Introduction to Perturbation Techniques

ALI HASAN NAYFEH

University Distinguished Professor
Virginia Polytechnic Institute and State University
Blacksburg, Virginia
and
Yarmouk University, Irbid, Jordan

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To my parents Hasan and Khadrah
my wife Samirah
and my children Mahir, Tariq, Samir, and Nader
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Preface

Many of the problems facing physicists, engineers, and applied mathematicians involve such difficulties as nonlinear governing equations, variable coefficients, and nonlinear boundary conditions at complex known or unknown boundaries that preclude solving them exactly. Consequently, solutions are approximated using numerical techniques, analytic techniques, and combinations of both. Foremost among the analytic techniques are the systematic methods of perturbations (asymptotic expansions) in terms of a small or a large parameter or coordinate. This book is concerned only with these perturbation techniques.

The author’s book Perturbation Methods presents in a unified way an account of most of the perturbation techniques, pointing out their similarities, differences, and advantages, as well as their limitations. Although the techniques are described by means of examples that start with simple ordinary equations that can be solved exactly and progress toward complex partial-differential equations, the material is concise and advanced and therefore is intended for researchers and advanced graduate students only. The purpose of this book, however, is to present the material in an elementary way that makes it easily accessible to advanced undergraduates and first-year graduate students in a wide variety of scientific and engineering fields. As a result of teaching perturbation methods for eight years to first-year and advanced graduate students at Virginia Polytechnic Institute and State University, I have selected a limited number of techniques and amplified their description considerably. Also I have attempted to answer the questions most frequently raised by my students. The techniques are described by means of simple examples that consist mainly of algebraic and ordinary-differential equations.

The material in Chapters 3 and 15 and Appendices A and B cannot be found in Perturbation Methods. Chapter 3 discusses asymptotic expansions of integrals. Chapter 15 is devoted to the determination of the adjoints of homogeneous linear equations (algebraic, ordinary-differential, partial-differential, and integral equations) and the solvability conditions of linear inhomogeneous problems. Appendix A summarizes trigonometric identities, and Appendix B summarizes the properties of linear ordinary-differential equations and describes the symbolic method of determining the solutions of homogeneous and inhomogeneous ordinary-differential equations with constant coefficients.
The reader should have a background in calculus and elementary ordinary-differential equations.

Each chapter contains a number of exercises. For more exercises, the reader is referred to *Perturbation Methods* by Nayfeh and *Nonlinear Oscillations* by Nayfeh and Mook. Since this book is elementary, only a list of the pertinent books is included in the bibliography without any attempt of citing them in the text.

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Ali Hasan Nayfeh

*Blackburg, Virginia*

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CHAPTER 1

Introduction

1.1. Dimensional Analysis

Exact solutions are rare in many branches of fluid mechanics, solid mechanics, motion, and physics because of nonlinearities, inhomogeneities, and general boundary conditions. Hence, engineers, physicists, and applied mathematicians are forced to determine approximate solutions of the problems they are facing. These approximations may be purely numerical, purely analytical, or a combination of numerical and analytical techniques. In this book, we concentrate on the purely analytical techniques, which, when combined with a numerical technique such as a finite-difference or a finite-element technique, yield very powerful and versatile techniques.

The key to solving modern problems is mathematical modeling. This process involves keeping certain elements, neglecting some, and approximating yet others. To accomplish this important step, one needs to decide the order of magnitude (i.e., smallness or largeness) of the different elements of the system by comparing them with each other as well as with the basic elements of the system. This process is called nondimensionalization or making the variables dimensionless. Consequently, one should always introduce dimensionless variables before attempting to make any approximations. For example, if an element has a length of one centimeter, would this element be large or small? One cannot answer this question without knowing the problem being considered. If the problem involves the motion of a satellite in an orbit around the earth, then one centimeter is very very small. On the other hand, if the problem involves intermolecular distances, then one centimeter is very very large. As a second example, is one gram small or large? Again one gram is very very small compared with the mass of a satellite but it is very very large compared with the mass of an electron. Therefore, expressing the equations in dimensionless form brings out the important dimensionless parameters that govern the behavior of the system. Even if one is not interested in approximations, it is recommended that one perform this important step before analyzing the system or presenting
experimental data. Next, we give a few examples illustrating the process of nondimensionalization.

