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## THE STABILITY OF MOTION IN A PERIODIC CUBIC FORCE FIELD

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The object of this research has been to investigate numerically the possible existence of closed, invariant regions of stability for solutions of a particular nonlinear differential equation. The Illiac II computer at the University of Illinois has been utilized for an accurate study of the stability properties of solutions of the "Cubic Equation". This nonautonomous equation has the form  $d^2X/dt^2 + p(t) X^3 = 0$ , where  $p(t)$  is a periodic square-wave function of time with magnitude  $p_0$  and period  $\tau$ . This equation can be integrated piecewise in terms of Jacobian elliptic functions, and hence is suitable for a high-speed digital computer. Not relying on numerical integration techniques allows us to perform many highly accurate computations over thousands of periods of  $\tau$ .

For fixed values of  $p_0$  and  $\tau$  we have studied the motions in a global neighborhood of a particular solution with period  $11\tau$  which has an elliptic fixed point on the X-axis. Initial exploratory calculations using numerical integration revealed a region associated with fixed points of period  $55\tau$  which apparently did not allow any of its interior points to escape to infinity. This

discovery suggested that we make a more careful study of solutions associated with the clusters of multiple points (fixed under multiples of  $11\tau$  periods) which surround the original elliptic fixed point. The solutions with  $22\tau$  and  $44\tau$  periods quickly carried points to infinity. The solutions with  $11\tau$  and  $33\tau$  periods possessed characteristic oscillating hyperbolic invariant curves, whereas the solutions with a  $55\tau$  period did not. It was found that the outer  $33\tau$  oscillating invariant curves intersected the outgoing  $11\tau$  oscillating invariant curves. However the inner  $33\tau$  oscillating invariant curves formed a tangled mesh of self-intersections. This "entanglement" effectively prevented them from intersecting the  $55\tau$  hyperbolic invariant curves. Consequently points in the domain of the  $55\tau$  solutions were "stable".

The final results clearly show (to the  $10^{-12}$  accuracy of the Illiac II computer) the existence of a closed invariant region, associated with fixed points of period  $55\tau$ , around the original  $11\tau$  elliptic fixed point. This region is bounded by the doubly-periodic union of the non-oscillating hyperbolic invariant curves associated with the  $55\tau$  solutions. It is surrounded by a jumbled global structure of oscillating and intersecting hyperbolic invariant curves belonging to mappings with lesser multiplicities. This intersection property provides the escape mechanism for points just outside the invariant region. Points initially inside the invariant region always map inside that region, even after an infinite number of repeated mappings of the phase-space. The points belonging to this region are said to be stroboscopically stable.

If we scale our flat nearly-elliptical twist mappings up to nearly-circular ones (by scaling  $p_0$ ), then the conditions of Moser's theorem on invariant curves for certain twist mappings apply. However they are not sufficient to locate the observed invariant region accurately. (Testing Moser's conditions involves as much work as actually searching for the non-oscillating invariant curves.) Thus the best method at present for obtaining practical knowledge about the precise location of invariant regions of stability is to use a digital computer. The most important features of our work should be reproducible on less accurate computers, even utilizing standard numerical integration techniques for the equations of motion. Our empirical methods would then afford some practical guidelines for a quantitative global understanding of the absolute stability of solutions of any particular nonautonomous differential equation.

## ACKNOWLEDGEMENT

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## I. INTRODUCTION

The stability of motion is a concept which is defined analogously to the stability of equilibrium, with which we are all familiar. Just as small displacements from an equilibrium position may result in a particle moving permanently away from that position (which is then said to be one of unstable equilibrium), so small initial displacements from a periodic orbit can result in steady motion away from that orbit, and the initial motion is termed unstable.

The planetary orbits appear to be "stable" since they have changed very little over several thousands of years. However, astronomical evolution occurs over billions of years, and it would be nice to know how stable the solar system is during such a time span. Evidence may be gained for this problem by observing the stability properties of charged particle trajectories in accelerators and magnetic bottles. These trajectories may involve millions of orbit repetitions in relatively short times, so that any instabilities will show up vividly. In fact the elimination of plasma instabilities may help in controlling nuclear fusion, so that precise knowledge of "long-period" stability is vital. To attain this, we need an adequate mathematical theory of stability of motion for dynamical systems in general, and for stellar systems, accelerators, and stellera-tors in particular.

In general, the dynamical equations are nonlinear, either explicitly time-dependent (i.e. non-autonomous) or else time-independent. The solutions of the equations may be found numerically, but then we have to examine their totality in order to make conclusions about stability. General treatises on the "stability of motion" (e.g. the survey by Sansone and

Conti<sup>1/</sup> or Bellman<sup>2/</sup>, with references) list many different definitions of "stability", depending partly on the approximations which are to be made. We shall begin by briefly discussing first-order stability. Let us consider a set of  $n$  coupled first-order equations of the motion

$$dX_i/dt = F_i(X_j) \quad \text{for } i, j = 1, 2, \dots, n \quad (1)$$

which possess periodic (closed) solutions of the form

$$X_i(t) = f_i(t) = f_i(t+T) \quad (2)$$

where  $T$  is the period of the motion. If we now look at some "nearby" solution whose  $n$  coordinates  $X_i'(t)$  at time  $t$  differ from  $X_i(t)$  by infinitesimal perturbations  $\delta X_i(t)$ , we can write

$$X_i'(t) = X_i(t) + \delta X_i(t) \quad (3)$$

Since this is also a solution of equation (1) by definition, we can obtain  $n$  variational equations for the  $\delta X_i(t)$  infinitesimals, namely

$$\frac{d}{dt} \delta X_i = \sum_{j=1}^n \frac{\partial F_i}{\partial X_j} \delta X_j \quad i = 1, 2, \dots, n \quad (4)$$

These linear differential equations (with coefficients periodic in  $t$ ) will possess solutions for the  $\delta X_i$  of the form

$$\delta X_i(t) = \sum_{k=1}^n e^{\alpha_k t} M_{ik} \quad (5)$$

where  $M_{ik}$  is an array of functions of period  $T$ , and  $\alpha_k$  are  $n$  constants known as characteristic exponents of the solution. If all the  $\alpha_k$  are purely imaginary, then the perturbations  $\delta X_i(t)$  themselves will be periodic and bounded. If, however, some of the  $\alpha_k$  are real and positive, then the

$\delta X_i(t)$  will grow exponentially with time. Discussions of the exact methods for determining the  $\alpha_k$  in general, as well as their implications, can be found in Whittaker<sup>3/</sup> or Minorsky<sup>4/</sup>. This variational approach, due to Poincaré, is a rather universal one, but is clearly only a first-order approximation to the effects of arbitrary (finite) changes in all the variables. Specific applications of these methods to the determination of the first-order stability of solutions to the restricted three-body problem<sup>5/</sup> are found in the work of Shearing<sup>6/</sup> and Henon<sup>7/</sup>.

The limitations to the usefulness of this first-order theory are apparent. As long as we consider solutions with only an infinitesimal separation from the known periodic solution, we can definitely say (provided, of course, that we can directly integrate the variational equations (4)) whether the nearby solution will grow in size with increasing time, or whether it will remain infinitesimally "close". But what about solutions which are a finite distance from the periodic solution? Clearly the variational theory is only a local theory, whereas what we need is a global study of solutions of the equations of motion.

A first step in this direction was made by Liapunov<sup>8/</sup> near the end of the nineteenth century. His so-called "second (or direct) method" yields conditions for "stability in the large", without the need to integrate the often difficult variational equations. His theory centers around the existence of certain algebraic functions (Liapunov V functions) of the  $n$  variables, which possess properties of definiteness of sign in different regions of phase-space. Liapunov's second theorem states that if there exists such a definite function  $V(X_i)$  whose Eulerian derivative (i.e. a time derivative taken along a trajectory) is either zero or of opposite sign to  $V$  in some domain  $D$ , then the differential system of equation (1)

possesses a stable equilibrium around the singular point  $X_i=0$ . Weakening of the hypotheses produces several modified theorems for other "types" of stability. In the general case, however, his theorems do not specify any method for determining the function  $V$ , nor do they prove its existence. If such a function can be generated somehow, then it can be used to demonstrate stability of motions in the large or at least in definite regions. Although Liapunov's theorems have found extensive applications during the last 70 years, they still present a number of analytical problems which can be as difficult as those of the variational theory. But they do begin to tackle the basic problem of determining the global properties of solutions, as opposed to the strictly local properties.

The same period of history saw the emergence of a fundamentally different point of view about stability -- the topological approach of H. Poincaré<sup>9/</sup>. He regarded the differential equations of motion as defining a transformation  $T^n$  of the phase-space onto itself. Specifically, the integration of the equations of motion takes initial coordinates on the manifold and maps them into final coordinates on the same manifold. It is not hard to imagine how simple regions of initial coordinates can become mapped onto quite complex, distorted regions on the manifold. Yet it is also conceivable that some particular points on the manifold could be transformed into themselves. Such points are called fixed points of the mapping  $T^n$ . In fact, Brouwer's remarkable fixed point theorem<sup>4/</sup> asserts that every continuous transformation of a closed ( $n$ -dimensional) disk onto itself has at least one such fixed point.

The works of Poincaré in the field of topological approaches to stability are very extensive, but we need to mention here only one of his

last theorems, the so-called "geometric theorem"<sup>10/</sup>. This theorem, which arose in connection with periodic solutions of the restricted three-body problem, concerns itself with any closed annular ring  $0 < a \leq r \leq b$  in the  $r\theta$ -plane ( $a$  and  $b$  are fixed radii), and asserts that if a given transformation  $T^n$  of the annular ring onto itself satisfies certain conditions, then at least two points of this ring will be fixed under  $T^n$ . Those necessary conditions on  $T^n$  are that it must be one-to-one, continuous, and area-preserving, and that it advances points on the (inner) circle  $r=a$  while regressing points on the (outer) circle  $r=b$ . (A later extension of this theorem by Birkhoff<sup>11/</sup> generalizes to the ring  $0 < a \leq r$ . If  $T^n$  advances points on the circle  $r=a$  and regresses all points  $r \geq R > a$  by at least some angle  $\theta_1 > 0$ , then there will exist at least two fixed points of  $T^n$  in the open ring  $a \leq r < R$ .) This extremely remarkable geometric theorem of Poincaré has set the stage for all later topological investigations of stability.

Having isolated at least two fixed points of the mapping  $T^n$ , we may logically wish to know the general properties of the mapping  $T^n$  for points in any neighborhood of the invariant points just found. If we choose a  $(U,V)$  coordinate system on the manifold with the fixed point in question at the origin  $U=V=0$ , then we can investigate transformations  $T^n$  of the general form

$$\left\{ \begin{array}{l} U_1 = aU_0 + bV_0 + \dots \\ V_1 = cU_0 + dV_0 + \dots \end{array} \right\} \quad (6)$$

where  $(U_1, V_1)$  is the map under  $T^n$  of some initial point  $(U_0, V_0)$  in a neighborhood surrounding the origin, and where  $a, b, c$ , and  $d$  are real constants satisfying  $ad-bc > 0$ . By studying the roots  $\rho_i$  of the

characteristic equation for the linear transformation above (e.g. see Minorsky):

$$\rho^2 - (a+d)\rho + (ad-bc) = 0$$

we find there are three basic types of such linear transformations. These three cases include

$$T_h^n: \left\{ \begin{array}{l} U_1 = \rho U_o \\ V_1 = U_o / \rho \end{array} \right\}, \quad (\rho \neq \pm 1) \quad (7)$$

$$T_e^n: \left\{ \begin{array}{l} U_1 = U_o \cos \theta - V_o \sin \theta \\ V_1 = U_o \sin \theta + V_o \cos \theta \end{array} \right\}, \quad (\rho = e^{i\theta}) \quad (8)$$

$$T_p^n: \left\{ \begin{array}{l} U_1 = \pm U_o \\ V_1 = \pm V_o + a U_o \end{array} \right\}, \quad \left( \begin{array}{l} \rho = \pm 1 \\ a \neq 0 \end{array} \right) \quad (9)$$

The symbols h, e, and p designate respectively hyperbolic, elliptic, and parabolic mappings. That is, under repeated applications of one of these mappings, points near the original fixed point ( $U=V=0$ ) move, or map, successively into other points whose locus is one of the three basic conic sections. Consequently we speak of hyperbolic (saddle), elliptic (vortex), and parabolic (nodal) fixed points under the general mapping  $T^n$ . (A fourth type of mapping exists if  $\rho_1$  is a complex root; the resulting motion is of a spiral nature, and the associated fixed point is called a focal point.)

In terms of our intuitive concept of stability we might guess the general stability properties of points "near" the original fixed point by reference to the basic type of  $T^n$  mapping under which that point is fixed. Thus elliptic points, plus spiral and parabolic points with nearby motions directed inwards to the point, ought to be relatively "stable", while all others ought to be basically "unstable". For example, the linear mapping

$T_e^n$  successively maps an initial (neighboring) point along the locus of an ellipse whose major and minor axes will depend on the exact initial coordinates.

But the  $T^n$  transformations discussed above are strictly linear mappings. In general, nonlinear differential equations will possess non-linear mappings  $T^n$ . Consequently any discussion of the possible stability or instability of a given fixed point of the mapping  $T^n$  will necessarily be of a more complex nature. Moreover, there are several additional theoretical and empirical results of the topological approaches of Poincaré and Birkhoff that complicate the general nonlinear problem. Birkhoff<sup>12/</sup> has proved that inside any neighborhood of a given alliptic-type fixed point in phase-space there exist clusters of multiple fixed points. These multiple points are fixed under some integer multiple  $m$  of the original fixed point mapping  $T^n$ , i.e. are fixed points of the mapping  $T^{mn}$  ( $m$  successive applications of  $T^n$ ). As we shall see later, for any particular value of the integer  $m$  there is an even number  $2m$  of these  $T^{mn}$  multiple points, and they alternate around the  $T^{mn}$  cluster between elliptic-type and hyperbolic-type fixed points. To complicate matters even more, Bartlett<sup>13/</sup> has shown, in a quantitative demonstration of Poincaré's theory of consequents<sup>9/</sup>, that the invariant curves issuing from and toward a  $T^n$  hyperbolic-type fixed point are, in general, oscillatory curves. (By invariant curve we mean a locus of points all of which are consequents, or maps under  $T^n$ , of themselves alone.) In fact, these hyperbolic invariant curves develop area-preserving oscillations of increasing amplitude as they approach the next hyperbolic-type point beyond the adjacent elliptic-type point.

It is this (denumerable) infinity of clusters of alternating  $T^{mn}$  fixed points and oscillating hyperbolic invariant curves which complicate the analysis of nonlinear mappings. In fact, this complex behavior suggests the possibility that all fixed points could be ultimately "unstable". For example, in any neighborhood of a given elliptic-type fixed point, no matter how small, one could always find regions containing hyperbolic-type fixed points of increasing multiplicity  $m$ . Points near these  $T^{mn}$  hyperbolic-type points ( $m$  of them) could be readily carried far away by the oscillating invariant curves issuing from (and coming into) the hyperbolic-type points. If successive mappings of these points under  $T^{mn}$  carry them out to infinity, then the original elliptic-type point would be "unstable".

The investigation of various forms of area-preserving mappings of annuli onto themselves has been a fruitful topological approach ever since the pioneering work of Poincaré. It is this same technique which forms the basis of several recent theorems on invariant curves and regions, published independently by Arnol'd<sup>14/</sup> and Moser<sup>15/</sup>. Moser's theorem is somewhat more general, and consequently we shall discuss his work. He first considers the annular region in the  $r\theta$ -plane defined by

$$0 < a \leq r \leq b \quad (10)$$

and the simple circular twist mapping:

$$\left\{ \begin{array}{l} \theta_1 = \theta_0 + \alpha(r_0) \\ r_1 = r_0 \end{array} \right\} \quad (11)$$

of that annular ring, where the function  $\alpha$  satisfies

$$d\alpha/dr_0 > 0 \quad (12)$$

This last requirement on  $\alpha(r_0)$  means that the amount of the angular twisting

of the mapping always increases with radial distance from the origin of the polar coordinates. (Note that the sense of the twist is opposite to that of Poincaré.) It is clear that this simple mapping preserves circles by simply twisting them around counter-clockwise through an angle  $\alpha(r_0)$ .

Moser next considers a "small perturbation" of this simple circular twist mapping, of the form

$$\left\{ \begin{array}{l} \theta_1 = \theta_0 + \alpha(r_0) + F(r_0, \theta_0) \\ r_1 = r_0 + G(r_0, \theta_0) \end{array} \right\} \quad (13)$$

where  $F$  and  $G$  are assumed to be small and of angular period  $2\pi$  radians. He then proves that for any  $\epsilon > 0$ , no matter how small, the above mapping possesses a smooth, closed invariant curve, lying inside an annular region, of the form

$$\left\{ \begin{array}{l} \theta = \theta^0 + p(\theta^0) \\ r = r_0 + q(\theta^0) \end{array} \right\} \quad (14)$$

where the functions  $p$  and  $q$  also have an angular period of  $2\pi$  radians and possess  $s$  continuous (partial) derivatives satisfying

$$|p|_s + |q|_s < \epsilon. \quad (15)$$

Moser also assumes that the following conditions hold:

- i) Every closed curve near a circle:

$$r = f(\theta) = f(\theta + 2\pi) = \text{a constant} \quad (16)$$

intersects its mapped (image) curve at least twice;

- ii) There exists a constant  $C_0 > 1$  such that

$$C_0^{-1} \leq d\alpha/dr_0 \leq C_0; \quad (17)$$

iii) There exists a  $\delta_0(\epsilon, s, C_0)$  and an  $l(s)$  such that

$$\text{and} \quad \left\{ \begin{array}{l} |F|_0 + |G|_0 < \delta_0 \\ |\alpha|_l + |F|_l + |G|_l < C_0 \end{array} \right\} \quad (18)$$

Finally he asserts that the mapping induced on the curve of equation (14) is given by

$$\theta'_1 = \theta' + \alpha(r_0) \quad . \quad (19)$$

Moser later develops a formula for  $l(s)$  and obtains the large value of  $l=333!$ . Thus the existence of only a (large!) finite number of partial derivatives of  $F, G$ , and  $\alpha$  are necessary (in contrast to the theorems of Arnol'd which require infinitely many derivatives to exist).

Moser's theorem has an important bearing on the stability of solutions of differential equations. For if its necessary conditions for the mapping are met, then the theorem asserts that there exists a smooth, closed invariant curve of the mapping. This statement in turn means that the invariant curve and the interior region it bounds must always be preserved, even after an infinite number of mappings of the region onto itself. Consequently, points which lie initially anywhere inside the invariant region will always map into other points somewhere inside this region and thus will be absolutely stable (i.e. for an infinite time). Points initially outside this region may map anywhere outside, in particular to infinity, so that we can make no general (absolute) predictions about such points.

One method of studying Moser's contentions (and the general question of stability) is to pick a simple nonlinear differential equation which can be integrated in some closed form. We can then use a high-speed digital

computer to investigate its periodic and "almost-periodic" solutions to a high degree of accuracy. This arithmetical approach has been adopted in our research in order to avoid the accuracy limitations of ordinary numerical integration techniques. (We will, however, employ the latter approach for certain preliminary calculations.)

The non-autonomous (time-dependent) equation chosen, the "Cubic equation", arises in connection with the motion of charged particles in accelerators<sup>16/</sup>. It has been studied numerically by Powell and Wright<sup>17/</sup>, and later by Bartlett<sup>13/</sup> using the Illiac I computer. The dimensionless form of this "Cubic equation" is

$$\frac{d^2x}{dt^2} + p(t) x^3 = 0 \quad (20)$$

where  $p(t)$  is here chosen to be a periodic square-wave function of time, with a constant magnitude  $p_0$  and a fixed period  $\tau$ , that is

$$\left\{ \begin{array}{ll} p(t) = p_0 > 0 & \text{for } -\tau/4 < t \leq \tau/4 \\ p(t) = -p_0 & \text{for } \tau/4 < t \leq 3\tau/4 \end{array} \right\} \quad (21)$$

This equation can be easily integrated in a piecewise manner (i.e. over each half-period  $\tau/2$  of the square-wave) to yield piecewise first integrals of the motion. These in turn can then be integrated once more to give a piecewise solution for  $X(t)$  and  $\dot{X}(t)$  in terms of Jacobian elliptic functions. Given initial phase-space coordinates  $(X_1, \dot{X}_1)$  we can thus obtain the coordinates  $(X_1, \dot{X}_1)$  at the end of the first quarter-period of the square-wave. Using these coordinates now as initial coordinates, we can calculate the coordinates  $(X_2, \dot{X}_2)$  at the end of the next half-period of motion, and so on in steps of  $\tau/2$  thereafter until we desire to stop. Thus a phase-space trajectory of period  $n\tau$  can be constructed, where  $n$  is

some integer multiple of  $\tau$  (i.e. an integral number of square-wave cycles). This trajectory, then, corresponds to a mapping  $T^n(\tau)$  of the phase-space onto itself. We shall study the properties of this mapping and its multiples. A brief analysis of the possible solutions to equation (20) together with detailed methods for numerically calculating the piecewise trajectories above (i.e. the mappings  $T^n(\tau)$ ) are given in the first appendix.

The object of this research will therefore be to analyze experimentally (i.e. by using computers) the mathematical contents of Moser's theorem for the above "Cubic equation". If we can assure ourselves of extreme precision in these computer calculations, then we ought to be able to locate and determine the precise shape and size of the closed invariant curves and regions, if they exist. It is hoped that these experimental methods will then be applicable to any nonlinear system of equations. If so, then a practical numerical understanding of the absolute stability of solutions for any particular nonlinear differential equations can be obtained easily and quickly.

## II. CALCULATIONS

### A. Exploratory mapping of the phase-plane

The analysis in the first appendix shows that we can utilize a digital computer for the rapid and accurate calculation of any desired solution of the original equation (20). That is, for any given initial coordinates  $(x_i, \dot{x}_i)$  in phase-space we can calculate the subsequent trajectory for as many periods of the periodic function  $p(t)$  as we desire. The accuracy of these solutions will be limited chiefly by the capacity of the computer. A program (listed for reference in the second appendix) has been written in the machine language (NICAP) of the University of Illinois' Illiac II computer in order to obtain maximum (12-place) accuracy. (The Illiac II computer utilizes 52-bit words in the form of four quarter-words, and has a speed of up to twice that of the IBM-7094 computer.)

[It has also been convenient to write a simpler Fortran program (also included in the second appendix) for the direct numerical integration of equation (20) using a modified Runge-Kutta-Gill integration routine. This program, run on both the IBM-7094 and the Illiac II, has been highly useful for general exploratory investigations of solutions. Although the calculations made with this program are limited to about 6-place accuracy, they are helpful in ascertaining just which solutions require the highly accurate NICAP program. The program has the additional virtues of being short, simple to write, and compatible with IBM computers. Time-wise, however, the Fortran program is about 3 to 4 times slower than the NICAP program, in spite of the latter program's length and complexity (1250 NICAP instructions compared with 110 Fortran statements). On the average the NICAP program can calculate about 2200 periods of  $p(t)$  per minute of

computer time, while the Fortran program can compute only about 660 periods of  $p(t)$  per minute.]

An exploratory scan of solutions (in the phase-plane) was made on the Illiac I computer by Bartlett, who found that all fixed points in the phase-plane are either on the two coordinate axes or else are maps of these points<sup>13/</sup>. A list of fixed points on the two axes for  $n=2,3,\dots,12$  is reproduced in Table 1 for reference. These values are accurate only to about 3 decimal places (which is adequate for initial studies). As  $n$  gets larger the corresponding points fixed under  $T^n(\tau)$  lie closer to the origin and closer to one another. We shall see that around any of these  $T^n(\tau)$  fixed points there lie clusters of "satellite" multiple points fixed under the transformation  $T^{mn}(\tau)$ , where  $m$  is an integer.

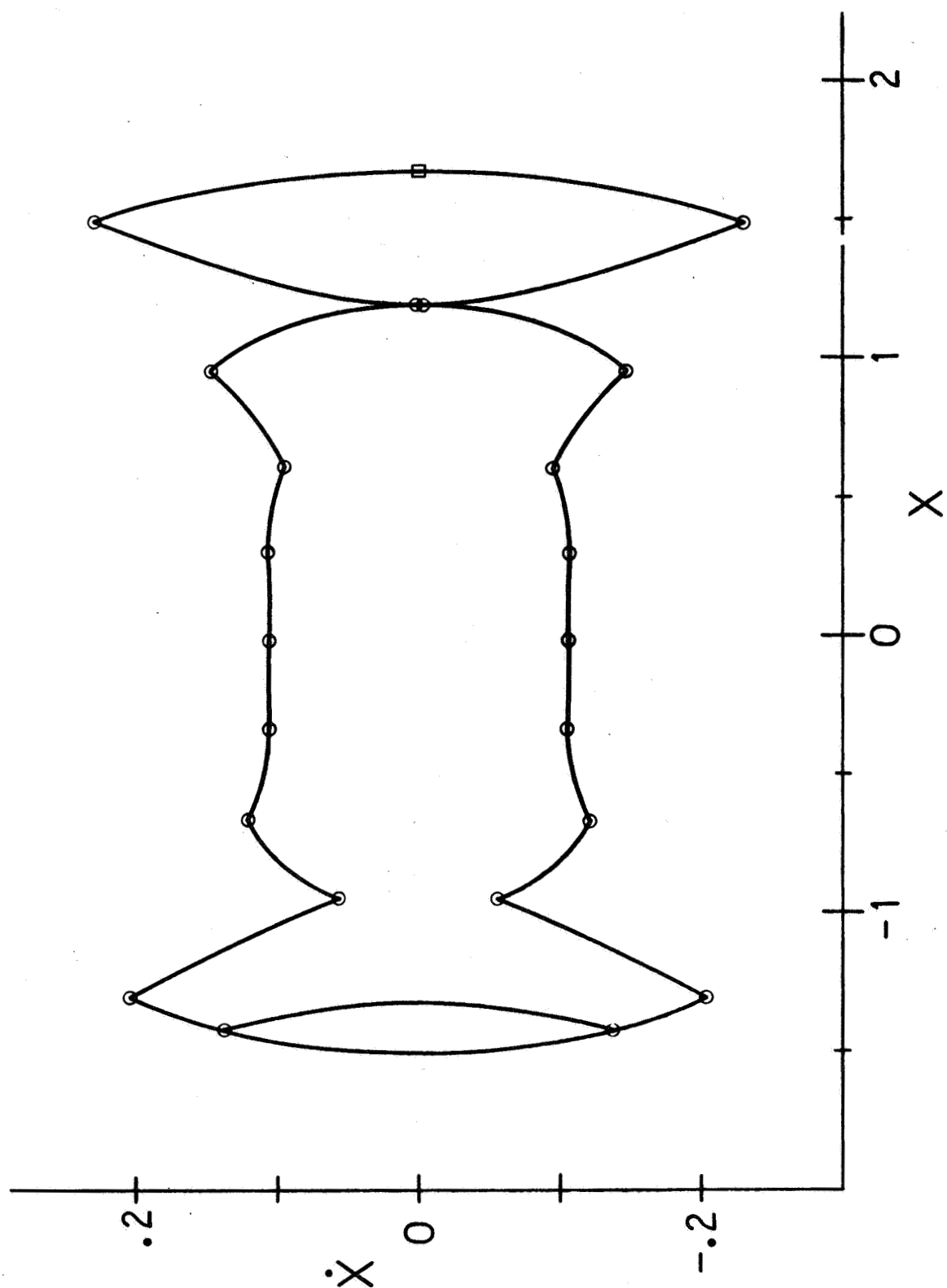
Knowing the location of the simpler fixed points, we can proceed to investigate some of the properties of their associated solutions. Preliminary repeated mappings under  $T^{11}(\tau)$  of points initially inside a relatively large neighborhood of the  $T^{11}$  elliptic fixed point on the +X-axis (at about  $X_1 = 1.667$ ) showed that the consequent (mapped) points generally return repeatedly to the same initial neighborhood. Few escape to infinity or wander very far from the initial neighborhood. Other  $T^n$  elliptic fixed points were similarly scanned, but none exhibited such an apparent "stability" of the repeated mappings. Thus the remainder of this work will concern itself with a detailed analysis of the region near the  $T^{11}$  elliptic fixed point. A plot of the periodic orbit in phase-space that is associated with this  $T^{11}$  fixed point is shown in Figure 1. The (22) encircled points mark the end of each half-period arc of motion, while the  $T^{11}$  fixed point on the +X-axis (the initial conditions for the orbit) is enclosed by a small square.

Table 1

Approximate axis locations of some fixed points  
of the mapping  $T^n(\tau)$ , for  $n=2,3,\dots,12$ .

<u>n</u>	<u>X-axis</u>	<u>X-axis</u>
2	3.4143	9.6012
3	3.1823	0.4591
4	2.387	0.3775
5	2.338	0.3579
6	2.1086	0.3088
7	2.058	0.1639
8	1.9106	0.1545
9	1.8506	0.1464
10	1.736	0.1259
11	1.667	0.1021
12	1.5777	0.09339

Figure 1. Periodic phase-space trajectory associated with the  $T^{11}$  elliptic fixed point on the X-axis.

