EXTENSION OF **GAUSS'** METHOD FOR THE SOLUTION OF KEPLER'S **EQUATION**

by

THOMAS J. FILL **B.S.,** Pennsylvania State University **(1975)**

SUBMITTED IN PARTIAL **FULFILLMENT**

OF THE REQUIREMENTS FOR THE

DEGREE OF MASTER OF **SCIENCE**

at the

MASSACHUSETTS INSTITUTE OF **TECHNOLOGY**

May, **1976**

Signature of Author Department oflAeronaut'ics and Astronautics May 21, **1976**

Certified **by**

Thesis Supervisor

Accepted **by F**Chairman, Departmental Graduate Committee

EXTENSION OF **GAUSS'** METHOD FOR THE **SOLUTION** OF KEPLER'S **EQUATION**

by

THOMAS J. FILL

Submitted to the Department of Aeronautics and Astronautics on

May 21, **1976,** in partial fulfillment of the requirements for the degree of

Master of Science.

ABSTRACT

The solution of a transcendental equation known as "Kepler's equation," which relates position in orbit with time, requires an iterative procedure for solution. **A** method is developed based on one presented **by** Gauss in his Theoria Motus dealing with the problem of position determination for time since pericenter passage in cases where elliptic and hyperbolic orbits approached very near unity. The problem of interest here is the more general one of determining final position and velocity from given initial conditions and a specified time interval. Kepler's equation is transformed to an equation which is of the form of a cubic and which provides the nucleus of an efficient iteration algorithm. The final algorithm is a general form valid for any orbit of any eccentricity and requires no knowledge of the nature of the orbit for application. Universal formulae are developed relating final position and velocity to initial values in terms of variables defined in the transformation. Finally, the method is tested over a wide range of orbits to observe its performance and comparison is made with the proposed Kepler subroutine for the **NASA** Space Shuttle orbiter vehicle.

> Thesis Supervisor: Richard H. Battin Title: Lecturer in Aeronautics and Astronautics

ACKNOWLEDGMENTS

The author is, first of all, indebted to Karl F. Gauss, whose work in and contributions to the field of astronomy have provided the foundation for the work presented in this thesis.

Special thanks to Dr. Richard H. Battin, thesis advisor, who suggested the topic presented here and whose advice, comments and insight have contributed greatly to its solution. Also, a note of appreciation is extended to Stanley W. Shepperd whose assistance and suggestions were so important during the course of the work.

TABLE OF **CONTENTS**

 $\label{eq:2.1} \frac{1}{\sqrt{2}}\int_{0}^{\infty}\frac{1}{\sqrt{2\pi}}\left(\frac{1}{\sqrt{2\pi}}\right)^{2}d\mu_{\rm{eff}}\,.$

 $\label{eq:2.1} \frac{1}{\sqrt{2}}\int_{\mathbb{R}^3}\frac{1}{\sqrt{2}}\left(\frac{1}{\sqrt{2}}\right)^2\frac{1}{\sqrt{2}}\left(\frac{1}{\sqrt{2}}\right)^2\frac{1}{\sqrt{2}}\left(\frac{1}{\sqrt{2}}\right)^2\frac{1}{\sqrt{2}}\left(\frac{1}{\sqrt{2}}\right)^2.$

References **90**

SYMBOLS

CHAPTER **1**

INTRODUCTION

The determination of the position and velocity in two-body orbits leads to the solution of a transcendental equation commonly referred to as "Kepler's equation" which relates the dependence of position in orbit with time.

In classical analysis, the shape of these two-body orbits is described through the use of conics and correspcnding to each conic Kepler's equation has a different form. **A** useful quantity in classifying conics is a constant e called the eccentricity. For a circle e **= 0,** for an ellipse e is between **0** and **1,** e equals **1** for a parabola, and is greater than **1** for the hyperbola. Also obtainable from elementary considerations, is the general polar equation for the conic which can be stated as

$$
r = \frac{h^2/\mu}{1 + e \cos f} = \frac{p}{1 + e \cos f}
$$
 (1.1)

where h is the massless angular momentum, μ is the product of the universal gravitational constant and the sum of the masses of the two bodies, **p** is the semi-latus rectum or parameter, and **f,** called the true anomaly, is the angle between the radius vector and the direction of pericenter or point of closest approach of the two bodies.

For the parabola, Kepler's equation is simply

$$
6\sqrt{\frac{\mu}{\beta^3}} (t - \tau) = \tan^3(f/2) + 3 \tan(f/2)
$$
 (1.2)

where **T** is the time of pericenter passage. Although **Eq.** (1.2) is a special form of Kepler's equation it is more commonly known as "Barker's formula." **A** graph of the parabola is shown in Fig. **1.1.**

Figure **1.1** Parabola **q**

In the case of the ellipse, use is made of an angle, denoted **by E,** called the eccentric anomaly, which is based on a reference circle referred to as the "auxiliary circle," and whose geometrical significance can be seen in Fig. 1.2.

Figure 1.2 Ellipse

In terms of **E,** Kepler's equation may be expressed as

$$
M = E - e \sin E \tag{1.3}
$$

$$
M = \sqrt{\frac{\mu}{a^3}} (t - \tau) \tag{1.4}
$$

^Mis the mean anomaly and a is the semi-major axis.

For the hyperbola, instead of an angle an area is employed as the auxiliary variable and is also based on a reference geometric shape referred to as the "equilateral hyperbola" as shown in Fig. **1.3.**

Then the appropriate variable H is defined as

$$
Area CAQ = \frac{a^2}{2} H
$$

so that Kepler's equation may be written as

$$
\sqrt{\frac{\mu}{a^3}} (t - \tau) = e \sinh H - H
$$
 (1.5)

Kepler's equation is transcendental, so that for a given time it cannot be solved algebraically for the position parameter. However, there is one and only one solution and for an analytic solution an iterative process must be employed.

In his Theoria Motus Gauss addressed the problem of determining the true anomaly from the time, for elliptic and hyperbolic orbits which are nearly parabolic. In such cases the conventional methods of solution could not give the precision required. As Gauss expressed it, "The methods above treated, both for the determination of the true anomaly from the time and for the determination of the time from the true anomaly, do not admit of all the precision that might be required in those conic sections of which the eccentricity differs but little from unity, that is, in ellipses and hyperbolas which approach very near to the parabola; indeed, unavoidable errors, increasing as the orbit tends to resemble the parabola, may at length exceed all limits. Larger tables, constructed to more than seven figures would undoubtedly diminish this uncertainty, but they would not remove it, nor would they prevent its surpassing all limits as soon as the orbit approached too near the parabola. Moreover, the methods given above become in this case very troublesome, since a part of them require the use of indirect trials frequently repeated, of which the tediousness

is even greater if we work with the larger tables. It certainly, therefore, will not be superflous, to furnish a peculiar method **by** means of which the uncertainty in this case may be avoided, and sufficient precision may be obtained with the help of the common tables."

Gauss' method of solution is applicable to orbits of any eccentricity. The required iterative scheme is a "Picard" type iteration, i.e. successive substitution, there being no need for trials or tests which are so characteristic of many iterative schemes. Furthermore the method is applicable to all conic orbits, the advantage here being that the type of conic encountered need not be known in order to apply the formulae. Also, continuity is maintained during transition from one conic to another while at the same time being free from ambiguities or indeterminant forms. As will be seen, the speed of convergence is quite rapid. Gauss' method is briefly outlined here for the elliptic orbit.

Rewriting **Eq.** (1.4) as

$$
M = \sqrt{\frac{\mu (1 - e)^3}{q^3} (t - \tau)}
$$
 (1.6)

q being the pericenter distance, Gauss then chose to replace **E** arid sin **E by** the quantities

$$
P = E - \sin E
$$

\n
$$
Q = \frac{9}{10} E + \frac{1}{10} \sin E
$$
 (1.7)

With these, **Eq. (1.3)** takes the form

$$
(1 - e) P + (\frac{1}{10} + \frac{9}{10} e) Q = M
$$
 (1.8)

then as long as **E** is a quantity of first order,

$$
P = \frac{1}{6} E^3 - \frac{1}{120} E^5 + \frac{1}{5040} E^7 - \dots
$$

is a quantity of the third order, while

$$
Q = E - \frac{1}{60} E^3 + \frac{1}{1200} E^5 - \dots
$$

is a quantity of the first order. Then defining

$$
Y = \frac{6}{Q} \qquad B^2 = \frac{Q^3}{6 P} = \frac{Q^2}{Y}
$$

$$
Y = E2 - \frac{1}{30} E4 - \frac{1}{5040} E6 - \dots
$$

is a quantity of second order, while

$$
B = 1 + \frac{3}{2800} E^{4} + \frac{1}{16800} E^{6} + \ldots
$$

is a quantity which differs from unity **by** a quantity of the fourth order. Finally, **Eq. (1.8)** becomes

$$
B\left\{ (1 - e) Y^{1/2} + \frac{1}{60} (1 + 9 e) Y^{3/2} \right\} = M \qquad (1.9)
$$

Hence it is readily seen that the choice of $\frac{1}{10}$ and $\frac{9}{10}$ in the definition of **Q** was to obtain B, which multiplies the entire left side of **Eq. (1.9),** as nearly constant as possible. **Eq. (1.9)** is essentially an algebraic equation of third order.

It is easy to see that B can be considered as a function of Y. Furthermore, in the words of Gauss, "Now, although B may be finally known from Y **by** means of our auxiliary table, nevertheless it can be foreseen, owing to its differing so little from unity, that if the divisor B were wholly neglected from the beginning, Y would be affected with a slight error only. Therefore, we will first determine roughly Y, putting B = **1;** with the approximate value of Y, we will find B in our auxiliary table, with which we will repeat more exactly the same calculations; most frequently, precisely the same value of B that had been found from the approximate value of Y will correspond to the value of Y thus corrected, so that a second repetition of the operation would be superfluous, those cases excepted in which the value of **E** may have been very considerable." The tables referred to **by** Gauss were constructed to further simplify the iteration process **by** reducing the amount of computation even more. Here corresponding values of B are listed for values of Y, which are in increments of .004 from **0** to 1.2. In this manner the tables provide a simple means of obtaining values of **E** up to 640 **7'.**

To further illustrate the speed of convergence, consider an elliptic orbit in the x-y plane where pericenter is given **by**

$$
r_p
$$
 = (2 x 10⁷ m) i_x $\frac{v}{p}$ = (6.2 x 10³ m/sec) i_y

then for the Earth as the central force and an arbitrary value of **5** hours for the time since pericenter passage

$$
e = .928735
$$
 M = .0764383

For an initial guess of 1 for B, Y is obtained from **Eq. (1.9)** and is found to be

$$
Y = .59995
$$

Then from Gauss' table,

log B **=** .0001734 or B **= 1.000399**

With this value of B, **Eq. (1.9)** gives

^Y= **.59998**

and again from the table

B **= 1.000399**

Hence we have converged to the solution after just one correction at which point **E** may be calculated from

$$
E = B \sqrt{Y} (1 + \frac{1}{60} Y)
$$

with a value of **33.70390.**

It is the purpose of this study to extend Gauss' method of solution of Kepler's equation in standard form to the general problem of determining final position and velocity vectors for a specified time interval from given initial conditions at any point in the orbit while at the same time preserving all or most of the qualities which were inherent in Gauss' method. Finally, to see just how practical this solution process is, comparison is made with the algorithm proposed for the on-board computer in the **NASA** Space Shuttle orbital vehicle.