**EXAMPLE 1**

We consider the motion of a particle of mass $m$ restrained by a linear spring having the constant $k$ and a viscous damper having the coefficient $\mu$, as shown in Figure 1-1. Using Newton's second law of motion, we have

$$m \frac{d^2u}{dt^2} + \mu \frac{du}{dt} + ku = 0$$

where $u$ is the displacement of the particle and $t$ is time. Let us assume that the particle was released from rest from the position $u_0$ so that the initial conditions are

$$u(0) = u_0, \quad \frac{du}{dt}(0) = 0$$

In this case, $u$ is the dependent variable and $t$ is the independent variable. They need to be made dimensionless by using a characteristic distance and a characteristic time of the system. The displacement $u$ can be made dimensionless by using the initial displacement $u_0$ as a characteristic distance, whereas the time $t$ can be made dimensionless by using the inverse of the system's natural frequency $\omega_0 = \sqrt{k/m}$. Thus, we put

$$u^* = \frac{u}{u_0}, \quad t^* = \omega_0 t$$

where the asterisk denotes dimensionless quantities. Then,

$$\frac{du}{dt} = \frac{d(u_0 u^*)}{dt^*}, \quad \frac{dt^*}{dt} = \omega_0 u_0 \frac{du^*}{dt^*}$$

$$\frac{d^2u}{dt^2} = \omega_0^2 u_0 \frac{d^2u^*}{dt^{*2}}$$

so that (1.1) becomes

![Figure 1-1. A mass restrained by a spring and a viscous damper.](image-url)
Hence,

\[ \frac{d^2 u^*}{dt^2} + \mu^* \frac{du^*}{dt} + \frac{k}{m\omega_0^2} u^* = 0 \]

or

\[ \frac{d^2 u^*}{dt^2} + \mu^* \frac{du^*}{dt} + u^* = 0 \]  

(1.3)

where

\[ \mu^* = \frac{\mu}{m\omega_0} \]  

(1.4)

In terms of the above dimensionless quantities, (1.2) becomes

\[ u^*(0) = 1 \quad \text{and} \quad \frac{du^*}{dt^*}(0) = 0 \]  

(1.5)

Thus, the solution to the present problem depends only on the single parameter \( \mu^* \), which represents the ratio of the damping force to the inertia force or the restoring force of the spring. If this ratio is small, then one can use the dimensionless quantity \( \mu^* \) as the small parameter in obtaining an approximate solution of the problem, and we speak of a lightly damped system. We should note that the system cannot be considered lightly damped just because \( \mu \) is small; \( \mu^* = \mu/m\omega_0 = \mu/\sqrt{k/m} \) must be small.

**EXAMPLE 2**

Let us assume that the spring force is a nonlinear function of \( u \) according to

\[ f_{\text{spring}} = ku + k_2 u^2 \]  

(1.6)

where \( k \) and \( k_2 \) are constants. Then, (1.1) becomes

\[ m \frac{d^2 u}{dt^2} + \mu \frac{du}{dt} + ku + k_2 u^2 = 0 \]  

(1.7)

Again, using the same dimensionless quantities as in the preceding example, we have

\[ m\omega_0^2 \frac{d^2 u^*}{dt^2} + \mu_0 \omega_0 \frac{du^*}{dt} + k_0 u^* + k_2 u_0^2 u^* = 0 \]

or
4 INTRODUCTION

\[
\frac{d^2 u^*}{dt^2} + \mu^* \frac{du^*}{dt} + u^* + \epsilon u^{*2} = 0
\]  

(1.8)

where

\[
\mu^* = \frac{\mu}{m\omega_0} \quad \text{and} \quad \epsilon = \frac{k_2 u_0}{k}
\]  

(1.9)

The initial conditions transform as in (1.5). Thus, the present problem is a function of the two dimensionless parameters \( \mu^* \) and \( \epsilon \). As before, \( \mu^* \) represents the ratio of the damping force to the inertia force or the linear restoring force. The parameter \( \epsilon \) represents the ratio of the nonlinear and linear restoring forces of the spring.

When we speak of a weakly nonlinear system, we mean that \( k_2 u_0/k \) is small. Even if \( k_2 \) is small compared with \( k \), the nonlinearity will not be small if \( u_0 \) is large compared with \( k/k_2 \). Thus, \( \epsilon \) is the parameter that characterizes the nonlinearity.

EXAMPLE 3

As a third example, we consider the motion of a spaceship of mass \( m \) that is moving in the gravitational field of two fixed mass-centers whose masses \( m_1 \) and \( m_2 \) are much much bigger than \( m \). With respect to the Cartesian coordinate system shown in Figure 1-2, the equations of motion are

\[
m \frac{d^2 x}{dt^2} = -\frac{mm_1 G x}{(x^2 + y^2)^{3/2}} - \frac{mm_2 G (x - L)}{[(x - L)^2 + y^2]^{3/2}}
\]  

(1.10)

\[
m \frac{d^2 y}{dt^2} = -\frac{mm_1 G y}{(x^2 + y^2)^{3/2}} - \frac{mm_2 G y}{[(x - L)^2 + y^2]^{3/2}}
\]  

(1.11)

where \( t \) is the time, \( G \) is the gravitational constant, and \( L \) is the distance between \( m_1 \) and \( m_2 \).