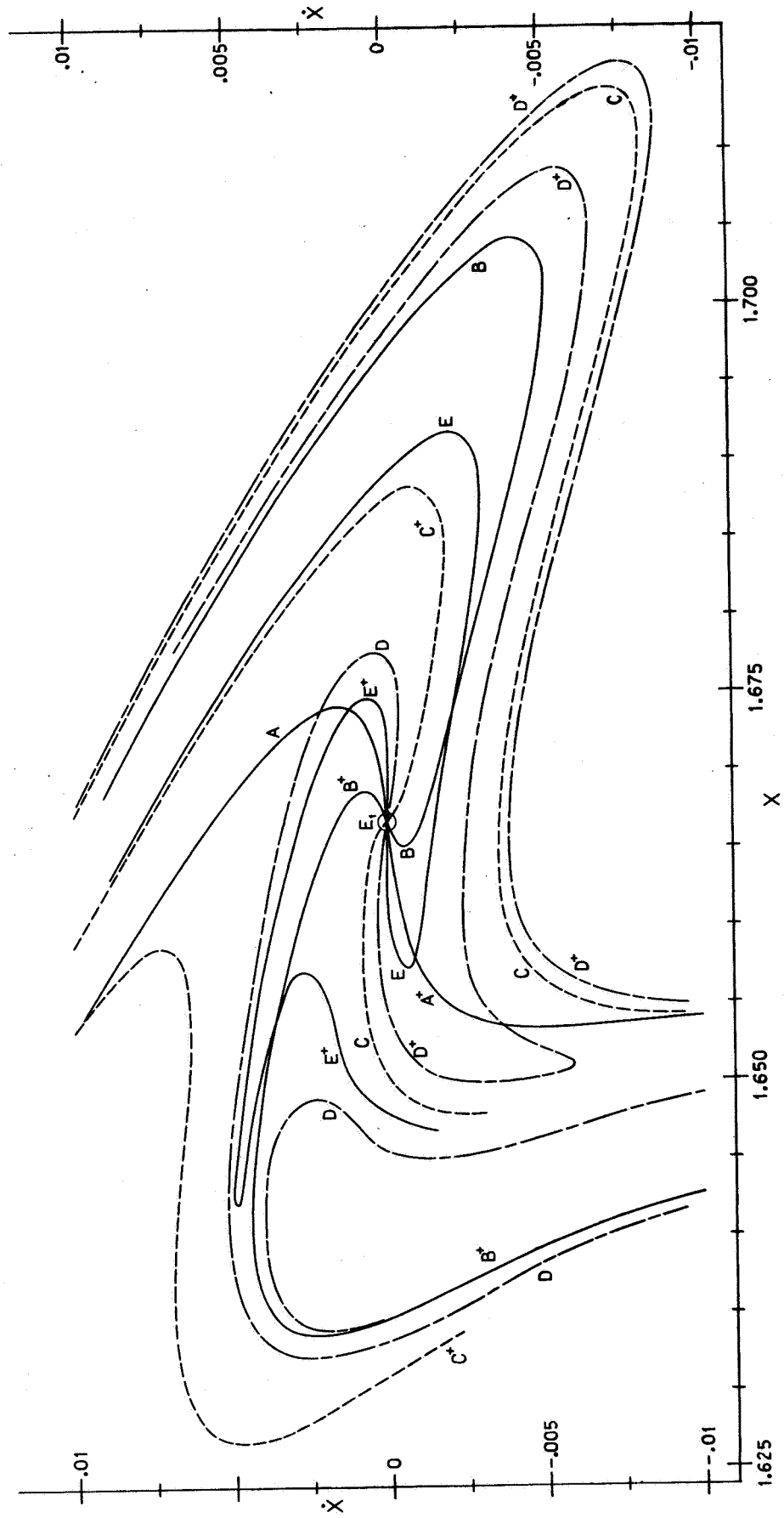


The X-axis near the point  $X_1 = 1.667$  was mapped for the multiples  $m=1,2,3,4$ , and 5 of the basic period  $n\tau=11\tau$ , and the results of the Runge-Kutta method are shown in Figure 2. Note the striking "twisted" character of these mappings, a feature that will play an increasingly important role. In the figure,  $E_1$  marks the  $T^{11}$  fixed point, and the (+) superscripts on the letters A,B,C,D, and E indicate the curve is a map of the +X-axis for  $m=1,2,3,4$ , and 5 respectively. Each successive mapping is a twisted (rotated) version of the previous one, so that even by  $m=3$  the plots become intricately twisted around each other. Sections get thrown outward and then return to intersect the axes again, more often the greater  $m$  is.

Yet a great deal of useful information can be obtained from this exploratory plot. First of all we observe that something in the structure of our nonlinear system is creating two wide "escape avenues" for points around the  $T^{11}$  fixed point. These begin in the neighborhoods of (1.66, 0.01) and (1.65, -0.01) and carry the mapped axis curves far away from the fixed point  $E_1$ . Already we have signs of what we would intuitively call "instability".

The other pieces of valuable information are the intersections of the various mapped axis curves with the X-axis and with each other. Curves  $E^+$  and  $C^+$  intersect the +X-axis and reveal fixed points for the  $T^{55}$  and  $T^{33}$  mappings respectively. Similarly E and C give fixed points for  $T^{55}$  and  $T^{33}$  on the -X-axis. Curves  $D^+$  and  $B^+$  intersect the -X-axis and locate points fixed under  $T^{88}$  and  $T^{44}$  respectively; the same applies for D and B intersecting the +X-axis. In addition, these curves cross each other. Curve  $A^+$  intersects  $D^+$  and so reveals a point fixed under  $T^{44-11} = T^{33}$ . Other  $T^{33}$  points arise from the following intersections: A with D,  $B^+$  with  $E^+$ ,

Figure 2.  $T^{11m}$  twist mappings of the X-axis around  
the  $T^{11}$  elliptic fixed point, for  
 $m=1,2,3,4$ , and 5.



and B with E. The intersections of A with  $E^+$  and  $A^+$  with E yield points fixed under  $T^{88}$ .

Already we see some of the multiple (fixed) point structure of the  $T^{11}$  fixed point revealed to us in this simple scan. So far we have the approximate coordinates of 6 points fixed under  $T^{33}$ , two under  $T^{44}$ , two under  $T^{55}$ , and 4 under  $T^{88}$ . The 6 points fixed under  $T^{33} = T^{3n}$  (for  $n=11$ ) represent the totality of  $T^{3n}$  points which have been found in the cluster around the  $T^n$  fixed point. In simple cases we can expect  $2m$  multiple fixed points of the mapping  $T^{mn}$  in a cluster around the  $T^n$  fixed point. Of these,  $m$  will be elliptic fixed points and  $m$  will be hyperbolic points, alternately spaced around the  $T^n$  fixed point.

#### B. Accurate mapping of the phase-plane

A finer scan, using the more accurate NICAP program for the Illiac II, resulted in accurate fixed point values for the  $T^{33}$  and  $T^{55}$  mappings. These values are listed in Table 2, along with the results of a later determination of the type of point (i.e. e or E stand for elliptic and h or H for hyperbolic-type fixed points). The points belonging to an e (or h) set for  $T^{33}$  and  $T^{55}$  are listed as the successive  $T^{11}$  transforms of the initial point on the X-axis. Approximate values for the two hyperbolic fixed points adjacent to the  $T^{11}$  elliptic fixed point on the X-axis are taken from the work of Bartlett<sup>13/</sup>.

To investigate the type of motion for points initially in a neighborhood of the basic  $T^n$  fixed point or its  $T^{mn}$  multiple points we can simply have the computer print out and plot the repeatedly mapped coordinates at intervals of  $n$  or  $mn$  periods. We can then draw a smooth curve through the locus of these successively mapped points and call this a

Table 2

Exact locations and designations of fixed points of the mapping  $T^{mn}$ , for  $n=11$  and  $m=1,3,5$ .

<u>n</u>	<u>X</u>	<u>X</u>	<u>Designation</u>
11	1.666735430	0.0	$E_1$
11	1.6395	0.01690	$H_1$
11	1.6395	- 0.01690	$H_2$
33	1.687355697	0.0	$e_1$
33	1.653501571	- 3.641919791 $\times 10^{-3}$	$e_2$
33	1.653501570	+ 3.641919722 $\times 10^{-3}$	$e_3$
33	1.650033472	0.0	$h_3$
33	1.673820978	+ 2.109859576 $\times 10^{-3}$	$h_1$
33	1.673820976	- 2.109859583 $\times 10^{-3}$	$h_2$
55	1.672374290	0.0	$e(1)$
55	1.661880306	- 3.119288766 $\times 10^{-4}$	$e(3)$
55	1.668547852	+ 5.276790971 $\times 10^{-4}$	$e(5)$
55	1.668547852	- 5.276791001 $\times 10^{-4}$	$e(2)$
55	1.661880306	+ 3.119288870 $\times 10^{-4}$	$e(4)$
55	1.660868636	0.0	$h(4)$
55	1.671135787	+ 3.399519824 $\times 10^{-4}$	$h(1)$
55	1.665062836	- 4.802075374 $\times 10^{-4}$	$h(3)$
55	1.665062836	+ 4.802075326 $\times 10^{-4}$	$h(5)$
55	1.671135787	- 3.399519730 $\times 10^{-4}$	$h(2)$

stroboscopic mapping of the initial point. This curve will be approximately elliptical and closed if the initial point is "close" to an elliptic fixed point. If, on the other hand, the fixed point is hyperbolic, then this curve will depart rapidly from it. In either case the resulting stroboscopic curve is an invariant one, since any point on it will map into some other point on the same curve (its "consequent") under  $T^n$ . Applying this technique to points near the  $T^{33}$  and  $T^{55}$  fixed points, we find that for  $m = 3, 5$  the  $T^{11m}$  elliptic or hyperbolic point on the  $+X$ -axis lies respectively to the right or left of the original  $T^{11}$  elliptic fixed point. These results are listed in Table 2.

As we move away from the nearly elliptical motions close to an elliptic fixed point (such as  $e_1$  of the  $T^{33}$  mapping), we find that many of the points sooner or later jump into regions which contain oscillating hyperbolic invariant curves and are then quickly lost to infinity. In making a stroboscopic scan of the  $X$ -axis, with  $\theta_1 = 0$  and  $r_1 \leq 0.040$ , where  $r_1$  is now the initial radial distance from the  $T^{11}$  elliptic fixed point  $E_1$  on the  $X$ -axis, we find that this sort of jumping behavior persists for all elliptic regions down to about  $r_1 = 0.0090$ . Now, as we approach the  $T^{55}$  elliptic point  $e(1)$  on the  $X$ -axis, the motion becomes approximately cyclic every 5 multiples of the base period ( $11\tau$ ) of the mappings.

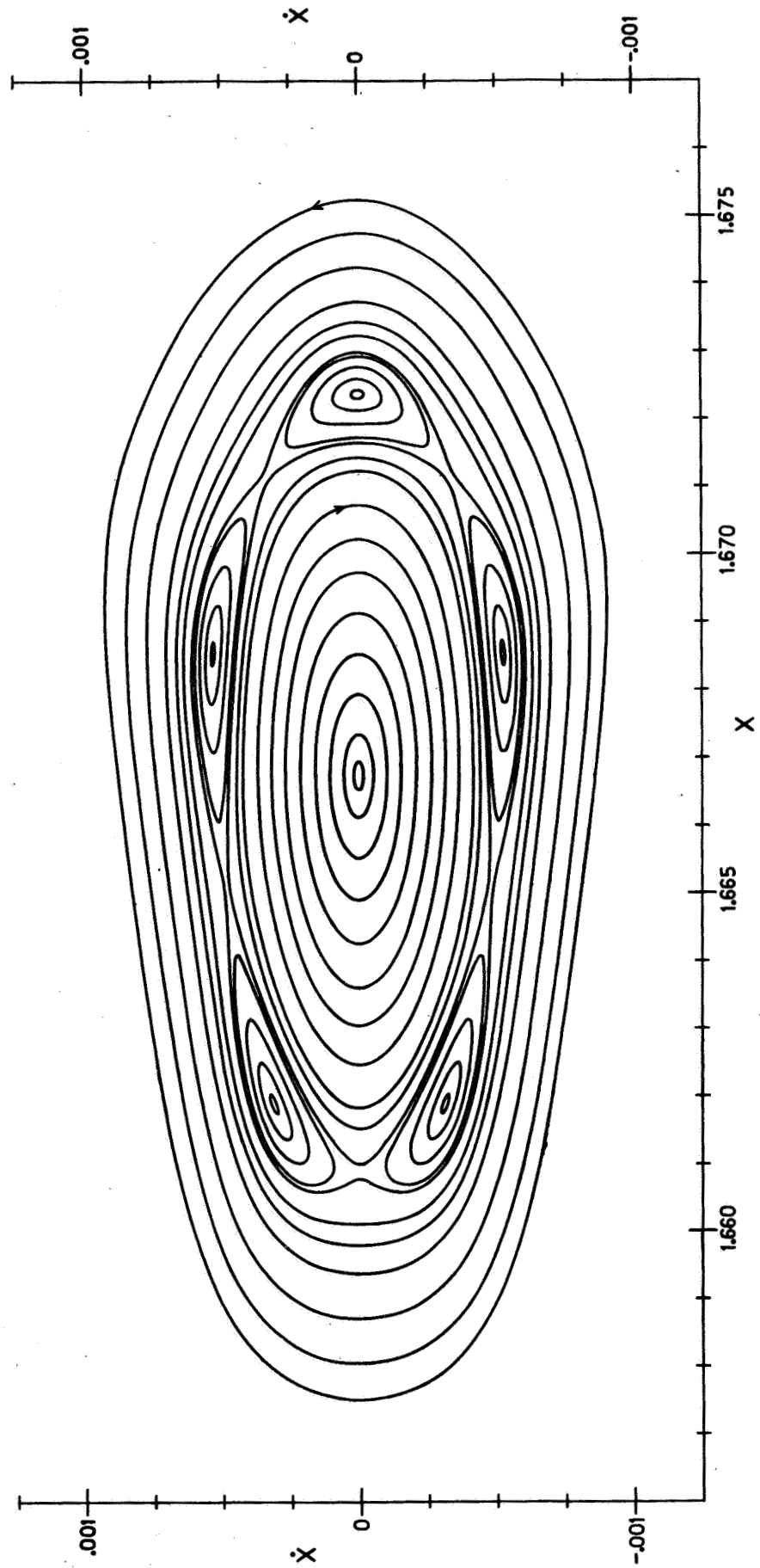
Plotting every fifth point we obtain counter-clockwise stroboscopic motion around the  $T^{11}$  elliptic fixed point  $E_1$ . The periods for the completion of one stroboscopic revolution increase to infinite values as we cross the outer invariant curves of the adjacent  $T^{55}$  hyperbolic points. Inside these invariant curves the periods decrease to zero when we arrive at the  $T^{55}$  elliptic point  $e(1)$  itself. As we move to the left of this

point, they quickly get large again, become infinite as we cross the inner invariant curves of the adjacent  $T^{55}$  hyperbolic points, and then start decreasing.

These approximately closed, elliptical stroboscopic motions under  $T^{55}$  continue in a clockwise sense now, no matter how close we go toward  $E_1$  (at least down to the distance  $r_i = 10^{-8}$ )! There seem to be no further hyperbolic "leakages" to infinity, even though we are crossing an increasingly denser set of clusters of alternating fixed points with multiplicities greater than 5. So it seems that we have found an invariant region of stability. For, once inside the invariant curves from the  $T^{55}$  hyperbolic points, at about  $r_i = 0.005$ , the endpoints of the  $T^{11}$  trajectories always remain inside this region, performing approximately closed, elliptical motions about  $E_1$ . A plot of these stroboscopic motions for initial X-axis points in the range  $0.0002 \leq r_i \leq 0.0090$  (for  $\theta_i = 0$  again) is shown in Figure 3. The five "islands" of locally elliptical motion under  $T^{55}$  separated by five "channels" of locally hyperbolic motion under  $T^{55}$  show up clearly, as does the transition from a counter-clockwise to a clockwise sense of stroboscopic revolutions. (As we move radially outward from the point  $E_1$  the periods for the plotted stroboscopic revolutions increase from 120 to 1040 multiples of  $11\tau$  for the inner clockwise motions; then from 1060 to 200 and back to 1060 multiples for the "island" regions; and finally from 1140 to 90 multiples for the outer counter-clockwise motions shown.)

Thus the motion near the  $T^{55}$  cluster is "stable" while the motion near the  $T^{33}$  cluster is "unstable". We need to explain this difference in behavior. As will be seen later, the invariant curves from the

Figure 3. Stroboscopic trajectories of the mapping  $T^{55}$  around the  $T^{11}$  elliptic fixed point.



hyperbolic points of  $T^{33}$  have the general oscillation property (first demonstrated by Bartlett<sup>13/</sup>). There do not appear to be any such oscillations of the  $T^{55}$  hyperbolic invariant curves in Figure 3. These hyperbolic invariant curves instead seem to join one another smoothly as they envelope the well-defined elliptic "islands". Perhaps, then, instability goes hand in hand with these oscillations, while stability is a result of the lack of such oscillations.

### C. The mapping $T^{11}$

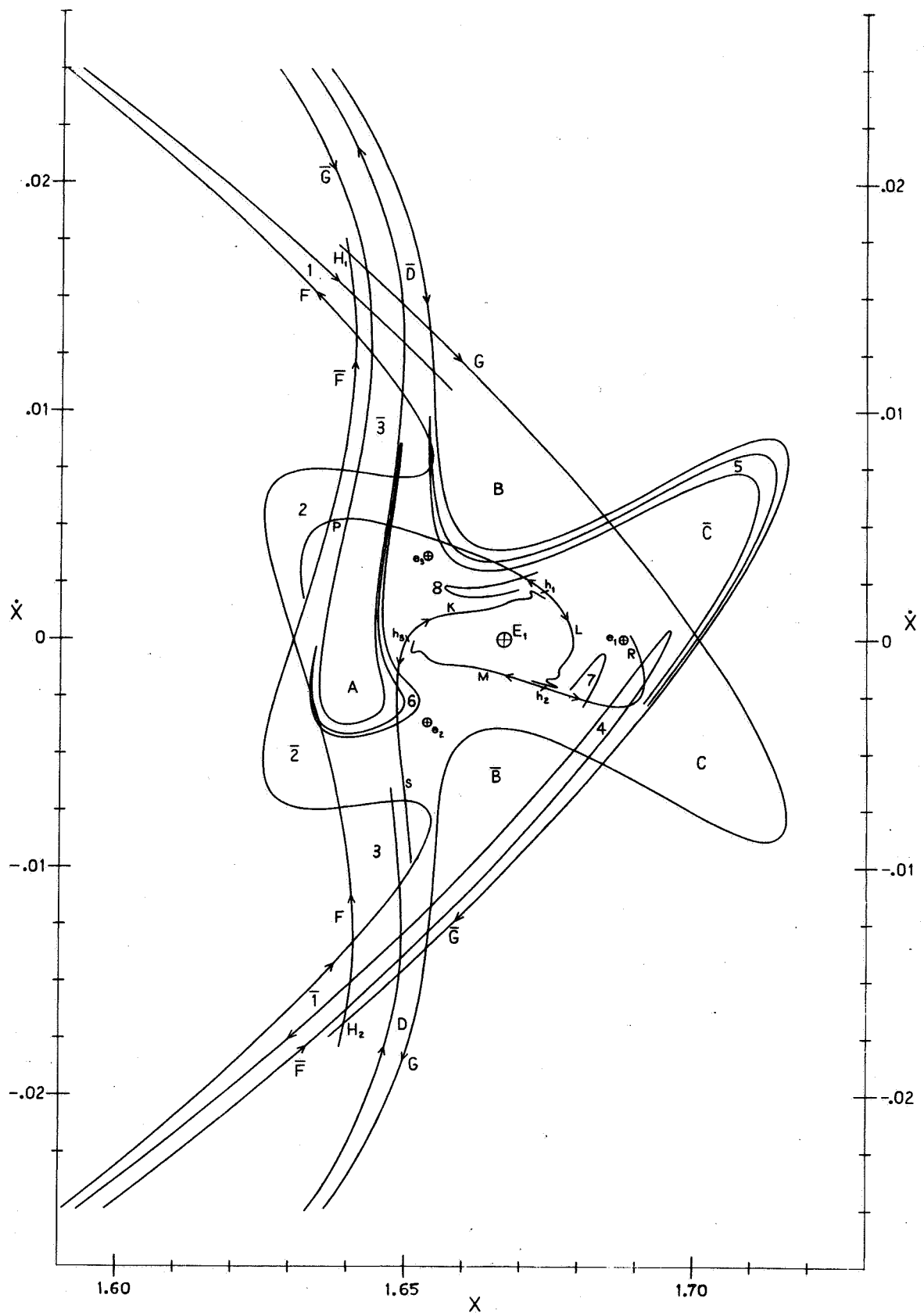
In order to see how the above apparently stable inner region coexists with the unstable outer regions, let us study the mappings for  $T^{11}$  in a region around the elliptic fixed point  $E_1$  and large enough to include the two adjacent hyperbolic points  $H_1$  and  $H_2$ . The invariant curves issuing from one hyperbolic point develop larger and larger oscillations as their overall motion carries them towards the nearest hyperbolic fixed points. This should hold true for both the inner and outer invariant curves, where we define "inner" as closer to the origin of our  $XX$  coordinate system, and "outer" as farther away from that origin. In our case we find that the inner invariant curves move from  $H_2$  toward  $H_1$  with increasing oscillations, while the outer invariant curves move from  $H_1$  to  $H_2$ . (A plot of similar hyperbolic invariant curves for the mapping  $T^{12}(\tau)$  is shown in Figure 3 of reference 13).

To obtain these curves we need only pick initial points rather close to  $H_1$  or  $H_2$  and then map them repeatedly. From earlier discussions of motion near hyperbolic points we know that these points will move rapidly towards the invariant curves and be asymptotic to them. The initial points can be taken to be on a segment of the straight line from  $E_1$  to  $H_2$ . At the end

of 11 periods the series of successive points will have been mapped into a new series of points forming the locus of the first mapped curve. This can be repeated for as many multiples  $m$  as desired, thus obtaining  $m$  mapped curves of the original straight line. If the range of initial radii is too small, these  $m$  curves will be disjoint. By choosing the initial interval just long enough, the  $m$  mapped curves will just overlap each other. We call one such barely overlapping mapped curve "one full oscillation" of the invariant curve. (Actually they will miss each other very slightly because two successive points on the initial straight line are at slightly different distances from the adjacent hyperbolic invariant curves.) These barely overlapping  $m$  curves can then be joined together smoothly to give a fairly accurate plot of the outgoing hyperbolic invariant curves.

In our calculations we use a section of the line between  $E_1$  and  $H_2$ , at an angle of  $211.85^\circ$  to the  $+X$ -axis, with radii in the range  $0.0300 \leq r_1 \leq 0.0316$ . Initially the scan can be in large increments of the  $r_1$  values (such as  $5 \times 10^{-5}$  or  $10^{-5}$ ), but eventually it becomes necessary to go to a finer grid with increments of  $2 \times 10^{-6}$  in  $r_1$ . The results of these mappings are plotted in Figure 4 as curve F, shown issuing from the point  $H_2$  and moving upwards toward  $H_1$ . This is the "inner" invariant curve, as it passes to the left of  $E_1$  (i.e. closer to the  $XX$  origin). The outer invariant curve, plotted as curve G in Figure 4, is obtained by using a section of the line between  $H_1$  and the point  $(1.716, 0.0)$ . The angle of this line is  $167.54^\circ$  to the  $+X$ -axis, and the interval of initial radii covers the range  $0.0750 \leq r_1 \leq 0.0780$  in increments as small as  $2 \times 10^{-6}$ . Figure 4 also contains the reflections in the  $X$ -axis of curves F and G, labeled  $\bar{F}$  and  $\bar{G}$  respectively. From the time-reversal invariance of our

Figure 4. Hyperbolic invariant curves of the mapping  $T^{11}$  and their relations to those of the mapping  $T^{33}$ .



original equation we see that these curves are solutions not only for the motion with time running backwards, but also for the reversed motion originating at an infinite negative time and continuing with time running forwards toward  $t=0$ . Thus, for example, curve  $\bar{F}$  shows a series of oscillations moving away from  $H_2$  and becoming smoother as it approaches and intersects  $H_1$  (as an inner invariant curve). Similarly  $\bar{G}$  smooths out its oscillations as it moves away from  $H_1$  to  $H_2$  as an outer invariant curve.

Because of the "wild" oscillations in these curves, the "forward-going" invariant curves ( $F$  and  $G$ ) do not join smoothly with their own reflections ( $\bar{F}$  and  $\bar{G}$ ), the "forward-coming" invariant curves. Instead they intersect one another ( $F$  with  $\bar{F}$ , and  $G$  with  $\bar{G}$ ) an infinite number of times. The regions between successive intersections of  $F$  with  $\bar{F}$  are labeled  $1, 2, 3, \dots, 8$  when  $F$  is on the outside (with respect to  $E_1$ ) of  $\bar{F}$ , and  $\bar{1}, \bar{2}, \bar{3}$  when  $F$  is on the inside of  $\bar{F}$  (these are reflections in the  $X$ -axis of the first set of areas). The regions between successive intersections of  $G$  with  $\bar{G}$  are labeled  $A, B, C$ , and  $D$  when  $G$  is on the outside of  $\bar{G}$ ; and  $\bar{A}, \bar{B}, \bar{C}$ , and  $\bar{D}$  when  $G$  is on the inside of  $\bar{G}$ . By "inside" we now mean closer to  $E_1$ , while "outside" means farther away from  $E_1$ .

As the oscillations of the invariant curve  $F$  issuing from  $H_2$  crowd in towards  $H_1$ , they become long and thin to conserve their phase-space area. By noting the directional arrows on the various sections of the curves  $F$  and  $\bar{F}$  in Figure 4 we can see that all the points in region  $\bar{1}$  will map into all of  $\bar{2}$ , then into all of  $\bar{3}$ , and so on into the inside of the hyperbolic invariant curves. But these points in  $\bar{1}$  come from a longer, thinner oscillation  $\bar{0}$  (not shown) that is closer to that incoming hyperbolic invariant curve which intersects  $H_2$  from a direction exactly

opposite to the incoming curve  $\overline{G}$ . That is to say, these points may come from regions quite distant from  $E_1$  and yet still get in close to  $E_1$ . Similarly we can see the fate of points within the "inner" region 6. They will map successively into regions 5,4,3,2, and 1 and then into a longer, thinner region which is closer to the outgoing hyperbolic invariant curve through  $H_1$  (i.e. into a region which is a simple reflection of region  $\overline{0}$  above, the "antecedent" of region  $\overline{1}$ ). Thus these points initially within an inside region will soon escape to regions quite distant from their original region. In fact, as Bartlett has shown<sup>13/</sup>, these oscillating curves around the  $T^{11}$  fixed point  $E_1$  will intersect similar invariant curves of the  $T^{10}$  mapping, so that some of the points in the region between these two sets of intersecting curves can be carried even further outward, eventually reaching infinity. The time-reversed motions will analogously bring points in from infinity to a region close to  $E_1$ .

This same discussion applies to the regions formed by the intersections of curves  $G$  and  $\overline{G}$ . Thus points from infinity can move into region  $\overline{D}$  and then map into  $\overline{C}, \overline{B}, \overline{A}$  (not shown), and so on into the inside of the invariant curves. Similarly, inside points within region  $A$  (or its unplotted antecedents) will map successively into regions  $B, C, D$  and then quickly move out directly to intersect the  $T^{10}$  hyperbolic oscillations. This is easy to visualize because region  $D$  lies outside that incoming hyperbolic invariant curve to  $H_2$  which forms an outer bound to the (incoming) oscillations of curve  $\overline{F}$ . This shows indirectly that the outgoing oscillations of curve  $F$  also reach infinity via intersections with  $T^{10}$  oscillations. To see why, consider the reflected curve  $\overline{F}$  coming in from the same distant regions reached by curve  $F$ . Points inside  $\overline{1}$  will

map into  $\overline{2}, \overline{3}$ , and so on. But notice that  $\overline{3}$  intersects region A (as  $\overline{4}$  will intersect B), so that points in the mutual intersection of regions  $\overline{3}$  and A will map into B, then C, D, and so on out directly to the  $T^{10}$  oscillations. This may sound like a tautology, but nonetheless the observed facts are that points do escape from inner regions around  $E_1$  to infinity.

We now have a mechanism whereby points may escape to infinity even when relatively close to an elliptic fixed point. Going back to Figure 4 we see that the inverse (time-reversed) maps of region 2 (i.e. regions 3, 4, 5, 6, ...) should eventually fill a great deal of the space inside the inner and outer invariant curves because of their property of preserving phase-space areas. In principle we should be able to keep mapping region 2 backward in time or  $\overline{2}$  forward in time, which in reality is how we actually make the maps (it is more convenient here to plot their reflections). Then in the limit perhaps we would end up with some envelope of the successively mapped regions which would coincide with some invariant inner region of stability.

Unfortunately we cannot do this in practice, for the mapped regions get progressively thinner and more convoluted. Region 5 is the last relatively complete region shown; under a time-reversed (inverse) mapping of  $T^{11}$  the section of region 5 lying inside region  $\overline{C}$  gets "squirted" out to region  $\overline{D}$ , the antecedent of  $\overline{C}$ , and only the tip of the mapped region 6 re-enters the inside of the invariant curves. For the next few inverse maps we can only plot very small pieces of the mapped regions such as regions 7 and 8 in Figure 4. These small "hairpin" regions correspond to just a short segment of the inside portion of region 6 in the neighborhood of the bump near the number 6! We will therefore have to resort to some other approach in order to continue the investigation.

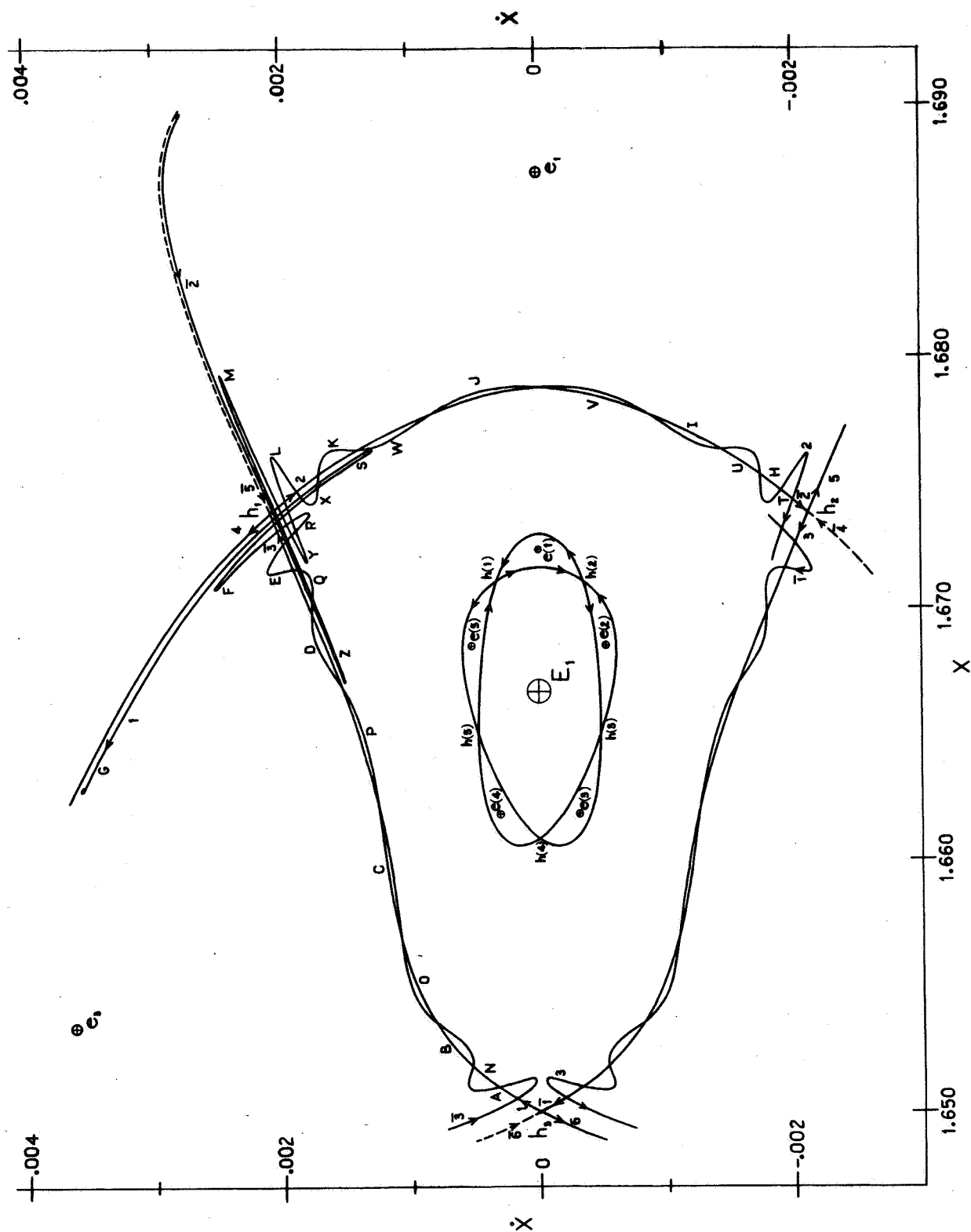
#### D. The mapping $T^{33}$

Since a knowledge of the oscillating invariant curves for  $T^{11}$  has helped us to explain the escape mechanism, let us next find these curves for  $T^{33}$ . (Referring to Figures 4 and 2 we note that the  $T^{11}$  oscillations already discussed have penetrated closer to  $E_1$  than the cluster of  $T^{44}$  points, so we may bypass them and instead study the next cluster of points closer to  $E_1$ .) For this purpose we draw the line between  $e_1$  and  $h_1$  at an angle of  $171.14^\circ$  to the X-axis, and scan the line segment with initial radii in the range  $0.0127 \leq r_1 \leq 0.0135$  using increments of about  $5 \times 10^{-6}$ . Successive mappings then yield the hyperbolic invariant curves for  $T^{33}$ .