CHAPTER 2

EXTENDED FORM OF **GAUSS'** METHOD

The extension of Gauss' method to the solution of Kepler's equation for some arbitrary interval of time to obtain the final position and velocity is presented here for the ellipse and hyperbola, respectively. The parabola is shown to be the limiting form of both the elliptic and hyperbolic solutions as the eccentricity tends to unity. Furthermore, universal formulae are derived which permit calculation of final position and velocity using the initial position and velocity without knowledge of the type of orbit encountered. Finally a generalized procedure of the iterative process is presented, the details being left as the topic for a later chapter.

2.1 The Ellipse

Kepler's equation for a time interval $t = t_f - t_0$ corresponding to an eccentric anomaly difference $\Delta E = E - E_0$, may be written as⁽²⁾

$$
\sqrt{\frac{\mu}{a^3}} t = \Delta E + \frac{\sigma_0}{\sqrt{a}} (1 - \cos \Delta E) - (1 - \frac{r_0}{a}) \sin \Delta E
$$
 (2.1)

where the quantity σ_{0} is defined as

$$
\sigma_0 = \frac{r_0 \cdot \underline{v}_0}{\sqrt{\mu}}
$$

Since $q = a(1 - e)$ we have

$$
\sqrt{\frac{\mu}{q^3}} t = (1 - e)^{-3/2} \{ \Delta E + \frac{\sigma_0}{\sqrt{q}} (1 - e)^{1/2} (1 - \cos \Delta E) - (1 - \frac{r_0}{q} (1 - e)) \sin \Delta E \}
$$
 (2.2)

Defining the variables

- P = **AE -** sin **AE (2.3)**
- $Q = \alpha \Delta E + \beta \sin \Delta E$ (2.4)
- $R = 1 \cos \Delta E$ **(2.5)**

$$
\gamma = \frac{1}{2} \left(1 - \alpha \frac{r_0}{q} \left(1 - e \right) \right)
$$
 (2.6)

where α and β are constants to be specified such that α + β = 1 . Then **Eq.** (2.2) becomes

$$
\sqrt{\frac{\mu}{q^{3}}} t = (1 - e)^{-3/2} \left\{ \frac{r_{0}}{q} (1 - e) \left(\alpha \Delta E + \beta \sin \Delta E \right) + \frac{\sigma_{0}}{\sqrt{q}} (1 - e)^{1/2} (1 - \cos \Delta E) + (1 - \alpha \frac{r_{0}}{q} (1 - e)) \left(\Delta E - \sin \Delta E \right) \right\}
$$

or

$$
\sqrt{\frac{\mu}{q^3}} \quad t = (1 - e)^{-3/2} \left(\frac{r_0}{q} (1 - e) \, q + \frac{\sigma_0}{\sqrt{q}} (1 - e)^{1/2} \, R + 2 \, \gamma \, P \right) \tag{2.7}
$$

Then, as did Gauss, defining the variables

$$
Y = \frac{6}{Q}P
$$
 $B^2 = \frac{Q^3}{6P}$ (2.8)

and also a new variable

$$
C = \frac{\sigma_0}{2 \sqrt{r_0}} C_0 \qquad C_0 = \sqrt{\frac{2}{3} \sqrt{p} \sqrt{q}} = \frac{2 R}{\gamma B} \qquad (2.9)
$$

it is easily verified that

$$
Q = B \sqrt{Y}
$$
 $P = \frac{1}{6} B Y^{3/2}$

Hence, we have in **Eq. (2.7)**

$$
\sqrt{\frac{\mu}{q^3}} \quad t = (1 - e)^{-3/2} B \left\{ \frac{r_0}{q} (1 - e) \gamma^{1/2} + \sqrt{\frac{r_0}{q}} (1 - e)^{1/2} C \gamma + \frac{1}{3} \gamma \gamma^{3/2} \right\}
$$

 $\overline{1}$

or, defining the variable **D** such that

$$
Y = \frac{r_0}{q} (1 - e) 0^2
$$
 (2.10)

then

$$
\sqrt{\frac{\mu}{r_0^3}} \quad t = B \quad (D + C \quad D^2 + \frac{1}{3} \quad \gamma \quad D^3)
$$
 (2.11)

Hence, we have succeeded in reducing Kepler's equation to a form which resembles a cubic equation, the solution of which will be discussed in a subsequent chapter. Notice also that B, C_O and Y are all functions of AE, so that we may regard B and C_O as functions of Y. This fact, as will be seen, is of great importance in the solution of **Eq.** (2.11). Also, in the definition of Q we leave α and β unspecified for the moment instead of using Gauss' choice of $\frac{9}{10}$ and $\frac{1}{10}$. The functions B and C_O are both dependent on Y and their sensitivity to particular values of α and β will be discussed in a later chapter.

2.2 The Hyperbola

Kepler's equation for the hyperbola may be treated in an analagous fashion. To the time interval $t = t_f - t_0$ corresponds the difference **AH =** H **- H0** and the relevant form of Kepler's equation is

$$
\sqrt{\frac{\mu}{a^3}} \quad t = -\Delta H + \frac{\sigma_0}{\sqrt{a}} \left(\cosh \Delta H - 1 \right) + \left(1 + \frac{r_0}{a} \right) \sinh \Delta H
$$

since $q = a (e - 1)$

$$
\sqrt{\frac{\mu}{q^3}} t = (e - 1)^{-3/2} \left[-\Delta H + \frac{\sigma_0}{\sqrt{q}} (e - 1)^{1/2} \left(\cosh \Delta H - 1 \right) + \left(1 + \frac{r_0}{q} (e - 1) \right) \sinh \Delta H \right]
$$
 (2.12)

Again defining the variables

$$
P = \sinh \Delta H - \Delta H \tag{2.13}
$$

$$
Q = \alpha \Delta H + \beta \sinh \Delta H \qquad (2.14)
$$

$$
R = \cosh \Delta H - 1 \tag{2.15}
$$

$$
\gamma = \frac{1}{2} \left(1 + \alpha \frac{r_0^*}{q} (e - 1) \right)
$$
 (2.16)

gives

$$
\sqrt{\frac{\mu}{q^3}} \quad t = (e - 1)^{-3/2} \left\{ \frac{r_0}{q} (e - 1) \quad Q + \frac{\sigma_0}{\sqrt{q}} (e - 1)^{1/2} \quad R + 2 \quad \gamma \quad P \right\}
$$

or, using Eqs. **(2.8)** and **(2.9)**

$$
\sqrt{\frac{\mu}{q^3}} \quad t = (e - 1)^{-3/2} B \left\{ \frac{r_0}{q} (e - 1) \gamma^{1/2} + \sqrt{\frac{r_0}{q}} (e - 1)^{1/2} C \gamma + \frac{1}{3} \gamma \gamma^{3/2} \right\}
$$
 (2.17)

By letting

$$
Y = \frac{r_0}{q} (e - 1) D^2
$$
 (2.18)

Eq. (2.17) becomes

$$
\frac{\mu}{r_0^3} \quad t = B \quad (D + C \quad D^2 + \frac{1}{3} \gamma \quad D^3)
$$
 (2.19)

The universality of this method becomes apparent here since, for different type orbits, the same resulting equation is obtained. Furthermore, from Eqs. **(2.18)** and (2.10), it is readily seen that Y may be used to classify conics in a similar fashion to the eccentricity. **If Eq.** (2.10) is accepted as the definition of **D** then for the ellipse Y is greater than **0,** Y is equal to **0** for the parabola, and is less than **0** for the hyperbola. Hence we now possess a general form for the solution of Kepler's equation for hyperbolic and elliptic orbits and will show subsequently that this general form is indeed also valid for parabolic orbits.

2.3 The Parabola

The parabola can be shown to be the limiting form of both the ellipse and hyperbola as e tends to unity. Here Kepler's equation is

$$
2 \sqrt{\frac{\mu}{p^3}} t = \{ \tan(f/2) - \tan(f_0/2) \} + \frac{1}{3} \{ \tan^3(f/2) - \tan^3(f_0/2) \}
$$
\n(2.20)

As e approaches unity from either the hyperbola or ellipse a approaches infinity and from the definition of **D,** Y tends to **0.** It is easily verified that both B and C_O approach unity. Hence we have from either **Eq.** (2.11) or **Eq. (2.19)**

$$
\sqrt{\frac{\mu}{r_0^3}} \quad t = D + \frac{\sigma_0}{2 \sqrt{r_0}} \quad D^2 + \frac{1}{6} \quad D^3 \tag{2.21}
$$

Using the fact that for parabolic motion

$$
\sin f_0 = \frac{\sqrt{p} \,^0 0}{r_0} \qquad \cos f_0 = \frac{p}{r_0} - 1 \qquad (2.22)
$$

and that the root of **Eq.** (2.21) is

$$
D = \frac{\sqrt{2} \sin \left(-\frac{f}{2}\right)}{\cos(f/2)}
$$
 (2.23)

then substitution of **Eq.** (2.22) and **Eq. (2.23)** does indeed lead to **Eq. (2.20).** As a case in point; at pericenter $f_0 = 0$; hence

D =
$$
\sqrt{2} \tan(f/2)
$$
 $\sigma_0 = 0$ $\frac{p}{r_0} = 2$

Therefore **Eq.** (2.21) becomes

$$
\sqrt{\frac{\mu}{r_0^3}} (t - \tau) = \sqrt{2} \tan(f/2) + \frac{\sqrt{2}}{3} \tan^3(f/2)
$$

$$
2 \sqrt{\frac{\mu}{p^3}} (t - \tau) = \tan(f/2) + \frac{1}{3} \tan^3(f/2)
$$

which is Barker's formula. Therefore as γ approaches 0 (e \rightarrow 1), the hyperbolic and elliptic forms reduce to the parabolic form.

2.4 Final Position and Velocity Vectors

Determination of the final position and velocity from given initial position and velocity vectors and a time interval may be done through the use of universal formulae expressed in terms of the variables Y, B, **D,** and C_{Ω} .