In this case, the dependent variables are \( x \) and \( y \) and the independent variable is \( t \). Clearly, a characteristic length of the problem is \( L \), the distance between the two mass centers. A characteristic time of the problem is not as obvious. Since

\[\text{Figure 1-2. A satellite in the gravitational field of two fixed mass centers.}\]
the motions of the masses \( m_1 \) and \( m_2 \) are assumed to be independent of that of the spaceship, \( m_1 \) and \( m_2 \) move about their center of mass in ellipses. The period of oscillation is

\[
T = \frac{2\pi L^{3/2}}{\sqrt{G(m_1 + m_2)}}
\]

so that the frequency of oscillation is

\[
\omega_0 = L^{-3/2} \sqrt{G(m_1 + m_2)}
\]  

Thus, we use the inverse of \( \omega_0 \) as a characteristic time. Then, we introduce dimensionless quantities defined by

\[
x^* = \frac{x}{L} \quad y^* = \frac{y}{L} \quad t^* = \omega_0 t
\]

so that

\[
\frac{dx}{dt} = \frac{d(x^*L)}{dt^*} = L \omega_0 \frac{dx^*}{dt^*} \quad \frac{d^2x}{dt^2} = L \omega_0^2 \frac{d^2x^*}{dt^*}
\]

\[
\frac{dy}{dt} = \frac{d(y^*L)}{dt^*} = L \omega_0 \frac{dy^*}{dt^*} \quad \frac{d^2y}{dt^2} = L \omega_0^2 \frac{d^2y^*}{dt^*}
\]

Hence, (1.10) and (1.11) become

\[
mL \omega_0^2 \frac{d^2x^*}{dt^*} = \frac{m_1 G L x^*}{L^2 (x^*^2 + y^*^2)^{3/2}} - \frac{m_2 G L (x^* - 1)}{L^2 (x^* - 1)^2 + L^2 y^*^2}^{3/2}
\]

\[
mL \omega_0^2 \frac{d^2y^*}{dt^*} = \frac{m_1 G L y^*}{L^2 (x^*^2 + y^*^2)^{3/2}} - \frac{m_2 G L y^*}{L^2 (x^* - 1)^2 + L^2 y^*^2}^{3/2}
\]

or

\[
\frac{d^2x^*}{dt^*} = \frac{m_1 G}{L^3 \omega_0^2} \frac{x^*}{(x^*^2 + y^*^2)^{3/2}} - \frac{m_2 G}{L^3 \omega_0^2} \frac{(x^* - 1)}{[(x^* - 1)^2 + y^*^2]^{3/2}}
\]

\[
\frac{d^2y^*}{dt^*} = \frac{m_1 G}{L^3 \omega_0^2} \frac{y^*}{(x^*^2 + y^*^2)^{3/2}} - \frac{m_2 G}{L^3 \omega_0^2} \frac{y^*}{[(x^* - 1)^2 + y^*^2]^{3/2}}
\]

Using (1.12), we have

\[
\frac{m_1 G}{L^3 \omega_0^2} = \frac{m_1}{m_1 + m_2} \quad \frac{m_2 G}{L^3 \omega_0^2} = \frac{m_2}{m_1 + m_2}
\]

Hence, if we put

\[
\frac{m_2}{m_1 + m_2} = \epsilon \quad \text{then} \quad \frac{m_1}{m_1 + m_2} = 1 - \epsilon
\]
and (1.14) and (1.15) become

\[
\frac{d^2 x^*}{dt^2} = -\frac{(1 - e)x^*}{(x^* + y^*)^{3/2}} - \frac{e(x^* - 1)}{[(x^* - 1)^2 + y^*]^3/2} \tag{1.17}
\]

\[
\frac{d^2 y^*}{dt^2} = -\frac{(1 - e)y^*}{(x^* + y^*)^{3/2}} - \frac{ey^*}{[(x^* - 1)^2 + y^*]^3/2} \tag{1.18}
\]

Therefore, the problem depends only on the parameter \( e \), which is usually called the reduced mass. If \( m_1 \) represents the mass of the earth and \( m_2 \) the mass of the moon, then

\[
e \approx \frac{80}{1 + \frac{1}{80}} = \frac{1}{81}
\]

which is small and can be used as a perturbation parameter in determining an approximate solution to the motion of a spacecraft in the gravitational field of the earth and the moon.