The invariant curve issuing from  $h_1$  develops oscillations as it moves toward  $h_2$ , and is partially plotted as curve 2 in Figure 5. This curve 2 is not an inner invariant curve with respect to  $e_1$  (i.e. is closer to our  $E_1$  origin than is  $e_1$ ), with its motion counter-clockwise about  $e_1$ . On the other hand the  $T^{11}$  invariant curve from  $H_1$  is an outer invariant curve with respect to  $E_1$  (i.e. is farther from the  $X\text{-}\dot{X}$  origin than is  $E_1$ ), with clockwise motion about  $E_1$ .

By getting computer printouts of every  $11\tau$  periods we can easily obtain the  $T^{11}$  and  $T^{22}$  maps of curve 2 by plotting the points for  $t=11\tau, 44\tau, 77\tau, \dots$  and  $t=22\tau, 55\tau, 88\tau, \dots$  (since curve 2 comes from points at  $t=0, 33\tau, 66\tau, \dots$ ). The  $T^{11}$  transform of curve 2 is plotted as curve 3, while the  $T^{22}$  transform of curve 2 is shown as the more detailed curve 1. These curves and their mirror reflections in the X-axis (labeled  $\bar{1}, \bar{2}$ , and  $\bar{3}$ ) are all the "inner"  $T^{33}$  hyperbolic invariant curves (i.e. they lie closer to  $E_1$  than  $e_1, e_2$ , or  $e_3$  and so we designate them as "inner" with respect to  $E_1$ ). To obtain the "outer" invariant curves (again outer with

Figure 5. Hyperbolic invariant curves of the mapping  $T^{33}$  and their relations to those of the mapping  $T^{55}$ .



respect to  $E_1$ , as they extend farther from  $E_1$  than  $e_1, e_2$ , or  $e_3$ ), the easiest method is to pick several points along a short segment of the "inner" outgoing curve 1 which is very close to the  $T^{33}$  hyperbolic point  $h_3$ . Extending this short segment of curve 1 through  $h_3$ , we can then map points on this new extended straight-line segment of the "outer" outgoing invariant curve. The results are partially plotted as curve 6 in Figure 5. The  $T^{11}$  and  $T^{22}$  maps of this invariant curve are plotted as curves 4 and 5 respectively. The reflections of these 3 "outer" outgoing invariant curves are plotted as curves  $\bar{4}, \bar{5}$ , and  $\bar{6}$ , and complete the picture for all the  $T^{33}$  "outer" invariant curves (about  $E_1$ ).

We shall discuss these intersecting oscillating invariant curves by comparing them to the picture for the  $T^{11}$  mapping. Looking at the intersections of curves 1 and  $\bar{3}$  we can predict that points initially between the invariant curves in region A will map successively into region B, C, D, E, F, G, ... outside the invariant curves along outgoing invariant curve 4. Similarly exterior points in region N will map into regions O, P, Q, R, S, ... and become interior points. The same phenomena occur for the intersections of curves 2 and  $\bar{2}$  or 3 and  $\bar{1}$ . For example, for the former, exterior points in region M (near the incoming invariant curve  $\bar{5}$ ) map into interior points in region H via the regions L, K, J, and I. Interior points in region Y become exterior points in region T via regions X, W, V, and U. Here we are again using "interior" to mean closer to  $E_1$ , and "exterior" to mean farther away from  $E_1$ . Thus the term "exterior" here would be called "interior" (to the regions around  $e_1, e_2$ , or  $e_3$ ) if we were investigating the stability of  $e_1, e_2$ , or  $e_3$  rather than  $E_1$ . (In effect we have already investigated  $e_1$  by observing that the  $T^{11}$  oscillations reach in close to it eventually.)

But now there are additional complications. Not only, for example, do curves 2 and  $\bar{2}$  intersect one another, but their (inner) regions of mutual intersection are themselves intersected by curves  $\bar{1}$  and 1 (issuing to and from the nearby  $T^{33}$  hyperbolic point  $h_3$ )! In the case of the  $T^{11}$  oscillations the nearest hyperbolic points to  $H_1$  and  $H_2$  are a good distance away along the cluster of 22 alternating  $T^{11}$  fixed points surrounding the  $XX$  origin. For  $T^{33}$  the nearest hyperbolic point to  $h_1$  and  $h_2$  is just the equivalent satellite multiple point  $h_3$ . This same pattern of mutual intersection of neighboring inner oscillations similarly occurs when curves 2 and  $\bar{3}$  intersect the regions bounded by curves 3 and  $\bar{1}$ , or when curves 3 and  $\bar{2}$  intersect the regions bounded by curves 1 and  $\bar{3}$ .

There are several important consequences of this mutual intersection property of the inner oscillations. One is that points from the 3 outer regions (bounded by the outer invariant curves) can get into the inner region, and vice versa. Consequently points from the outer region surrounding  $e_3$ , for example, can reach the inner region around  $E_1$  (via curve  $\bar{3}$ ) and then pass out to the outer region surrounding  $e_1$  (via curve 2). These points might eventually reach the region of  $e_2$ , or even  $e_3$  again by the same device. A second result is that the areas of the outer oscillations are considerably smaller (in proportion to the total "outer" area) than those for the  $T^{11}$  mapping. It thus seems unlikely that these oscillations will fill up the space around, for example,  $e_1$  to the extent that the  $T^{11}$  oscillations did around  $E_1$ . But these outer oscillations have the same area as the inner oscillations, and so we are led to suspect empirically that the inner region around  $E_1$  will be relatively free of inner oscillations from the  $T^{33}$  mapping crowding in toward  $E_1$  (as they did for the  $T^{11}$  maps).

A third consequence is that some points will move around the periphery of the inner region for several entire stroboscopic revolutions before escaping to the outer regions. To see this let us consider points inside area Q of curve 1. Those points inside the rectangular area which is the mathematical intersection of regions Q and Z will map into the transforms of this area, i.e. the regions  $R \cap Y$ ,  $S \cap X$ , and so on. But since these points lie inside region X, they will map subsequently into regions W,V,U,T, and so on into the outer regions. But let us now consider points inside the tip of region Q but not inside  $Q \cap Z$ . These will map into corresponding points in the inner tip of R,S, and so on. But the tip of S does not intersect W, so that the tip of the  $T^{11}$  map of S (call it  $S'$ ) will not intersect region V. (However a portion of  $S'$  will intersect region W, namely the  $T^{11}$  map of the region  $S \cap X$ , since X transforms into W.) The tip of each successive map of  $S'$  will move towards  $h_2$  in this fashion, not intersecting regions U or T (which are transforms of V), and thus will remain on the inside. By the same token the tip will not intersect any of the inner regions I,H,... because these are transforms of regions L,K,J,... which in turn were not intersected by region S or its antecedents. (These calculations were in fact carried out, but are not shown in Figure 5 because of the difficulty of getting enough mapped points for a continuous curve.)

So the successively mapped tip will travel on toward  $h_2$  on the inside, paralleling the inner oscillations of curve 2. Once close to  $h_2$ , the mapped tips enter the domain of tight oscillations of curve  $\bar{1}$  as it moves left towards  $h_3$ . Now there are several alternatives. Since the tip of  $S'$  (too long and thin to be shown) extends below the X-axis at an angle  $\neq 90^\circ$ , this tip of  $S'$  will intersect its own mirror image in the X-axis,  $\bar{S}'$  (which is just the same inner oscillation of curve  $\bar{1}$  as  $S'$

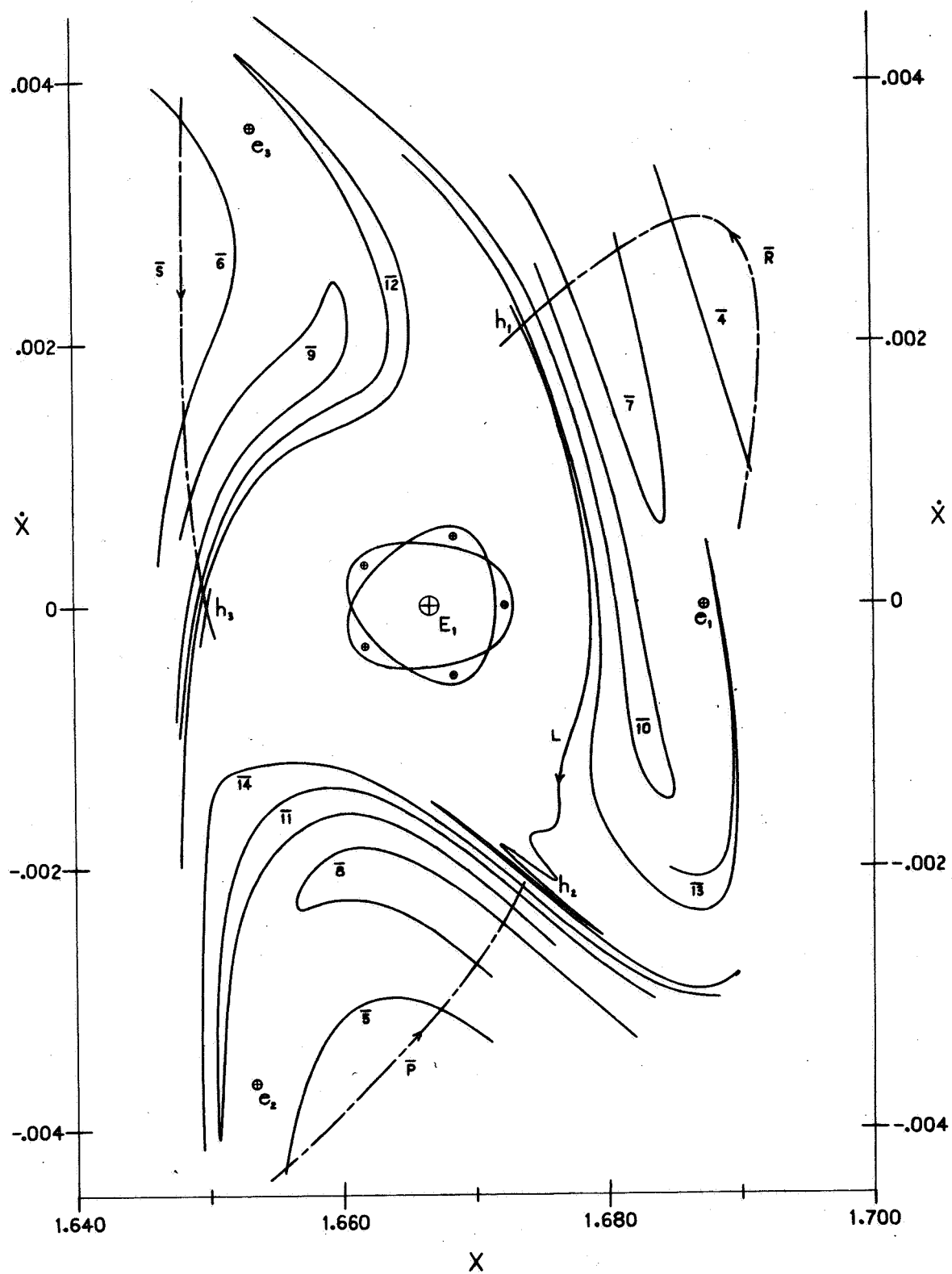
is of curve 1). Now  $\bar{S}'$  is the antecedent of  $\bar{S}$ , which in turn maps successively into regions  $\bar{R}, \bar{Q}, \bar{P}, \bar{O}, \bar{N}$ , and so on to the outer region. Thus if our mapped points in the tip of  $S'$  lie inside the region defined by  $S' \cap \bar{S}'$ , they will eventually go to the outer region along outgoing invariant curve 6. But if they still lie in the tip of  $S'$  and not inside  $S' \cap \bar{S}'$ , they will move from  $h_2$  to  $h_3$  on the inside, paralleling curve  $\bar{1}$  but not intersecting it. So these points will then move around to the neighborhood of  $h_3$  and into the influence of curve  $\bar{3}$ . There the same alternatives and processes will take place as above. Thus we can see that some points originally inside the region  $Q$  can move all the way around the inside periphery of the invariant curves, while others will gradually leak out near  $h_2$  or  $h_3$  (or  $h_1$ ,  $h_2$ , or  $h_3$  on the next revolution, etc.). Again because of the great stretching and twisting of each successive map of the inner oscillations, it is impractical to determine exactly which points (in a given initial region like  $Q$ ) leak out and when, and which points circulate around the inside periphery and for how long. In principle it can be done.

But our main concern is with the vast majority of interior points that eventually leak to the outer regions around  $e_1, e_2$ , or  $e_3$ . To investigate this let us study both the  $T^{33}$  and the  $T^{11}$  mappings simultaneously. Referring back to Figure 4, the curves labeled  $K, L$ , and  $M$  are the 3 inner invariant curves of  $T^{33}$  plotted in Figure 5 as curves 1, 2, and 3 respectively. The curves  $P, R$ , and  $S$  in Figure 4 are more complete plots of the 3 outer invariant curves of  $T^{33}$ , partially plotted in Figure 5 as curves 4, 5, and 6 respectively. We recall from Figure 5 that the interior points which do escape to the outer regions do so by moving outward with the (outer) oscillations of curves 1, 2, and 3, that is, along the general direction of the outer invariant curves 4, 5, and 6 respectively. Looking

at Figure 4 we see the connecting link to infinity: the outer invariant curves P, R, and S (i.e. curves 4, 5, and 6) intersect the time-reversed oscillations of the  $T^{11}$  mapping! Thus, in particular, interior points which escape (via curves 2 and 5 of Figure 5) along outer invariant curve R in Figure 4 will quickly intersect region 7 of the  $T^{11}$  mapping. In one more map of  $T^{33}$  (three more under  $T^{11}$ ) these points will be inside region 4, and will then move successively under  $T^{11}$  through regions 3, 2, 1 to infinity. The same applies to points moving out along invariant curve S to intersect region 6, or those moving out along invariant curve P to intersect region 8. (Note the double intersection of R with 4, corresponding to similar double intersections of S with 6 and P with 5.)

A more detailed plot of the relationship between the  $T^{33}$  and  $T^{11}$  mappings is shown in Figure 6. Curve L is one of the oscillating  $T^{33}$  inner invariant curves (plotted in Figure 5 as curve 2). Curves  $\bar{P}, \bar{R}$ , and  $\bar{S}$  are incoming outer invariant curves (labeled  $\bar{4}, \bar{5}$ , and  $\bar{6}$  in Figure 5) which are simply reflections of the outgoing outer invariant curves P, R, and S of Figure 4. The directional arrows indicate the sense of the motion under  $T^{33}$ . Also plotted are the reflected regions  $\bar{4}, \bar{5}, \dots, \bar{14}$  of portions of the time-reversed  $T^{11}$  oscillations previously discussed. For example,  $\bar{7}$  and  $\bar{8}$  are reflections of regions 7 and 8 in Figure 4. Since these are reflections of time-reversed mappings themselves (for which 8 maps into 7 and so on as time goes forward), then the order of these mappings is clearly  $\bar{4}, \bar{5}, \dots, \bar{14}$ . (In fact, these are the original forward-going mappings obtained for  $T^{11}$ , and the regions plotted as 1, 2, 3, ..., 8 in Figure 4 are in reality themselves the mirror reflections of these original mappings.) The plot shows vividly how the inner  $T^{11}$  oscillations wrap themselves around the  $T^{33}$

Figure 6. Hyperbolic invariant curves of the mapping  $T^{33}$  and their relations to those of the mappings  $T^{11}$  and  $T^{55}$ .



elliptic fixed points in cycles of 3 (e.g. note the similarity and proximity of regions  $\overline{5,8,11}$ , and  $\overline{14}$  as they wrap around  $e_2$ ). It also shows that in the limit as  $t \rightarrow \infty$  these  $T^{11}$  oscillations crowd in along  $\overline{P}, \overline{R}$ , and  $\overline{S}$  to the invariant curves of the  $T^{33}$  mapping (curves 1-6 in Figure 5, or curves K, L, M, P, R, and S in Figure 4). Finally, the plot clearly shows the 3 channels (along invariant curves P, R, and S) through which points can escape from the  $T^{33}$  mapping to the  $T^{11}$  mapping.

We therefore can understand the complete mechanisms for the escape of points from inside the  $T^{33}$  hyperbolic invariant curves out to infinity. The first link in the chain of mechanisms is the intersection of the hyperbolic invariant curves of the  $T^{33}$  and  $T^{11}$  mappings. The subsequent links are the intersections of the  $T^{11}$  mapping with the mapping of  $T^{10}$ , that of  $T^{10}$  with  $T^9$ , and so on out to  $T^1$  and infinity. The key word is intersection.

#### E. The mapping $T^{55}$

We can now ask if the same mechanism allows escape from inside the  $T^{55}$  hyperbolic invariant curves. Our original stroboscopic calculations seemed to show smooth (non-oscillatory) behavior for  $T^{55}$ . The first strong indication that this behavior does occur has already been noted: the apparent inability of the  $T^{33}$  hyperbolic oscillations to fill up the region inside the  $T^{33}$  hyperbolic invariant curves and around  $E_1$  (as do the  $T^{11}$  curves). Apparently the inner oscillations of the  $T^{33}$  invariant curves do not come close to the cluster of 10 fixed points of the  $T^{55}$  mapping, but instead stretch out thinner and thinner along the inside periphery of their mutual outer boundary, far away from the  $T^{55}$  invariant curves.

But to be positive of this apparent non-intersection of the  $T^{33}$  and  $T^{55}$  curves (just as for the  $T^{11}$  curves, we are here limited in practice

from mapping very many oscillations of the  $T^{33}$  curves) it is necessary to study the  $T^{55}$  invariant curves to a high degree of accuracy. The simplest approach is to study the inner invariant curve issuing from the  $T^{55}$  hyperbolic fixed point on the X-axis ( $h(4)$  in Table 2), analogous to curve 1 in Figure 5. The first step is to choose initial points on the X-axis very close to that hyperbolic fixed point as the origin of a polar coordinate system (with  $r_i \leq 10^{-6}$  and  $\theta_i = 0$ ). We can then make a stroboscopic plot of their subsequent motions under repeated mappings of  $T^{55}$ . These points will rapidly approach the outgoing hyperbolic curve, which is itself asymptotic to a short straight-line segment through  $h(4)$ . Knowing the angle this asymptote makes with the X-axis, we can then map a segment of this radial line for many successive applications of  $T^{55}$ .

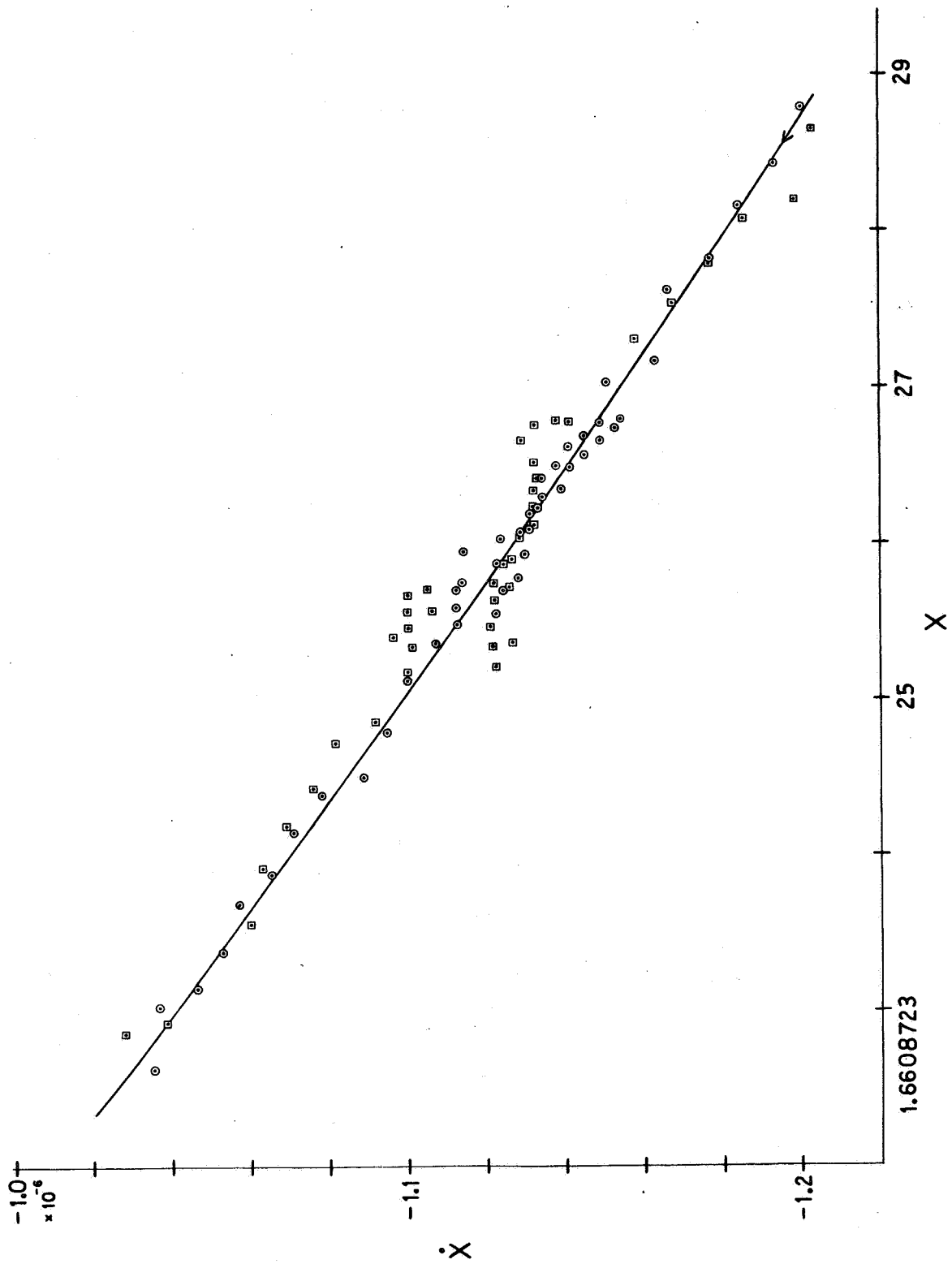
In our particular case the asymptote makes an angle of about  $15.9^\circ$  with the X-axis, and the range of initial radii used is  $7.00 \times 10^{-7} \leq r_i \leq 8.15 \times 10^{-7}$  with increments as small as  $2 \times 10^{-9}$ . Once again it is useful to print out our coordinates every 11 periods, so that we can plot the invariant curve issuing from the original hyperbolic point  $h(4)$  or from any of the other 4 hyperbolic points which are transforms of it (see Table 2). In particular the choice of the invariant curve issuing from  $h(3)$  is convenient because its oscillations (if they exist) would crowd into the asymptote to  $h(4)$  which has already been determined above.

The results of these calculations differ strikingly from those for the  $T^{11}$  or  $T^{33}$  invariant curves. First of all there is a pronounced increase in the stroboscopic period of the motion. Whereas points on the  $T^{11}$  or  $T^{33}$  invariant curves cover most of the distance between alternate hyperbolic points in 4 or 6 multiples of the basic mapping ( $11\tau$  or  $33\tau$ ),

points on the  $T^{55}$  invariant curves require more than 100 multiples of the basic mapping of  $55\tau$  periods. Secondly, as points in the original radial segment near  $h(3)$  slowly approach  $h(4)$ , they do develop a small irregular "scatter", undoubtedly due to the accumulation of truncation and inversion errors after so many successive mappings by the computer. [Note that 100 maps under  $T^{55}$  correspond to 5500 periods of  $\tau$ , or 11,000 half-period calculations. Now each of these calculations involves 9 or 10 computational steps such as basic arithmetic operations, square roots, and elliptic function interpolations and inversions. These latter computations in turn involve separate subprograms with many calculational steps themselves (on the order of 10 to 100). Thus each final plotted point may require some  $10^6$  or  $10^7$  arithmetic operations! If the computer truncation error is in the thirteenth decimal place (as it is on the Illiac II computer), then this error might accumulate constructively for  $10^6$  operations or more and hence show up in even the seventh decimal place. In our calculations the scatter shows up in the eighth decimal place, entirely consistent with the predictable errors.]

However, there is no tendency of the  $T^{55}$  invariant curves to develop any regular pattern of oscillations analogous to the  $T^{11}$  or  $T^{33}$  invariant curves. The 46 circled points in Figure 7 are results of the 442<sup>nd</sup> successive application of the  $T^{11}$  mapping to each of 46 points on the original radial segment through  $h(4)$  (with  $\theta_1 = 15.9^\circ$  and  $7.00 \times 10^{-7} \leq r_1 \leq 8.15 \times 10^{-7}$ ). These circled points are the results using an elliptic function inversion-routine error of  $E = 10^{-12}$ , while those points inside squares are the results for  $E = 10^{-11}$ . The smooth continuous curve is a quadratic, least-squares computer fit of the 46 circled points for

Figure 7. Computer fluctuations over the 442<sup>nd</sup>  
"oscillation" of one hyperbolic  
invariant curve of the mapping  $T^{55}$ .



$\epsilon = 10^{-12}$  (the best possible accuracy for the inversion routines). This plot clearly shows the nature and magnitude of the random scatter of mapped points. Also, this mapped section does indeed correspond to a "full period of an oscillation" of the  $T^{55}$  invariant curve. For, the first circled point in the lower-right corner of Figure 7, which is a map of the point with  $r_i = 7.00 \times 10^{-7}$ , transforms under  $T^{55}$  into a point lying below and to the right of the last circled point in the upper-left corner of the plot, a map of the point with  $r_i = 8.15 \times 10^{-7}$ . That is, the stroboscopic mapping of the original radial segment at the end of the 447<sup>th</sup> successive application of  $T^{11}$  overlaps the 442<sup>nd</sup> mapping of the same original radial segment. This overlapping of the ends of the mapped segment with the next  $T^{55}$  map of that segment is what we mean by a "full period of one oscillation" of the  $T^{55}$  invariant curve.

The most convincing contrast between the behavior of this  $T^{55}$  invariant curve and that of the  $T^{11}$  or  $T^{33}$  invariant curves is the phase-space area occupied by any one oscillation of the mapping. Using Figures 4 and 5 we can measure the phase-space area enclosed by one of the  $T^{11}$  or  $T^{33}$  oscillations, for example, region 3 in Figure 4 and region E in Figure 5. For the randomly scattered points of the  $T^{55}$  invariant curve in Figure 7 we can simply connect all points in a piecewise (discontinuous) fashion and then measure the area of the largest areal fluctuation. The results are

$$\left\{ \begin{array}{ll} T^{11}: & A_{11} \approx 5.65 \times 10^{-5} \quad (\text{XX units}) \\ T^{33}: & A_{33} \approx 1.4 \times 10^{-7} \quad (\text{XX units}) \\ T^{55}: & A_{55} \approx 8 \times 10^{-17} \quad (\text{XX units}) . \end{array} \right\} \quad (22)$$

Now the ratio  $(A_{11}/A_{33}) \approx 400$ , so that if we use this same ratio for  $(A_{33}/A_{55})$  we might expect  $A_{55} \approx 10^{-10}$  (XX units). What we actually find is that  $A_{55}$  is  $10^6$  times smaller than this area, suggesting that  $A_{55}$  is almost, if not exactly, zero. In fact, recalling that computer truncation errors may build up to  $10^{-8}$  in both  $X$  and  $\dot{X}$  after long periods of time, we see that the product of the predictable errors in  $X$  and  $\dot{X}$  agrees with the largest areal fluctuation found for the  $T^{55}$  mapping.

Having thus determined, to the  $10^{-12}$  accuracy of the Illiac II computer, that there exist no oscillations in the  $T^{55}$  hyperbolic invariant curves, we can now look at the composite picture of all of the  $T^{55}$  hyperbolic invariant curves. The (inner) invariant curve issuing from  $h(3)$  approaches  $h(4)$  via the radial asymptote to  $h(4)$  which makes an angle of  $360^\circ - 15.9^\circ = 344.1^\circ$  to the  $X$ -axis. In what follows we assume that our point lies an infinitesimal distance from the true mathematical invariant curve, and that we can perform our  $T^{55}$  mappings with extreme precision. After repeated applications of  $T^{55}$ , our point slowly emerges near the outer invariant curve at an angle around  $180^\circ - 15.9^\circ = 164.1^\circ$ . Referring to the stroboscopic map of the region shown in Figure 3 we see that the point will now move around the outside of the fixed point  $e(4)$  and will then proceed into fixed point  $h(5)$ . But now it will move towards the  $X$ -axis and along an inner invariant curve. It will next go around  $e(5)$  on the inside and then head out into  $h(1)$ . Next the point will cross the  $X$ -axis with an intercept  $(X_r, 0)$  to the right of  $e(1)$  and then proceed to  $h(2)$ . In like manner the point will move inside  $e(2)$ , through  $h(3)$ , outside  $e(3)$ , through  $h(4)$ , inside  $e(4)$ , through  $h(5)$ , outside  $e(5)$ , through  $h(1)$ , then cross the  $X$ -axis with an intercept

$(X_\ell, 0)$  to the left of  $e(1)$ , next pass through  $h(2)$ , and finally move outside  $e(2)$  and return to  $h(3)$  along the same initial invariant curve. This motion represents one complete stroboscopic revolution.

