is integrated to give

$$\int_{\theta_0}^{\theta} \pm \left[\frac{2g}{l} (\cos\theta - \cos\alpha) \right]^{-1/2} d\theta = t - t_0.$$

The period of oscillation is [cf. Eq. (3)]

$$\begin{aligned} P &= 2 \int_{-\alpha}^{\alpha} \left[\frac{2g}{l} (\cos\theta - \cos\alpha) \right]^{-1/2} d\theta \\ &= (2l/g)^{1/2} \int_{-\alpha}^{\alpha} (\cos\theta - \cos\alpha)^{-1/2} d\theta \end{aligned} \quad (9)$$

Obviously, P depends on the angular amplitude α . Let us assume, however, that this amplitude is very small: $\alpha \ll 1$. We may then make the approximations

$$\cos\theta \approx 1 - \theta^2/2 \quad \text{and} \quad \cos\alpha \approx 1 - \alpha^2/2.$$

Furthermore, we set $\theta/\alpha = u \Leftrightarrow \theta = \alpha u$. It is then a straightforward exercise to show that (9) reduces to

$$\begin{aligned} P &= 2(l/g)^{1/2} \int_{-1}^1 (1-u^2)^{-1/2} du = 2(l/g)^{1/2} [\arcsin u]_{-1}^1 \\ &= 2(l/g)^{1/2} \left[\frac{\pi}{2} - \left(-\frac{\pi}{2} \right) \right] \Rightarrow \\ P &= 2\pi(l/g)^{1/2}, \end{aligned}$$

which is the familiar expression for the period of oscillation of a pendulum executing simple harmonic motion for small angles of deflection from the vertical. Once again, the SHM is seen to be the only one-dimensional periodic motion in which the period does not depend on the amplitude of oscillation.

Reference

- [1] C. J. Papachristou, *Introduction to Mechanics of Particles and Systems* (Springer, 2020).⁴

⁴ Manuscript: <http://metapublishing.org/index.php/MP/catalog/book/68>