**EXAMPLE 4**

As a fourth example, we consider the vibration of a clamped circular plate of radius \( a \) under the influence of a uniform radial load. If \( w \) is the transverse displacement of the plate, then the linear vibrations of the plate are governed by

\[
D \nabla^4 w - P \nabla^2 w - \rho \frac{\partial^2 w}{\partial t^2} = 0 \tag{1.19}
\]

where \( t \) is the time, \( D \) is the plate rigidity, \( P \) is the uniform radial load, and \( \rho \) is the plate density per unit area. The boundary conditions are

\[
w = 0 \quad \frac{\partial w}{\partial r} = 0 \quad \text{at } r = a
\]

\[
w < \infty \quad \text{at } r = 0 \tag{1.20}
\]

In this case, \( w \) is the dependent variable and \( t \) and \( r \) are the independent variables. Clearly, \( a \) is a characteristic length of the problem. The characteristic time is assumed to be \( T \) and it is specified below. Then, we define dimensionless variables according to

\[
w^* = \frac{w}{a} \quad r^* = \frac{r}{a} \quad t^* = \frac{t}{T}
\]

Hence,
\[ \begin{align*}
\frac{\partial w}{\partial r} &= \frac{\partial (aw^*)}{\partial r^*} \frac{dr^*}{dr} = \frac{\partial w^*}{\partial r^*} \\
\frac{\partial w}{\partial \theta} &= \frac{\partial (aw^*)}{\partial \theta} = a \frac{\partial w^*}{\partial \theta} \\
\frac{\partial w}{\partial t} &= \frac{\partial (aw^*)}{\partial t^*} \frac{dt^*}{dt} = a \frac{\partial w^*}{\partial t^*} \\
\end{align*} \]

Since

\[ \nabla^2 = \frac{\partial^2}{\partial r^2} + \frac{1}{r} \frac{\partial}{\partial r} + \frac{1}{r^2} \frac{\partial^2}{\partial \theta^2} \]

(1.19) becomes

\[ \begin{align*}
\frac{D}{a^3} \left( \frac{\partial^2}{\partial r^*^2} + \frac{1}{r^*} \frac{\partial}{\partial r^*} + \frac{1}{r^*^2} \frac{\partial^2}{\partial \theta^*^2} \right)^2 w^* - \frac{P}{a} \left( \frac{\partial^2}{\partial r^*^2} + \frac{1}{r^*} \frac{\partial}{\partial r^*} + \frac{1}{r^*^2} \frac{\partial^2}{\partial \theta^*^2} \right) w^* \\
- \frac{\rho a^2}{T^2} \frac{\partial^2 w^*}{\partial t^*^2} = 0
\end{align*} \]

or

\[ \frac{D}{a^2 P} \nabla^*^4 w^* - \nabla^*^2 w^* - \frac{\rho a^2}{PT^2} \frac{\partial^2 w^*}{\partial t^*^2} = 0 \]

(1.21)

We can choose \( T \) to make the coefficient of \( \frac{\partial^2 w^*}{\partial t^*^2} \) equal to 1, that is, \( T = \frac{a}{\sqrt{\rho P}} \). Then, (1.21) becomes

\[ \epsilon \nabla^*^4 w^* - \nabla^*^2 w^* - \frac{\partial^2 w^*}{\partial t^*^2} = 0 \]

(1.22)

where

\[ \epsilon = \frac{D}{a^2 P} \]

(1.23)

In terms of dimensionless quantities, the boundary conditions (1.20) become

\[ \begin{align*}
\frac{\partial w^*}{\partial r^*} &= 0 \quad \text{at} \quad r^* = 1 \\
\frac{\partial w^*}{\partial r^*} &= 0 \quad \text{at} \quad r^* = 0 \\
\end{align*} \]

(1.24)

Therefore, the problem depends on the single dimensionless parameter \( \epsilon \). If the radial load is large compared with \( \frac{D}{a^2} \), then \( \epsilon \) is small and can be used as a perturbation parameter.
EXAMPLE 5

As a final example, we consider steady incompressible flow past a flat plate. The problem is governed by

\[ \frac{\partial u}{\partial x} + \frac{\partial v}{\partial y} = 0 \]  

\[ \rho \left( \frac{\partial u}{\partial x} + v \frac{\partial u}{\partial y} \right) = -\frac{\partial p}{\partial x} + \mu \left( \frac{\partial^2 u}{\partial x^2} + \frac{\partial^2 u}{\partial y^2} \right) \]  

\[ \rho \left( \frac{\partial v}{\partial x} + v \frac{\partial v}{\partial y} \right) = -\frac{\partial p}{\partial y} + \mu \left( \frac{\partial^2 v}{\partial x^2} + \frac{\partial^2 v}{\partial y^2} \right) \]  

\[ u = v = 0 \quad \text{at} \quad y = 0 \]  

\[ u \to U_\infty, \quad v \to 0 \quad \text{as} \quad x \to -\infty \]  

where \( u \) and \( v \) are the velocity components in the \( x \) and \( y \) directions, respectively, \( p \) is the pressure, \( \rho \) is the density, and \( \mu \) is the coefficient of viscosity.