Therefore the union of the forward-going and the forward-coming (reflected) invariant curves is in reality one continuous, smooth curve with an angular period of  $4\pi$  radians. It simultaneously bounds the inner region as well as the 5 outer "islands" of  $T^{55}$  elliptic fixed points. This curve, then, is what we may properly call the continuous doubly-periodic invariant curve of the  $T^{55}$  mapping.

Determining curves which lie very close to this  $T^{55}$  invariant curve is an easy matter on the computer. We now make long-period stroboscopic mappings under  $T^{55}$  of a series of points in an interval of the X-axis which brackets the inner invariant curve intercept  $(X_\ell, 0)$ . From our initial stroboscopic scan of the X-axis (see Figure 3) we see that the region in question has a radius  $r_i \approx 5 \times 10^{-3}$  for  $\theta_i = 0$  (again using  $E_1$  as origin of our polar coordinates). From Figure 3 it is clear that initial points with  $X_i < X_\ell$  will move stroboscopically in a clockwise sense about  $E_1$ , always sticking to the inner periphery of the envelope of the  $T^{55}$  invariant curves. On the other hand, points with  $X_i > X_\ell$  will move stroboscopically in a counter-clockwise sense around  $e(1)$ , defining the "island" curves about the elliptic point. (Note that if we print out every  $T^{11}$  map of these points with  $X_i > X_\ell$  we will simultaneously obtain the other 4 "island" curves as transforms of the original one around  $e(1)$ .)

We can now take successively finer grids of initial points around the suspected  $X_\ell$  value and calculate for about 300 mappings of  $T^{11}$  (about all that is needed before "island" motions separate from inner-region motions).

It is then an easy matter to narrow the interval down as close to  $X_\ell$  as the errors in computation will permit. The result of this scan is the following bracketed interval for  $X_\ell$ :

$$1.671665069 = X_{\ell 1} < X_\ell < X_{\ell 2} = 1.671665070 . \quad (23)$$

Setting  $X_1 = X_{\ell 1}$  and applying  $T^{55}$  successively, we obtain an approximately closed curve very close to the envelope of the inner  $T^{55}$  invariant curves, as suggested by its long stroboscopic period of 2610 multiples of  $11\tau$ . Similarly by setting  $X_1 = X_{\ell 2}$  and now applying  $T^{11}$  repeatedly, we can obtain the 5 approximately closed curves around the 5 "islands" (we find each has a stroboscopic period of 1060 multiples of  $11\tau$ ). We might note that from the mapping of the initial point  $X_1 = X_{\ell 2}$  (after about 530 multiples of  $11\tau$ ) we can obtain an approximate value for the outer invariant curve intercept  $X_r$  of the  $X$ -axis. We find this value to be  $X_r \approx 1.672981475$ . (In polar coordinates about  $E_1$ , with  $\theta=0$ , this outer intercept has  $r_r \approx 0.00624605$ ; likewise the inner intercept  $X_\ell$  has  $r_\ell \approx 0.00492964$ .)

The resulting 6 disjoint "closed" curves, lying very close to the  $T^{55}$  hyperbolic invariant curves, are plotted in Figure 5, with directional arrows for the motion indicated on a few. The five  $T^{55}$  elliptic fixed points are plotted and labeled, while the labeled positions of the five  $T^{55}$  hyperbolic points show up clearly between the five elliptic points. The same curves are plotted in Figure 6 as one continuous closed curve, i.e. as the closed invariant curve of the  $T^{55}$  mapping, to an accuracy better than the width of the lines in the figure. This closed curve thus bounds the invariant region of stability which we have been seeking.

### III. CONCLUSIONS

#### A. Summary of results; an application

Let us summarize the results of our detailed calculations. There does exist to the  $10^{-12}$  accuracy of the Illiac II computer a closed region of phase-space which surrounds the  $T^{11}$  elliptic fixed point ( $E_1$ ) on the X-axis and which is invariant under successive  $T^{55}$  mappings of the phase-space onto itself. Many points initially outside this region will escape to infinity after a finite number of  $T^{55}$  mappings. Some points are stable for a long time but eventually escape to infinity also, so that they can be considered to lie in a "twilight zone" that surrounds the invariant region. The outer  $T^{33}$  invariant curves oscillate and intersect the inner  $T^{11}$  oscillating invariant curves, so that points can move from near the  $T^{33}$  fixed points about  $E_1$  into the  $T^{11}$  oscillations and escape to infinity.

However in our particular case the mutual intersection of oscillating invariant curves which correspond to multiples of the basic ( $T^{11}$ ) mapping terminates at the  $T^{33}$  level. The inner oscillations of the  $T^{33}$  mapping do not intersect the outer hyperbolic invariant curves of the  $T^{55}$  mapping which apparently have no oscillations. The union of all the  $T^{55}$  hyperbolic invariant curves is actually one continuous closed curve with an angular period of  $4\pi$  radians. This curve bounds a stable region around the  $T^{11}$  elliptic fixed point  $E_1$ .

This region of stable motions consists of 6 cells: a large inner region, plus 5 smaller "islands" which surround the 5 elliptic fixed points of the mapping  $T^{55}$ . The inner region has a generally flat elliptical shape but has 5 distinct bulges toward the 5 hyperbolic points of  $T^{55}$ . Points initially inside this inner region always remain inside the

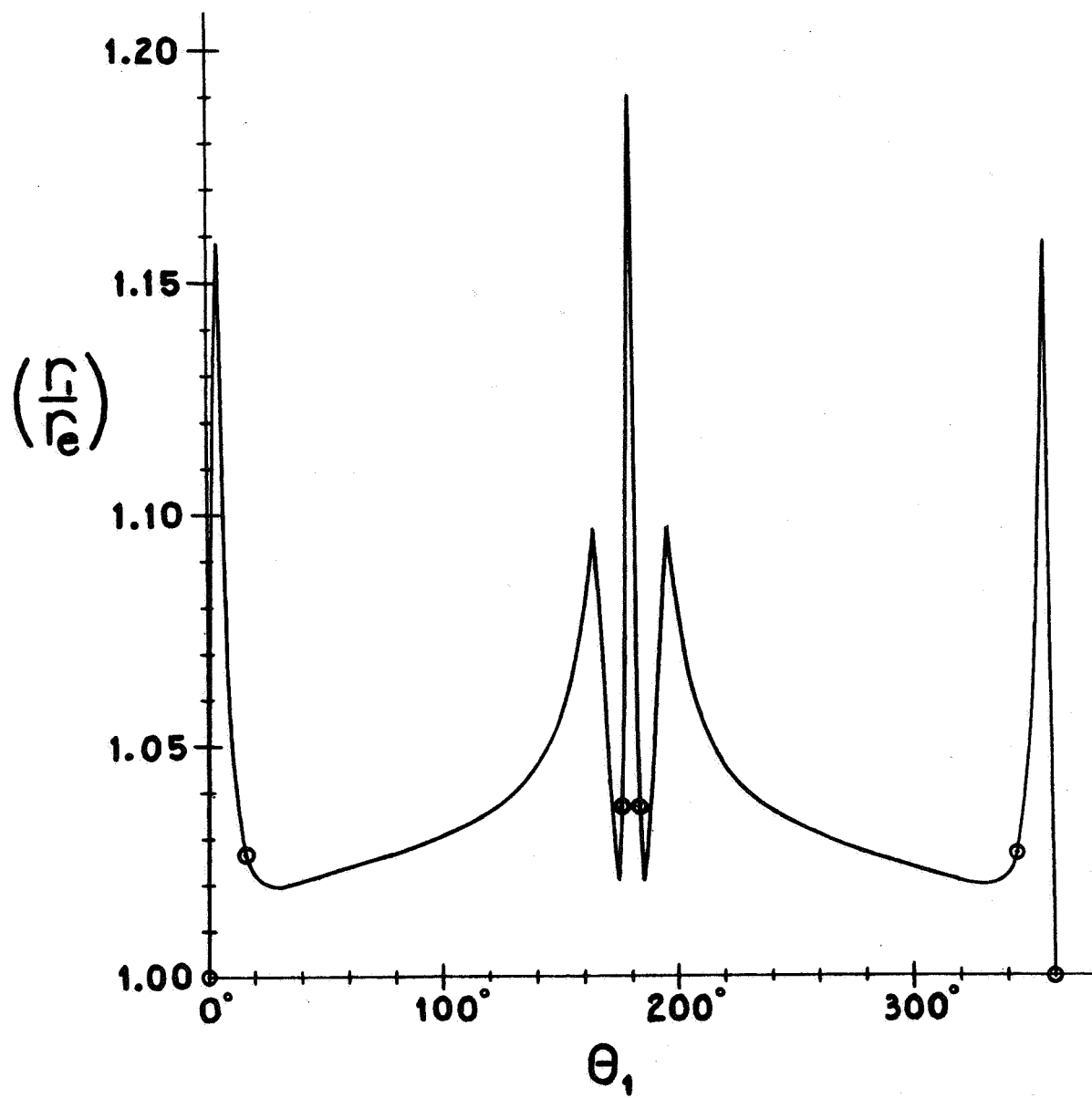
region under successive applications of the mapping  $T^{11}$ . They execute clockwise, almost-elliptical motions under successive mappings of  $T^{55}$ .

For the purpose of comparison let us arbitrarily select the inner, almost-elliptical motion generated by successive  $T^{55}$  mappings of the initial point with  $\theta_i = 0$ ,  $r_i = 0.0045$ , which is just inside the inner hyperbolic invariant curves of the  $T^{55}$  mapping. Next let us match an ellipse to this motion at two points on the curve, for example at  $0^\circ$  and  $90^\circ$ . Identifying the radii at these two angles with the semi-major and semi-minor axes of the ellipse respectively, we calculate the corresponding "standard" eccentricity to be  $e_s = 0.995635$ . The actual deviations of the mapped motions from a "standard" ellipse which has an eccentricity  $e_s$  are obtained by calculating the ratio of the mapped radius  $r_1$  at some angle  $\theta_1$  to the polar radius  $r_e$  of a standard ellipse for the same angle  $\theta_1$ . This standard ellipse is completely specified by  $e_o = e_s$  and  $a_o$  values. We then use the initial radius  $r_i$  at  $\theta_i = 0$  for the semi-major axis  $a_o$ . The value of  $r_e$  at the mapped angle  $\theta_1$  is given by the equation

$$r_e(\theta_1) = a_o \sqrt{\frac{1 - e_o^2}{1 - e_o^2 \cos^2 \theta_1}} \quad (24)$$

We now stroboscopically map the initial point  $\theta_i = 0$ ,  $r_i = r_\ell = 0.004929639$  in order to obtain a curve  $r_1(\theta_1)$  very close to the inner invariant curves of  $T^{55}$  which bound the inner cell. A plot of the results for the ratio  $r_1/r_e$  as a function of  $\theta_1$  around this mapped curve (using  $a_o = r_i = r_\ell$ ) is shown in Figure 8. This plot shows a large 19% deviation in the direction of  $h(4)$ , 15.8% deviations toward  $h(1)$  and  $h(2)$ , and 9.7% deviations toward  $h(3)$  and  $h(5)$ . Minimum deviations are about 2% in the directions

Figure 8. Deviation from a standard ellipse of a curve very close to the union of inner hyperbolic invariant curves of the mapping  $T^{55}$ .



of  $e(2)$ ,  $e(3)$ ,  $e(4)$ , and  $e(5)$ . (The values of  $r_1/r_e$  for these  $T^{55}$  elliptic points are indicated by small circles.) Thus, very close to the outer edge of the inner invariant cell the mapped motions deviate from our very flat standard ellipse for all non-zero angles. Nevertheless the concept of a standard ellipse does constitute a good way of comparing a wide range of mapped motions around  $E_1$ .

Points lying initially inside any one of the 5 "island" regions will always remain inside that particular "island" for successive mappings of  $T^{55}$ , but will jump discontinuously through all 5 "island" regions under the mapping  $T^{11}$ . Under  $T^{55}$  the mapped motions will be counter-clockwise and nearly elliptical in each of the 5 "islands". It is therefore clear that we can break up our original invariant region into 6 cells possessing separate invariance properties and predictable stroboscopic motions under  $T^{55}$ . But because these same 6 cells are also bounded by one continuous closed invariant curve, we can say that there is one invariant region with a doubly-periodic boundary.

These results can be applied to a totally different type of problem, the numerical search for additional integrals of motion (see, for example, the work of Bozis<sup>18/</sup> on the restricted three-body problem). Hénon and Heiles<sup>19/</sup>, in an earlier general search, plotted disjoint curves which in reality are stroboscopic trajectories around various  $T^{mn}$  satellite fixed points (e.g. their Figure 5 is exactly analogous to our Figure 3). They studied these disjoint curves for a wide range of values of the energy  $E$  (they considered motions under a general nonlinear potential). Interpreted in the light of our results, their work shows that the particular  $T^{mn}$  cluster of satellite multiple points which bounds the invariant region(s) depends on the energy  $E$ . Their plots for lower energies demonstrate that

the  $T^{3n}$  satellite fixed points bound the invariant region(s), whereas at higher energies the  $T^{5n}$  points form the boundary. At still higher energies the invariant regions get smaller and smaller, undoubtedly involving  $T^{mn}$  points with very large  $m$ . (One of the figures in the paper of Bozis shows a multiplicity of  $m = 37$ !)

Hénon and Heiles refer to their closed curves (i.e. the stroboscopic trajectories) around the central point and around the outlying "islands" as disjoint "isolating third integrals of the motion". Points, with puzzling random (or jumping) behaviors, which lie outside these curves are referred to as points on the "ergodic trajectories in the sea between the islands" (i.e. where the third integral is non-isolating). Our work shows that these points are following completely predictable trajectories, as long as the global  $T^{mn}$  mappings are known. It is also clear that the invariant region(s) we have found are just those regions where an "isolating third integral of the motion" exists. In fact, the bounding invariant curve of our 6 invariant cells is a graphic example of such an "isolating integral of the motion", for it exhibits a functional relationship between  $X$ ,  $\dot{X}$ , and  $t$  for successive stroboscopic (integrated) motions. Thus the application of our methods not only will yield the invariant curves and regions of stability (for a given equation with given parameter values), but also will exhibit simultaneously the regions of phase-space in which "isolating integrals of the motion" exist. In reality, these two concepts are equivalent.

#### B. Applicability of Moser's theorem to the $T^{55}$ mapping

Since the delineation of a locally stable invariant region for our equation has involved much computer time, it would be desirable to utilize

any general mathematical methods available for the reduction of the time required. To this end we shall discuss the applicability of a recent theorem of Moser<sup>15/</sup>. Our exploratory mappings of the X-axis (a few of them are shown in Figure 2) are highly suggestive of properties which might be necessary conditions for stroboscopic stability. The most obvious one is the property of twist mappings of the phase-plane, where the angles of the mapped points increase monotonically with radial distance from the given fixed point. In Figure 2 we see that this property holds for all the mappings  $T^{11m}$  (for  $m=1,2,3,4,5$ ) as long as we stay moderately close to  $E_1$  (within, roughly, an elliptical region with  $e_0 = e_s$  and semi-major axis  $a_0 \leq 0.0100$ ). As we get farther away from  $E_1$  the twisting gets much more complicated, so that a monotonically increasing  $T^{11m}$  twist mapping may no longer exist. For example, for the  $T^{44}$  mapping of the X-axis to the right of  $E_1$  (curve  $D^+$ ) we see that the condition of equation (12) is not satisfied for a small region around (1.652, -0.005). This region lies just outside the  $T^{33}$  elliptic fixed point  $e_2$  (which is located by the intersection of curves  $D^+$  and  $A^+$ ). Therefore, it appears that some type of a twist mapping property does exist in the interior of an elongated irregular region bounded, roughly, by the 6 fixed points of the  $T^{33}$  mapping.

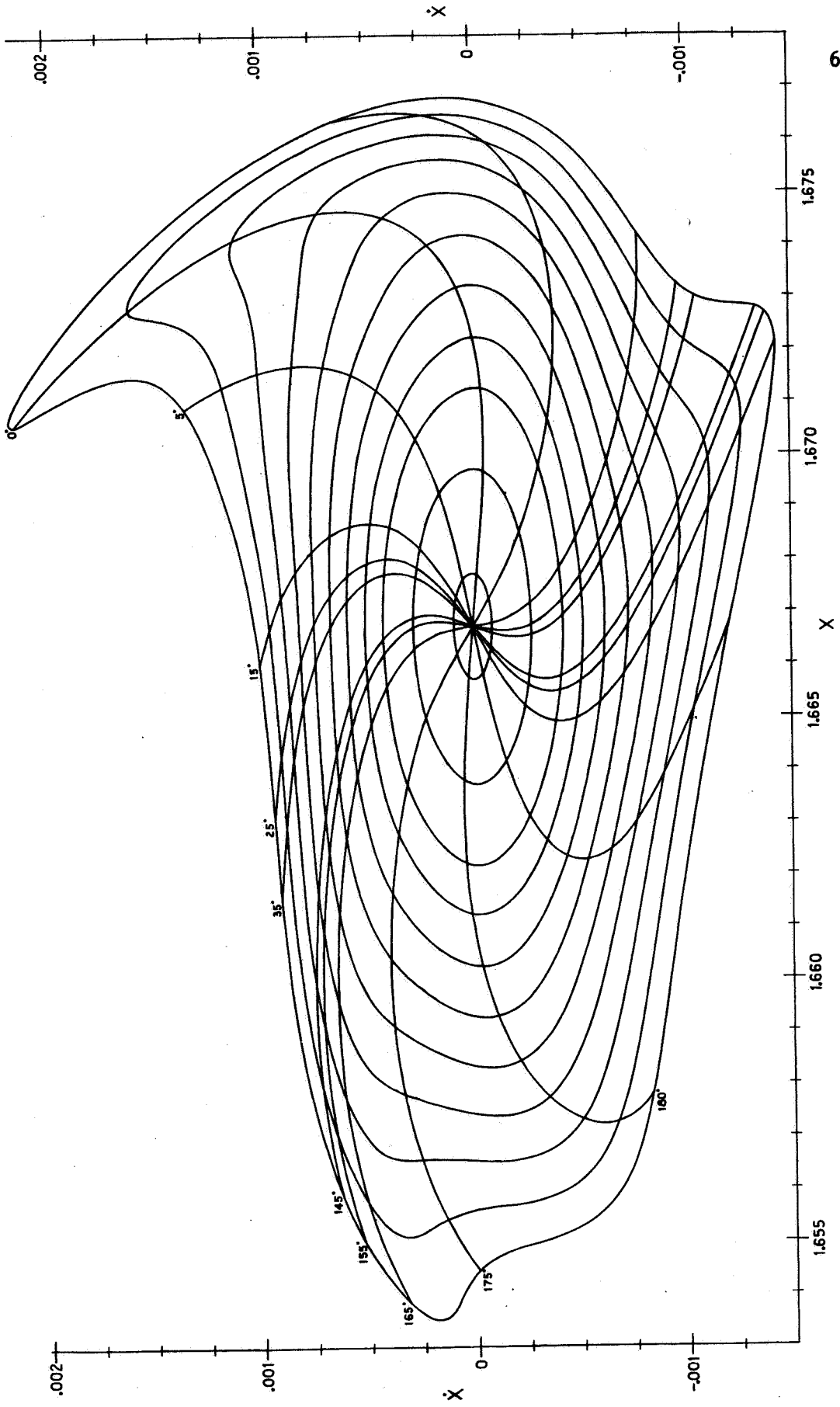
The maps of Figure 2 indicate the general regions in which a  $T^{11m}$  twist mapping appears to exist, but these regions are much more extensive (in area) than the known invariant regions of Figures 5 and 6. There seems to be no clue here that the  $T^{55}$  invariant curves form the boundary of the invariant region. So let us examine in detail the general region of the  $T^{55}$  fixed points, where we know that an invariant region exists. Since the inner stroboscopic motions are clockwise and almost-elliptical

(see Figure 3), let us make a  $T^{55}$  mapping of a set of concentric ellipses and radiating straight lines about and through the point  $E_1$ . The maps of these initial curves and lines are shown in Figure 9. The eccentricity of the initial ellipses is  $e_o = e_s = 0.995635$ . The values of the initial semi-major axes  $a_o$  are  $10 \times 10^{-4}$ ,  $30 \times 10^{-4}$ , and  $45 \times 10^{-4}$  to  $125 \times 10^{-4}$  in steps of  $10 \times 10^{-4}$ . The angles  $\theta_o$  of the initial straight lines are indicated near their  $T^{55}$  maps. Evidently a  $T^{55}$  twist mapping condition does hold throughout this range of initial semi-major axis values. Again, however, there does not appear to be any readily visible indication of the actual existence of an invariant region for  $a_o \approx 50 \times 10^{-4}$ .

For  $a_o \approx 100 \times 10^{-4}$  the mapped curves deviate considerably from their initial elliptical shapes, especially for the initial angles around  $0^\circ$ ,  $165^\circ$ , and  $215^\circ$ . Referring to the plots of the  $T^{33}$  mappings (Figures 5 and 6), we see that the three prominent directions for "bulging" of the  $T^{55}$  twist maps in Figure 9 are toward the 3 hyperbolic "escape channels" of the  $T^{33}$  mapping (i.e. along its 3 outgoing outer invariant curves). These bulges, signs of impending escape to infinity via the  $T^{33}$  and  $T^{11}$  oscillating invariant curves, begin to show up faintly in the mapped curves of Figure 9 at about  $a_o = 85 \times 10^{-4}$ , but not noticeably before that. So once again we find no striking behaviors of the mappings which might indicate the precise location of the invariant region around  $E_1$ .

Moser<sup>15/</sup> has demonstrated the theoretical existence of invariant curves and regions for mappings which are small perturbations of the circular twist mapping. Our mappings are very flat elliptical twist mappings, but we can change our scale so that the minor axis of our elliptical motions approximately equals the major axis (e.g. we can adjust

Figure 9.  $T^{55}$  mappings of initially concentric standard ellipses and radiating straight lines around the  $T^{11}$  fixed point.



the value of  $p_0$  in equation (20)). An equivalent procedure is to choose a particular standard ellipse, map its points, and compare the resulting locus (i.e. the image curve) with the original ellipse. This procedure is precisely the one which we followed in order to obtain Figure 9.

Now let us see if Moser's conditions apply to the  $T^{55}$  mapping. The computations for Figure 9 show that the image curve and the initial curve (a standard ellipse) generally intersect each other four times. Thus Moser's first condition (equation (16)) is satisfied. In order to check his other conditions (equations (17) and (18)) we would first have to exhibit his functions  $\alpha(r_0)$ ,  $F(r_0, \theta_0)$ , and  $G(r_0, \theta_0)$ . From equations (13) we see that this can be done directly by calculating the increments

$$\Delta\theta = (\theta_1 - \theta_0) = \alpha(r_0) + F(r_0, \theta_0) \quad (25)$$

$$\text{and} \quad \Delta r = (r_1 - r_0) = G(r_0, \theta_0), \quad (26)$$

where  $(r_0, \theta_0)$  is the initial point and  $(r_1, \theta_1)$  is the mapped point, for a whole series of initially concentric ellipses. (Recall from equation (24) that  $e_s$ ,  $a_0$ , and  $\theta_0$  uniquely determine the  $r_0 = r_e$  coordinates of points on the initial standard ellipse.) By plotting the  $\Delta\theta$  increment versus  $\theta_0$  we obtain the angular dependence of the  $F$  function (for one  $a_0$  value) plus one value of the  $\alpha$  function. Similarly a plot of  $\Delta r$  versus  $\theta_0$  yields the angular dependence of the  $G$  function for one  $a_0$  value. By repeating this process for many initial ellipses with different initial semi-major axes  $a_0$  (keeping  $e_s$  constant) we can obtain  $\alpha$ ,  $F$ , and  $G$  as functions of  $a_0$  for various values of  $\theta_0$ .

Before examining these functional dependences in detail, let us write down a more convenient elliptical analog of Moser's twist mapping,

using different functional symbols:

$$\theta_1 = T^n(\theta_0) = \theta_0 + \beta(a_0, e_0) + R(a_0, e_0, \theta_0) \quad (27)$$

$$\text{and} \quad r_1 = T^n(r_0) = r_0 + S(a_0, e_0, \theta_0) \quad (28)$$

where  $a_0$  and  $e_0$  are the semi-major axis and the eccentricity respectively of an ellipse close to its image curve which lies within some elliptical annular region. Note that the variable  $r_0$  (or  $r_1$ ) can be expressed in terms of the basic variables  $a_0$  (or  $a_1$ ) and  $e_0$  by using equation (24).

Because of the three variables now ( $a_0, e_0$ , and  $\theta_0$ ), the task of exhibiting the functions  $\beta$ ,  $R$ , and  $S$  is greatly complicated. If, however, we set  $e_0 = e_s$  for convenience, then we have the simpler problem of mapping concentric initial ellipses and radial straight lines and determining the functional dependences as before. Using the data from the plots in Figure 9 we can calculate the basic  $T^{55}$  mapping increments

$$\Delta\theta = (\theta_1 - \theta_0) = \beta(a_0, e_s) + R(a_0, e_s, \theta_0) \quad (29)$$

$$\text{and} \quad \Delta r = (r_1 - r_0) = S(a_0, e_s, \theta_0) \quad (30)$$

as functions of  $a_0$  or  $\theta_0$  (for  $e_0 = e_s$ ). Now in contrast to the previous discussion of methods for obtaining Moser's functions ( $\alpha, F$ , and  $G$ ), we find it more convenient to first calculate  $\Delta\theta$  and  $\Delta r$  as functions of  $a_0$  for various fixed initial values of  $\theta_0$  (i.e. for a series of different initial straight lines).

The results of these calculations for  $\Delta\theta$  and  $\Delta r$  are plotted in Figures 10 and 11 respectively. The dashed curves in both figures (for  $\theta_0 = 15^\circ$  and  $95^\circ$ ) are plots of  $\Delta\theta$  or  $\Delta r$  versus  $r_0$ . The solid curves (for

Figure 10. Angular  $T^{55}$  mapping increment  $\Delta\theta$  as a function of the semi-major axis  $a_0$  of an initial standard ellipse, for several initial angles  $\theta_0$ .

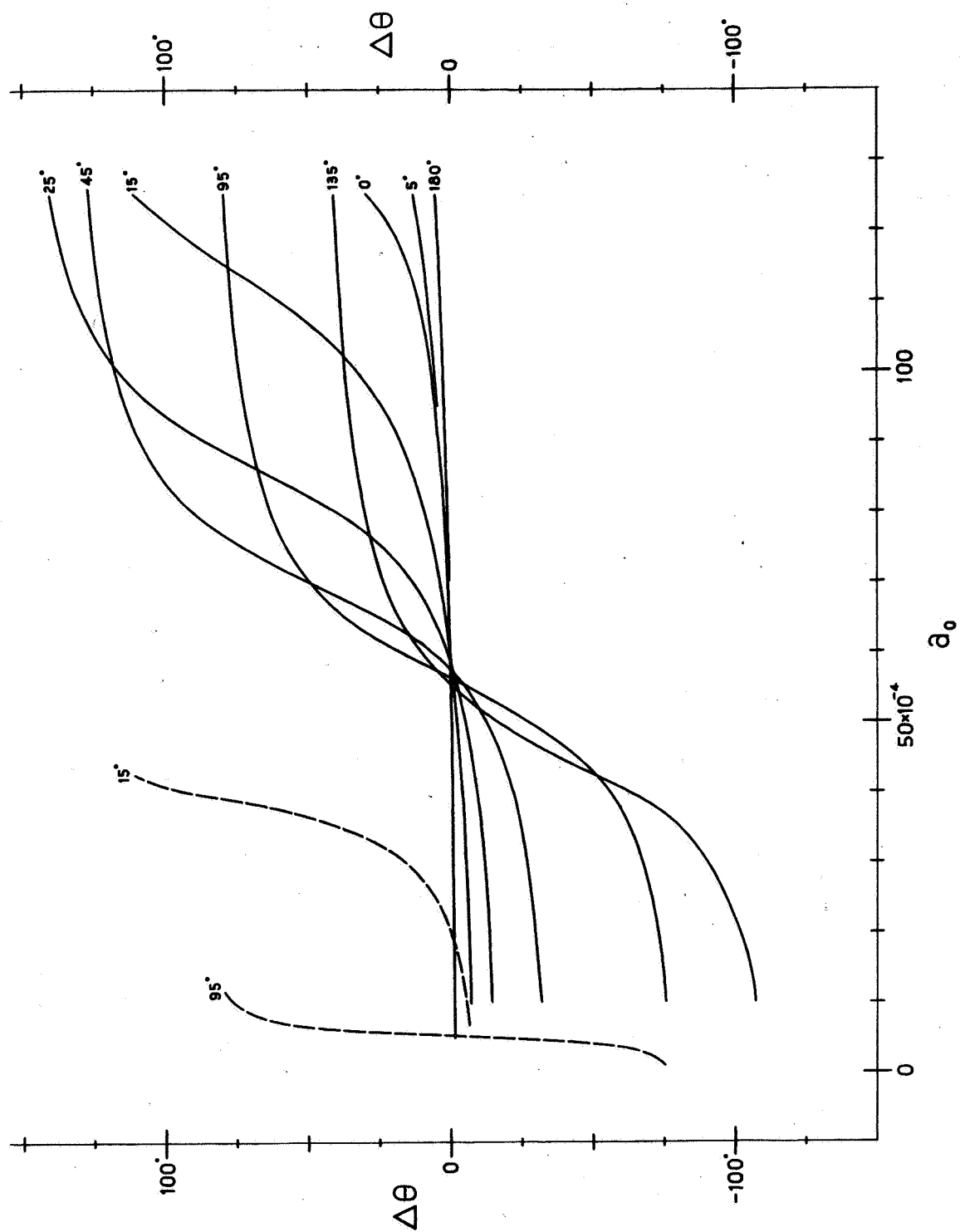
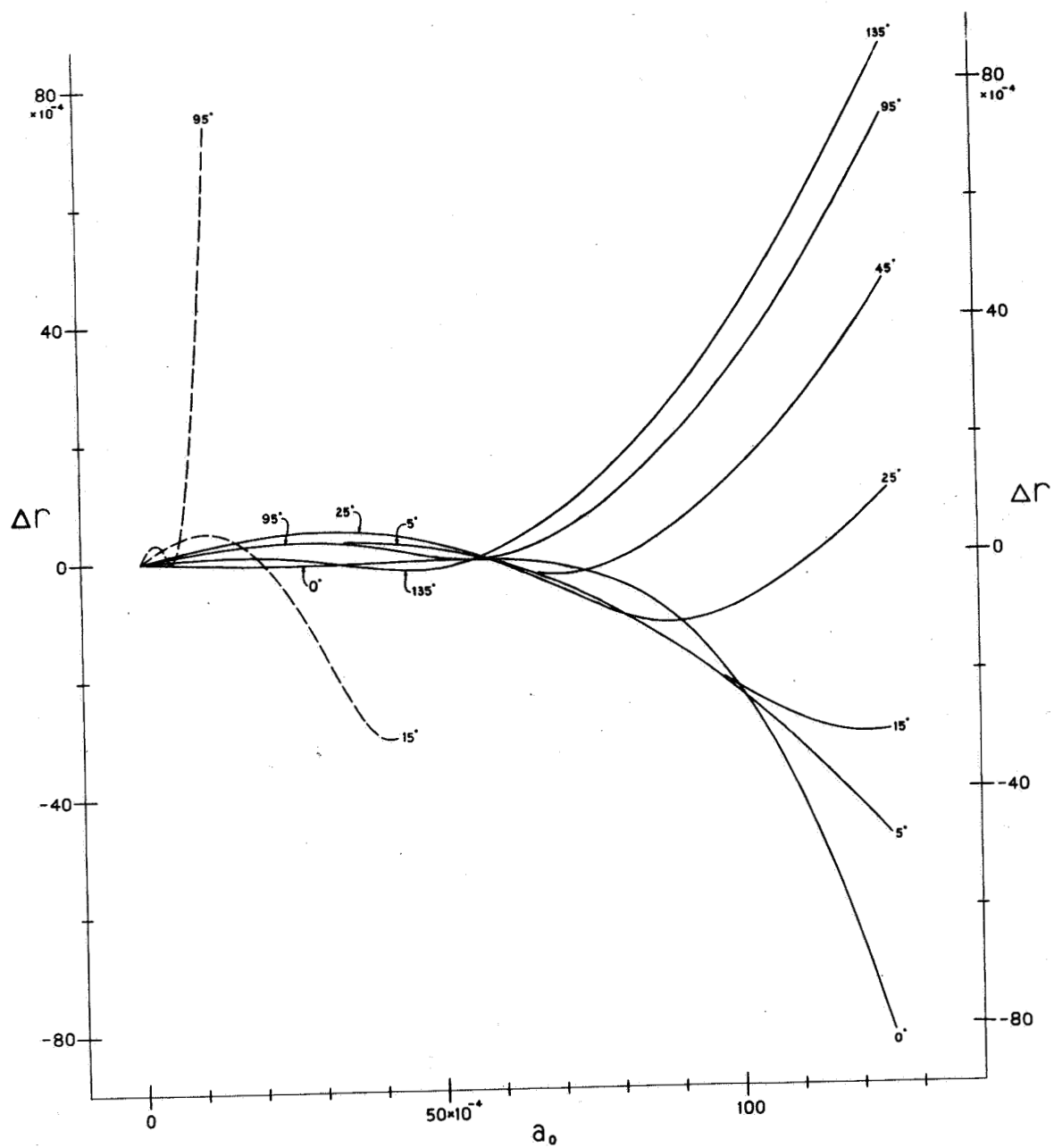


Figure 11. Radial  $T^{55}$  mapping increment  $\Delta r$  as a function of the semi-major axis  $a_0$  of an initial standard ellipse, for several initial angles  $\theta_0$ .