In general, the final position and velocity vectors may be written in terms of the Lagrange F and **G** functions as

$$
\underline{v} = F_{\underline{r}_0} + G_{\underline{v}_0}
$$
\n
$$
\underline{v} = F_{\underline{r}_0} + G_{\underline{v}_0}
$$
\n(2.24)

For elliptic orbits F and **G** are found to be

$$
F = 1 - \frac{a}{r_0} (1 - \cos \Delta E)
$$

$$
F_t = - \frac{\sqrt{\mu a}}{r r_0} \sin \Delta E
$$

or

$$
G = t - \sqrt{\frac{a^3}{\mu}} (\Delta E - \sin \Delta E)
$$

$$
G_t = 1 - \frac{a}{r} (1 - \cos \Delta E)
$$

making use of Eq. (2.9) and Eq. (2.5)

$$
1 - \cos \Delta E = \frac{1}{2} Y B C_0
$$

Also from Eq. (2.3) and Eq. (2.4)

$$
\sin E = Q - \alpha P = B \sqrt{Y} - \frac{\alpha}{6} B Y^{3/2}
$$

$$
= B \sqrt{Y} (1 - \frac{\alpha}{6} Y)
$$

but since $Y = \frac{r_0}{a} D^2$ sin ΔE = B D $\sqrt{\frac{r_0}{a}}$ (1 - $\frac{\alpha}{6}$ Y)

and

$$
P = \Delta E - \sin \Delta E = \frac{1}{6} B Y^{3/2} = \frac{D^3 B}{6} \sqrt{\frac{r_0^3}{a^3}}
$$

Therefore, for an ellipse

$$
F = 1 - \frac{1}{2} D^{2} B C_{0}
$$
\n
$$
G = t - \frac{B D^{3}}{6} \sqrt{\frac{r_{0}^{3}}{\mu}}
$$
\n
$$
F_{t} = - \frac{B D}{r} \sqrt{\frac{\mu}{r_{0}}} (1 - \frac{\alpha}{\delta} Y) \qquad G_{t} = 1 - \frac{r_{0}}{r} (\frac{1}{2} D^{2} B C_{0})
$$
\n(2.25)

For hyperbolic orbits F and G are

$$
F = 1 - \frac{a}{r_0} (\cosh \Delta H - 1)
$$

\n
$$
F_t = - \frac{\sqrt{\mu a}}{r r_0} \sinh \Delta H
$$

\n
$$
G = t - \sqrt{\frac{a^3}{\mu}} (\sinh \Delta H - \Delta H)
$$

\n
$$
G_t = 1 - \frac{a}{r} (\cosh \Delta H - 1)
$$

making use of Eq. (2.9) and Eq. (2.15)

$$
\cosh \Delta H - 1 = \frac{1}{2} Y B C_0 = \frac{r_0}{a} (\frac{1}{2} D^2 B C_0)
$$

Also from Eq. (2.13) and Eq. (2.14)

$$
\sinh \Delta H = Q + \alpha P = B \sqrt{Y} (1 + \frac{\alpha}{6} Y)
$$

$$
= B D \sqrt{\frac{r_0}{a} (1 + \frac{\alpha}{6} Y)}
$$

and

$$
P = \sinh \Delta H - \Delta H = \frac{1}{6} B Y^{3/2} = \frac{B D^3}{6} \sqrt{\frac{r_0^3}{a^3}}
$$

Therefore, for the hyperbola

$$
F = 1 - \frac{1}{2} D^{2} B C_{0}
$$

\n
$$
G = t - \frac{B D^{3}}{6} \sqrt{\frac{r_{0}^{3}}{\mu}}
$$

\n
$$
F_{t} = - \frac{B D}{r} \sqrt{\frac{\mu}{r_{0}}} (1 + \frac{\alpha}{6} Y) \qquad G_{t} = 1 - \frac{r_{0}}{r} (\frac{1}{2} D^{2} B C_{0})
$$

\n(2.26)

The only difference between **Eq. (2.25)** and **Eq. (2.26)** is **in** the expression for F_t. This sign difference may be avoided by the convention that Y<O for the hyperbola, Y=O for the parabola, and Y>O for the ellipse. Thus, the final position and velocity vectors may be determined through the use of **Eq.** (2.24) and **Eq. (2.25),** where Y determines the conic. Notice that nowhere in **Eq. (2.25)** is there any need to know the nature of the conic.

One further point to be made here is that another means of calculating **y** and also Y will be needed if the new form of Kepler's equation is to be valid for rectilinear orbits as well. This may be accomplished using the fact that

$$
q = \frac{p}{(1+e)} \qquad e^2 = \left(\frac{p}{r_0} - 1\right)^2 + \frac{p \sigma_0^2}{r_0} \qquad (2.27)
$$

Recall that

$$
p = \frac{h^2}{\mu} = \frac{(r_0 \times v_0) \cdot (r_0 \times v_0)}{\mu}
$$

Using the identity

 $(\underline{A} \times \underline{B}) \cdot (\underline{C} \times \underline{D}) = (\underline{A} \cdot \underline{B}) (\underline{B} \cdot \underline{D}) - (\underline{A} \cdot \underline{D}) (\underline{B} \cdot \underline{C})$ 20

we have \sim 2

$$
p = \frac{(r_0 v_0)^2}{\mu} - \sigma_0^2
$$

Now

$$
(1 - e)^2 = \frac{2 p}{r_0} - \frac{p (p + \frac{\sigma_0^2}{2})}{r_0^2}
$$

and
\n
$$
\frac{(1-e)}{q} = \frac{(1-e^2)}{p} = \frac{2}{r_0} - \frac{p+{}^{0}0}{r_0}^{2} = \frac{2}{r_0} - \frac{{v_0}^{2}}{r}
$$

Hence **Eq. (2.6)** and **Eq. (2.16)** become

$$
\gamma = \frac{1}{2} \left(1 - \alpha \left(2 - \frac{r_0 v_0^2}{\mu} \right) \right) \tag{2.28}
$$

$$
\text{while} \\
$$

while
\n
$$
Y = (2 - \frac{r_0 v_0^2}{\mu}) p^2
$$
\n(2.29)

2.5 General Summary of the Iterative Method

The purpose for the following general discussion of the method of solution is to summarize the work presented in this chapter and also to give purpose to the material presented in the following chapters.

Given an initial position, velocity and a time interval one may find the final position and velocity **by** performing the following sequence of steps:

1. Calculate
$$
r_0
$$
, v_0 , σ_0 , and γ .

- 2. For an initial guess set B and **C0** equal to unity.
- **3.** Substitute these values into the cubic

$$
\sqrt{\frac{\mu}{r_0^3}} \quad t = B \quad (D + C \quad D^2 + \frac{1}{3} \quad \gamma \quad D^3)
$$

and solve for **D.**

- 4. With this value of **D,** calculate Y using **Eq. (2.29).**
- **5.** From B and **C0** expressed as functions of Y, calculate new values for B and C_0 .
- **6.** Test for convergence by using the present values of B, C₀ and **^D**in the cubic and recomputing t. Check the error between this value and the given time interval. If the error is larger than some specified value, return to step **3.** If the error is smaller, then calculate the final position and velocity using **Eq.** (2.24) and **Eq. (2.25).**

As in Gauss' solution there must be a limit to the magnitude of **AE** or **AH** in the ellipse and hyperbola and hence on the time interval or else the number of iterations does not remain small. Therefore it is necessary to account for larger intervals of time. This aspect along with determining B and **C0** as functions of Y, and determining the root to the cubic, which in this case may not always possess one real root as was the case with Gauss' solution since **C** and **y** are not constant, will be discussed in the subsequent chapters.

CHAPTER **3**

SERIES EXPANSIONS

Probably the most convenient method of obtaining B and **C0** as functions of Y is **by** power series representation. The means of obtaining these series is discussed along with the selection of the constants α and **. A** procedure to economize the series is then presented and a potential means **by** which further reduction in the amount of computation is obtained.

3.1 Selection of the Constants α **and** β **.**

Consideration in the selection of the constants α and β , in this instance, must be given not only to B but also to c_0 . For elliptic motion, if B and **C0** are expanded as powers of **AE** in terms of cx and 3 the resulting first terms are

$$
B = 1 + \frac{1}{4} \left(\frac{1}{10} - \beta \right) (\Delta E)^{2} + \dots
$$

$$
C_0 = 1 + \frac{1}{12} (\beta - \frac{7}{10}) (\Delta E)^{2} + \dots
$$

Clearly selection of $\alpha = \frac{9}{10}$, $\beta = \frac{1}{10}$ makes B a value which differs from $\frac{3}{2}$ $\frac{2}{3}$ unity by a quantity of fourth order in Δ E. $\,$ But if α = $\overline{10}$, β = $\overline{10}$ are selected then C_O will differ from unity by a quantity of fourth order in ΔE . On the other hand if $\alpha = \frac{3}{4}$, $\beta = \frac{1}{4}$ are chosen, both B and C_n differ from unity **by** the same quantity of second order in **AE.**

It is possible to observe both B and C_0 as functions of Y for these selections **by** evaluating B, **C0** and Y for different values of **AE (AH)** and then plotting B vs. Y and C_0 vs. Y. The resulting curves are shown in Figs. **3.1** through 3.4. It is easy to see, for both the hyperbola and ellipse, that to obtain either B or C_0 as nearly flat as possible the other varies significantly. Also notice that for $\alpha = \frac{3}{10}$, C_0 remains much flatter than B does for $\alpha = \frac{9}{10}$. These two choices of α will be compared in tests to see which gives better results.

The following procedure to obtain B and C₀ as series in powers of Y will be explained with $\alpha = \frac{9}{10}$ but results for both $\alpha = \frac{3}{10}$ and $\frac{1}{4}$ are also presented.