In this case, \( u, v, \) and \( p \) are the dependent variables and \( x \) and \( y \) are the independent variables. To make the equations dimensionless, we use \( L \) as a characteristic length, where \( L \) is the distance from the leading edge to a specified point on the plate as shown in Figure 1-3, and use \( U_\infty \) as a characteristic velocity. We take \( \rho U_\infty^2 \) as a characteristic pressure. Thus, we define dimensionless quantities according to

\[ u^* = \frac{u}{U_\infty}, \quad v^* = \frac{v}{U_\infty}, \quad p^* = \frac{p}{\rho U_\infty^2}, \quad x^* = \frac{x}{L}, \quad y^* = \frac{y}{L} \]

Then,

\[ \frac{\partial u}{\partial x} = \frac{\partial (U_\infty u^*)}{\partial x^*}, \quad \frac{\partial u}{\partial y} = \frac{U_\infty}{L} \frac{\partial u^*}{\partial y^*}, \quad \frac{\partial^2 u}{\partial x^2} = \frac{U_\infty}{L^2} \frac{\partial^2 u^*}{\partial x^{*2}}, \quad \frac{\partial^2 u}{\partial y^2} = \frac{\partial^2 u^*}{\partial y^{*2}} \]

\[ \frac{U_\infty}{L} \]

\[ U_\infty \]

\[ \to \]

\[ x, y \]

\[ \to \]

\[ x, u \]
\[
\begin{align*}
\frac{\partial v}{\partial x} &= \frac{U_\infty}{L} \frac{\partial v^*}{\partial x^*} \quad \frac{\partial v}{\partial y} = \frac{U_\infty}{L} \frac{\partial v^*}{\partial y^*} \\
\frac{\partial^2 v}{\partial x^2} &= \frac{U_\infty}{L} \frac{\partial^2 v^*}{\partial x^{*2}} \\
\frac{\partial^2 v}{\partial y^2} &= \frac{U_\infty}{L} \frac{\partial^2 v^*}{\partial y^{*2}} \\
\frac{\partial p}{\partial x} &= \rho U_\infty^2 \frac{\partial p^*}{\partial x^*} \quad \frac{\partial p}{\partial y} = \rho U_\infty^2 \frac{\partial p^*}{\partial y^*} \\
\end{align*}
\]

Hence, (1.25) through (1.28) become
\[
\begin{align*}
\frac{U_\infty}{L} \frac{\partial u^*}{\partial x^*} + \frac{U_\infty}{L} \frac{\partial v^*}{\partial y^*} &= 0 \\
\frac{\rho U_\infty^2}{L} \frac{\partial u^*}{\partial x^*} + \frac{\rho U_\infty^2}{L} \frac{\partial v^*}{\partial y^*} &= -\frac{\rho U_\infty^2}{L} \frac{\partial p^*}{\partial x^*} + \frac{\mu U_\infty}{L} \frac{\partial^2 u^*}{\partial x^{*2}} + \frac{\partial^2 u^*}{\partial y^{*2}} \\
\frac{\rho U_\infty^2}{L} \frac{\partial u^*}{\partial x^*} + \frac{\rho U_\infty^2}{L} \frac{\partial v^*}{\partial y^*} &= -\frac{\rho U_\infty^2}{L} \frac{\partial p^*}{\partial y^*} + \frac{\mu U_\infty}{L} \frac{\partial^2 v^*}{\partial x^{*2}} + \frac{\partial^2 v^*}{\partial y^{*2}} \\
\end{align*}
\]

\[
\begin{align*}
u^* &= v^* = 0 \quad \text{at } y^* = 0 \\
U_\infty u^* &\to U_\infty, v^* \to 0 \quad \text{as } x^* \to -\infty
\end{align*}
\]

Equations (1.29) through (1.32) can be rewritten as
\[
\begin{align*}
\frac{\partial u^*}{\partial x^*} + \frac{\partial v^*}{\partial y^*} &= 0 \\
u^* \frac{\partial u^*}{\partial x^*} + \nu^* \frac{\partial u^*}{\partial y^*} &= -\frac{\partial p^*}{\partial x^*} + \frac{1}{R} \left( \frac{\partial^2 u^*}{\partial x^{*2}} + \frac{\partial^2 u^*}{\partial y^{*2}} \right) \\
u^* \frac{\partial v^*}{\partial x^*} + \nu^* \frac{\partial v^*}{\partial y^*} &= -\frac{\partial p^*}{\partial y^*} + \frac{1}{R} \left( \frac{\partial^2 v^*}{\partial x^{*2}} + \frac{\partial^2 v^*}{\partial y^{*2}} \right) \\
u^* &= v^* = 0 \quad \text{at } y^* = 0 \\
u^* &\to 1 \quad v^* \to 0 \quad \text{as } x^* \to -\infty
\end{align*}
\]

where
\[
R = \frac{\rho U_\infty L}{\mu}
\]

is called the Reynolds number.