$\theta_0 = 0^\circ, 5^\circ, 15^\circ, 25^\circ, 45^\circ, 95^\circ, 135^\circ$ , and  $180^\circ$ ) are plots of  $\Delta\theta$  or  $\Delta r$  versus  $a_0$ . We immediately see that plotting the increments  $\Delta\theta$  or  $\Delta r$  versus  $a_0$  gives a spectrum of curves that can be easily compared, whereas plotting versus  $r_0$  does not.

Referring now to the solid curves in Figure 10, we notice that the  $\Delta\theta$  function is smooth and differentiable over the full range  $0 < a_0 < 125 \times 10^{-4}$ , well beyond the known invariant region at  $a_0 \approx 50 \times 10^{-4}$ . The  $\Delta\theta$  function is always negative for the inner region  $0 < a_0 < 55 \times 10^{-4}$ , always positive for the outer region  $60 \times 10^{-4} < a_0 < 125 \times 10^{-4}$ , and is monotonically increasing over the full range. The first two features correspond to the known (opposite) directions of the stroboscopic motions for the two regions. The small region between the inner and outer regions corresponds to the cluster of  $T^{55}$  fixed points, where the sense of the motion changes sign at different  $a_0$  values depending on the angle  $\theta_0$  (e.g. compare a straight line through e(1) with one through h(1)).

The function  $\beta$  is ordinarily obtained by plotting  $\Delta\theta$  versus  $\theta_0$  for various fixed initial values of  $a_0$ , and is the angularly independent part of each of these curves. Now the  $\Delta\theta$  function is approximately symmetric about  $\theta_0 = 180^\circ$ , so we might use the  $\Delta\theta$  intercepts at  $\theta_0 = 0^\circ$  or  $180^\circ$  as the value of  $\beta(a_0)$  for each curve, provided they are identical ( $\beta$  is just like an additive constant for all  $\theta_0$  values). If they are not identical, then the function  $\beta$  contains some angular dependence (and there is no exact symmetry about  $\theta_0 = 180^\circ$ ). One way to check this possibility is to compare the  $\theta_0 = 0^\circ$  and  $180^\circ$  curves in Figure 10 for the full range of  $a_0$  values. For the inner region we find that these two curves always differ from one another by less than  $0.1^\circ$  in the ordinate  $\Delta\theta$ . Both curves

increase almost linearly from about  $\Delta\theta \approx -1.25^\circ$  for  $a_0 = 10 \times 10^{-4}$  to about  $\Delta\theta \approx -0.3^\circ$  for  $a_0 = 50.5 \times 10^{-4}$ . For the outer region we find that the two curves differ widely from one another by as much as  $24^\circ$  for  $a_0 = 125 \times 10^{-4}$ . They still increase monotonically.

Hence we may conclude that for the inner region the function  $\beta$  is a small, negative, almost linearly increasing function of  $a_0$ , with approximately no angular dependence. For the outer region  $\beta$  is a rather large, positive, increasing function of  $a_0$ , with very definite angular dependence.

The function  $R(a_0, e_s, \theta_0)$  for the inner region is a negative, monotonically increasing function of  $a_0$  for all angles  $\theta_0$ . For this inner region we see that for fixed  $a_0$  values,  $\Delta\theta$  first decreases with increasing  $\theta_0$  and then increases again. In other words the second derivative of  $\Delta\theta$  with respect to  $\theta_0$  is positive for all  $\theta_0$  and all  $a_0$  in this region. For the outer region the conclusions are not so definite because of the angular dependence of  $\beta$  (which decreases as  $\theta_0$  increases). Here  $R$  is a positive function of  $a_0$  which increases with increasing  $a_0$ , except for  $\theta_0 \gtrsim 15^\circ$  when the reverse is true. Scanning the angular dependence we observe that the second derivative of  $\Delta\theta$  with respect to  $\theta_0$  is now negative for all  $\theta_0$  when  $a_0 < 80 \times 10^{-4}$ . When  $a_0 > 80 \times 10^{-4}$  this derivative is still negative for most values of  $\theta_0$ , but it does develop a pronounced positive region for  $0 < \theta_0 \lesssim 15^\circ$ .

Let us now turn to the solid curves for  $\Delta r$  versus  $a_0$  in Figure 11. Analogous to the  $\Delta\theta$  function, the  $\Delta r = S(a_0, e_s, \theta_0)$  function appears to be smooth and differentiable over the full range  $0 < a_0 < 125 \times 10^{-4}$ , again well beyond the known invariant region. Now, however, the function takes on both positive and negative values inside each region, depending on

the  $\theta_o$  value. There is also no striking pattern of monotonicity of the function over the two regions. We can only say that the function  $S(a_o, e_s, \theta_o)$  has a varied (continuous) dependence on  $a_o$  in both regions. For the inner region we notice that the function  $\Delta r$  is generally positive and rather small for all angles  $\theta_o$ . In fact the values of  $\Delta r$  range between  $+5.1 \times 10^{-4}$  and  $-1.3 \times 10^{-4}$  for all values of  $\theta_o$  (when  $a_o$  is in the inner region). For the outer region the fluctuation of  $\Delta r$  with varying  $\theta_o$  values is markedly greater. For example, for  $a_o = 125 \times 10^{-4}$  the range of values for  $\Delta r$  is  $+86 \times 10^{-4}$  to  $-81 \times 10^{-4}$ . Thus we can conclude that for the inner region the function  $S$  is nearly independent of the angle  $\theta_o$ , whereas for the outer region it is strongly dependent on both  $a_o$  and  $\theta_o$ .

All of these results for the behavior of our  $\beta$ ,  $R$ , and  $S$  functions (or Moser's  $\alpha$ ,  $F$ , and  $G$  functions if we convert to the polar coordinates  $r, \theta$  exclusively) can be summarized in the following relations for the two regions discussed above:

#### Inner (Invariant) Region

$$T^{55}: \left\{ \begin{array}{l} \theta_1 = \theta_o + \beta(a_o, e_s) + R(a_o, e_s, \theta_o) \\ r_1 = r_o + S(a_o, e_s, \theta_o) \end{array} \right\} \quad (31)$$

$$\left. \begin{array}{l} \text{where } \beta < 0, \quad R < 0, \quad |\beta| \ll |R| \\ \text{with } d\beta/da_o > 0, \quad d^2\beta/da_o^2 \approx 0 \\ \text{and } \partial R/\partial a_o > 0, \quad \partial^2 R/\partial \theta_o^2 > 0, \quad \partial^2 S/\partial \theta_o^2 \approx 0; \end{array} \right\} \quad (32)$$

Outer (Non-invariant) Region

$$T^{55}: \left\{ \begin{array}{l} \theta_1 = \theta_0 + \beta(a_0, e_s, \theta_0) + R(a_0, e_s, \theta_0) \\ r_1 = r_0 + S(a_0, e_s, \theta_0) \end{array} \right\} \quad (33)$$

$$\left. \begin{array}{l} \text{where } \beta > 0, \quad \partial\beta/\partial\theta_0 < 0, \quad \partial\beta/\partial a_0 > 0 \\ \text{and } R > 0, \quad \partial R/\partial a_0 > 0, \quad \partial^2 R/\partial\theta_0^2 < 0 \text{ for } \theta_0 \gtrsim 15^\circ \end{array} \right\} \quad (34)$$

We are now in a position to check the applicability of Moser's second and third conditions (equations (17) and (18)) to our  $T^{55}$  mapping. First of all, it is clear that  $\beta$  (or  $\alpha$ ) is a monotonically increasing function of  $a_0$  over the full range of  $a_0$  values plotted. The value of  $d\beta/da_0$  is positive and approximately constant in the inner region, while  $\partial\beta/\partial a_0$  is positive, increasing and bounded in the outer region. Therefore if we choose any number greater than the upper bound of  $\partial\beta/\partial a_0$  as the constant  $C_0$  in equation (17), then Moser's second condition will be satisfied by our  $T^{55}$  mapping over the open interval  $0 < a_0 < 125 \times 10^{-4}$ .

In order to check Moser's third condition we first observe that both the  $R$  and  $S$  functions are bounded and continuous over the open interval  $0 < a_0 < 125 \times 10^{-4}$ . Thus we can choose some finite value of  $\delta_0$  which is greater than the sum of the absolute bounds of the  $R$  and  $S$  functions. From the smooth, continuous behavior of the  $\beta$ ,  $R$ , and  $S$  functions we can estimate the sum of the absolute values of the  $l = 333^{\text{rd}}$  partial derivatives of these three functions as approximately zero, and in any case less than the value of  $C_0$  chosen above. Hence Moser's third condition is satisfied by our  $T^{55}$  mapping over the open interval  $0 < a_0 < 125 \times 10^{-4}$ .

Because all of Moser's conditions are satisfied, we may conclude that invariant closed curves and regions will exist somewhere within the

standard ellipse having  $a_0 = 125 \times 10^{-4}$  and  $e_s = 0.995635$ . Indeed we have already found an invariant region which is bounded by the union of the hyperbolic invariant curves of the  $T^{55}$  mapping and which is located approximately within the standard ellipse having  $a_0 = 50 \times 10^{-4}$ . However we have also seen that points outside the latter ellipse but still inside the larger ellipse having  $a_0 = 125 \times 10^{-4}$  do escape to infinity in a finite time. Since Moser's conditions apply to this outer, unstable (non-invariant) region as well, we conclude that Moser's conditions are not sufficient to locate the invariant region precisely.

The empirical properties summarized in equations (31) through (34) exhibit several distinctions between the inner invariant region and the outer non-invariant region for the  $T^{55}$  mapping. For example, the transition of  $\beta$  from a linear to a nonlinear function and the abrupt change of magnitude and functional dependences of  $S$  may be possible theoretical clues to the exact location of the invariant region. At this point it should be noted that exhibiting the  $\beta$ ,  $R$ , and  $S$  functions involves as much work as actually locating the invariant region accurately. But an adequate theorem would definitely be valuable for rigorously proving the actual invariance of the region.

### C. Summary of methods

Having shown that invariant regions of stability do exist, and that at present they can be located most readily by using a computer, let us now summarize our practical, empirical methods for exhibiting the existence and form of the invariant curves.

The first step is to find the fixed points by scanning the phase-plane in some systematic fashion, such as repeatedly mapping a grid or mesh

of points. If these mappings are superimposed on the original grid, their intersections will yield all the fixed points in the phase-space. For this work a computer program based on numerical integration (e.g. a Runge-Kutta-Gill routine) should be quite adequate.

Next we can look for invariant regions around any particular one of these fixed points. If  $T^n(\tau)$  refers to the mapping under which the given point is fixed, then we can apply this mapping a number of times  $m$  in succession to each of a series of points along an axis through the fixed point (e.g.  $m=5$  or 10 times). This initial scan can be rather coarse, but should cover a reasonable neighborhood around the fixed point (for example, a radial range of perhaps  $10^{-8}$  to  $10^{-2}$  or more). Now for each value of  $m$  plot the locus of that particular  $m^{\text{th}}$  map of the entire series of points along the axis. If the original point is strongly unstable, then the maps of points along the axis will jump about discontinuously and may quickly reach infinity (even for  $m=2$  or 3). But if the mapped loci of the axis points seem to possess some form of a general elliptical twist mapping of that axis, then look for any repeating patterns in the successive  $T^n$  mappings of each initial axis point by itself. For example, do the radii or especially the angles of the mapped points seem to exhibit similarities every  $m$  cycles of the mapping  $T^n$ ? If so, for what regions? If not, make a finer grid and scan closer to the original fixed point; increase the number  $m$  of successive  $T^n$  mappings of each point.

If a pattern eventually emerges, it will probably be indicative of precisely which  $m^{\text{th}}$  cluster of  $T^{mn}(\tau)$  satellite multiple points of the original fixed point bounds the invariant region (and which hyperbolic invariant curves make up the doubly-periodic invariant curve around that

region). If no pattern ever emerges, even though the mappings are quasi-stable, then the original point is probably an elliptic point which is nevertheless unstable. Again we should stress that the appearance of an elliptical twist mapping around an elliptic fixed point seems to be a necessary (but not necessarily sufficient) condition for the existence of an invariant region. The appearance of a definite pattern in cycles of every  $M$  mappings of  $T^n$  seems to be the chief indicator of the particular mapping  $T^{Mn}(\tau)$  which exhibits the seemingly necessary twist mapping conditions.

To obtain the invariant region more accurately, we can first scan all our initial  $T^{Mn}$  twist maps of the axis for all mutual intersections to see whether or not we already may have found any of the  $2M$  fixed points of the  $T^{Mn}$  mapping of the axis. If not, we will have to make a crude mapping of the axis for points whose radii coincide approximately with the outer bounds of the region which exhibits the  $M$ -fold mapping patterns. In any case, once approximate values for the two  $T^{Mn}$  fixed points on the axis are found, their exact location can be determined as accurately as desired by taking successively finer and finer grids of initial points. By making repeated  $T^n$  mappings of both of these points for a total of  $(M-1)$  times each, the entire cluster of  $2M$  alternating hyperbolic and elliptic  $T^{Mn}$  satellite fixed points can be located. These will bound the invariant region.

We now make some preliminary short-period  $T^{Mn}$  stroboscopic mappings (perhaps for 50 to 100 strobe periods) of a series of points on the axis and inside the cluster of  $T^{Mn}$  satellite points. These mappings will show the elliptical character of the inner motions, plus the sense of

their stroboscopic rotations (which reverse outside the region). If, in fact, we map a coarse grid of inner axis points which extend towards the  $T^{Mn}$  elliptic fixed point on that axis, we will find that those grid points farthest from the  $T^{Mn}$  fixed point will move nearly elliptically around the inner region in one direction. Grid points closer to the  $T^{Mn}$  fixed point will move outward around that elliptic fixed point itself, in an opposite sense. By taking finer and finer grids we can locate (to the accuracy of the computer) the approximate borderline between these motions. Now, longer-period stroboscopic mappings (e.g. for 500 to 1000 strobe periods of  $T^{Mn}$ ) of two points on opposite sides of this borderline will trace out both the inner union of hyperbolic invariant curves (if we use  $T^{Mn}$  printouts), plus all the outer hyperbolic invariant curves surrounding the  $M$  elliptical cells or "islands" (if we use  $T^n$  printouts). The union of these  $(M+1)$  curves will be the one continuous doubly-periodic invariant curve of  $T^{Mn}$  which bounds the invariant region, in the limit that our two initial points approach the inner invariant curve.

All of the preceding steps can be carried out quickly and efficiently utilizing only a moderately accurate computer and a numerical integration routine (such as our Fortran program included in the second appendix). In special cases some numerical analysis may be necessary to obtain sufficient accuracy for the longer-period stroboscopic mappings (needed to generate the actual closed invariant curve). But in general the preceding methods should enable one to easily and relatively accurately determine the invariant regions and hence the stroboscopic stability of any particular solution (fixed point) of a particular differential equation.

If more detailed proof of the stroboscopic stability of a particular

solution is desired, then we must study carefully the hyperbolic invariant curves of all  $T^{mn}$  mappings (at least up to  $m=M$ ) of a global region around our particular fixed point. By mapping the inner and outer hyperbolic invariant curves "layer by layer" (i.e. for consecutive values of  $m$ ), we can check for intersections of these consecutive  $T^{mn}$  invariant curves. Thus we can gradually exhibit the mechanisms or "channels" whereby points can escape to infinity from regions close to our particular fixed point. After "peeling off" each successive "layer" until we reach one layer whose  $T^{Mn}$  hyperbolic invariant curves do not intersect those of the next outer layer, we can conclude that the escape mechanism no longer holds for this layer. (That is, the hyperbolic invariant curves of the  $T^{Mn}$  mapping develop no oscillations with which to carry points out to infinity.)

To explicitly demonstrate the non-oscillatory behavior of these  $T^{Mn}$  hyperbolic invariant curves would require more accuracy still. However, this seems not to be important for our conclusion that the union of the hyperbolic invariant curves of the  $T^{Mn}$  mapping isolated above will therefore form the continuous, closed, doubly-periodic boundary of an invariant region of absolute stroboscopic stability.

## IV. APPENDIX -- ANALYSIS

A. Integration of the equations of motion

The first integral of equation (20), when  $p=p_0$  and  $X_a$  is an intercept on the X-axis, is simply

$$(dX/dt)^2 + \frac{1}{2}p_0 X^4 = \frac{1}{2}p_0 X_a^4 > 0 . \quad (35)$$

A half-period of motion described by this equation in phase-space is an arc of a quasi-ellipse which intersects the X-axis at  $X_a$  and is concave inward. Similarly when  $p=-p_0$  and  $X_b$  is another X-axis intercept, the first integral of equation (20) is

$$(dX/dt)^2 - \frac{1}{2}p_0 X^4 = -\frac{1}{2}p_0 X_b^4 < 0 . \quad (36)$$

A half-period of motion described by this equation in phase-space is an arc of a quasi-hyperbola which intersects the X-axis at  $X_b$  and is concave outward. Finally when  $p=-p_0$  and  $\dot{X}_c$  is an intercept on the  $\dot{X}$ -axis, the first integral of equation (20) is

$$(dX/dt)^2 - \frac{1}{2}p_0 X^4 = \dot{X}_c^2 = \frac{1}{8}p_0 c^4 > 0 \quad (37)$$

where  $c$  is a real constant. Again a half-period of the motion is an arc of a quasi-hyperbola which now intersects the  $\dot{X}$ -axis at  $\dot{X}_c$  and is also concave outward. Notice that a transition must occur between the last two types of quasi-hyperbolic motion whenever  $\dot{X}^2 - \frac{1}{2}p_0 X^4 = 0$ , i.e. when

$$\dot{X} = \pm \sqrt{p_0/2} X^2 . \quad (38)$$

The motion described by this last equation consists of two parabolas through the origin in phase-space. It reveals that the origin is a parabolic (nodal) fixed point of equation (20).

Equations (35), (36), and (37) above can be reduced to a standard form

$$\dot{q}^2 = \frac{1}{2} p_0 C_i (1 - q^4) \quad (39)$$

where  $C_i$  is  $X_a$ ,  $X_b$ , or  $c$  respectively. This standard form can be integrated in terms of the Jacobian elliptic functions snu, cnu, and dnu. This integration has been done by Bartlett<sup>13/</sup>, and the resulting solutions for the three basic types of motion considered above are reproduced here.

Quasi-Elliptic Case (QE Motion):

$$X = X_a \text{ cnu} \quad (40a)$$

$$\dot{X} = -\sqrt{p_0} X_a^2 \text{ snu dnu} \quad (40b)$$

$$u = \sqrt{p_0} X_a t \quad (40c)$$

Quasi-Hyperbolic Case I (QH I Motion):

$$X = X_b / \text{cnu} \quad (41a)$$

$$\dot{X} = \sqrt{p_0} X_b^2 \text{ snu dnu} / \text{cn}^2 u \quad (41b)$$

$$u = \sqrt{p_0} X_b t \quad (41c)$$

Quasi-Hyperbolic Case II (QH II Motion):

$$X = \sqrt{c^2/2} \text{ snu} / (1 + \text{cnu}) \quad (42a)$$

$$\dot{X} = \sqrt{p_0/2} c^2 \text{ dnu} / (1 + \text{cnu}) \quad (42b)$$

$$u = \sqrt{p_0} c t \quad (42c)$$

A number of theoretical predictions based on these results can be made about the general properties of the motion<sup>13/</sup>, but we will not dwell on them here. Instead we will go directly to the more specific problem of numerically calculating many highly accurate solutions  $(X(t), \dot{X}(t))$  from any arbitrary set of initial conditions  $(X_i, \dot{X}_i)$ .

## B. Jacobian elliptic functions

The first obvious problem in doing our numerical calculations is one of obtaining accurate tables of the Jacobian elliptic functions  $\text{snu}$ ,  $\text{cnu}$ , and  $\text{dnu}$ . Each of these functions can be written as products and quotients of the four Jacobian Theta functions, which in turn are expressible as infinite series. The resulting formulas, taken from the Smithsonian Elliptic Function Tables<sup>20/</sup>, are reproduced here for reference.

$$\text{snu} = \frac{\theta_3(0,q)}{\theta_2(0,q)} \times \frac{\theta_1(v,q)}{\theta_0(v,q)} \quad (43)$$

$$\text{cnu} = \frac{\theta_0(0,q)}{\theta_2(0,q)} \times \frac{\theta_2(v,q)}{\theta_0(v,q)} \quad (44)$$

$$\text{dnu} = \frac{\theta_0(0,q)}{\theta_3(0,q)} \times \frac{\theta_3(v,q)}{\theta_0(v,q)} \quad (45)$$

where

$$\theta_0(v,q) = 1 - 2q \cos 2\pi v + 2q^4 \cos 4\pi v - 2q^9 \cos 6\pi v + \dots \quad (46)$$

$$\theta_1(v,q) = \sin \pi v - q^2 \sin 3\pi v + q^6 \sin 5\pi v - q^{12} \sin 7\pi v + \dots \quad (47)$$

$$\theta_2(v,q) = \cos \pi v + q^2 \cos 3\pi v + q^6 \cos 5\pi v + q^{12} \cos 7\pi v + \dots \quad (48)$$

$$\theta_3(v,q) = 1 + 2q \cos 2\pi v + 2q^4 \cos 4\pi v + 2q^9 \cos 6\pi v + \dots \quad (49)$$

with

$$v = u/2K \quad (50)$$

and

$$\left\{ \begin{array}{l} K = 1.854074677301372\dots \\ q = 0.04321391826377225\dots \\ \pi = 3.14159265358979323846\dots \end{array} \right\} \quad (51)$$

Since highly accurate routines exist for the calculation of the circular functions used above, these formulas are admirably suited to a high-speed, highly accurate digital computer such as the Illiac II.

In our work, tables for the elliptic functions  $\text{snu}$ ,  $\text{cnu}$ , and  $\text{dnu}$  are calculated and stored by the computer in terms of the "angular" argument  $b$ , where

$$u = K(b/90) \quad (52)$$

so that

$$v = b/180 \quad (53)$$

The tables are calculated for the first quadrant values of  $b$  (i.e. for  $0.0 \leq b \leq 90.0$ ) in steps of 0.1 in  $b$ . Using the periodic properties of the elliptic functions (e.g.  $\text{snu}$  and  $\text{cnu}$  are analogous to  $\sin u$  and  $\cos u$ , with  $K$  playing the role of  $\pi/2$ ) and using the first quadrant tables calculated above, we can easily obtain any of the elliptic functional values for the full range  $0 \leq u \leq 4K$ . For example, it will be necessary to calculate and store a table of values of the function

$$S(u) = \text{snu}/(1 + \text{cnu}) \quad (54)$$

for the range  $0 \leq u \leq 3K/2$ . For the range  $0 \leq u \leq K$  it is a trivial calculation. For the range  $K < u \leq 3K/2$  the table is assembled by calculating the function  $\text{snu}/(1-\text{cnu})$  starting with the  $u=K$  end of the first quadrant tables and going backwards toward  $u=K/2$  (i.e.  $b=45.0$ ). These results are then stored as  $S(u)$  starting at  $u=K$  and going forwards toward  $u=3K/2$  (i.e.  $b=145.0$ ).

At this point we should remark that if only a few values of the elliptic functions are needed, it would be most efficient to use equations (43) through (53) each time the need arose. However for an inordinately

large number of required values (e.g. as many as  $10^5$  per trajectory, as in the present work), it becomes more efficient time-wise to construct relatively small, accurate tables and then use an accurate (nonlinear) interpolation routine. In this work we make use of a Lagrange six-point (equal interval) interpolation routine<sup>21/</sup>, which has an average execution time of 375 microseconds per entry on the IBM-7094 computer. This method utilizes the three nearest table entries on each side of the input value of  $u$  and then interpolates the (elliptic) functional value corresponding to  $u$  from a sixth-order fit of the 6 tabulated functional values. The results are then nearly as accurate as if they had been calculated from the original infinite series expressions. It should be noted that this method allows us to use relatively small tables (our tables contained 900 entries for the range  $0 \leq u \leq K$ ) and still obtain highly accurate interpolations. For a small computer memory core this efficiency is essential.

### C. Inversion of elliptic functions

A second major problem which we will face in the actual calculations is that of obtaining accurate values for the inverted Jacobian elliptic functions. Instead of using the obvious elliptic integrals (i.e. the formal inversions of the elliptic functions) to calculate still more tables, we will consider several alternatives which are more efficient with regard to saving time and memory core. One alternative is to adapt a Lagrangian six-point (unequal interval) interpolation routine to the calculations. This approach is the desirable one for highest accuracy and overall efficiency. In the present work, however, it has been more expedient to follow a different approach, one of approximation and iteration. From equations (40a,b,c), (41a,b,c), and (42a,b,c) for the three different types of motion

it is seen that we only have to invert  $\text{cnu}$  (for both QE and QH I Motions) and the  $S(u) = \text{snu}/(1+\text{cnu})$  function (for QH II Motion only).

Given some value of  $\text{cnu}_i$  which we wish to invert in order to obtain  $u_i$ , let us quickly scan our  $\text{cnu}$  table and determine between which two entries our value of  $\text{cnu}_i$  lies. Then let us choose the one with the smaller  $b$  value, call it  $b_o = 90u_o/K$  (where  $u_o < u_i$  is the nearest smaller value of  $u$ , with  $\text{cnu}_o > \text{cnu}_i$  the nearest larger value of  $\text{cnu}$ ). Now let us write a Taylor expansion for  $\text{cnu}_i$  about the point  $\text{cnu}_o$ . Using the properties of the elliptic functions<sup>22/</sup> we obtain

$$\text{cnu}_i = f_o - g_o \Delta - 1/2 f_o^3 \Delta^2 + 1/2 g_o f_o^2 \Delta^3 + 1/8 (f_o^5 - 2f_o g_o^2) \Delta^4 + \dots \quad (55)$$

where

$$\left\{ \begin{array}{l} f_o = \text{cnu}_o \\ g_o = \text{snu}_o \, d\text{nu}_o \\ \Delta = (u_i - u_o) \end{array} \right\} \quad (56)$$

For the first approximation we neglect all nonlinear terms and obtain

$$\Delta_1 = (f_o - \text{cnu}_i)/g_o \quad (57)$$

Neglecting terms beyond  $\Delta^2$  and using  $\Delta_1$  in the second-order term, we obtain the second-order approximation

$$\Delta_2 = \Delta_1 - (f_o^3 \Delta_1^2)/2g_o \quad (58)$$

Similarly we obtain the third-order approximation

$$\Delta_3 = \Delta_2 + \frac{1}{2} f_o^2 \Delta_2^3 - (f_o^6 \Delta_1^3)/2g_o^2 + (f_o^5 \Delta_1 \Delta_2^3)/2g_o + \dots \quad (59)$$

Given  $\text{cnu}_i$  and having chosen the nearby smaller base value  $u_o$  we can readily calculate first  $\Delta_1$ , then  $\Delta_2$ , and finally  $\Delta_3$ . The third-order approximation

for  $u_i$  is then

$$u_{i3} = u_o + \Delta_3 . \quad (60)$$

In the second case of inversion, if we are given a value for the function  $S(u_i) = \text{sn}u_i / (1 + \text{cn}u_i)$  and we wish to invert it to obtain  $u_i$ , we can proceed as follows. We first make a Taylor expansion about the point  $u_o$  (chosen as before with  $u_o < u_i$ , but now with  $S(u_o) < S(u_i)$ ) and obtain

$$S(u_i) = f_o + g_o \Delta + \frac{1}{2} f_o^3 \Delta^2 + \frac{1}{2} f_o^2 g_o \Delta^3 + (3f_o g_o^2 + 3f_o^5/4) \Delta^4/24 + \dots \quad (61)$$

where now

$$\left\{ \begin{array}{l} f_o = \text{sn}u_o / (1 + \text{cn}u_o) = S(u_o) \\ g_o = \text{dn}u_o / (1 + \text{cn}u_o) \\ \Delta = (u_i - u_o) . \end{array} \right\} \quad (62)$$

Analogous to the case of the  $\text{cn}u_i$  inversion, we can invert this expansion by means of successive approximations to obtain

$$\Delta_1 = (S(u_i) - f_o) / g_o \quad (63)$$

$$\Delta_2 = \Delta_1 - (f_o^3 \Delta_1^2) / 4g_o \quad (64)$$

$$\Delta_3 = \Delta_2 - \frac{1}{2} f_o^2 \Delta_2^3 + (f_o^6 \Delta_1^3) / 8g_o^2 + (f_o^5 \Delta_1 \Delta_2^3) / 8g_o + \dots . \quad (65)$$

Given  $S(u_i) = \text{sn}u_i / (1 + \text{cn}u_i)$  we can readily calculate  $\Delta_1$ , then  $\Delta_2$ , and finally  $\Delta_3$ , yielding the third-order approximation for  $u_i$  as in equation (60).