3.2 Methods of Derivation

Two methods of obtaining B and C₀ as a power series in Y were tried. The first uses the technique of series reversion. The procedure is to obtain Y as a series in **AE (AH)** and revert the series such that now **AE (AH)** is a series in powers of Y. **A** convenient expression can be obtained relating B to **AE (AH)** and Y in which case substitution of the reverted series gives the final result. The series for C_O can then be obtained **by** making use of the reverted series and the B power series. This method may prove more applicable if the computation is to be done **by**

Figure 3.2 B vs. Y, hyperbola

hand, but with the aid of a computer the second method proved to be much quicker and easier. For this reason the reversion method is explained in detail in Appendix **A.**

The second method -vas performed on MACSYMA (Project MAC's SYmbolic MAnipulation system), a large computer programming system used for performing symbolic as well as numerical mathematical manipulations and developed **by** the Mathlab Group Project **MAC,** of MIT. The procedure was to obtain the power series in **AE,** in the case of the ellipse, about **AE 0,** for Y, B and C₀ from their definitions using the algebra for power series. Hence,

$$
Y = \Delta E^2 - \frac{\Delta E^2}{30} - \frac{\Delta E^4}{5040} + \frac{\Delta E^6}{36000} - \frac{79 \Delta E^8}{498960000} - \frac{6469 \Delta E^{10}}{340540200000} + \cdots
$$

$$
E_0^2 = 1 + \frac{3 \Delta E^4}{2800} - \frac{\Delta E^6}{84000} + \frac{71 \Delta E^8}{258720000} - \frac{527 \Delta E^{10}}{100900800000} + \dots
$$

$$
C_0 = 1 - \frac{\Delta E^2}{20} + \frac{\Delta E^4}{4200} + \frac{11 \Delta E^6}{504000} - \frac{43 \Delta E^8}{194040000} - \frac{8747 \Delta E^{10}}{908107200000} + \dots
$$

Now, if a series for B is assumed in the form

$$
B = \sum_{0}^{\infty} A_{n} Y^{n}
$$
 (3.1)

substitution of the power series for Y then gives B as a power series in **AE** which must be equal to the one above. Therefore, the coefficients of like powers must be equivalent. The end result is a system of simultaneous

equations in the coefficients which may be solved for the coefficients of Eq. (3.1) . The same procedure can be used to obtain C_0 as a power series in Y. Values for the coefficients of the B and C₀ power series in Y for the cases when $\alpha = \frac{9}{10}$, $\frac{3}{10}$ and $\frac{3}{4}$ are listed in tables 3.1 through 3.3.

Expansions for the hyperbola were found to differ from those for the ellipse only in the sign of the coefficients of the odd powers of Y. Therefore if the convention that Y<0 for the hyperbola is invoked, then the coefficients for the B and C₀ power series in Y for the hyperbola become identical to those listed in tables **3.1, 3.2,** and **3.3.**

In the remainder of this report this convention will be assumed unless otherwise stated.

3.3 Limits on Y

Before presenting a procedure to determine the number of terms needed in the series presented above, some discussion is needed on the range of Y. Clearly before determining the number of terms needed in the series, it is important to know the magnitude of Y.

Figs. **3.5** and **3.6** illustrate graphs of Y versus **AE** and **AH,** respectively. Clearly, for the ellipse, Y is not single-valued;therefore **AE** maximum must be limited such that no ambiguities arise in the series. Furthermore, from consideration of the figures presented earlier, Y must be limited in such a way that the maximum difference of B or **C0** from unity is fairly small.

	Table 3	

Coefficients of B and C_O series for $\alpha = \frac{9}{10}$

Table 3.3

Coefficients of B and C₀ series for $\alpha = \frac{3}{10}$

For comparison purposes, two ranges of Y were selected on the basis of the figures, to illustrate both the number of terms needed in the B and C₀ series and the number of iterations required in the basic algorithm for certain test cases to be presented subsequently. These two ranges are $-2 \le Y \le 2$ and $-1 \le Y \le 1$.

9 In the case of an ellipse; for **a** = , a Y of 2 corresponds to a **AE** of approximately 84' while for a Y of **1, AE** is approximately **580.** 3 These values are roughly the same for **a** =

3.4 Series Economization

A simple procedure to predict the effect of a term in a power series is through the use of Chebyshev polynomials (1) If a function is expanded in a series of Chebyshev polynomials,then the contribution of the n $^\mathrm{th}$ Chebyshev (T_n) term will never be greater than the magnitude of its coefficient since the magnitude of any Chebyshev polynomial does not exceed unity. Specifically, for a given function

$$
f(x) = \sum_{0}^{m} a_{n} x^{n} \quad \text{where } -1 \leq x \leq 1
$$
 (3.2)

let it be required to find a function

$$
g(x) = \sum_{n=0}^{K} b_n x^n
$$
 (3.3)

with K as small as possible, such that

$$
|f(x) - g(x)| < \rho
$$

where **p** is some maximum permissable error. This may be done **by** expressing f(x) in terms of Chebyshev polynomials

$$
f(x) = \sum_{0}^{m} c_n T_n(x)
$$

and since $|T_n(x)| \le 1$, for $-1 \le x \le 1$, then

$$
g(x) = \sum_{n=0}^{K} c_n T_n(x) \qquad (3.4)
$$

within the desired accuracy provided that

$$
\sum_{n=K+1}^{m} |c_n| < \rho
$$

After **Eq.** (3.4) is obtained, then **by** expressing the Chebyshev polynomials in terms of the powers of x, we obtain **Eq. (3.3).**

Economizations were performed to obtain **10** decimal place accuracy, for reasons which will become apparent later, for the two ranges of ^Y discussed earlier. For Y between ± 2 , when $\alpha = \frac{9}{10}$, the B series required **3** terms and the C₀ series 9 terms. When $\alpha = \frac{3}{10}$ 11 terms and the C_O series 7 terms. For Y between ± 1 ; when $\alpha = \frac{9}{10}$, the B series required 6 terms and the C₀ series 7 terms. When $\alpha = \frac{3}{10}$, the B series required 8 terms and the C_O series 5 terms. The new

coefficients for the four cases noted here are listed in Tables 3.4 through **3.7.** Out of curiosity, the B series was economized for $|Y| < \pm \frac{1}{2}$ and α = $\frac{10}{10}$. In this case only 4 terms were required.

3.5 Series About an Arbitrary Point ^Y

If a greater accuracy is required for the values of B and **C0** and the resulting number of terms required in the above series, expanded about ^Y**= 0,** is too large or still further reduction in the amount of computation in the B and **C0** series is desired, this may be had **by** expanding B and **C0** about some arbitrary point Y_0 . In this way the range of Y may be divided into smaller ranges. For example, B and C_O might be expanded about the points $Y_0=0$, ± 1 , ± 2 . Then, if Y varies from -2 to ± 2 , this range may be subdivided, where the different expansions would cover the ranges

$$
-\frac{5}{2} \le (Y + 2) \le -\frac{3}{2}
$$

$$
-\frac{3}{2} \le (Y + 1) \le -\frac{1}{2}
$$

$$
-\frac{1}{2} \le (Y - 0) \le \frac{1}{2}
$$

$$
\frac{1}{2} \le (Y - 1) \le \frac{3}{2}
$$

$$
\frac{3}{2} \le (Y - 2) \le \frac{5}{2}
$$

Clearly the symmetry of the expansions is lost here in that different series are needed for negative values of Y, (i.e. for hyperbolic motion,

Economized coefficients of B and C_O series for $|Y| < 2$ and $\alpha = \frac{9}{10}$

I _ _ _ _ _ _ _ __ _ _ _ _ _ _ _ _ _ **_L I__** _ _ _ _ _ _ _ _ _ _ __ _ _ _ _ _ _ _ _ _ _ _

Table 3.5

Economized coefficients of B and C₀ series for $|Y| < 2$ and $\alpha = \frac{3}{10}$

Table 3.6

Economized coefficients of B and C_O series for $|Y| < 1$ and $\alpha = \frac{9}{10}$

Table 3.7

Economized coefficients of B and C_O series for $|Y| < 1$ and $\alpha = \frac{3}{10}$

 \bar{z}

B and C₀ must be expanded about $|Y_0|$ using the hyperbolic definitions and then the sign of the coefficients for the odd powers of Y changed.).

The methods presented earlier are no longer of practical use in this case. The method of obtaining these coefficients and their values are presented in Appendix B.

CHAPTER 4

ROOTS OF THE **CUBIC**

The cubic equation in Gauss' solution is monotonic in that the coefficients are such that only one real root exists. This is not the case here because of the existence of **C** and **y** which are functions of initial position and velocity. It is now possible that three real roots may exist and clearly only one is valid. Two methods are presented here for obtaining the one correct solution of the cubic.

 ~ 100

4.1 Method **A**

Equation (2.11) is

$$
\frac{1}{3} \gamma \, 0^3 \div C \, 0^2 \div D = \frac{t}{B} \sqrt{\frac{v}{r_0^3}}
$$

which may be written in the form

$$
\frac{1}{3} (\gamma D)^{3} + C (\gamma D)^{2} + \gamma (\gamma D) = \frac{\gamma^{2} \pm}{B} \sqrt{\frac{\mu}{r_{0}^{3}}}
$$

Then **by** defining

 γ D = x - C (4.1) $F = \frac{t}{B} \sqrt{\frac{\mu}{r_0^3}}$ (4.2) \mathcal{A}^{\prime}

we have

$$
x^3 - 3 \varepsilon x = 2 b \tag{4.3}
$$

where

$$
\varepsilon = C^2 - \gamma \tag{4.4}
$$

$$
2 b = 3 \gamma^2 F + C^3 - 3 \epsilon C \qquad (4.5)
$$

From elementry algebra, the criterion for one real root in **Eq.** (4.3) is

$$
b^2 - \varepsilon^3 > 0
$$

In this case it is easily shown that

$$
x = (b + \sqrt{b^2 - \epsilon^3})^{1/3} + (b - \sqrt{b^2 - \epsilon^3})^{1/3}
$$
 (4.6)

is the real root. This may be written in the form

$$
x = \frac{2 b m_0}{(m_0 - \varepsilon)^2 + m_0 \varepsilon} \tag{4.7}
$$

where

$$
m_0 = (|b| + \sqrt{b^2 - \epsilon^3})^{1/3}
$$
 (4.8)

The absolute value of **b** results from taking the positive square root of **b2** Eqs. (4.7) and (4.8) are more desirable than **Eq.** (4.6) in computing x due to the fact that only one and not two cube roots is required and

that m_0 is obtained by adding two positive numbers in contrast to the subtraction required in **Eq.** (4.6).

When $b^2 - \epsilon^3 \leq 0$, three real roots exist, two of which are equal when the equality holds. Here the three roots are obtained **by** calculating

$$
3 \theta = \arccos(\frac{\sqrt{b^2}}{\sqrt{\epsilon^3}})
$$
 (4.9)

so that

$$
x_1 = \pm 2 \sqrt{\epsilon} \cos \theta
$$

\n
$$
x_2 = \pm 2 \sqrt{\epsilon} \cos(\theta + \frac{2\pi}{3})
$$

\n
$$
x_3 = \pm 2 \sqrt{\epsilon} \cos(\theta + \frac{4\pi}{3})
$$

\n
$$
x_4 = \pm 2 \sqrt{\epsilon} \cos(\theta + \frac{4\pi}{3})
$$

are the three real roots.