Equations (1.33) through (1.37) show that the problem depends only on the dimensionless parameter \( R \). For the case of small viscosity, namely \( \mu \) small compared with \( \rho U_\infty L \), \( R \) is large and its inverse can be used as a perturbation parameter to determine an approximate solution of the present problem. This process leads to the widely used boundary-layer equations of fluid mechanics. When the flow is slow, namely \( \rho U_\infty L \) is small compared with \( \mu \), \( R \) is small and it can be
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used as a perturbation parameter to construct an approximate solution of the present problem. This process leads to the Stokes-Oseen flow.

1.2 Expansions

In determining approximate solutions of algebraic, differential, and integral equations or evaluating integrals, we need to expand quantities in power series of a parameter or a variable. These power series expansions are usually obtained either as binomial expansions or Taylor series. These are explained next.

BINOMIAL THEOREM

Using straight multiplication, we have

\[(a + b)^2 = a^2 + 2ab + b^2\]

\[(a + b)^3 = a^3 + 3a^2b + 3ab^2 + b^3\]

\[(a + b)^4 = a^4 + 4a^3b + 6a^2b^2 + 4ab^3 + b^4\]

The process can be generalized for general \(n\) as

\[(a + b)^n = a^n + na^{n-1}b + \frac{n(n-1)}{2!} a^{n-2}b^2 + \frac{n(n-1)(n-2)}{3!} a^{n-3}b^3 + \cdots\]

(1.39a)

which can be rewritten as

\[\begin{align*}
(a + b)^n &= \sum_{m=0}^{n} \frac{n!}{m!(n-m)!} a^{n-m}b^m \\
&= \sum_{m=0}^{n} ^nC_m a^{n-m}b^m \quad \text{where } ^nC_m = \frac{n!}{m!(n-m)!} \\
\end{align*}\]

(1.39b)

or

\[(a + b)^n = \sum_{m=0}^{n} ^nC_m a^{n-m}b^m \quad \text{where } ^nC_m = \frac{n!}{m!(n-m)!} \quad (1.39c)\]

It turns out that (1.39a) terminates and hence it is valid when \(n\) is a positive integer. If it does not terminate, it is valid for any positive or negative number \(n\) provided that \(|b/a|\) is less than 1; otherwise, the series diverges because

\[
\lim_{m \to \infty} \frac{m^{\text{th term}}}{(m-1)^{\text{th term}}} = \lim_{m \to \infty} \frac{(m-1)!n(n-1)(n-2)\cdots(n-m+1)a^{n-m}b^m}{m!n(n-1)(n-2)\cdots(n-m+2)a^{n-m+1}b^{m-1}}
\]

\[
= \lim_{m \to \infty} \frac{(n-m+1)b}{ma} = -\frac{b}{a}
\]

For example,
EXPANSIONS

\[(a + b)^5 = \sum_{m=0}^{5} \frac{5!}{m!(5-m)!} a^{5-m} b^m\]
\[= a^5 + 5a^4b + 10a^3b^2 + 10a^2b^3 + 5ab^4 + b^5\]

\[(a + b)^6 = \sum_{m=0}^{6} \frac{6!}{m!(6-m)!} a^{6-m} b^m\]
\[= a^6 + 6a^5b + 15a^4b^2 + 20a^3b^3 + 15a^2b^4 + 6ab^5 + b^6\]

\[(a + b)^{1/2} = a^{1/2} + \frac{1}{2} a^{-1/2} b + \frac{\left(\frac{1}{2}\right)\left(-\frac{1}{2}\right)}{2!} a^{-3/2} b^2 + \frac{\left(\frac{1}{2}\right)\left(-\frac{1}{2}\right)\left(-\frac{3}{2}\right)}{3!} a^{-5/2} b^3 + \cdots\]
\[= a^{1/2} + \frac{1}{2} a^{-1/2} b - \frac{1}{8} a^{-3/2} b^2 + \frac{1}{16} a^{-5/2} b^3 + \cdots\]

\[(a + b)^{-1} = a^{-1} - a^{-2} b + \frac{(-1)(-2)}{2!} a^{-3} b^2 + \frac{(-1)(-2)(-3)}{3!} a^{-4} b^3 + \cdots\]
\[= a^{-1} - a^{-2} b + a^{-3} b^2 - a^{-4} b^3 + \cdots\]

We note that the first two series corresponding to \(n = 5\) and 6 terminate. The last two series corresponding to \(n = \frac{1}{2}\) and -1 do not terminate, and hence, they are valid only when \(|b| < |a|\).