In the special case that  $u_o = 0$  or  $u_o = K$  (the ends of the table) we can simplify our original expansions for  $\text{cn}u_i$  and  $S(u_i)$  given in equations (55) and (61). Recalling that  $\text{sn}0=0$  and  $\text{cn}0=\text{dn}0=1$ , we obtain the third-order expressions

$$u_{i3} = \sqrt{2(1 - \text{cn}u_i)} \quad \text{for } u_o = 0 , \quad (66a)$$

and 
$$u_{i3} = K - \sqrt{2} \text{cnu}_i \quad \text{for } u_0 = K ; \quad (66b)$$

also 
$$u_{i3} = 2S(u_i) \quad \text{for } u_0 = 0 , \quad (67a)$$

and 
$$u_{i3} = K - \sqrt{2} (1 - \sqrt{2S(u_i)-1}) \quad \text{for } u_0 = K . \quad (67b)$$

Before we start the calculations we can see that our above successive approximations for the  $\Delta_i$  in the case of the  $\text{cnu}_i$  inversion do not converge as  $u_0 \rightarrow 0$ . This divergence is due to the presence of  $\text{snu}_0$  in the denominators of all of the expressions. Since the approximations are increasingly accurate toward the middle of the tables, it is clear that we must devise an accurate (converging) approximation for the case  $u_0 \rightarrow 0$  in the  $\text{cnu}_i$  inversion. To this end we make use of a direct expansion for  $\text{cnu}$  in powers of  $u^{22/}$ :

$$\text{cnu} = 1 - \frac{1}{2}u^2 + u^4/8 - 3u^6/80 + \dots \quad (68)$$

Forming the quotient

$$R(u) = \frac{1 - \text{cnu}}{1 + \text{cnu}} = \frac{\frac{1}{2}u^2 - u^4/8 + 3u^6/80}{2 - \frac{1}{2}u^2 + u^4/8 - 3u^6/80} \quad (69)$$

and performing the indicated long-division, we obtain

$$R(u) = \frac{1 - \text{cnu}}{1 + \text{cnu}} = \frac{u^2}{4} + \frac{u^6}{320} + \frac{7u^8}{1280} + \dots \quad (70)$$

If we drop all terms higher than  $u^2$  we get the first-order approximation for  $u_i^2$  to be

$$u_{i1}^2 = 4(1 - \text{cnu}_i)/(1 + \text{cnu}_i) . \quad (71)$$

Inserting this value in the two higher-order terms in equation (70) and solving for the resulting approximation, we obtain

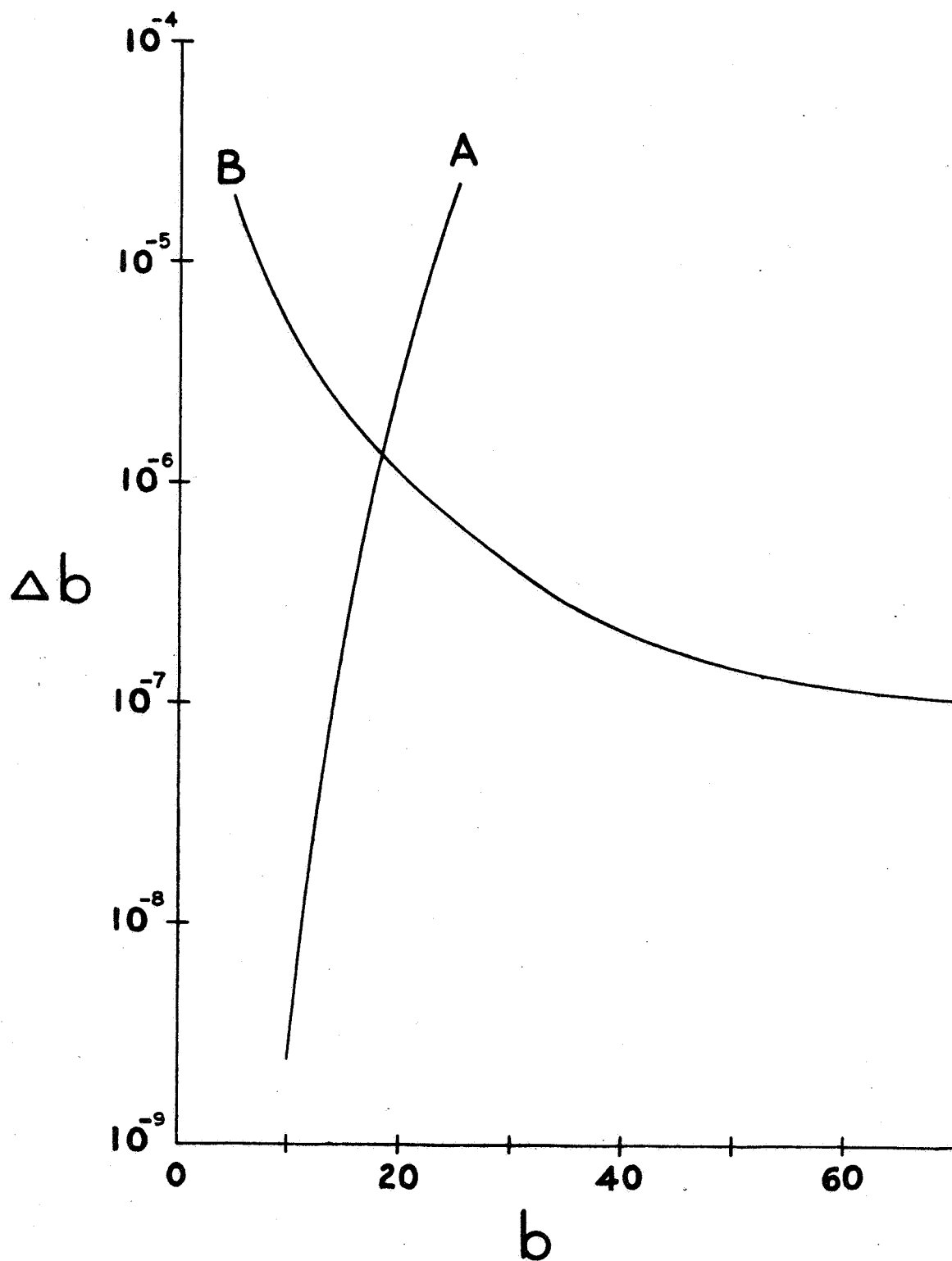
$$\frac{1}{2}u_{i2}^2 = R(u_i) - R^3(u_i)/5 - 7R^4(u_i)/5 - \dots \quad (72)$$

When  $u_i \rightarrow 0$  both sides of the last equation converge to zero, so that we have a valid expression for small  $u_i$  in terms of the given  $cnu_i$  value. When  $u_i$  gets larger, this expression will become decreasingly accurate. On the other hand, our original equations for  $cnu_i$  inversion (equations (56) to (60)) will become increasingly accurate. Clearly there will be some crossover point between the accuracies of the two inversion methods, and we will need to know this point if we are to use both methods for the full range of the table values. To find this crossover point we must do some hand calculations (faster than writing a small computer program in this case). We can select, for example, an integral value of  $b_i$  and then take the accurate table value for the associated  $cnu_i = cn(Kb_i/90)$ . Next we can invert this  $cnu_i$  value using both inversion methods. These results for  $u_{i3}$  (or  $u_{i2}$ ) can then be compared with the original value of  $u_i = (Kb_i/90)$ . Thus we can tabulate the error  $\Delta b = |b_i - b_{i3}|$  as a function of the original  $b_i$  for both inversion methods. A semi-log plot of these results is shown in Figure 12. Curve B is the error plot for the original inversion method, while curve A is for the second method for small angles. From this plot we find the crossover value of  $b_i$  as

$$b_c = 18.2 \quad (73)$$

In our calculations we use the original  $cnu_i$  inversion method (equations (56) to (60)), for  $b_i > 18.2$  and the second inversion method of equation (72) for  $b_i \leq 18.2$ . Note that the maximum error in  $u_{i3}$  will then be about  $10^{-6}$  when  $b_i = b_c$ .

Figure 12. Numerical errors arising in two methods  
for the inversion of Jacobian elliptic  
functions.



In order to obtain more accurate inversions than the above  $u_{i3}$ , we now add an iteration process. After inverting  $cnu_i$  or  $S(u_i)$  to obtain the approximate value of  $u_{i3}$  (or  $u_{i2}$ ) we turn around and calculate the corresponding value of  $cnu_{i3}$  or  $S(u_{i3})$  using the accurate Lagrangian interpolation routine and our original tables. The resulting value can then be compared with the given value of  $cnu_i$  or  $S(u_i)$  that we initially tried to invert. If  $|cnu_{i3} - cnu_i| \leq E$  or  $|S(u_{i3}) - S(u_i)| \leq E$ , where  $E$  is some preset error limitation, then we can regard  $u_{i3}$  (or  $u_{i2}$ ) as close enough to  $u_i$  for our purposes. In our work we eventually set  $E = 10^{-12}$ .

However, if the error is greater than  $E$  we can make a new estimate for  $u_i$  by expanding  $cnu_i$  or  $S(u_i)$  in a Taylor series about  $u_{i3}$  or  $u_{i2}$ , namely

$$cnu_i = cnu_{i3} - snu_{i3} dnu_{i3} (u_i - u_{i3}) + \dots \quad (74)$$

and

$$S(u_i) = S(u_{i3}) + \left( \frac{dnu_{i3}}{1 + cnu_{i3}} \right) (u_i - u_{i3}) + \dots \quad (75)$$

Solving each of these equation for  $u_i$  we obtain

$$u_{i4} = u_{i3} + \frac{(cnu_{i3} - cnu_i)}{snu_{i3} dnu_{i3}} + \dots \quad (76)$$

and

$$u_{i4} = u_{i3} + \frac{(S(u_i) - S(u_{i3})) (1 + cnu_{i3})}{dnu_{i3}} + \dots \quad (77)$$

as the next approximations for the inversions of  $cnu_i$  and  $S(u_i)$  respectively (provided that  $u_{i3}$  or  $u_{i2}$  are not zero). This procedure can be repeated over and over again until the desired accuracy is achieved. In practice it has not been necessary to use these corrections more than once.

Although the previous analysis should be complete and contain no sources for computational errors (greater than  $10^{-12}$ ), one final error did arise. This error manifested itself in the eighth decimal place after about  $10^5$  quarter-period calculations, but only for a few random initial conditions. Its effects were observed as random discontinuities in certain plots that should have been smooth ones. After many careful hand-calculations, the error was tracked down to the random truncation error in the 13<sup>th</sup> decimal place by the Illiac II computer! This seemingly insignificant computer "fact-of-life" has overwhelming effects on the  $\text{cnu}_i$  inversion process when  $\text{cnu}_i$  is very close to one. If under QE or QH I motion we have  $X_i \approx X_a$  or  $X_b$  within about one part in  $10^{13}$ , then the ratio  $\text{cnu}_i = X_i/X_a$  or  $\text{cnu}_i = X_b/X_i$  will be about 1.0 to the same accuracy. But the computer would sometimes call this exactly 1.000..., while at other times it would truncate it to  $0.999... = 1 - 10^{-13}$  (this is at the limit of its accuracy). Again this may seem inconsequential, but when we look at the inversions of these slightly different  $\text{cnu}_i$  values we find that  $\text{cn}^{-1}(1.0) = 0.00$  while  $\text{cn}^{-1}(1.0 - 10^{-13}) \approx 10^{-7}$ . So suddenly we see the cause of the random discontinuities!

To overcome this problem we refer back to equations (35) and (36) for  $X_a$  and  $X_b$  in terms of the initial coordinates  $X_i$  and  $\dot{X}_i$ . Factoring these expressions and expanding them by means of the binomial theorem, we obtain respectively

$$X_a = X_i(1+\delta)^{\frac{1}{2}} = X_i(1 + \frac{1}{2}\delta - 3\delta^2/32 + \dots) \quad (78)$$

and

$$X_b = X_i(1-\delta)^{\frac{1}{2}} = X_i(1 - \frac{1}{2}\delta + 5\delta^2/32 - \dots) , \quad (79)$$

where

$$\delta = 2\dot{X}_i^2/p_o X_i^4 \ll 1 . \quad (80)$$

From equations (40a) and (41a) we obtain

$$\text{cnu} = X_i/X_a = (1 - \frac{1}{4}\delta + 3\delta^2/32 - \dots) \quad \text{for QE} \quad (81)$$

and

$$\text{cnu} = X_b/X_i = (1 - \frac{1}{4}\delta + 5\delta^2/32 - \dots) \quad \text{for QH I} \quad (82)$$

Thus to first-order in  $\delta$  (fourth-order in  $X_i$ ) we obtain for both QE and QH I motions

$$\text{cnu} \approx 1 - \frac{1}{4}\delta \quad (83)$$

Since our given value of  $\text{cnu}_i$  differed from 1.0 by about  $10^{-13}$ , the above approximation to first-order in  $\delta$  is justified. Note that this value for  $\text{cnu}_i$  requires  $u_0 = 0$  in our  $\text{cnu}_i$  inversion routines. Then using equation (66a) for  $u_{i3}$  when  $u_0 = 0$  we obtain

$$u_{i3} = \sqrt{2(\frac{1}{4}\delta_i)} = \sqrt{\delta_i/2} \quad (84)$$

or

$$u_{i3} = \pm \hat{X}_i / (\sqrt{p_0} X_i^2) \quad (\text{for } u_0 = 0). \quad (85)$$

Thus whenever

$$1.0 \geq \text{cnu}_i \geq \text{cn}(K/900) \quad (86)$$

(where  $\text{cn}(K/900)$  is the first table entry for  $b = 0.1$ ), the value of  $u_{i3}$  would be calculated directly from  $(X_i, \hat{X}_i)$  using equation (85), thus bypassing the regular inversion routines.

#### D. Piecewise calculations

Let us now turn to the actual procedures for the numerical calculations. In order to limit the amount of computations involved, we first choose fixed numerical values for our two parameters  $p_0$  and  $\tau$ , namely

$$p_0 = 0.037 \quad \text{and} \quad \tau = 6.0 . \quad (87)$$

These are arbitrarily chosen values, but they will generate solutions which can be scaled (both in phase-space and in time) in order to yield solutions for other values of  $p_0$  and  $\tau$ . For more general equations with non-scaling parameters, or for non-integrable equations which possess some constants of the motion (e.g. the energy), we would have to use different values for those parameters or constants of the motion and calculate all the corresponding solutions. A complete representation of the structure of periodic solutions to the equations could then be obtained by plotting the phase-space trajectories for many different values of the parameters or constants of the motion. Treating these parameters then as continuous variables, we could thus generate sets of "eigensurfaces" of solutions in phase-parameter space<sup>23/</sup>. But for our equation (20) we need select only one value for the parameters.

Since our numerical calculations will begin at  $t=0$ , it is clear from our definition of  $p(t)$  in equation (21) that for the first quarter-period ( $\tau/4 = 1.5$ ) we will have  $p(t) = p_0 > 0$  and hence a quasi-elliptical (QE) motion. Then for the next half-period ( $\tau/2 = 3.0$ ) we will have  $p(t) = -p_0 < 0$  and hence a quasi-hyperbolic (QH I or QH II) motion. Next will follow another half-period of QE motion, and so on for as long as we wish to calculate. A simple index counter can be employed to keep track of the sign of  $p(t)$  as a function of the number of quarter-periods of elapsed time. Note that we could have defined  $p(t)$  so that  $p(0) = -p_0$  instead of  $+p_0$  simply by shifting our  $p(t)$  by a half-period in time. Then the motion would have begun as QH I or QH II motion for the first quarter-period. We have arbitrarily chosen  $p(t)$  here so that we will begin with QE motion at  $t=0$ .

Before presenting a summary of the program logic for the computer calculations, it will be to our advantage to discuss in detail a few typical solutions for the possible motions. These will then illustrate the various features of the computational techniques. Let us first choose some set of initial coordinates  $(X_i, \dot{X}_i)$  in our two-dimensional phase-space, for  $t=0$ . (In our program we make use of the polar coordinates  $r_i$  and  $\theta_i$  as input variables, where  $X_i = r_i \cos\theta_i$  and  $\dot{X}_i = r_i \sin\theta_i$ .) Specifically, let us choose this initial phase-space point to lie somewhere in the first quadrant and near the positive X-axis. This choice is not necessary in the general calculations.

Using the initial coordinates  $(X_i, \dot{X}_i)$  we can calculate the X-axis intercept  $X_a$  from the first integral for QE motion, given in equation (35). The result for  $X_i \neq 0$  is

$$X_a = +(2\dot{X}_i^2/p_o + X_i^4)^{1/4} \geq X_i . \quad (88)$$

(The case when  $X_i = 0$  is discussed in the next section.) Next, using equations (40c) and then (40a) we can compute the time it would take to go from  $(X_i, \dot{X}_i)$  to  $(X_a, 0)$ , namely

$$t_a = \left( \sqrt{p_o} X_a \right)^{-1} \text{cn}^{-1}(X_i/X_a) \geq 0 . \quad (89)$$

Now if  $t_a = 1.5$  then the initial QE motion would just be completed at  $(X_a, 0)$ . If, however,  $t_a < 1.5$  then the motion will continue on beyond (below) the X-axis for a time interval  $t_1 = (1.5 - t_a) > 0$ . To find the final coordinates (when  $t = 1.5$ ) we simply use (40c) and then (40a) and (40b). Thus

$$X_1 = X_a \text{cnu}_1 > 0 \quad (90a)$$

$$\dot{X}_1 = - \sqrt{p_o} X_a^2 \text{snu}_1 \text{dnu}_1 < 0 \quad (90b)$$

where

$$u_1 = \sqrt{p_0} X_a t_1 > 0 . \quad (90c)$$

If instead  $t_a > 1.5$  then the QE motion will fall short of (or above) the X-axis by the time interval  $t_1 = (1.5 - t_a) < 0$ . With  $t_1 < 0$  we clearly have  $u_1 < 0$  also. From the symmetry properties of the elliptic functions ( $\text{sn}(-u) = -\text{sn}u$ ;  $\text{cn}(-u) = \text{cn}u$ ;  $\text{dn}(-u) = \text{dn}u$ ) we see that for  $u_1 < 0$  we have  $X_1 > 0$  and also  $\dot{X}_1 > 0$ . The magnitudes of  $X_1$  and  $\dot{X}_1$  are found by using equations (90a) and (90b). These coordinates, then, are simply a reflection in the X-axis of the coordinates of the motion for an equal but positive interval of  $t_1$ . Indeed, inspection of our original equation (20) readily reveals that it is invariant under the time-reversed operation  $t \rightarrow -t$ , as well as under the space reflection operation  $X \rightarrow -X$ , or under both simultaneously. This invariance follows from equation (20) plus the symmetry property  $p(t) = p(-t)$  built into our periodic square-wave function in equation (21).

Now that we have determined our phase-space coordinates  $(X_1, \dot{X}_1)$  at the end of the first quarter-period, we can calculate the next half-period of quasi-hyperbolic motion. From equations (36) and (37) we see that it will be either QH I or QH II motion depending on our new initial coordinates  $(X_1, \dot{X}_1)$  satisfying the conditions:

$$|\dot{X}_1| < \sqrt{p_0/2} X_1^2 \quad \text{for QH I motion} \quad (91)$$

or

$$|\dot{X}_1| > \sqrt{p_0/2} X_1^2 \quad \text{for QH II motion.} \quad (92)$$

Note that if  $|\dot{X}_1| = \sqrt{p_0/2} X_1^2$  then the motion is parabolic (see the discussion of equation (38)).

Let us assume first that the motion is QH I. Paralleling the QE calculations we first find the intercept  $X_b$  on the X-axis

$$X_b = +(X_1^4 - 2\dot{X}_1^2/p_o)^{1/2} \leq X_1 . \quad (93)$$

The time to go from  $(X_1, \dot{X}_1)$  to  $(X_b, 0)$  would be

$$t_b = \left( \sqrt{p_o} X_b \right)^{-1} \text{cn}^{-1}(X_b/X_1) \geq 0 . \quad (94)$$

If  $t_b = 3.0$  then the QH I motion would be completed at the point  $(X_b, 0)$ .

But if  $t_b < 3.0$  then the motion will go beyond (above) the X-axis for a time interval  $t_2 = (3.0 - t_b) > 0$ . The final coordinates (when  $t = 1.5 + 3.0 = 4.5$ ) will thus be

$$X_2 = X_b / \text{cnu}_2 > 0 \quad (95a)$$

$$\dot{X}_2 = \sqrt{p_o} X_b^2 \text{snu}_2 \text{dnu}_2 / \text{cn}^2 u_2 > 0 \quad (95b)$$

where

$$u_2 = \sqrt{p_o} X_b t_2 . \quad (95c)$$

Similarly if  $t_b > 3.0$  the motion will fall short of (or below) the X-axis for a negative time interval  $t_2$ . The final coordinates will still be calculated from the above equations, but now  $\dot{X}_2 < 0$ .

For the next half-period the motion will be QE again. The calculations for  $(X_3, \dot{X}_3)$  will take the same form as equations (88), (89), and (90) except that  $(X_1, \dot{X}_1)$  will be replaced by  $(X_2, \dot{X}_2)$  and that (from now on) we must subtract  $t_a$  from 3.0 instead of 1.5. Thus in the present case we must calculate  $t_3 = (3.0 - t_a)$ . If  $(X_2, \dot{X}_2)$  had been in the first quadrant, then we would obtain  $X_3 > 0$ ,  $\dot{X}_3 < 0$  for  $t_3 > 0$  (or  $\dot{X}_3 > 0$  for  $t_3 < 0$ ). If  $(X_2, \dot{X}_2)$  had been in the fourth quadrant, a different result would be obtained:

$(X_3, \dot{X}_3)$  would be in either the fourth or the third quadrant. We will return to this case in the next section of this appendix.

But first let us consider the next half-period of quasi-hyperbolic motion to be QH II motion starting from initial coordinates  $(X_3, \dot{X}_3)$ . Let us for the moment assume  $X_3 > 0$  and  $\dot{X}_3 < 0$  (which required  $t_3 > 0$  plus the initial requirement that  $X_2 > 0, \dot{X}_2 > 0$ ). We proceed in a manner similar to the QE and QH I cases already discussed above. First we calculate the constant  $c^2$  from the first integral for QH II motion, equation (37):

$$c^2 = \pm (8\dot{X}_3^2/p_0 - 4X_3^4)^{\frac{1}{2}}. \quad (96)$$

Notice that the absolute value of the  $\dot{X}$ -axis intercept  $\dot{X}_c$  will be given by

$$|\dot{X}_c| = \sqrt{p_0 c^4/8} < |\dot{X}_3|. \quad (97)$$

The time to go from  $(X_3, \dot{X}_3)$  to  $(0, -|\dot{X}_c|)$  would be

$$t_c = \left( \sqrt{p_0} c \right)^{-1} f^{-1} \left( \sqrt{2} X_3/c \right) > 0 \quad (98)$$

where as before

$$f(u) = \text{snu}/(1 + \text{cnu}).$$

(As mentioned before it is convenient to tabulate the function  $f(u)$  for  $0 \leq u \leq 3K/2$  in order to facilitate taking its inverse.) If  $t_c = 3.0$  then the QH II motion would be completed at  $(0, -|\dot{X}_c|)$ . If however  $t_c < 3.0$  then the motion will go beyond (or to the left of) the  $\dot{X}$ -axis by a time interval  $t_4 = (3.0 - t_c)$  which is positive. The final coordinates for this motion will be

$$X_4 = - \sqrt{c^2/2} \operatorname{sn} u_4 / (1 + \operatorname{cn} u_4) < 0 \quad (99a)$$

$$\dot{X}_4 = - \sqrt{p_0/2} c^2 \operatorname{dn} u_4 / (1 + \operatorname{cn} u_4) < 0 \quad (99b)$$

where

$$u_4 = \sqrt{p_0} c t_4 > 0 \quad (99c)$$

and where the minus signs have been chosen to agree with QH II motion originating in the fourth quadrant (see the following section of this appendix). If  $t_c > 3.0$  the motion will fall short (or to the right) of the  $\dot{X}$ -axis for a negative time interval  $t_4$ . The final coordinates will still be calculated from the above equations, but now  $X_4 > 0$ .

#### E. Final coordinate signs

Now that we have discussed several typical motions for a particular set of initial coordinates, it is clear that one problem remains to be solved. That problem is the unambiguous determination of signs for the final coordinates of the QE, QH I, or QH II motions originating in any arbitrary quadrant of the phase-plane.

The simplest way to calculate the final coordinates from any initial coordinates is first to calculate the absolute values of all quantities (e.g.  $|\dot{X}_c|$ ,  $|t_c|$ ,  $|t_4|$ ,  $|X_4|$ , and  $|\dot{X}_4|$  for the last case discussed). Then we can attach the appropriate signs. These signs can be determined a priori from a consideration of the general properties of motion. Let us look at QE motion first. This motion is described in clockwise quasi-elliptic arcs about the origin, concave inward (see equation (88)). Thus if  $(X_1, \dot{X}_1)$  is in the first or third quadrant (i.e. if  $X_1 \dot{X}_1 > 0$ ) then the motion will be towards the X-axis, with an intercept  $X_a$ . Thus the results outlined in

equations (88) to (90) will hold. That is, the sign of  $X_1$  will equal the sign of  $X_i$ , while the sign of  $\dot{X}_1$  will be equal or opposite to the sign of  $\dot{X}_i$  according as  $|t_a| > 3.0$  (or 1.5 in the first step only) or  $|t_a| < 3.0$ . In the special cases where  $(X_i, \dot{X}_i)$  is on the  $\pm X$ -axis, then we will have  $\dot{X}_1 < 0$  for  $X_i > 0$ , and  $\dot{X}_1 > 0$  for  $X_i < 0$ .

If  $(X_i, \dot{X}_i)$  is in the second or fourth quadrants (i.e. if  $X_i \dot{X}_i < 0$ ) then the motion will be in clockwise arcs away from the  $X$ -axis, with an  $X$ -axis intercept  $X_a$  now in past time. Thus  $|t_a|$  is just the time it would have taken to go from  $(X_a, 0)$  on the axis to  $(X_i, \dot{X}_i)$ . Therefore the total time from the  $X$ -axis to the final point  $(X_i, \dot{X}_i)$  will be the sum  $t_1 = (3.0 + |t_a|)$ . To get the time it would take to go from  $(X_i, \dot{X}_i)$  to the  $\dot{X}$ -axis (i.e. the axis towards which the motion is directed), we can subtract  $|t_a|$  from the total time  $T_a$  it would take for a point to go from the  $X$ -axis initially to the  $\dot{X}$ -axis finally under QE motion. To find  $T_a$  we use equation (90a) and see that for  $X_1 = 0$  (on the  $\dot{X}$ -axis finally) we must have  $cnu_1 = 0$  or  $u_1 = \sqrt{p_0} X_a T_a = K$ . From equation (88) it is clear that for  $\dot{X}_i = 0$  (on the  $X$ -axis initially) we have  $X_a = X_i$ . Thus the maximum time  $T_a$  between axes under QE motion will be

$$T_a = K / (\sqrt{p_0} X_i) \quad (100)$$

We can therefore conclude that the sign of  $X_1$  will be equal or opposite to the sign of  $X_i$  corresponding to  $(T_a - |t_a|) > 3.0$  or  $(T_a - |t_a|) < 3.0$ . The sign of  $\dot{X}_1$  will generally remain the same as the sign of  $\dot{X}_i$ , unless the motion goes beyond the next quadrant, i.e. unless  $(2T_a - |t_a|) < 3.0$ . In the special cases that  $(X_i, \dot{X}_i)$  is on the  $\pm \dot{X}$ -axis we see that  $|t_a| = T_a$  above, so that  $t_1 = (3.0 + T_a)$ . Then for  $\dot{X}_i < 0$  we will generally have  $X_1 < 0$  (unless  $2T_a < 3.0$ , in which case  $X_1 > 0$ ), while the sign of  $\dot{X}_1$  will be equal or opposite to the sign of  $\dot{X}_i$  corresponding to  $T_a > 3.0$  or  $T_a < 3.0$ . For  $\dot{X}_i > 0$

the results for the sign of  $\dot{X}_1$  will still hold, but in general we now will have  $X_1 > 0$  (unless  $2T_a < 3.0$ , in which case  $X_1 < 0$ ).

Next let us consider QH I motion. This motion is described in counter-clockwise quasi-hyperbolic arcs about the origin, concave outward (see equation (93)). Thus if  $(X_1, \dot{X}_1)$  is in the second or fourth quadrants (i.e. if  $X_1 \dot{X}_1 < 0$ ) then the motion will be towards the X-axis, with an intercept  $X_b$ . Thus the results for the final coordinates will be those given in equations (93) to (95), with the following signs. The sign of  $X_2$  will equal the sign of  $X_1$ , while the sign of  $\dot{X}_2$  will be equal or opposite to the sign of  $\dot{X}_1$  corresponding to  $|t_b| > 3.0$  or  $|t_b| < 3.0$ . There is no possible sign change for  $X_2$  (as there is for QE motion) since the QH I motion is outward and asymptotic to the characteristic straight lines of the quasi-hyperbola. On the other hand, if  $(X_1, \dot{X}_1)$  is in the first or third quadrants (i.e. if  $X_1 \dot{X}_1 > 0$ ) or if it is on the  $\pm X$ -axis (the  $\pm \dot{X}$ -axis is excluded by equation (91)) then the motion will be away from the X-axis, with the intercept  $X_b$  now in past time. That is,  $|t_b|$  will be the time it would take to go from  $(X_b, 0)$  on the X-axis to  $(X_1, \dot{X}_1)$ . Therefore the total time to move from the X-axis to the final point  $(X_2, \dot{X}_2)$  will be  $t_2 = (3.0 + |t_b|)$ . Clearly the signs of  $X_2$  and  $\dot{X}_2$  will be equal to those of  $X_1$  and  $\dot{X}_1$  respectively (if  $\dot{X}_1 = 0$ , then  $\dot{X}_2 > 0$  for  $X_1 > 0$ , while  $\dot{X}_2 < 0$  for  $X_1 < 0$ ).