Since ϵ must be greater than zero and b lies between $\pm \sqrt{\epsilon^3}$ for for these three real roots to exist, selection of the proper root can be deduced, and easily verified, **by** plotting the three roots as functions of **b** and c as is shown in Fig. 4.1. Using **Eq.** (4.1), if **C** is positive then this equation is equivalent to translation of the **y D** axis to the right of the origin whereas if **C** is negative, translation is to the left. Rewriting **Eq.** (4.4) in the form

$$
|C| = \sqrt{1 + \frac{\gamma}{\epsilon}} \sqrt{\epsilon}
$$

Figure 4.1 Graph of three real roots, b vs. x

it can be seen that if **y** is negative then absolute **C** will always be less than $\sqrt{\epsilon}$. Hence transiation of the axis will be such that x_3 is the correct choice of the root. If γ is positive then selection of the roots x_1 or x_2 will be based on the signs of **b** and **C**.

A simple criterion for selection of the roots can be obtained **by** letting **6** run from **0** to **R,** starting at pt. **A** in Fig. 4.1 and moving along the curve_instead of oscillating between 0⁰ and 30⁰ as is seen in **Eq.** (4.9). Therefore the roots may be computed using the following criteria: for $b^2 \leq \epsilon^3$,

$$
3 \theta = \arccos(\frac{b}{\sqrt{\epsilon^3}})
$$

then for $\gamma < 0$, $120^{\circ} - \theta \rightarrow \theta$

for $\gamma > 0$ and $C < 0$, $\theta + 120^{\circ} \rightarrow \theta$

otherwise, θ remains as calculated (always in the first quadrant) and

$$
x = 2 \sqrt{\epsilon} \cos \theta
$$

From a continuity standpoint, since small changes in t and hence in F and **b ,** must correspond to small changes in the root **D,** then depending on the values of γ and C there are only certain ranges of t for which this solution method is valid. For example, consider **C** positive with **b** and e being such that three real roots exist. **If b** is decreased **(by** changing t), some value will be reached such that $b = -\sqrt{\epsilon^3}$ at which point the value of $x = \sqrt{\epsilon}$ is the solution (pt. B in Fig. 4.1). Any further

decrease in **b** results in the existence of one real root and Fig. 4.1 illustrates that a jump from pt. B to point **D** exists. Hence a discontinuity in x exists at the point where $b = -\sqrt{\epsilon^3}$ when $c > 0$. Since this is not physically possible, t must be constrained, in this case, such that $b \geq -\sqrt{\epsilon^3}$. (we know that requiring that **b** be less than $-\sqrt{\epsilon^3}$ is incorrect from the simple fact that three real roots exist when $t = 0$ and hence continuity must be maintained from this point). In a similar fashion it is seen that when C is negative, t must be constrained such that $b \leq \sqrt{\epsilon^3}$. Hence several tests must be incorporated in the algorithm to maintain these requirements. This is, of course, not desirable.

1 Depending on the selection of α , in particular if α > $\frac{1}{2}$, it is possible for **y** to be equal to or less than zero. **If y** equals zero then **Eq.** (2.11) reduces to a quadratic with solution

$$
D = \frac{-1 \pm \sqrt{1 + 4 \, \text{CF}}}{2 \, \text{C}} \tag{4.10}
$$

Clearly the (+) sign is the correct choice from the simple fact that **^D** must equal zero when F is zero.

When **y** is in the vicinity of zero, computation of the root using **Eq.** (4.1) leads to an indeterminant form. In this case, an asymptotic expansion of the form

$$
D = D_0 + D_1 \left(\frac{Y}{3}\right) + D_2 \left(\frac{Y}{3}\right)^2 + \dots \tag{4.11}
$$

is appropriate.⁽⁴⁾ Substitution into Eq. (2.11) and equating powers of γ gives a set of expressions for the coefficients of **Eq.** (4.11) which may be written in the form

$$
A = \sqrt{1 + 4 c F}
$$
\n
$$
D_0 = \frac{2 F}{1 + A}
$$
\n
$$
A_1 = C D_1 + 3 D_0^2
$$
\n
$$
D_1 = -\frac{D_0^3}{A}
$$
\n
$$
A_2 = C D_2 + 3 D_0 D_2
$$
\n
$$
D_2 = -\frac{D_1 A_1}{A}
$$
\n
$$
A_3 = C D_3 + 3 D_0 D_2 + D_1^2
$$
\n
$$
D_3 = -\frac{D_2 A_1 + D_1 A_2}{A}
$$
\n
$$
D_4 = -\frac{D_3 A_1 + D_2 A_2 + D_1 A_3}{A}
$$

A more convenient form which eliminates all of the indirect computation may be arrived at after some inspection of the above set of coefficients. Rewriting **Eq.** (4.11) in the form

$$
D = D_0 \left(1 + \sum_{1}^{\infty} W_k \left(\frac{D_0^2 \gamma}{A^2} \right)^k \right) = D_0 H \qquad (4.12)
$$

and expressing **C** and F in the forms

$$
F = \frac{(1 + A) D_0}{2} \qquad C = \frac{A^2 - 1}{4F} = \frac{A - 1}{2 D_0}
$$

we substitute into **Eq.** (2.11) to obtain

$$
\frac{1}{3} A^2 \left(\frac{D_0^2 \gamma}{A^2}\right) H^3 + \frac{(A-1)}{2} H^2 + H = \frac{A+1}{2}
$$

The series definition of H then provides expressions for the coefficients W_k in terms of A. Table 4.1 provides a list of several of these coefficients.

The efficiency and practical use of the asymptotic series both in the amount of computation required and in the accuracy leaves much to be desired. The coefficients do alternate in sign so that the truncation error is smaller in absolute value than the first term omitted. On the other hand, no precise criterion is apparent to decide when the asymptotic series must be used or when obtaining the root using **Eq.** (4.1) is no longer valid. Also, when **C** is negative (or t is negative in which case F is also negative), then for **y** equal zero, t must be such that the radical in **Eq** (4.10) is nonnegative. Furthermore, when **y** is small, t must be such that **A** is not in the vicinity of zero, since **A** is present in the denominator of the asymptotic series.

4.2 Method B

The second method involves a change in variable from **D** to x according to

$$
D = \frac{3 F}{1 + x}
$$
 (4.13)

T

Substituting into **Eq. (2.11)** gives

$$
x^3 - 3 \varepsilon x = 2 b \tag{4.14}
$$

where now

$$
\varepsilon = 1 + 3 \quad \text{C} \quad \text{F} \tag{4.15}
$$

$$
2 b = 2 + 9 C F + 9 \gamma F^{2}
$$
 (4.16)

Note that **Eq.** (4.14) is identical to **Eq.** (4.3) with different definitions of the coefficients. Thus, if $b^2 - \varepsilon^3$ is positive, then the root is computed using **Eq.** (4.7).

When $b^2 - \epsilon^3 < 0$, selection of the proper root can be made by plotting the three roots as was previously done in Fig. 4.1. When $F = 0$ 2 (t = 0), the roots of Eq. (4.14) are x_1 = 2, x_2 = x_3 = -1 and b² = e^{3} Hence point **A** is the solution point for all values of **C** and **y.** From a continuity standpoint, small changes in F, and hence in **b,** correspond to small changes in the root. Therefore if **b** decreases then the correct root will always be located on that portion of the curve from points **A** to B. Hence for three real roots, the correct root is simply calculated from

$$
3 \theta = \arccos(\frac{b}{\sqrt{\epsilon^3}}) \tag{4.17}
$$

and

$$
x = 2\sqrt{\epsilon} \cos \theta \tag{4.18}
$$

To avoid any discontinuities, the requirement that $b \geq -\sqrt{\epsilon^3}$ must be met. This requirement avoids obtaining a value of x in the vicinity of $x = -1$ for which case **Eq.** (4.13) has a singularity.

The most important characteristic of this method is the elimination of division **by y** and hence the need for an asymptotic expansion is avoided. Also computation of the real root has been simplified along with the number of criteria to maintain continuity. Another quality evident is the effect of errors resulting in the calculation of x. Small errors in x using method A are magnified by a factor of $\frac{1}{\gamma}$ which is not the case in this **Y** method.

CHAPTER **5**

FINAL ALGORITHM

Incorporating any iterative procedure into an algorithm requires the use of tests along with certain logic operations to account for cases that may arise which could pose problems if precautions are not taken. In this solution of Kepler's equation only two such cases arise.

The first deals with time intervals which require a larger transfer angle than that which is permitted **by** limitations on the range of Y. This can be handled simply **by** computing the time interval corresponding to the maximum value of Y and comparing with the desired time interval. **If** the time interval corresponding to Y_{max} is greater than the desired time interval then the iteration process is initiated immediately to obtain the final position and velocity using the desired time interval. **If,** on the other hand, the opposite is true then a transfer angle step corresponding to Ymax is taken. This is done **by** computing a position and velocity **by** means of the universal formulae derived previously using the values of B, C_O, D, t and Y corresponding to Y_{max}. Then the time corresponding to Y_{max} is subtracted from the desired transfer time interval. This is continued until the time interval corresponding to Y_{max} is greater than the present desired time interval. This procedure will be referred to subsequently as "time stepping."

For the second case, recall that one requirement in the solution of

Eq. (4.14) was that $b \ge -\sqrt{\epsilon^3}$ (when ϵ is positive), to maintain continuity and also to avoid the singularity at x **= -1** in computing **D.** Note, **b** and c are both functions of **y, C,** and F while **y** and **C,** for the most part, are functions of the initial position and velocity. Thus, **b** and ϵ vary with F and hence t. When e is positive and a value of **b** which is less than $-\sqrt{\epsilon^3}$ is encountered, a simple relation to determine that value of F such that $\mathbf{b} = -\sqrt{\varepsilon^3}$ does not seem to exist.

The value of F where $b = -\sqrt{\epsilon^3}$ is determined from the quadratic

$$
b^{2} - \varepsilon^{3} = 9 F^{2} (\frac{9}{4} \gamma^{2} F^{2} + \frac{9}{2} C (\gamma - \frac{2}{3} C^{2}) F + (\gamma - \frac{3}{4} C^{2})) = 0
$$

and, unfortunately, division **by y** is required. **If y** is very nearly zero, problems will arise. Furthermore, when **c** is negative, there being always only one real root in this case, there is continuity for all values of **b** and ϵ but now b must be restricted to only positive values to positively avoid the singularity at x **= -1.** Hence the first test is to see if **b** is positive or negative. If positive, there is no problem. If negative, then **E** must be checked to see if it is positive or negative. If negative, then a value of F must be determined which will make **b** at least zero. Here **b** is a quadratic itself in F and also requires division **by y.** If **C** is positive then $b \geq -\sqrt{\epsilon^3}$ must be tested. If it is true, there is no problem. If it is false then the above quadratic must be solved for that value of F where $b = -\sqrt{\epsilon^3}$.