TAYLOR SERIES EXPANSIONS

If a function \(f(x)\) is infinitely differentiable at \(x = x_0\), we express it in a power series of \((x - x_0)\) as

\[f(x) = a_0 + a_1(x - x_0) + a_2(x - x_0)^2 + a_3(x - x_0)^3 + \cdots\]

\[= \sum_{n=0}^{\infty} a_n (x - x_0)^n\]  \hspace{1cm} (1.41)

where the \(a_n\) are constants related to \(f\) and its derivatives at \(x = x_0\). Putting \(x = x_0\) in (1.41), we find that \(a_0 = f(x_0)\). Differentiating (1.41) with respect to \(x\), we have

\[f'(x) = a_1 + 2a_2(x - x_0) + 3a_3(x - x_0)^2 + 4a_4(x - x_0)^3 + \cdots\]  \hspace{1cm} (1.42)

which, upon putting \(x = x_0\), yields \(a_1 = f'(x_0)\). Differentiating (1.42) with respect to \(x\), we have

\[f''(x) = 2!a_2 + 3!a_3(x - x_0) + 4 \cdot 3a_4(x - x_0)^2 + \cdots\]  \hspace{1cm} (1.43)

which, upon putting \(x = x_0\), yields \(a_2 = (1/2!) f''(x_0)\). Differentiating (1.43) with respect to \(x\) gives

\[f'''(x) = 3!a_3 + 4!a_4(x - x_0) + \cdots\]  \hspace{1cm} (1.44)
which, upon putting \( x = x_0 \), yields \( a_3 = (1/3!) f'''(x_0) \). Continuing the process, we obtain

\[
a_n = \frac{1}{n!} f^{(n)}(x_0) \quad f^{(n)} = \frac{d^n f}{dx^n}
\]

(1.45)

and \( f^{(0)} = f(x_0) \). Therefore, (1.41) can be rewritten as

\[
f(x) = \sum_{n=0}^{\infty} \frac{f^{(n)}(x_0)}{n!} (x - x_0)^n
\]

(1.46)

which is called the Taylor series expansion of \( f(x) \) about \( x = x_0 \).

Since

\[
\frac{d}{dx} (\sin x) = \cos x \quad \text{and} \quad \frac{d}{dx} (\cos x) = -\sin x
\]

we have

\[
\sin x = x - \frac{x^3}{3!} + \frac{x^5}{5!} - \frac{x^7}{7!} + \cdots = \sum_{n=0}^{\infty} \frac{(-1)^n x^{2n+1}}{(2n+1)!}
\]

(1.47)

\[
\cos x = 1 - \frac{x^2}{2!} + \frac{x^4}{4!} - \frac{x^6}{6!} + \cdots = \sum_{n=0}^{\infty} \frac{(-1)^n x^{2n}}{(2n)!}
\]

(1.48)

Since

\[
\frac{d}{dx} (e^x) = e^x
\]

\[
e^x = 1 + \frac{x}{1!} + \frac{x^2}{2!} + \frac{x^3}{3!} + \cdots = \sum_{n=0}^{\infty} \frac{x^n}{n!}
\]

(1.49)

Since

\[
\frac{d}{dx} [\ln(1 + x)] = (1 + x)^{-1} \quad \frac{d}{dx} (1 + x)^{-n} = -n(1 + x)^{-n-1}
\]

\[
\ln(1 + x) = x - \frac{x^2}{2} + \frac{x^3}{3} - \frac{x^4}{4} + \cdots = \sum_{n=1}^{\infty} \frac{(-1)^{n+1} x^n}{n}
\]

(1.50)

The above Taylor series expansions are frequently used in subsequent chapters.

1.3. Gauge Functions

In this book, we are interested in the limit of functions such as \( f(\epsilon) \) as \( \epsilon \) tends to zero, denoted by \( \epsilon \to 0 \). This limit might depend on whether \( \epsilon \) tends to
zero from below, denoted by $e^0_{\downarrow}$, or from above, denoted by $e^0_{\uparrow}$. For example,
\[
\lim_{e^0_{\downarrow}} e^{-1/e} = 0 \quad \lim_{e^0_{\uparrow}} e^{-1/e} = \infty
\]

In what follows, we assume that the parameters have been normalized so that $e \geq 0$. If the limit of $f(e)$ exists (i.e., it does not have an essential singularity at $e = 0$ such as $\sin e^{-1}$), then there are three possibilities
\[
\begin{align*}
& f(e) \to 0 \\
& f(e) \to A \\
& f(e) \to \infty
\end{align*}
\]

as $e \to 0$, $0 < A < \infty$ \hspace{1cm} (1.51)