Finally let us consider QH II motion. This motion is described in clockwise arcs about the origin, concave outward (see equation (97)), with its intercept  $(\dot{X}_c)$  now on the  $\dot{X}$ -axis. Thus for  $(X_1, \dot{X}_1)$  in the second or fourth quadrants the motion is towards the  $\dot{X}$ -axis. Analogous to the  $X_1 \dot{X}_1 < 0$  case for the QH I motion, the sign of  $\dot{X}_4$  will equal the sign of  $\dot{X}_1$ , while the sign of  $X_4$  will be equal or opposite to the sign of  $X_1$  corresponding to

$|t_c| > 3.0$  or  $|t_c| < 3.0$ . Again it is impossible for the sign of  $\dot{X}_4$  to differ from the sign of  $\dot{X}_1$  because the QH II motion arcs outward towards infinity. If now  $(X_1, \dot{X}_1)$  is in the first or third quadrants or if it is on the  $\pm\dot{X}$ -axis (the  $\pm X$ -axis is excluded by equation (92)), the motion will be away from the  $\dot{X}$ -axis, and so  $|t_c|$  will be the time to go from  $(0, \dot{X}_c)$  on the  $\dot{X}$ -axis to  $(X_1, \dot{X}_1)$ . Thus the total time to move from the  $\dot{X}$ -axis to the final point  $(X_4, \dot{X}_4)$  will be  $t_4 = (3.0 + |t_c|)$ . Again, the signs of  $X_4$  and  $\dot{X}_4$  will be equal to those of  $X_1$  and  $\dot{X}_1$  respectively; if  $X_1 = 0$ , then  $X_4 > 0$  if  $\dot{X}_1 > 0$  while  $X_4 < 0$  if  $\dot{X}_1 < 0$ .

In conclusion, it can be seen that the calculation of the final coordinates of the motion is best done by computing only the magnitudes of the necessary quantities. Then a logical sign routine based on the above results can be constructed. A simple test of the product  $(X_1 \dot{X}_1)$  plus a determination of the type of motion involved will then simply and easily fix the signs of the final coordinates. (Recall that the motion is QE if  $p(t) > 0$ , and QH I or QH II according to equations (91) and (92) when  $p(t) < 0$ .) A summary of all these sign results together with the piecewise computational steps previously discussed are given in the beginning of the next appendix.

## V. APPENDIX -- PROGRAMS

A. Summary of piecewise calculations and sign routines

Initial Tests: $p(t) = +p_0$ QE Motion
--

$$|x_a| = + (2\dot{x}_i^2/p_0 + x_i^4)^{1/2}$$

$$|t_a| = (|x_a| \sqrt{p_0})^{-1} \operatorname{cn}^{-1} (|x_i|/|x_a|)$$

$$|u_1| = \sqrt{p_0} |x_a| |t_1|$$

$$|x_1| = |x_a| \operatorname{cn}|u_1|$$

$$|\dot{x}_1| = \sqrt{p_0} x_a^2 \operatorname{sn}|u_1| \operatorname{dn}|u_1|$$

$$T_a = K/(|x_i| \sqrt{p_0})$$

Case I:  $(x_i \dot{x}_i) > 0$

Use:  $t_1 = (3.0^* - |t_a|)$  in third equation above;

Set:  $\operatorname{Sign} x_1 = \operatorname{Sign} x_i$  if  $t_1 < T_a$ ,

or  $\operatorname{Sign} x_1 = -\operatorname{Sign} x_i$  if  $t_1 > T_a$ ;

Set:  $\operatorname{Sign} \dot{x}_1 = \operatorname{Sign} \dot{x}_i$  if  $t_1 < 0$  or  $t_1 > 2T_a$ ,

or  $\operatorname{Sign} \dot{x}_1 = -\operatorname{Sign} \dot{x}_i$  if  $0 < t_1 < 2T_a$ .

Case II:  $\dot{X}_i = 0$

Use:  $t_1 = 3.0^*$  in third equation above;

Set:  $\text{Sign } X_1 = \text{Sign } X_i$  if  $T_a > 3.0^*$ ,  
 or  $\text{Sign } X_1 = -\text{Sign } X_i$  if  $T_a < 3.0^*$ .

Subcase 1:  $X_i > 0$

Set:  $\dot{X}_1 < 0$  if  $2T_a > 3.0^*$ ,  
 or  $\dot{X}_1 > 0$  if  $2T_a < 3.0^*$ .

Subcase 2:  $X_i < 0$

Set:  $\dot{X}_1 > 0$  if  $2T_a > 3.0^*$ ,  
 or  $\dot{X}_1 < 0$  if  $2T_a < 3.0^*$ .

Case III:  $X_i = 0$

Use:  $t_1 = (3.0^* + T_a)$  in third equation above;

Set:  $\text{Sign } \dot{X}_1 = \text{Sign } \dot{X}_i$  if  $T_a > 3.0^*$ ,  
 or  $\text{Sign } \dot{X}_1 = -\text{Sign } \dot{X}_i$  if  $T_a < 3.0^*$ .

Subcase 1:  $\dot{X}_i < 0$

Set:  $X_1 < 0$  if  $2T_a > 3.0^*$ ,  
 or  $X_1 > 0$  if  $2T_a < 3.0^*$ .

Subcase 2:  $\dot{X}_i > 0$

Set:  $X_1 > 0$  if  $2T_a > 3.0^*$ ,  
 or  $X_1 < 0$  if  $2T_a < 3.0^*$ .

Case IV:  $(X_i \dot{X}_i) < 0$

Use:  $t_1 = (3.0^* + |t_a|)$  in third equation above ;

Set:  $\text{Sign } X_1 = \text{Sign } X_i$  if  $t_1 < T_a$  ,

or  $\text{Sign } X_1 = -\text{Sign } X_i$  if  $t_1 > T_a$  ;

Set:  $\text{Sign } \dot{X}_1 = \text{Sign } \dot{X}_i$  if  $t_1 < 2T_a$  ,

or  $\text{Sign } \dot{X}_1 = -\text{Sign } \dot{X}_i$  if  $t_1 > 2T_a$  .

Return to Initial Tests

\* (Use 1.5 instead of 3.0 for the first quarter-period of QE motion, plus the last quarter-period of the motion if the total period of the solution is to be an integer multiple of the period  $\tau$ ).

Initial Tests:  $\left\{ \begin{array}{l} p(t) = -p_o \\ |\dot{X}_i| < \sqrt{p_o/2} X_i^2 \end{array} \right\} \text{ QH I Motion}$

$$|X_b| = + (X_i^4 - 2\dot{X}_i^2/p_o)^{1/2}$$

$$|t_b| = (|X_b| \sqrt{p_o})^{-1} \text{cn}^{-1}(|X_b|/|X_i|)$$

$$|u_2| = \sqrt{p_o} |X_b| |t_2|$$

$$|X_2| = |X_b|/\text{cn}|u_2|$$

$$|\dot{X}_2| = \sqrt{p_o} X_b^2 (\text{sn}|u_2| \text{dn}|u_2|)/\text{cn}^2|u_2|$$

Case I:  $(X_i \dot{X}_i) > 0$

Use:  $t_2 = (3.0 + |t_b|)$  in third equation above ;

Set:  $\text{Sign } X_2 = \text{Sign } X_i$

and  $\text{Sign } \dot{X}_2 = \text{Sign } \dot{X}_i$  .

Case II:  $\dot{X}_1 = 0$

Use:  $t_2 = 3.0$  in third equation above ;

Set:  $\text{Sign } X_2 = \text{Sign } X_1$

and  $\dot{X}_2 > 0$  if  $X_1 > 0$  ,

or  $\dot{X}_2 < 0$  if  $X_1 < 0$  .

Case III:  $X_1 = 0$

Excluded (QH II Motion Only) .

Case IV:  $(X_1 \dot{X}_1) < 0$

Use:  $t_2 = (3.0 - |t_b|)$  in third equation above ;

Set:  $\text{Sign } X_2 = \text{Sign } X_1$

and  $\text{Sign } \dot{X}_2 = -\text{Sign } \dot{X}_1$  if  $t_2 > 0$  ,

or  $\text{Sign } \dot{X}_2 = \text{Sign } \dot{X}_1$  if  $t_2 < 0$  .

**Return to Initial Tests**

Initial Tests:  $\left\{ \begin{array}{l} p(t) = -p_o \\ |\dot{X}_1| > \sqrt{p_o/2} X_1^2 \end{array} \right\}$  QH II Motion

$$|c| = + (8\dot{X}_1^2/p_o - 4X_1^4)^{\frac{1}{2}}$$

$$|t_c| = (|c| \sqrt{p_o})^{-1} f^{-1} (\sqrt{2} |X_1|/|c|)$$

where  $f(u) \equiv \text{sn}u/(1 + \text{cn}u)$

$$|u_4| = \sqrt{p_o} |c| |t_4|$$

$$|X_4| = \sqrt{c^2/2} \text{sn}|u_4|/(1 + \text{cn}|u_4|)$$

$$|\dot{X}_4| = \sqrt{p_o/2} c^2 \text{dn}|u_4|/(1 + \text{cn}|u_4|)$$

Case I:  $(X_i \dot{X}_i) > 0$

Use:  $t_4 = (3.0 + |t_c|)$  in fourth equation above;

Set:  $\text{Sign } \dot{X}_4 = \text{Sign } \dot{X}_i$

and  $\text{Sign } X_4 = \text{Sign } X_i$ .

Case II:  $\dot{X}_i = 0$

Excluded (QH I Motion Only).

Case III:  $X_i = 0$

Use:  $t_4 = 3.0$  in fourth equation above;

Set:  $\text{Sign } \dot{X}_4 = \text{Sign } \dot{X}_i$

and  $X_4 > 0$  if  $\dot{X}_i > 0$ ,

or  $X_4 < 0$  if  $\dot{X}_i < 0$ .

Case IV:  $(X_i \dot{X}_i) < 0$

Use:  $t_4 = (3.0 - |t_c|)$  in fourth equation above;

Set:  $\text{Sign } \dot{X}_4 = \text{Sign } \dot{X}_i$

and  $\text{Sign } X_4 = -\text{Sign } X_i$  if  $t_4 > 0$ ,

or  $\text{Sign } X_4 = \text{Sign } X_i$  if  $t_4 < 0$ .

**Return to Initial Tests**

## B. NICAP program for the Illiac II

### Opening routines (only one entry per run)

Calling of library subroutines for input/output, Lagrangian 6-point interpolation, arctangent, square root, cosine and sine functions;

SYSETC, TLOOP, QUAD1, QUAD2: Calculation and storage of Jacobian elliptic function tables ( $snu = SNARRY$ ,  $cnu = CNARRY$ ,  $dnv = DNARRY$ , and  $f(u) = snu/(1 + cnu) = SBYC$ ).

### Main program (EFPR3)

Input/output controls, various counters;

EFPR3: Conversion from input polar to cartesian coordinates.

ELOOP: Preparation for entry into QE Motion Loop.

H1LOOP: Preparation for entry into QH I Motion Loop.

H2LOOP: Preparation for entry into QH II Motion Loop.

OP: Output options with final coordinates, e.g. conversion back to polar coordinates, calculation of special functions such as  $\Delta\theta$ ,  $\Delta r$ , and the ratio  $r_1/r_e = RUCRIN$ .

### Subroutines

ELLIP: Calculations for QE Motion.

HYP1: Calculations for QH I Motion.

HYP2: Calculations for QH II Motion.

NVERS1:  $cnu = CNARRY$  inversion (both methods).

NVERS2:  $f(u) = SBYC$  inversion.

LOOKUP, LKAUX, AUXLK: Preparation of data for using the library subroutine for the six-point Lagrangian interpolation of the tables (i.e. for the direct lookup of off-table values).

SIGNEA: Final coordinate sign routine for QE Motion.

SIGN1A: Final coordinate sign routine for QH I Motion.

SIGN2A: Final coordinate sign routine for QH II Motion.

STXCAL: Special subroutine to calculate the statistical properties of certain functions (e.g. RUCRIN) for large numbers of solutions.

RDUM: Counter routine to allow recycling of calculations for multiples  $m$  of the basic mapping period  $n$ .

#### Input parameters

XPERM: The  $X$ -coordinate of the point to be used as the origin of the polar coordinates.

XDPERM: The  $\dot{X}$ -coordinate of the point to be used as the origin of the polar coordinates.

AMJOUT, AMJIN: The semi-major axes of two ellipses which form the outer and inner boundaries respectively of the (elliptical) annular region under study.

ECOUT, ECIN: The eccentricities of these two (outer and inner, respectively) bounding ellipses.

AMJOR: The semi-major axis  $a_0$  corresponding to the initial point in phase-space being mapped.

DELRAD: An increment in  $a_0$  if a regular series of  $a_0$  values is desired.

RADFIN: The final  $a_0$  in the desired series of  $a_0$  values (may equal AMJOR).

PHINP: The original angle  $\theta_0$  corresponding to the initial point in phase-space being mapped.

DELPHI: An increment in  $\theta_0$  if a regular series of  $\theta_0$  values is desired.

ECCEN: The fixed eccentricity  $e_0$  corresponding to the initial point in phase-space being mapped.

- XMAX:** The maximum value of  $X$  which a final coordinate is allowed to attain.
- N:** The integer number of periods of  $\tau$  of the basic mapping (printout every  $N$  periods).
- M:** The integer number of multiples of the basic period  $N$  ( $M$  lines of printout of  $N$  periods apiece).
- P:** The number of incremented values for  $\theta_0$  (the number of times DELPHI is added for a given  $a_0$ ; DELRAD is controlled by RADFIN).

#### Output parameters

- M1:** The running counter for the number of multiples  $m$  of the basic period  $n$ .
- RUCRIN:** The ratio of the mapped radius  $r_1$  at the angle  $\theta_1$  to the radius  $r_e$  of a standard ellipse (specified by AMJOR and ECCEN) at the same angle  $\theta_1$ .
- LARGE:** The ratio of the mapped radius  $r_1$  to the original radius  $r_0$ .
- F:** Moser's  $\Delta\theta$  function =  $(\theta_1 - \theta_0)$ .
- SMALL:** Moser's  $\Delta r$  function =  $(r_1 - r_0)$ .
- RAD:** The mapped radius  $r_1$  about (XPERM, XDPERM).
- PHI:** The mapped angle  $\theta_1$  about (XPERM, XDPERM).
- XF:** The mapped  $X$ -coordinate  $X_1$ .
- XDF:** The mapped  $\dot{X}$ -coordinate  $\dot{X}_1$ .
- Q:** The running counter for the number of increments  $P$  in the initial angle  $\theta_0$ .

```

PRINT
READ
LAG6
ATAN1
SQR1
COS1
SIN1
SYSETC
*
*      M0      PAGE COUNT.
*      M2      'LOOPS'.
*      M4      FIXED POINT NUMBER, 'N'.
*      M5      ORBITS, 'M'.
*      M6      'P'.
*      M7      'Q'.
*      M8      USED ONCE ONLY, TO START OFF-AXIS IN ELLIP.
*      M9      SAME AS M5 (M), BUT USED IN CALCULATIONS.
*      M10     SAME AS M6 (P), BUT USED IN CALCULATIONS.
*      M11     OUTPUT COUNT.
*
      CALL      SYSETC
      DECQL     ,,IGNORE,
      CAD       K
      DIV       900.
      STR       SCALE
      CAD       15,3,
      STU       SNARRY          TABLE OF SN, CN, AND DN.

      ADD       1.
      STR       CNARRY          SET INITIAL TABLE VALUES

      STR       DNARRY
      CAD       PSZ

      CALL      SQR1
      STR       RPSZ
*
*
*      RED      RED      RED      RED      RED      RED      RED      RED
      CSM       0,900          LOOP COUNTER.
      LFR       2,Q           F2 = Q = QFAST
      CAD       QFAST          1
      MPY       F1
      STU       F3
      SFR       3,Q2
      MPY       F3
      STU       Q4             Q**4 = Q4
      MPY       F1
      STU       Q6             Q**6 = Q6

```

	MPY	F1	
	MPY	QFAST	
	STU	Q9	$Q \times 9 = Q9$
	MPY	F1	
	MPY	F3	
	STU	Q12	$Q \times 12 = Q12$
	ADD	Q6	
	ADD	Q2	
	ADD	1.	$Q12 + Q6 + Q2 + 1$
	STR	CQ	GOES TO CQ
	CAD	Q9	
	ADD	Q4	
	ADD	QFAST	
	ADD	10,3,2048	$Q9 + Q4 + Q + 1/2$
	STR	AQ	GOES TO AQ
TLOOP	LFR	3,RINCR	
	CAT	RFAST	
	CALL	COS1	
	MPY	Q	$\times Q$
	STU	F5	TO F5
	CAD	RFAST	
	MPY	6.	
	CALL	COS1	
	MPY	Q9	$\times Q9$
	STU	F6	TO F6
	CAD	RFAST	
	MPY	4.	
	CALL	COS1	
	MPY	Q4	$\times Q4$
	ADD	10,3,2048	$1/2$
	SUB	F5	$-Q \cos(RINCR) - Q9 \times \cos(6 \times RINCR)$
	SUB	F6	
	MPY	CQ	
	STR	BNQCQ	TO BNQCQ
	CAD	RFAST	
	MPY	5.	
	CALL	SIN1	
	MPY	Q6	$\times Q6$
	STR	F5	TO F5
	CAD	RFAST	
	MPY	3.	

	CALL	SIN1	
	MPY	Q2	* Q2
	STU	F6	TO F6
	CAD	RFAST	
	MPY	7.	
	CALL	SIN1	
	MPY	Q12	*Q12
	STU	F7	TO F7
	CAD	RFAST	
	CALL	SIN1	
	ADD	F5	$Q6\sin(5 \cdot \text{RINCR})$
	SUB	F6	$-Q2\sin(3 \cdot \text{RINCR})$
	SUB	F7	$-Q12\sin(7 \cdot \text{RINCR})$
	MPY	AQ	$\cdot (Q9+Q4+Q+1/2)$
	DIV	BNQCQ	BY $Q4\cos(4 \cdot \text{RINCR})+1/2-$
	STR	SNARRY+M0+901	ELLIPTIC SINE.
	MPY	F0	$1 - (\text{SNU} \cdot \sin(\text{RINCR}))$
	STR	F5	$-Q6\sin(5 \cdot \text{RINCR})$ ROOTED
	CAD	1.	
	SUB	F5	
	CALL	SQR1	
	STR	CNARRY+M0+901	ELLIPTIC COSINE.
	CAD	2.	$2 - Q6\sin(5 \cdot \text{RINCR}) \cdot 1/2$
	SUB	F5	ROOTED
	DIV	2.	
	CALL	SQR1	
	STR	DNARRY+M0+901	ELLIPTIC TANGENT.
	CAD	RINCR	
	ASC	RFAST	INCREMENT R...
	CJU	M0,TLOOP	AND LOOP.
	CSM	0,901	
	CAM	1,SNARRY	
	CAM	2,CNARRY	
QUAD1	CAD	2,1,	
	ADD	1.	
	VID	1,1,	
	STR	SBYC+M0+901	
	CJU	0,QUAD1	
	CSM	0,801	
	CAM	1,SNARRY+899	

	CAM	2,CNARRY+899	
QUAD2	CSB	2,,	
	ADD	1.	
	VID	1,,	
	STR	SBYC+M0+1702	
	SBM	1,1	
	SBM	2,1	
	CJU	0,QUAD2	
	FIL		
* MAIN PROGRAM---(E)LLIPTIC (F)IXED (P)OINT, (R)ED---			
*			
* LOOPS = MODIFIER 2			
*			
EFPR3	CAM	1,READIN	
	CALL	READ	
	CAM	1,CTOFF	
	CALL	READ	
	CAM	1,PRINTN	
	CALL	PRINT	
	CAM	1,CTOFFP	
	CALL	PRINT	
	CSM	0,41	PAGE COUNTER.
	LFR	5,SETUP	
	CRN	4,12	
	CSM	4	FIXED POINT NUMBER.
	CSM	5,M5	TOTAL NUMBER OF ORBITS.
	CAM	7,1	SET INITIAL VALUE OF Q.
	SFR	5,SETUP	SAVE INPUT COUNTERS.
	CSM	9,M5	
	CAM	10,M6	
	CAM	11	SET INITIAL VALUE OF M(6)
	SFR	6,COUNT	3,COUNT = 0.
	CAD	15,3,	
	STU	AVRUC	
	STU	RMRUC	
	STU	AVG	
	STU	RMG	
EFPRB	CAM	1,HEAD	
	CALL	PRINT	

	CAD	PHINP	
	STR	PHIIN	
EFPRA	CALL	RDUM	
	DECQL	AMJOR,ECCEN,PHIIN,RADIN	
	CAD	PHIIN	
	CIV	180.	
	STR	F3	
	CALL	COS1	
	MPY	RADIN	
	ADD	XPERM	
	STR	XINP	
	STR	XIN	
	STR	XINS	
	STR	XF	
	CAD	F3	
	CALL	SIN1	
	MPY	RADIN	
	ADD	XDPERM	
	STR	XDINP	
	STR	XDIN	
	STR	XDINS	
	STR	XDF	
	CAD	15,3,	
	STU	F	
	STU	SMALLG	
	STU	LARGE G	
	ADD	1.	
	STR	RUCRIN	
	CAD	RADIN	
	STR	RAD	
	CAD	PHIIN	
	STR	PHI	
	CAN	1,OUT	
	CALL	PRINT	
	CSM	LOOPS,1	
	CSM	8,1	FIRST ELLIPTIC ENTRY.
ELOOP	CALL	ELLIPA	
	CALL	OP	
	CJU	LOOPS,ELOOP	
	CSM	LOOPS,2	
	CAN	8	
	CAD	X	
	CAV	XMAX	

	TP	CXMAX	PRINT X, XD, THEN READ IN NEW
	CAD	X	
	MPY	F1	
	STR	F3	
	CAD	PSZ	
	DIV	2.	
	CALL	SQR1	
	MPY	F3	
	CAV	XD	
	TN	H2LOOP	
H1LOOP	CALL	HYPIA	
	CALL	OP	
*			
	CJU	LOOPS,H1LOOP	
	CSM	LOOPS,2	
	TRA	ELOOP	
H2LOOP	CALL	HYP2A	
	CALL	OP	
*			
	CJU	LCOPS,H2LOOP	
	CSM	LOOPS,2	
	TRA	ELOOP	
OP	SFR	4,OPBSS	SAVE RETURN.
	CJZ	LOOPS,OPC	JUMP IF ON SECOND LOOP TO OUTP
	CAD	XZERO	
	STR	XIN	
	CAD	XCZERO	
	STR	XDIN	
	CAD	TBY4	
	ASC	T	
	CJU	4,OPX	JUMP IF NO CUTPUT IS DESIRED.
	LDM	4,SETUP	RESET OUTPUT COUNTER.
	CAD	X	
	SUB	XPERM	
	STR	DELX	
	CAD	XD	
	SUB	XDPERM	
	STR	DELXD	
	CAD	RAC	
	STR	RADOLD	
	CAD	PHI	
	STR	PHIOLD	

CAD	DELX	
MPY	F1	
STR	F0	
CAD	DELXD	
MPY	F1	
ADD	F0	
CALL	SQR1	
STR	RAD	
CAD	DELXD	
DIV	DELX	
CALL	ATAN1	
MPY	180.	
DIV	PI	
STR	PHI	
CALL	PHIST	
LFR	6,COUNT	
ADM	11,1	INCREMENT RUNNING COUNTER (M11)
SFR	6,COUNT	
CALL	RDUM	
DECQL	AMJOR,ECCEN,PHI,RADNEW	
CAD	RAD	
CIV	RADNEW	
STR	RUCRIN	
ASC	AVRUC	
CAD	RUCRIN	
MPY	F1	
ASC	RMRUC	
CALL	RDUM	
DECQL	AMJOUT,ECOUT,PHI,ROUT	
CALL	RDUM	
DECQL	AMJIN,ECIN,PHI,RIN	
CAD	RAD	
SUB	ROUT	
TU	OPG	
CAD	15,3,	
STU	LARGE G	
OPG	TRA	OPF
	CAD	RAD
SUB	RCUT	
DAV	15,3,	
STR	F2	
CAD	ROUT	

	SUB	RIN
	CAV	15.3.
	STR	F3
	MPY	F2
	MPY	2.
	STR	F7
	CAD	RAD
	SUB	RIN
	CAV	15.3.
	ADD	F2
	SUB	F3
	CIV	F7
	STR	F0
	CAD	RAD
	SUB	ROUT
	MPY	F0
	STR	LARGE
OPF	CAD	PHI
	SUB	PHIOLD
	TP	OPD
	ADD	360.
OPD	STR	F
	CAD	RAD
	SUB	RADOLD
	STR	SMALL
	CAD	X
	STR	XF
	CAD	XC
	STR	XDF
	CAD	LARGE
	ASC	AVG
	CAD	LARGE
	MPY	F1
	ASC	RMG
	JNM	0,OPE
	CAM	1,HEAD
	CALL	PRINT

HEAD TOP OF PAGE.

OPE	CSM	0,44	OUTPUT PAGE COUNTER.
	CAM	1,OUT	
	CALL	PRINT	
	ADM	0,2	UPDATE PAGE COUNTER.
	CAD	X	
	CAV	XMAX	
	TP	CXMAX	PRINT X, XD, THEN READ IN NEW C
	CJZ	5,OPA	
OPX	LDM	3,OPBSS	RESTORE RETURN.
	LDM	LOOPS,OPBSS	RESTORE LOOP COUNTER.
OPA	JLH	3,,	EXIT.
	LFR	6,COUNT	
	CAC	AVRUC	
	DIV	M9.	
	STR	AVGR	
	MPY	F0	
	STR	F3	
	CAD	RMRUC	
	DIV	M9.	
	CALL	SQR1	
	STR	RMSR	
	MPY	F0	
	SUB	F3	
	CALL	SQR1	
	STR	SIGMAR	
	CAD	AVG	
	DIV	M9.	
	STR	AVGG	
	MPY	F0	
	STR	F3	
	CAD	RMG	
	DIV	M9.	
	CALL	SQR1	
	STR	RMSG	
	MPY	F0	
	SUB	F3	
	CALL	SQR1	
	STR	SIGMAG	
	CAM	1,SMOUT1	
	CALL	PRINT	
	ADM	0,8	INCREMENT PAGE COUNTER.
	CAD	15,3,	
	STU	AVRUC	

	STU	RMRUC	
	STU	AVG	
<hr/>			
	STU	RMG	
	CAM	8,M7-M6	
<hr/>			
	JNM	8,OPB	
	TRA	STXCAL	
OPB	CAD	DEL RAD	
	MPY	DELPHI	
	TU	OPH	
	CAD	DEL RAD	
	ASC	AMJOR	
OPH	CAD	DELPHI	
	ASC	PHI IN	
	LFR	5,SETUP	
	ADM	7,1	
	SFR	5,SETUP	
	LFR	6,COUNT	
	CAM	11	
	SFR	6,COUNT	
	CAM	1,SKIP	
<hr/>			
	CALL	PRINT	
	ADM	0,3	INCREMENT PAGE COUNTER.
	TRA	EFRA	
OPC	CAD	X	
	STR	XIN	
	STR	XINS	
	CAD	XD	
	STR	XDIN	
	STR	XCINS	
	TRA	OPX	
OPBSS	BSS	1	
<hr/>			
•			
	ELLIPA	SFR 4,ELLBSS	SAVE RETURN.
		SFR 5,ELLBSS+1	
		CJZ 8,ELLIPG	
		CJZ LOOP6,ELLIPC	
ELLIPG	CAD	XIN	
	MPY	XDIN	
	STR	F0	
	CAV	15,3,	
	VID	F0	
	TOR	ELLIPE	

TRA	ELLIPF	
ELLIPE CSB	1.	
ELLIPF STR	MCTION	
CAD	XIN	
MPY	F1	
MPY	F1	
MPY	F1	
STR	F7	
ELLIPB CAD	XDIN	
MPY	F1	
MPY	2.	
DIV	PSZ	
STR	DELONE	
DIV	F7	
STR	DELTWO	
SUB	RERROR	
TP	ELLIPQ	
CAD	DELTWO	APPROXS. FOR WHEN XZERO IS NEAR
DIV	4.	
ADD	1.	
MPY	XIN	
STR	XZERO	
CALL	SIGNEA	
CAD	DELTWO	
NDV	4.	
ADD	1.	
STR	TSFCNU	
CAD	DELTWO	
DIV	2.	
CALL	SQR1	
STR	U	
TRA	ELLIPR	
ELLIPQ CAD	DELONE	
ADD	F7	
CALL	SQR1	
CALL	SQR1	
STR	XZERO	X AT ABCISSA
CALL	SIGNEA	

	CAD	XZERO	
	VID	XIN	
	STR	TSTCNU	
*			
	CALL	NVERS1	
*			
•			
ELLIPR	CAD	15,3,	
	STU	XDZERO	
	CAD	RPSZ	
	MPY	XZERO	
	VID	U	
	CAV	15,3,	
	STR	TZERO	TIME TO ABCISSA
*			
	MPY	MCTION	
	SUB	TBY4	
	STN	T	NEW TIME
*			
*			
ELLIPC	CAD	T	
	MPY	XZERO	
	STR	F7	
	CAD	RPSZ	
	MPY	F7	
	STR	TESTU	
*			
	CALL	LCOKUP	
	CAD	CNU	
•			
	MPY	XZERO	
	STR	X	NEW X
*			
	CALL	SIGNEB	
*			
*			
*			
*			
	CAD	RPSZ	
	MPY	XZERO	
	MPY	F1	
	MPY	SNU	
	MPY	DNU	
	STN	XD	NEW XD