Clearly all this testing is very tedious and detracts from the appeal and potential of this method of solution. It was found that for all orbits tested that a root less than zero was never encountered and for the most part was in fact quite large. Thus, the need for the tests described above can be avoided **by** the simple expediency of assuring only positive roots. If a negative root is encountered then a prescribed fraction of the time interval is subtracted and the iteration is reinitiated. This fraction would be preferred to be near unity **(** say **.95)** since fairly small changes in t have a significant effect on the root. Also, a decrease, rather than an increase, in F (t) is taken due to the fact that as F approaches zero both **b** and **c** approach unity while the root, x, approaches a value of 2. Maintaining x positive not only satisfies the continuity requirement when **c** is positive but avoids the possibility of having to subtract a small quantity from unity when computing **D.**

5.1 Procedure

Before presenting the final algorithm the following quantities are defined:

 t_{max} = time interval corresponding to Y_{max} t **=** input time interval t_T = sum of time steps (if any) T **=** time interval computed in the iteration for convergence test

Figure **5.1** illustrates the flow chart diagram used in the tests. In the following discussion, an iteration will be referred to as executing blocks 12 to 20 in the flow diagram.

5.2 Tests and Data

All tests described below were made on a Hewlett Packard model **9820A** calculator which employs 12 significant figures internally and displays **10.**

The first set of **28** tests consist of a series of orbits which comprise the test package for the Kepler subroutine in the Apollo project. The characteristics of these orbits are listed in Table **5.1.** These tests were first run to examine different ranges of Y. Two ranges were examined; $-1 \leq Y \leq 1$ and $-2 \leq Y \leq 2$, both for $\alpha = \frac{9}{10}$. The results of the number of time steps and iterations needed along with the results of the Kepler subroutine proposed for the **NASA** Space Shuttle orbiter vehicle are listed in Table **5.2.**

The second set of tests was used to examine the performance of this method for orbits of very high eccentricities. These cases evolved from modelling the Earth as a series of point masses, resulting in orbits of very high eccentricities. The orbits and results are listed in Table **5.3.**

Finally, we note that several cases were run for different values

Table **5.1**

Characteristics of test orbits in Apollo test package

 \bar{z}

Table **5.2** Number of iterations for Apollo test package

 $\ddot{}$

Table **5.3**

High Eccintricity Test Orbits resulting from use of Point Masses for a Circular Orbit around Earth

Characteristics

Time interval equals a quarter orbit Point masses scaled to $(1/\mu_e)$

Results

of α . The value of $\frac{9}{10}$ was superior to an α of $\frac{3}{10}$ for reasonable values of transfer angle although they were comparable for very small transfer angles. Hence, it was concluded that the value of $\frac{9}{10}$ was the proper choice, just as Gauss had determined for the more elementary problem.

5.3 Discussion of Results

The preceeding results clearly illustrate the application of this method to the solution of transfer problems for all types of orbits and for a wide range of eccentricity.

In the comparison of the two ranges of Y, the smaller range did tend to reduce the number of iterations in several instances. And, as would be expected, the smaller range increases the number of time steps but the amount of computation required in taking a time step is small. Also, the smaller range of Y requires fewer terms in the B and **C0** power series. Selection of the range of Y is largely up to the user although evidence points towards a smaller range being more benificial. Clearly there is a point of diminishing returns in reducing the allowable range of Y. Furthermore reduction of the range of Y below $\frac{1}{2}$ yields very little in the further reduction of the B and C_O power series. It is felt that a range of Y between ± 1 or possibly $\pm \frac{1}{2}$ would be the best selection.

5.4 Comparison with the Proposed NASA Shuttle Kepler Subroutine

The comparison with the proposed **NASA** Shuttle Shuttle subroutine shows a considerable decrease in the number of iterations in most cases.

It should be kept in mind that the amount of computation in an iteration for both methods is different. Interesting enough, this method shows a substantial decrease in both the amount of logic and computation in one iteration.

The convergence test in the Kepler subroutine was based on a relative error of **10-13** between the computed time in the iteration and the desired time. This error criteria was unobtainable due to the limitations of the HP **9820A** model. Hence a convergence test was selected based on an absolute difference of **10~4** seconds. Cases were run where this error was tightened whenever possible and little or no increase was evidenced in the number of iterations. Hence the comparison of these results could be made with little or no hesitation that these results would differ significantly if the same criteria were used.

Of equal importance to the number of iterations is the accuracy of the final position and velocity vectors. Here again the calculator posed limitations. For instance, the desired time interval could not be posed to full accuracy stice only **10** significant figures could be used without calculator round off. Nevertheless, with a convergence criteria of 10^{-4} seconds, accuracy of the final position and velocity vectors was at least **8** significant figures and sometimes even **10.** Hence, this method is not only quicker but also quite accurate. The reason for this accuracy is due to an added feature in this algorithm. Referring to the flow chart of Fig. **5.1,** the use of T, the computed time in the iteration for the convergence test, in the statement numbers 22 and **23** has the
net effect, upon reaching statement **27** , of taking all the errors which remain in statement 21, on the final iteration of any number of steps, and iterating on that time interval error. Almost always this error is so small that only one iteration is required to obtain a good approximation of Y, B, and **C0** . Here again, in most cases, the limits of the calculator prevented the effectiveness of this feature from being realized since Y was quite small and the calculator set B and C_O exactly to 1. In several cases, though, this was the reason for the increased accuracy of the final position and velocity.

CHAPTER **6**

CONCLUSIONS

The method developed here for the solution of Kepler's equation for the general problem of determining final position and velocity vectors from given initial conditions for a specified time interval through the extension of Gauss' method in the standard form of position determination for' time since pericenter passage also has resulted in a Picard type iteration, requiring only successive substitution. Furthernore, the form is general, thereby being applicable to all conics without knowledge of the conic encountered and at the same time is continuous during transition from one conic to another and is free from ambiguities or indeterminant forms. Also, it has proven itself to be applicable to both rectilinear motion and to all orbits of any eccentricity even in the cases where the eccentricity is very large in which case motion is nearly rectilinear. And, the resulting universal formulae, relating final position and velocity to initial values, are not only simple but are also expressed in terms of variables which have already been computed in the iteration process.

The final algorithm exhibits both simplicity and strong convergence and at the same time has very good accuracy in determining final position and velocity which results from internal correction of errors in the time interval accumulated in the iteration process.

66

Finally, comparison with the Apollo version of solving Kepler's equation showed not only a decrease in the number of iterations but also a reduction in the amount of logic and computation in any one iteration thereby further illustrating its simplicity and potential of finding a wide range of application in computer oriented problems and also presents itself as a simple method for hand or calculator computation.

Most important is the basic concept behind the method which was to transform Keplers equation to an equation which is nearly cubic and hence is solvable through algebraic methods, the result of which is a simple, straightforward and expedious means of obtaining the final position and velocity.

It is recommended that if increased accuracy is desired, say to **16** significant figures in which case better accuracy will be required for B and C_0 , then serious consideration should be given to using the B and **C0** power series about points other than Y **= 0,** since the number **of** terms needed in the series about Y = **0** could get quite large depending on the range of Y selected.

67

APPENDIX **A**

l,

SERIES REVERSION

A procedure for obtaining the coefficients of the expansions of B and C₀ in powers of Y, which is useful if the coefficients are to be calculated **by** hand, is through the use of series reversion and the algebra of power series. The algebraic relations used here are:⁽¹⁾

If
$$
S_1 = 1 + a_1 x + a_2 x^2 + a_3 x^3 + \dots
$$

\n
$$
S_2 = 1 + b_1 x + b_2 x^2 + b_3 x^3 + \dots
$$
\n
$$
S_3 = 1 + c_1 x + c_2 x^2 + c_3 x^3 + \dots
$$

then for

$$
S_{3} = S_{1} S_{2} \t C_{n} = b_{n} + a_{n} + \sum_{0}^{(n-1)} a_{k} b_{n-k}
$$

$$
S_{3} = S_{1} / S_{2} \t C_{n} = a_{n} - b_{n} - \sum_{0}^{(n-1)} c_{k} b_{n-k}
$$

$$
S_{1}^{2} = S_{3} \t a_{n} = \frac{1}{2} c_{n} - \frac{1}{2} \sum_{0}^{(n-1)} a_{k} a_{n-k}
$$

where $c_0 = b_0 = a_0$

The variable Y may be expanded in powers of **AE** using **Eq. (2.8)**

$$
Y = \frac{6 P}{Q} = \frac{60 (\Delta E - \sin \Delta E)}{9 \Delta E + \sin \Delta E}
$$

$$
\frac{y}{\Delta E^2} = \frac{1 + a_1 \Delta E^2 + a_2 \Delta E^4 + a_3 \Delta E^6 + \dots}{1 + b_1 \Delta E^2 + b_2 \Delta E^4 + b_3 \Delta E^6 + \dots}
$$

where

$$
a_n = \frac{6 (-1)^n}{(2 n + 3)!} \qquad b_n = \frac{(-1)^n}{10 (2 n + 1)!} \qquad (1)
$$

Then dividing the two series yields

$$
\frac{\gamma}{\Delta E^2} = 1 + c_1 \Delta E^2 + c_2 \Delta E^4 + c_3 \Delta E^6 + \dots
$$
 (2)

where

$$
c_n = a_n - b_n - \sum_{0}^{(n-1)} c_k b_{n-k}
$$
 (3)

The reversion theorem for a power series states that given an expansion of the form of **Eq.** (2) which is convergent in some interval, then if $c_1 \neq 0$ there exists one and only one function which can be expanded in the form

$$
\frac{\Delta E^2}{\gamma} = 1 + d_1 \gamma + d_2 \gamma^2 + d_3 \gamma^3 + d_4 \gamma^4 + \dots \tag{4}
$$

Clearly one method of obtaining the coefficients to **Eq.** (4) is **by** substituting **Eq.** (2) into **Eq.** (4) and equating powers. Unfortunately, no general expression for the n^{th} coefficient is obtainable through this

process. An alternate procedure is based on the fact that

$$
d_{n-1} = \frac{d^n n(\Delta E^2)}{n!} \tag{5}
$$

Differentiating **Eq.** (2) with respect to **AE2** and inverting gives

$$
\frac{d}{d\gamma}(\Delta E^2) = 1 + a_1 \Delta E^2 + a_2 \Delta E^4 + \dots
$$
 (6)

where in general

$$
a_n = -(n + 1) c_n - \sum_{0}^{(n-1)} (k + 1) c_k a_{n-k}
$$
 (7)

Now the n^{th} coefficient in Eq. (4) may be obtained as follows: **1.** Compute **k** coefficients in **Eq. (6).**

2. Starting with n **=** 2, compute

$$
\frac{d^n}{d\gamma^n}(\Delta E^2) = \frac{d}{d\gamma}(\Delta E^2) \frac{d}{d(\Delta E^2)} \left[\frac{d^{(n-1)}}{d\gamma^{(n-1)}}(\Delta E^2) \right]
$$