Most often, the above classification is not very useful because there are infinitely many functions that tend to zero as $e \to 0$. For example,
\[
\begin{align*}
\lim_{e \to 0} \sin e &= 0 \\
\lim_{e \to 0} (1 - \cos e) &= 0 \\
\lim_{e \to 0} (e - \sin e) &= 0 \\
\lim_{e \to 0} [\ln(1 + e)]^4 &= 0 \\
\lim_{e \to 0} e^{-1/e} &= 0
\end{align*}
\]

Also, there are infinitely many functions that tend to $\infty$ as $e \to 0$. For example,
\[
\begin{align*}
\lim_{e \to 0} \frac{1}{\sin e} &= \infty \\
\lim_{e \to 0} \frac{1}{1 - \frac{1}{2} e^2 \cos e} &= \infty \\
\lim_{e \to 0} e^{1/e} \ln \frac{1}{e} &= \infty
\end{align*}
\]

Therefore, to narrow down the above classification, we subdivide each class according to the rate at which they tend to zero or infinity. To accomplish this, we compare the rate at which these functions tend to zero and infinity with the rate at which known functions tend to zero and infinity. These comparison functions are called \textit{gauge functions}. The simplest and most useful of these are the powers of $e$
\[
1, e, e^2, e^3, \ldots
\]

and the inverse powers of $e$
\[
e^{-1}, e^{-2}, e^{-3}, e^{-4}, \ldots
\]

For small $e$, we know that
\[
1 > e > e^2 > e^3 > e^4 > \cdots
\]

and
\[
e^{-1} < e^{-2} < e^{-3} < e^{-4} < \cdots
\]
Let us determine the rate at which the preceding functions tend to zero or infinity. Using the Taylor series expansion (1.47), we have

\[ \sin \epsilon = \epsilon - \frac{\epsilon^3}{3!} + \frac{\epsilon^5}{5!} - \frac{\epsilon^7}{7!} + \cdots \]

so that \( \sin \epsilon \to 0 \) as \( \epsilon \to 0 \) because

\[ \lim_{\epsilon \to 0} \frac{\sin \epsilon}{\epsilon} = \lim_{\epsilon \to 0} \left(1 - \frac{\epsilon^2}{3!} + \frac{\epsilon^4}{4!} + \cdots \right) = 1 \]

Using (1.48), we have

\[ 1 - \cos \epsilon = \frac{\epsilon^2}{2!} - \frac{\epsilon^4}{4!} + \cdots \]

so that \( 1 - \cos \epsilon \to 0 \) as \( \epsilon^2 \to 0 \) because

\[ \lim_{\epsilon \to 0} \frac{1 - \cos \epsilon}{\epsilon^2} = \lim_{\epsilon \to 0} \left(1 - \frac{\epsilon^2}{2!} + \frac{\epsilon^4}{4!} + \cdots \right) = \frac{1}{2!} \]

Using (1.47), we have

\[ \epsilon - \sin \epsilon = \frac{\epsilon^3}{3!} - \frac{\epsilon^5}{5!} + \cdots \]

so that \( \epsilon - \sin \epsilon \to 0 \) as \( \epsilon^3 \to 0 \) because

\[ \lim_{\epsilon \to 0} \frac{\epsilon - \sin \epsilon}{\epsilon^3} = \lim_{\epsilon \to 0} \left(\frac{1}{3!} - \frac{\epsilon^2}{5!} + \cdots \right) = \frac{1}{3!} \]

Using (1.50), we have

\[ [\ln(1 + \epsilon)]^4 = \left(\epsilon - \frac{\epsilon^2}{2} + \frac{\epsilon^3}{3} + \cdots \right)^4 \]

so that \( [\ln(1 + \epsilon)]^4 \to 0 \) as \( \epsilon^4 \to 0 \) because

\[ \lim_{\epsilon \to 0} \frac{[\ln(1 + \epsilon)]^4}{\epsilon^4} = \lim_{\epsilon \to 0} \left(1 - \frac{\epsilon}{2} + \frac{\epsilon^2}{3} + \cdots \right)^4 = 1 \]

To determine the rate at which \( \exp \left(- \frac{1}{\epsilon}\right) \to 0 \) as \( \epsilon \to 0 \), we attempt to expand it in a Taylor series for small \( \epsilon \). To accomplish this, we need the derivatives of \( \epsilon \) at \( \epsilon = 0 \). But

\[ f'(\epsilon) = \frac{d(e^{-1/\epsilon})}{d\epsilon} = \frac{1}{\epsilon^2} e^{-1/\epsilon} \quad (1.54a) \]

which, at \( \epsilon = 0 \), gives 0 over 0. Hence, we need to use l'Hospital's rule to determine its limit as \( \epsilon \to 0 \). Thus,