For example, with n **=** 2,

 \sim σ

$$
\frac{d^{2}}{dY^{2}}(\Delta E^{2}) = \frac{d}{dY}(\Delta E^{2}) - \frac{d}{d(\Delta E^{2})}\left[\frac{d}{dY}(\Delta E^{2})\right]
$$

= $(1 + a_{1} \Delta E^{2} + ... + a_{k-1} \Delta E^{2k-2})(a_{1} + ... + k a_{k} \Delta E^{2k-2})$

$$
= a_1 (1 + b_1^2 \Delta E^2 + b_2^2 \Delta E^4 + \ldots + b_i^2 \Delta E^2 i + \ldots)
$$

where

$$
b_{i}^{2} = \frac{1}{a_{1}} \Big[(i+1) a_{i+1} + i a_{i} a_{1} + (i-1) a_{i-1} a_{2} + \dots + a_{1} a_{i} \Big] \qquad \qquad 1 \leq i \leq (k-1)
$$

With
$$
n = 3
$$
,
\n
$$
\frac{d^{3} \Delta E^{2}}{dY^{3}} = a_{1}(1 + a_{1} \Delta E^{2} + ... + a_{k-2} \Delta E^{2k-4})(b_{1}^{2} + 2 b_{2}^{2} \Delta E^{2} + ... + (k-1) b_{k-1}^{2} \Delta E^{2k-4})
$$
\n
$$
= a_{1} b_{1}^{2} (1 + b_{1}^{3} \Delta E^{2} + ... + b_{i}^{3} \Delta E^{2i} + ...)
$$

where

$$
b_{\hat{1}}^3 = \frac{1}{b_1^2} \left[(i+1) b_{\hat{1}+1}^2 + \sum_{1}^{i} m b_m^2 a_{\hat{1}-m+1} \right]
$$
 $1 \leq i \leq (k-2)$

and in general

$$
\frac{d^{n}\Delta E^{2}}{dY^{n}} = a_{1} b_{1}^{2} b_{1}^{3} ... b_{1}^{k} (1 + b_{1} \Delta E^{2} + ... + b_{i} \Delta E^{2} \dagger + ...)
$$
\n
$$
1 \leq i \leq (k + 1 - n)
$$

where

$$
b_{i} = \frac{1}{b_{1}^{n}} \left[(i+1) b_{i+1}^{n} + \sum_{1}^{n} m b_{m}^{n} a_{i-m+1} \right]
$$

It can be seen that $1 \le n \le (k+1)$. At the same time

$$
\frac{d^{n}\Delta E^{2}}{dY^{n}}\Big|_{Y=0} = a_{1} b_{1}^{2} b_{1}^{3}...b_{1}^{3}
$$

$$
= \left[\frac{d^{n-1}\Delta E^{2}}{dY^{n-1}}\Big|_{Y=0}\right]_{Y=0}^{n-1}
$$

Hence, if **k** coefficients are calculated in **Eq. (6),** then **k** coefficients are obtainable in **Eq.** (4). With the series in **Eq.** (4) determined, B as a power series in powers of Y may be obtained since

$$
B^2 = \frac{\Delta E^2}{Y} (1 + \frac{Y}{60})^{-2}
$$

Using the binomial expansion

$$
(1 + \frac{\gamma}{60})^{-2} = 1 + a_1 Y + a_2 Y^2 + \dots
$$

where

$$
a_n = \frac{(-1)^n (n + 1)}{(60)^n}
$$

then

$$
B2 = (1 + d1 Y + d2 Y2 + ...) (1 + a1 Y + a2 Y2 + ...)
$$

= 1 + b₁ Y + b₂ Y² + b₃ Y³ + ...

where

$$
b_n = a_n + d_n + \sum_{0}^{n-1} a_i d_{n-1}
$$

Hence

$$
B = 1 + A_1 Y + A_2 Y^2 + A_3 Y^3 + \dots
$$

with

$$
A_n = \frac{1}{2} b_n - \frac{1}{2} \sum_{0}^{n-1} b_i b_{n-i}
$$

To obtain the C_0 power series in Y;

$$
2R = 2(1 - \cos \Delta E) = \Delta E^2 (1 - \frac{2 \Delta E^2}{4!} + \ldots + \frac{2 (-1)^n \Delta E^{2n}}{(2 n + 2)!} + \ldots)
$$

and after substitution of **Eq.** (4) for **AE2** into the above series and expanding the result is

$$
2R = \Delta E^{2} (1 + b_{1} Y + b_{2} Y^{2} + b_{3} Y^{3} + ...)
$$

where the b's are the result of raising **Eq.** (4) to the corresponding power of ΔE^2 in the above series and adding the coefficients of similar powers in Y. Now using the definition of **C0**

$$
C_0 = (\frac{1}{B}) (\frac{\Delta E^2}{Y}) (1 + b_1 Y + b_2 Y^2 + b_3 Y^3 + ...)
$$

or

$$
C_0 = (1 + a_1 Y + a_2 Y^2 + ...) (1 + d_1 Y + ...) (1 + b_1 Y + ...)
$$

where

$$
a_n = -A_n - \sum_{0}^{n-1} a_k A_{n-k}
$$

Hence, performing the required multiplication

$$
c_0 = (1 + w_1 Y + w_2 Y^2 + ...) (1 + b_1 Y + b_2 Y^2 + ...)
$$

where

$$
w_n = a_n + d_n + \sum_{0}^{n-1} a_k d_{n-k}
$$

and finally

$$
c_0 = 1 + A_1 Y + A_2 Y^2 + A_2 Y^3 + \dots
$$

where

$$
A_n = w_n + b_n + \sum_{0}^{n-1} w_k b_{n-k}
$$

 $\sim 10^{11}$

APPENDIX B

SERIES EXPANSION ABOUT AN ARBITRARY POINT YO

The methods presented earlier for obtaining B and **C0** as power series in Y about Y **⁼0** are no longer practical if the expansion is required about some arbitrary point Y₀. An alternate procedure can be used in this instance to obtain numerical values of the coefficients for these two series which makes use of the simplicity of the functions P, Q, and R and their derivatives and Leibnitz's formula for the differentiation of a product, which states

$$
\frac{d^{n}}{dx^{n}} (u \ v) = \frac{n}{0} {n \choose k} u^{n-k} v^{k} \qquad \text{where} \qquad {n \choose k} = \frac{n!}{k! (n-k)!}
$$

For convenience, let

$$
P^{(i)} = \frac{d^{i} P}{d E^{i}} E = E_0
$$

be the convention for all functions, other than B and C₀, in which case

$$
B^{(i,j)} = \frac{d^{i+j}}{d E^{i} dY^{j}} B
$$

will be the convention

Leibnitz's formula for computing the $n-$ derivative of a function assumes that the values of the (n-1) derivatives are known. To compute

$$
\gamma^{(n)}
$$
; from Eq. (2.8)
Q Y = 6 P

Differentiating both sides n times gives

$$
\sum_{0}^{n} \binom{n}{k} q^{(n-k)} \gamma^{(k)} = 6 p^{(n)} \qquad n = 0, 1, ...
$$

or solving for $\gamma^{(n)}$

$$
Q Y^{(n)} = 6 P^{(n)} - \sum_{0}^{n-1} {n \choose k} Q^{(n-k)} Y^{(k)} \qquad n = 1, 2, ... (1)
$$

For
$$
\frac{d^n B}{d\Delta E^n} \Big|_{\Delta E_0} = B^{(n,0)}
$$
; If we let
 $x = z^2 = B^2 Y^2 = 6 P Q$ (2)

Then differentiating $x = z^2$ gives

$$
x^{(n)} = \sum_{0}^{n} {n \choose k} z^{(n-k)} z^{(k)}
$$
 n = 0, 1, ... (3)

but x **= 6** P **Q,** hence

$$
x^{(n)} = 6 \sum_{0}^{n} {n \choose k} p^{(n-k)} q^{(k)} \qquad n = 0, 1, ... \qquad (4)
$$

Equating Eqs. **(3)** and (4) gives, when n = 1

$$
2 z z(1) = 6 (P(1) Q + P Q(1))
$$
 (5)

and for $n \ge 2$,

$$
2 z z^{(n)} = 6 \sum_{0}^{n} {n \choose k} p^{(n-k)} q^{(k)} - \sum_{1}^{n-1} {n \choose k} z^{(n-k)} z^{(k)}
$$
(6)

And since z **=** B Y, differentiating with respect to **E** gives

$$
z^{(n)} = \sum_{0}^{n} {n \choose k} B^{(n-k,0)} \gamma^{(k)}
$$

= $\sum_{1}^{n} {n \choose k} B^{(n-k,0)} \gamma^{(k)} + B^{(n,0)} \gamma$

or

$$
\gamma B^{(n,0)} = z^{(n)} - \frac{n}{2} {n \choose k} B^{(n-k,0)} \gamma^{(k)} \qquad n = 1, 2, ... \qquad (7)
$$

Now $B^{0,0,1} = \frac{1}{B}$ dY n **AE0** be determined using the chain rule

$$
B^{(1,j)} = \frac{d}{d\Delta E} \frac{d^{j}B}{d\gamma^{j}} \bigg|_{\Delta E_{0}} = \gamma^{(1)} B^{(0,j+1)}
$$
 (8)

Differentiating i times with respect to **AE** gives

$$
B^{(i,j)} = \int_{0}^{i-1} {i-1 \choose k} \gamma^{(i-k)} B^{(k,j+1)} \qquad i = 2, 3, ...
$$

$$
\gamma^{(1)} B^{(i-1,j+1)} = B^{(i,j)} - \int_{0}^{i-2} {i-1 \choose k} \gamma^{(i-k)} B^{(i,j+1)} \qquad (9)
$$

or

If **j** is replaced **by** n-1 in **Eq. (8)** then

$$
B^{(1,n-1)} = \gamma^{(1)} B^{(0,n)}
$$
 (10)

hence, $B^{(1)}$, $B^{(2)}$ must be determined. The procedure for doing this can be seen **by** expanding **Eq. (9)** for several values of n:

For $n = 1$, $B^{(1,0)}$ is determined from Eq. (7)

from Eq. (10) $B^{(1,0)} = \gamma^{(1)} B^{(0)}$

For n **=** 2, **B(2,0)** is determined from **Eq. (7)**

In **Eq. (9)** For i=2; j=0 $\gamma^{(1)}$ B^(1,1) = B^(2,0) - Y⁽²⁾ From **Eq. (10)** $R(1,1) = \gamma(1) R(0,2)$

For n **= 3, B(3,0)** is determined from **Eq. (7)**

In **Eq. (9)** For i=3; j=0 $\gamma^{(1)}$ B^(2,1) = B^(3,0) - $\gamma^{(3)}$ B^(0,1) - 2 $\gamma^{(2)}$ B^(1,1) For i=2; j=1 $\gamma^{(1)}$ $\beta^{(1,2)}$ = $\beta^{(2,1)}$ ₋ $\gamma^{(2)}$ $\beta^{(0,2)}$ From **Eq. (10)** $B^{(1,2)} = Y^{(1)} B^{(0,2)}$

Hence in general for $n \ge 2$, in Eq. (9)

$$
\gamma^{(1)} \; \mathbf{B}^{(n-j-1, j+1)} = \mathbf{B}^{(n-j,j)} - \sum_{0}^{n-j-2} \; \mathbf{B}^{(n-j-1)} \; \gamma^{(n-j-k)} \; \mathbf{B}^{(k, j+1)}
$$
\n
$$
j = 0, 1, ..., n-2 \qquad (11)
$$

To determine
$$
\frac{d^n c_0}{d\Delta E^n}\Big|_{\Delta E_0}
$$
,

using **Eq. (2.9)** and **Eq.** (2)

$$
z C_0 = 2 R \tag{12}
$$

Differentiating n times with respect to **AE** gives

$$
\sum_{0}^{n} \binom{n}{k} C_{0}^{(k,0)} z^{(n-k)} = 2 R^{(n)}
$$

or

$$
z C_0^{(n,0)} = 2 R^{(n)} - \sum_{0}^{n-1} {n \choose k} C_0^{(k,0)} z^{(n-k)} \qquad n = 1, 2, ... \qquad (13)
$$

l''C and now **-0** is obtained using Eqs. **(8), (10)** and **(11)** where B dYn **AE0** is simply replaced by c_0 .

To summarize the procedure of determining the coefficients of the B and C_O power series in Y about some arbitrary point Y_O; To calculate m coefficients to these series the sequence is as follows

- **1.** Evaluate the variables P, **Q** and R and their m derivatives at the point ΔE_0 (ΔH_0).
- 2. Evaluate

$$
\gamma^{(0)} = \frac{6 \ P^{(0)}}{Q^{(0)}}
$$
 $B^{(0,0)} = \frac{Q^{(0)}}{\gamma^{(0)}}$ $C_0^{(0,0)} = \frac{2 \ R^{(0)}}{\gamma^{(0)} \ B^{(0,0)}}$

3. For n **= 1,** evaluate

$$
\gamma^{(1)} = (6 \rho^{(1)} - q^{(1)} \gamma^{(0)})/q^{(0)}
$$

\n
$$
z^{(1)} = 3 (\rho^{(1)} q^{(0)} + \rho^{(0)} q^{(1)})/z^{(0)}
$$

\n
$$
B^{(1,0)} = (z^{(1)} - B^{(0,0)} \gamma^{(1)})/Y^{(0)}
$$

\n
$$
B^{(0,1)} = B^{(1,0)}/Y^{(1)}
$$

\n
$$
C_0^{(1,0)} = (2 R^{(1)} - C_0^{(0,0)} z^{(1)})/z^{(0)}
$$

\n
$$
C_0^{(0,1)} = C_0^{(1,0)}/Y^{(1)}
$$

4. For n **= 2, 3,** 4, .. **.,** in

$$
Q^{(0)} Y^{(n)} = 6 P^{(n)} - \sum_{0}^{n-1} {n \choose k} Q^{(n-k)} Y^{(k)}
$$

\n
$$
2 Z^{(0)} Z^{(n)} = 6 \sum_{0}^{n} {n \choose k} P^{(n-k)} Q^{(k)} - \sum_{1}^{n-1} {n \choose k} Z^{(n-k)} Z^{(k)}
$$

\n
$$
Y^{(0)} B^{(n,0)} = Z^{(n)} - \sum_{1}^{n} {n \choose k} B^{(n-k,0)} Y^{(k)}
$$

\n
$$
Z^{(0)} C^{(n,0)}_{0} = 2 R^{(n)} - \sum_{0}^{n-1} {n \choose k} C^{(k,0)} Z^{(n-k)}
$$

 \sim

For j **= 0, 1, ... ,** (n-2) i^{-2} , n-j-1, $v(n-j-k)$ $R(k,j+1)$ 1) **b** $B(n-1-1, 3+1) = B(n-1, 3) - \sum_{k=0}^{n} {n-1-1 \choose k} \lambda_{k} (n-1-k) B(k, 3)$

For
$$
j = 0, 1, ..., (n-2)
$$

\n
$$
\gamma^{(1)} C_0^{(n-j-1)} = C_0^{(n-j,j)} - \sum_{0}^{n-j-2} {n-j-1 \choose k} \gamma^{(n-j-k)} C_0^{(k,j+1)}
$$
\n
$$
\gamma^{(1)} B^{(0,n)} = B^{(1,n-1)}
$$
\n
$$
\gamma^{(1)} C_0^{(0,n)} = C_0^{(1,n-1)}
$$
\n
$$
a_n = B^{(0,n)}/(n!) \qquad b_n = C_0^{(0,n)}/(n!)
$$

where

$$
B = 1 + \sum_{1}^{m} a_{n} (Y - Y_{0})^{n} \qquad C_{0} = 1 + \sum_{1}^{m} b_{n} (Y - Y_{0})^{n}
$$

Table B.1 lists the coefficients of the B and C₀ power series expanded about the points **.5, 1.0, 1.5,** 2.0. Table B.2 lists the coefficients for the points **-.5, -1.0, -1.5,** -2.0. In both tables the number of significant figures drops from **29** for the zeroth coefficient to 15 for the $15\frac{th}{ }$. It will be hence safe to assume that the coefficients from 1 to **5** are correct to **25** figures; coefficients **6** through **10** are correct to 20 figures and those from **11** through **15** correct to **15** figures. Finally, it is important to take note of the fact that the sign of Y has not been accounted for in Table B.2. Hence for Y **< 0,** the sign of the odd powered coefficients must be changed.

Table **B.1**

B and C_O series coefficients for Y > 0

y = **4.99920793653689037165097540983p-1 ⁰**

 \sim \sim

Table B.1 Continued

 $b8$

 $0.4809182228992886082126522226010^{*} = \frac{0}{0}$

 98

$S.B.B.B.I$

 $0 > Y$ not zametoittes coefficients for $Y < 0$

$$
\begin{array}{cccc}\n\angle T - 371120281889007926269109906216^2 & - & 91 \\
& 91 - 916690999117091892186166668828^2 & 61 \\
& 91 - 9126182022228895110962922920222^2 & - & 21 \\
& 61 - 912618225228895110962922920222^2 & - & 21 \\
& 61 - 926622522886662252065222^2 & - & 11 \\
& 61 - 92662298256880012821762821762^2 & - & 11 \\
& 61 - 921266229825688001288717622^2 & - & 6 \\
& 61 - 915660922821266118622126022562^2 & - & 6 \\
& 61 - 9156609228212666158100556668^2 & 8 \\
& 61992201222221836892608828012^2 & - & 2 \\
& 616085088283666919652256662989^2 & - & 6 \\
& 6198292012222218368892608828012^2 & - & 6 \\
& 6160850882166222218838892608828012^2 & - & 6 \\
& 61608508216622221836889260838012^2 & - & 6 \\
& 61608508216622221836662989^2 & - & 6 \\
& 6160209221262626219829051180106^2 & - & 6 \\
& 6160209221262626219829051180106^2 & - & 6 \\
& 6160209221262626225666298^2 & - & 6 \\
& 61602
$$

 \mathcal{A}^{max}

LT-8ZZEt166tZ0089TZEZttZ9tZ08TSTt*Z ST

 $\sim 10^{-10}$

$$
\mathbf{0} = -\mathbf{r} + \mathbf{0} \mathbf{a} \mathbf{c} + \mathbf{c} \mathbf{c} \mathbf{c} \mathbf{c} + \mathbf{c} \mathbf{c} \mathbf{c} \mathbf{c} \mathbf{c} \mathbf{c}
$$

 $\overline{18}$

 $\langle \cdot \rangle$

ا في المساوية المساوية المساوي التي تتم التي تتم التي يعني ورياد السيورية المساوية المستقلة المستقلة المساوية

 0_{λ} $0.45828207500765376088097266660$

- \overline{u}_{p} $\overline{\mathsf{u}}$
- $\mathbf 0$ 0858684865866878816695899186888800*1
- Σ =872202365364335463403845544345455485 Ţ
- V-BSZ886CZIZE8V6V9CEZ86SEIZ96ZE9'8 Z
- 9-HEZZ928122220GZGVZ00Z821VI6VV96*Z 5-872921992912566026268872256288918 £
- $\mathbf G$ $\mathbf b$
- 6-81022984919626112999988048686561 - ℓ 8-H9ZbETTSOTTTBS90Zbb0b0bZ9Tb6Z'T 9
- OT-860254662TB4TS86208928T680T042*T 8
- ΣΤ-ΩΦΟΣΦΖΙΩΟΦΖΒΙΦΟΖ9ΒΣΒ9ΖΣΖΒΣΒΟΣΖ'6 ΟΤ TT-HOVOVTSO9ZIV8ZTO98V0S0SSSSTX0T*T-- $\overline{\mathcal{A}}$
- VI-HIEIEBIB6IS944929420IBE9VV8269'8 II
- 91-8ZVSZZS8Z969Z9S£19SZ90S9VISOSI*Z £1 ST-HVE60ZS08Z0ZZTIZETZE0ZZE2S8V8'Z ZT
- ZI-HOVIZIZZVOVESEEEZEVIBVOZIVS9919 VI
- 81-41450854900650951287451640690'9 51

81-46006928050198555645664442660'S ST ZT-H92ZI8OS62VI628OSSZIV9S2V9V2ZS'S - VT 91-Hb8SZZEE60ZI9E86b0bSI896IIIZb1'9 21

 $BT = H996Zb bS9ZBZBZbZSZb b90b900b9b0b b8b b7Z - 9T = 1$ $LT = H6bZZQZQZBbSLZSbLSBbRBOZ9SbBQbSLSZ = bL$ 91-89162261690221090169088092692*£ \$1

 71

- T-T

51-89209099822/1969/082829902128'9 -

VI-BIBESZEVV9EZZ8V9IZS600ZZEZVZS'V

 Σ T-876970288811966618009687166676'S - 01

REFERENCES

- 1. Abramowitz, M., and I. A. Stegun, Handbook of Mathematical Functions, Dover Pub., Inc., **1965.**
- 2. Battin, R. H., Astronautical Guidance, McGraw-Hill, 1964.
- **3.** Gauss, K. F., Theory of the Motion of the Heavenly Bodies Moving about the Sun in conic Sections, a translation of Theoria Motus, Dover Pub., Inc., **1963.**
- 4. Knopp, K., Theory and Application of Infinite Series, Blackie and Son, **1951.**