	CALL	SIGNEC	
	CAD	EQCT	
	SAM	MCTION	
	JZM	8, ELLIPH	
	JUM	LCOPS, ELLIPD	
ELLIPH	CAD	XINS	
	MPY	XDINS	
	TP	ELLIPD	
	CAD	PSI	
*			
	STR	T	
*			
*			
ELLIPD	LFR	4, ELLBSS	DELETE ELLIPD LABEL FOR PH (PRIN
	LFR	5, ELLBSS+1	
	JLH	3,,	EXIT
ELLBSS	BSS	2	
*			
		ASSIGN DELONE, DELTWO	
RERROR	DEC	1.0E-07	
HYP1A	SFR	4, HY1BSS	SAVE RETURN.
	SFR	5, HY1BSS+1	
	CJZ	LOOPS, HYP1C	
	CAD	XIN	
	MPY	XDIN	
	STR	F0	
	CAV	15, 3,	
	VID	F0	
	STR	MOTION	
	CAD	XIN	
	MPY	F1	
	MPY	F1	
	MPY	F1	
	STR	F7	
HYP1B	CAD	XCIN	
	MPY	F1	
	MPY	2.	
	NDV	PSZ	
	STR	DELONE	
	NDV	F7	
	STR	DELTWO	
	SUB	RERROR	
	TP	HYP1Q	

	CAD	DELTWO	APPROXS. FOR WHEN XZERO IS NEAR
	NDV	4.	
	ADD	1.	
	STR	TSTCNU	
	MPY	XIN	
	STR	XZERO	
	CALL	SIGN1A	
	CAD	DELTWO	
	DIV	2.	
	CALL	SQR1	
	STR	U	
	TRA	HYP1R	
HYP1Q	CAD	DELONE	
	ADD	F7	
	CAV	15,3,	
	CALL	SQR1	
	CALL	SQR1	
	STR	XZERO	X AT ABCISSA
*			
	CALL	SIGN1A	
	CAD	XZERO	
	DIV	XIN	
	STR	TSTCNU	
*			
	CALL	NVERS1	
*			
HYP1R	CAD	15,3,	
	STU	XDZERO	
	CAD	RPSZ	
	MPY	XZERO	
	VID	U	
	CAV	15,3,	
	STR	TZERO	TIME TO ABCISSA
*			
	MPY	MCTION	
	ADD	TBY4	
	STR	T	
HYP1C	CAD	T	
*			
	MPY	XZERO	
	STR	F7	
	CAD	RPSZ	

	MPY	F7	
	STR	TESTU	
*			
	CALL	LOOKUP	
*			
•			
	CAD	XZERO	
	DIV	CNU	
	STR	X	NEW X
*			
	CALL	SIGN1B	
*			
•			
*			
	CAD	XZERO	
	MPY	F1	
	MPY	SNU	
	MPY	DNU	
	STR	F7	
	CAD	RPSZ	
	MPY	F7	
	DIV	CNU	
	DIV	F1	
	STR	XD	NEW XD
*			
	CALL	SIGN1C	
	CAD	H10CT	
	SAM	MOTION	
*			
•			
	LFR	4, HY1BSS	RESTORE RETURN.
	LFR	5, HY1BSS+1	
	JLH	3,,	EXIT
	HY1BSS BSS	2	
*			
HYP2A	SFR	4, HY2BSS	SAVE RETURN.
	SFR	5, HY2BSS+1	
	CJZ	LOOPS, HYP2C	
	CAD	XIN	
	MPY	XDIN	
	STR	F0	
	CAV	15, 3,	
	VID	F0	
	STR	MOTION	
	CAD	XIN	

MPY	F1	
MPY	F1	
MPY	F1	
MPY	4.	
STN	F7	
CAD	XDIN	
MPY	F1	
MPY	8.	
CIV	PSZ	
ADD	F7	
CAV	15,3,	
CALL	SQR1	
CALL	SQR1	
STR	F6	F6 = A
CAD	F7	
MPY	PSZ	
CIV	8.	
STR	F7	
HYP2B	CAD	XDIN
MPY	F1	
ADD	F7	
CAV	15,3,	
CALL	SQR1	
STR	XDZERO	XD AT ORDINATE
*		
CALL	SIGN2A	
CAD	2.	
CALL	SQR1	
MPY	XIN	
CIV	F6	
STR	TSTSBC	
CAV	15,3,	
CALL	NVERS2	
*		
*		
CAD	15,3,	
STU	XZERO	
CAD	RPSZ	
VID	U	
CIV	F6	
STR	TZERO	TIME TO ORDINATE
•		
MPY	MOTION	
ADD	TBY4	
STR	T	

*	HYP2C	CAD	RPSZ	
		MPY	T	
		MPY	F6	
		STR	TESTU	
•		CALL	LCOKUP	
		CAD	CNU	
*		ADD	1.	
		STR	F7	
		CAD	2.	
		CALL	SQR1	
		MPY	F7	
		STR	F7	
		CAD	SNU	
		MPY	F6	
		CIV	F7	
		STR	X	NEW X
*		CALL	SIGN2B	
*				
*				
•				
*				
		CAD	PSZ	
		CIV	2.	
		CALL	SQR1	
		MPY	DNU	
		MPY	F6	
		MPY	F6	
		STR	F7	
		CAD	CNU	
		ADD	1.	
		VID	F7	
		STR	XC	NEW XD
*		CALL	SIGN2C	
		CAD	H2OCT	
		SAM	MCTION	
		LFR	4,HY2BSS	RESTORE RETURN.
		LER	5,HY2BSS+1	
		JLH	3,,	EXIT
	HY2BSS	BSS	2	

NVERS1	SFR	4,NV1BSS	SAVE RETURN.
	STR	F2	F2 = TEST CNU
	CAM	0	
	CAM	1,CNARRY.	M1 = CNARRY
	LDM	4,COVAL	
	CSM	4,M4	
N1A	CAD	F2	
	CAV	1,1,	
	TP	N1B	
	ADM	0,1	
	TRA	N1A	
N1B	SBM	0,183	ACCURACY CROSSOVER OCCURS AT ENT
	JPM	0,N1C	
	CAD	1.	BEGIN APPROX. FOR SMALL TABLE VA
	SUB	F2	
	STR	F3	
	CAD	1.	
	ADD	F2	
	VID	F3	
	STR	F3	
	MPY	F3	
	MPY	F3	
	NDV	5.	
	ADD	F3	
	ADE	1	
	CALL	SQR1	
	TRA	N1D	
N1C	ADM	0,183	2ND ORDER ITERATION FOR REST OF
	SBM	0,900	
	JUM	0,N1E	
	ADM	0,900	
	CAD	F2	
	STR	F3	
	CAD	2.	
	CALL	SQR1	
	MPY	F3	
	SUB	K	
	STN	F0	
	CAD	F0	
	TRA	N1D	
N1E	ADM	0,900	
	CAD	M0-1.	
	CIV	900.	
	MPY	K	
	STR	U0	
	CAD	CNARRY+M0-1	
	MPY	F1	
	MPY	F1	

	STR	F3	
	CAC	SNARRY+M0-1	
	MPY	DNARRY+M0-1	
	STR	F2	F2 = SN(I) * DN(I)
	CAD	TSTCNU	
	SUB	CNARRY+M0-1	
	NDV	F2	
	STR	F0	
	MPY	F0	
	MPY	F3	
	NDV	F2	
	MPY	1C,3,2048	
	ADD	F0	
	ADD	U0	
N1D	STR	U	
	STR	TESTU	
	CALL	LOOKUP	
	CAD	CNU	
	SUB	TSTCNU	
	STR	DELCNU	
	CAV	15,3,	
	SUB	ERROR	
	TZN	N1X	
	CJZ	4,EFPR3	
	CAD	DELCNU	
	DIV	SNU	
	DIV	DNU	
	ADD	U	
	TRA	N1D	
N1X	LFR	4,NV1BSS	
	JLH	3,,	EXIT
NV1BSS	BSS	1	
*			
ASSIGN DELCNU,DELSBC			
ERROR	DEC	1.0E-12	
NVERS2	SFR	4,NV2BSS	SAVE RETURN.
	SFR	5,NV2BSS+1	SAVE F5
	SFR	6,NV2BSS+2	SAVE F6
	SFR	7,NV2BSS+3	SAVE F7
	STR	F5	F5 = TEST SN/1+CN
	CAV	SBYC+901	
	TP	N2B	IF SBYC IS IN SECOND QUADRANT.

	CAM	0	
	CAM	1,SBYC	SBYC IS PRIME VALUED.
N2C	CAD	1,1,	
	CAV	F5	
	TP	N2G	
	ADM	0,1	
	TRA	N2C	
N2G	CAM	2,M0	
	SBM	0,183	
	JPM	0,N2E	
	CAD	F5	BEGIN APPROX. FOR SMALL TABLE VA
	STR	F2	
	MPY	F2	
	MPY	F2	
	MPY	F2	
	MPY	F2	
	NDV	10.	
	ADD	F2	
	MPY	2.	
	TRA	N2X	
N2E	ADM	0,183	2ND ORDER ITERATION FOR REST OF
	ANM	0,8190	
	CRM	0,1	
	EOM	0,450	
	JUM	0,N2F	
	CAM	0,M2	
	CAD	F5	
	MPY	2.	
	SUB	1.	
	CALL	SQR1	
	SUB	1.	
	STR	F3	
	CAD	2.	
	CALL	SQR1	
	MPY	F3	
	ADD	K	
	TRA	N2X	
N2F	CAM	0,M2	
	CAD	M1-2	
	STR	F2	
	CAV	F5	

	STN	F3
	CAD	MC-1.
	CIV	9C0.
	MPY	K
	STR	UC
	CAD	1.
	ADD	CNARRY+M0-1
	STR	F0
	CAC	1.
	SUB	CNARRY+M0-1
	CIV	F0
	STR	F6
	CAD	CNARRY+M0-1
	ADD	1.
	VID	DNARRY+M0-1
N2B	TRA	N2A
	CAM	1,SBYC+901
	CAM	0
N2D	ADM	0,1
	CAC	1,1,
	CAV	F5
	TN	N2D
	CAC	M1-2
	STR	F2
	CAV	F5
	STN	F3
	CAD	M0-1.
	CIV	9C0.
	MPY	K
	ADD	K
	STR	UC
	CAD	1.
	SUB	CNARRY-M0+901
	STR	F0
	CAD	1.
	ADD	CNARRY-M0+901
	CIV	F0
	STR	F6
	CSB	CNARRY-M0+901

	ADD	1.	
	VID	DNARRY-M0+901	
N2A	STR	F7	
	VID	F3	
	STR	F3	
	MPY	F3	
	MPY	F6	
	MPY	F2	
	CIV	F7	
	NDV	4.	
	ADD	F3	
	ADD	U0	
N2X	STR	U	
	STR	TESTU	
	CALL	LCOKUP	
	CAD	CNU	
	ADD	1.	
	VID	SNU	
	CAV	15,3,	
	SUB	F5	
	STN	DELSBC	
	CAV	15,3,	
	SUB	ERROR	
	TZN	N2XT	
	CAD	CNU	
	ADD	1.	
	MPY	DELSBC	
	CIV	DNU	
	ADD	U	
	TRA	N2X	
N2XT	LFR	4,NV2BSS	
	LFR	5,NV2BSS+1	RESTORE F5
	LFR	6,NV2BSS+2	RESTORE F6
	LFR	7,NV2BSS+3	RESTORE F7
	JLH	3,,	EXIT
*NV2BSS	BSS	4	
LOOKUP	SFR	4,LKPBSS	SAVE RETURN.
	CAM	0	MO = 0
	I7P	LKP6	TRANSFER IF TESTU .GE. 0
LKP1	ADD	K4	TESTU .L. 0
	TN	LKP1	

LKP2	SUB TN	K LKP3	
	ADM TRA	0,1 LKP2	
LKP3	ADD	K	
LKP4	CRM	0,1	
	JNM CIV	0,LKP5 SCALE	JUMP IF MO IS ODD.
	STR TRA	F2 LKAUX	MO IS EVEN. TRANSFER TO LKAUX.
LKP5	SUB CIV	K SCALE	MO IS ODD.
	STN TRA	F2 LKAUX	TRANSFER TO LKAUX.
LKP6	SUB	K4	TESTU .GE. 0
	TP ADD	LKP6 K	
LKP7	ADM TZN	0,1 LKP7	
	EOM SBM	0,7 0,3	MASK MO FOR U .G. 0 ... AND SET MO.
	TRA LKP BSS	LKP4 1	
*			
LKAUX	SFR SIA SBM	5,LKXBSS 4 4,2	SAVE F5. M4 = TABLE VALUE. OK IF M4 .GE. 2
	JNM SBM	4,LKX4 4,897	JUMP IF M4 .L. 2 OK IF M4 .L. 898
	JNM CAM	4,LKX3 4,SNARRY+894	JUMP IF M4 IS IN RANGE. M4 .GE. 898, END OF TABLE.
	CAM CAM	5,CNARRY+894 6,DNARRY+894	
	CSM CAD	7,6 4,1,	SET END OF TABLES VALUES.
LKX1			

	STR	SNARRY+900-M7	
	CAD	5,1,	
	STN	CNARRY+900-M7	
	CAD	6,1,	
	STR	DNARRY+900-M7	
	CJU	7,LKX1	
	CAD	F2	
	SUB	894.	
LKX2	CAM	1,ELKAUX	
	CALL	LAC6	
	TRA	ALXLK	
LKX3	CAD	F2	
	CAM	1,CWLAG6	
LKX4	TRA	LKX2	TO LAG6.
	CAM	4,SNARRY+1	M4 .L. 2, BEGINNING OF TABLE.
	CAM	5,CNARRY+1	
	CAM	6,DNARRY+1	
LKX5	CSM	7,6	
	CAD	4,1,	SET BEGINNING OF TABLE VALUES.
	STN	SNARRY-7-M7	
	CAD	5,1,	
	STR	CNARRY-7-M7	
	CAD	6,1,	
	STR	DNARRY-7-M7	
	CJU	7,LKX5	
	CAD	F2	
	ADD	6.	
	CAM	1,BLKAUX	
	TRA	LKX2	TO LAG6.
*LKXBSS	BSS	1	
AUXLK	JPM	0,AXLK3	JUMP IF M0 IS EVEN.
	CRM	0,1	M0 IS ODD.
	JNM	0,AXLK2	JUMP IF M0 = 3 = 7 = ...
	CAD	CNU	M0 = 1 = 5 = ...
	STN	CNU	
AXLK1	LFR	4,LKPBSS	RESTORE RETURN.
	LFR	5,LKXBSS	RESTORE F5
	JLH	3,,	EXIT.
AXLK2	CAD	SNU	M0 = 3 = 7 = ...
	STN	SNU	
	TRA	AXLK1	TRANSFER TO EXIT.

AXLK3	CRM	0,1	MO IS EVEN.
	JPM	0,AXLK1	JUMP TO EXIT IF MO = 0 = 4 =.
	CAD	SNU	MO = 2 = 6 = ...
	STN	SNU	
	CAD	CNU	
	STN	CNU	
SIGNEA	TRA	AXLK1	TRANSFER TO EXIT.
	MPY	XINS	
	TP	SIGNEA+2	
	CAD	XZERO	
	STN	XZERO	
	JLH	3,,	
SIGNEB	CAD	XINS	
	MPY	XDINS	
	TN	SEBA	
	CAD	X	
	MPY	XINS	
	TP	SEBBX	
	CAD	X	
	STN	X	
SEBBX	JLH	3,,	EXIT.
SEBA	CAD	K	
	DIV	XZERO	
	DIV	RPSZ	
	DAV	15,3,	
	SUB	TZERO	
	STR	PSI	
	JZM	LCOPS,SEBC	
	DAV	TBY4	FIRST LOOP.
SEBF	TP	SEBD	
	CAD	XINS	
	MPY	X	
	TN	SEBBX	
SEBE	CAD	X	
	STN	X	
	TRA	SEBBX	
SEBD	CAD	XINS	
	MPY	X	
	TP	SEBBX	
	TRA	SEBE	
SEBC	DAV	3.	SECOND LOOP.
	TRA	SEBF	

SIGNEC	CAD	XINS	
	MPY	XDINS	
	TN	AUXE	
	CAD	TZERO	
	JZM	LCOPS,SECA	
	CAV	TBY4	
SECE	TN	SECB	
AUXE	CAD	XD	
	MPY	XDINS	
	TP	SECCX	
SECD	CAD	XD	
	STN	XD	
SECCX	JLH	3,,	EXIT.
SECB	CAD	XD	
	MPY	XDINS	
	TZN	SECCX	
	TRA	SECD	
SECA	CAV	3.	
	TRA	SECE	
SIGN1A	EQU	SIGNEA	
SIGN1B	MPY	XINS	
	TP	SIGN1B+2	
	CAD	X	
	STN	X	
	JLH	3,,	EXIT.
SIGNIC	CAD	XINS	
	MPY	XDINS	
	TP	AUX1	
	CAD	TZERO	
	JZM	LOOPS,S1A	
	CAV	TBY4	
S1E	TN	S1B	
AUX1	CAD	XD	
	MPY	XDINS	
	TP	S1CX	
S1D	CAD	XD	
	STN	XD	
S1CX	JLH	3,,	
S1B	CAD	XD	
	MPY	XDINS	
	TZN	S1CX	
	TRA	S1D	

S1A	CAV TRA	3. S1E	
SIGN2A	MPY TP	XDINS SIGN2A+2	
	CAD STN	XDZERO XDZERO	
SIGN2C	JLH MPY TP	3.,, XDINS SIGN2C+2	
	CAD STN	XD XD	
SIGN2B	JLH CAD	3.,, XINS	
	MPY TP	XDINS AUX2	
	CAD JZM	TZERO LCOPS,S2A	
S2E	CAV TN	TBY4 S2B	
AUX2	CAD MPY	X XINS	
	TP	S2CX	
S2D	CAD	X	
S2CX	STN	X	
S2B	JLH	3.,,	EXIT.
	CAD	X	
	MPY	XINS	
	TZN	S2CX	
	TRA	S2D	
S2A	CAV TRA	3. S2E	
STXCAL	CAD	DELRAD	
	MPY	DELPHI	
	TZ	EFPR3	
	CAD	AMJOR	
	SUB	RADFIN	
	TZP	EFPR3	
	CAD	DELRAD	
	ASC	AMJOR	

	LFR	5, SETUP	
	CAM	7, 1	
	SFR	5, SETUP	
	LFR	6, COUNT	
	CAM	11	
	SFR	6, COUNT	
	CAC	15, 3,	
	TRA	EFPRB	
*	CALL	RDUM	
*	DECQL	ADUM, ECDUM, PHIDUM, RDUM	
RDUM	SFR	2, RDMBSS	SAVE F2.
	SFR	3, RDMBSS+1	AND F3.
	SFR	4, RDMBSS+2	SAVE RETURN.
	SFR	5, RDMBSS+3	AND F5.
	LFR	5, M3	M4 = ADUM, M5 = ECDUM, M6 = PHID
	CAD	M6	
	DIV	180.	
	CALL	COS1	
	MPY	M5	
	STR	F2	
	MPY	F2	
	STR	F2	
	CAD	1.	
	SUB	F2	
	CALL	SQR1	
	STR	F3	
	CAD	M5	
	MPY	F1	
	STR	F2	
	CAD	1.	
	SUB	F2	
	CALL	SQR1	
	DIV	F3	
	MPY	M4	
	STR	M7	
	LFR	2, RDMBSS	
	LFR	3, RDMBSS+1	
	LFR	4, RDMBSS+2	
	LFR	5, RDMBSS+3	
	JLH	M3+1	EXIT.
RDMBSS	BSS	4	

---

ASSIGN SCALE,RPSZ,Q2,Q4,Q6,Q9,Q12,XINP,XDINP,XINS,XDINS,PSI

---

ASSIGN AVRUC,RMRUC,AVG,RMG,PHIIN,RADIN,ROUT,RIN

---

ASSIGN XIN,XDIN,TZERO,XZERO,XDZERO,MOTION,T,X,XD,DELX,DELXD

---

ASSIGN AQ,CQ,BNQCQ,TESTU,TSTSBC,TSTCNU,U,UO,SNU,CNU,DNU,RADNEW

---

ASSIGN COUNT,RUCRIN,LARGE,F,SMALLG,RAD,PHI,XF,XDF,RADOLD,PHIOLD

---

ASSIGN DELRAD,RADFIN,PHINP,DELPHI,ECCEN,XMAX,SETUP

---

ASSIGN XPERM,XDPERM,AMJOUT,AMJIN,ECOUT,ECIN,AMJOR

---

ASSIGN AVGR,RMSR,SIGMAR,AVGG,RMSG,SIGMAG

---

READIN DECQL ,XPERM,16,INREAD

---

INREAD CHR 30,4F15.0,2F10.0/6F10.0,F5.0,3D5\*

---

PRINTN DECQL ,XPERM,16,NPRINT

---

NPRINT CHR 48,8H1XPERM =1PE24.15//9H XDPERM =1PE24.15//9H AMJO

---

CHR 48,UT =1PE24.15//8H AMJIN =1PE24.15//8H ECOUT =1PE1

---

CHR 48,9.10//7H ECIN =1PE19.10//8H AMJOR =1PE19.10//9H

---

CHR 48,DELRAD =1PE19.10//9H RADFIN =1PE19.10//8H PHINP

---

CHR 48,=2PE19.10//9H DELPHI =2PE19.10//8H ECCEN =1PE19

CHR 48,.10//7H XMAX =F10.2//19H FIXED POINT NUMBERD5//2

CHR 48,0H NUMBER OF ORBITS ISD5//21H NUMBER OF SECTORS

CHR 5,ISD5\*

ASSIGN COVAL

CTOFF DECQL ,COVAL,1,OFFCT

OFFCT CHR 3,D5\*

CTOFFP DECQL ,COVAL,1,OFFPCT

OFFPCT CHR 24,//16H CUTOFF VALUE ISD5\*

SMOUT1 DECQL ,AVGR,6,OUTSM1

OUTSM1 CHR 48,//17H AVERAGE RUCRIN =1PE14.6,15H RMS RUCRIN =

CHR 48,1PE14.6,21H RUCRIN DEVIATION =1PE14.6//14H A

CHR 48,VERAGE G =1PE14.6,10H RMS G =1PE14.6,16H G D

CHR 18,EVIATION =1PE14.6\*

HEAD DECQL ,,,DAEH

DAEH CHR 48,1H1,X,3HM11,11X,6HRUCRIN,11X,6HLARGE G,10X,1HF,11

CHR		48,X,6HSMALLG,14X,3HRAD,8X,3HPHI,15X,2HXF,14X,3HXDF
CHR		10,,2X,1HQ//*
OUT	DECQL	4096+3,COUNT,9,TUD
	DECQL	3,SETUP,1,
TUD	CHR	46,D5,1P2E17.9,F11.5,1P2E17.9,F11.5,1P2E17.9,D3/*
SKIP	DECQL	,,,PIKS
PIKS	CHR	5,3(/)*
*BUFF1	BSS	6
*SNARRY	BSS	901
*BUFF2	BSS	6
*CNARRY	BSS	901
*BUFF3	BSS	6
*DNARRY	BSS	901
*BUFF4	BSS	6
*SBYC	BSS	1702
CWLAG6	DECQL	SNARRY,1,SNU,1
	DECQ	CNARRY,1,CNU,1
	DECQ	DNARRY,1,DNU,
ELKAUX	DECQL	SNARRY+894,1,SNU,1
	DECQ	CNARRY+894,1,CNU,1
	DECQ	DNARRY+894,1,DNU,
BLKAUX	DECQL	SNARRY-6,1,SNU,1
	DECQ	CNARRY-6,1,CNU,1
	DECQ	DNARRY-6,1,DNU,
IGNORE	EQU	4096
LOOPS	EQU	2
RFAST	EQU	3
QFAST	EQU	2
TBY4	DEC	1.5
K	DEC	1.854074677301372
K4	DEC	7.416298709205488
RINCR	DEC	.55555555555555E-03
Q	DEC	.43213918263772E-01
PSZ	DEC	.037
PI	DEC	3.141592653589793
EOCT	CCTQL	.65,4343,7147
H1OCT	OCTQL	,,7030,4701
H2OCT	OCTQL	,,7030,4702

```

PHIST  CAD  DELX      SUBROUTIN TO CORRECT 'PHI'.
        TP  PHISTA
        TZ  PHISTB
        CAD 180.
        ASC PHI
        TRA PHISTX
PHISTA  CAD  DELXD
        TP  PHISTX
        CAD 360.
        ASC PHI
        TRA PHISTX
PHISTB  CAD  DELXD
        TP  PHISTC
        CAD 270.
        STR PHI
        TRA PHISTX
PHISTC  CAD  90.
        STR PHI
PHISTX  JLH  3,,      EXIT.
CXMAX   SFR  4,CBSS    SAVE RETURN.
        CAM  1,PXMAX

        CALL PRINT
        LFR  4,CBSS
        TRA  STXCAL
PXMAX   DECQL ,X,2,XMAXP
XMAXP   CHR  48,1H1///24H ***** XMAX EXCEEDED *****//3H X=1PE24.
        CHR  17,///4H XD=1PE24.15*

*CBSS   BSS   1
*
*
*
*
*
GO

```

### C. Fortran program for the IBM-7094

#### Input parameters

- NSTART:** An integer (1,2,3, or 4) representing the quarter-period of  $p(t)$  which initiates the motion.
- MLTPLS:** The number of integer multiples  $m$  of the basic mapping period.
- PERS:** The basic mapping period  $n$ .
- VINC:** The numerical value of the period  $\tau$  of  $p(t)$ .
- PZERO:** The numerical value of the magnitude  $p_0$  of  $p(t)$ .
- XMAX:** The maximum value of  $X$  which a final coordinate is allowed to attain.
- XPERM:** The  $X$ -coordinate of the point to be used as the origin of the polar coordinates.
- XDPERM:** The  $\dot{X}$ -coordinate of the point to be used as the origin of the polar coordinates.
- ECCEN:** The fixed eccentricity  $e_0$  corresponding to the initial point in phase-space being mapped.
- AMJOR:** The semi-major axis  $a_0$  corresponding to the initial point in phase-space being mapped.
- PHIIN:** The original angle  $\theta_0$  corresponding to the initial point in phase-space being mapped.
- PHIINC:** An increment in  $\theta_0$  if a regular series of  $\theta_0$  values is desired.
- PHIFIN:** The final  $\theta_0$  in the desired series of  $\theta_0$  values (may equal PHIIN).

#### Output parameters

- NCYCLE:** The running counter for the number of multiples  $m$  of the basic period  $n$ .
- V(1):** The running value of time  $t_1$ .

- V(2): The mapped X-coordinate  $X_1(t_1)$ .
- V(3): The mapped  $\dot{X}$ -coordinate  $\dot{X}_1(t_1)$ .
- RAD: The mapped radius  $r_1$  about (XPERM, XDPERM).
- PHI: The mapped angle  $\theta_1$  about (XPERM, XDPERM).
- RUERIN: The ratio of the mapped radius  $r_1$  at the angle  $\theta_1$  to the radius  $r_e$  of a standard ellipse (specified by AMJOR and ECCEN) at the same angle  $\theta_1$ . [This is called RUCRIN in the NICAP program.]
- FMOSE: Moser's  $\Delta\theta$  function =  $(\theta_1 - \theta_0)$ .
- GMOSE: Moser's  $\Delta r$  function =  $(r_1 - r_0)$ .

```

DIMENSION D(4),E(4),F(4),VIP(3),VD(3),V(3)
COMMON V,VD,VINC,VIP,D,E,F,PZERO,S
1 READ INPUT TAPE 7,2,NSTART,MLTPLS,PERS,VINC,PZERO,XMAX,
  XPERM,XDPERM
2 FORMAT ( I2,I6,2F6.0,2F10.0,2F15.0 )
  D(1) = 0.5
  D(2) = 0.29289322
  D(3) = 1.7071068
  D(4) = 0.16666667
  E(1) = 2.0
  E(2) = 1.0
  E(3) = 1.0
  E(4) = 2.0
  F(1) = -0.5
  F(2) = -0.29289322
  F(3) = -1.7071068
  F(4) = -0.5
  CONV = 0.0174532925
  NTOTAL = 6.0*PERS/VINC
3 READ INPUT TAPE 7,5,ECCEN,AMJOR,PHIIN,PHIINC,PHIFIN
5 FORMAT (5F15.0)
12 WRITE OUTPUT TAPE 6,15
15 FORMAT (18H1 FIXED INPUT DATA ////)
18 WRITE OUTPUT TAPE 6,20,ECCEN,AMJOR,PHIIN,PHIINC,PHIFIN
20 FORMAT (10H ECCEN=1PE16.9, 8H AMJOR=1PE16.9,8H PHIIN=1PE16.9,
  1 9H PHIINC=1PE16.9, 9H PHIFIN=1PE16.9 //)
23 WRITE OUTPUT TAPE 6,25,NSTART,MLTPLS,PERS,VINC
25 FORMAT(11H NSTART= 12, 11H MLTPLS= 16, 9H PERS= F6.2,
  1 9H VINC= 1PE13.6 //)
26 WRITE OUTPUT TAPE 6,27,PZERO,XMAX,XPERM,XDPERM
27 FORMAT(10H PZERO= 1PE13.6, 9H XMAX= 1PE13.6,
  1 10H XPERM= 1PE20.12, 11H XDPERM= 1PE20.12 ////)
28 WRITE OUTPUT TAPE 6,30
30 FORMAT (7HONCYCLE,6X,2H T,13X,2H X,12X,5H XDOT,10X,4H RAD,11X,
  1 4H PHI,10X,7H RUERIN,8X,7H FMOSER,8X,7H GMOSER ////)
50 TOP = (1.0-ECCEN**2)*AMJOR*AMJOR
51 BOT = 1.0-(ECCEN*COS(PHIIN*CONV))**2

```

```

52 RADIN = SQRT(TOP/BOT)
53 RUERIN = 1.0
55 FMOSER = 0.0
56 GMOSER = 0.0
60 V(1) = 0.0
61 V(2) = XPERM+RADIN*COS(PHIIN*CONV)
62 V(3) = XDPERM+RADIN*SIN(PHIIN*CONV)
65 NCYCLE = 0
70 WRITE OUTPUT TAPE 6, 75,NCYCLE,V,RADIN,PHIIN,RUERIN,FMOSER,GMOSER
75 FORMAT (1H 14,2X, 8(1H 1PE14.7 ) //)
100 GO TO (105,115,110,120), NSTART
105 S = 1.0
106 T = 1.5
107 GO TO 150
110 S = -1.0
111 T = 1.5
112 GO TO 150
115 S = -1.0
116 T = 3.0
117 GO TO 150
120 S = 1.0
121 T = 3.0
150 NTRACE = 0
207 NCOUNT = 0
210 VIP(1) = 0.0
211 VIP(2) = 0.0
2111 VIP(3) = 0.0
212 IF (ABSF(V(2)) - XMAX) 216, 213, 213
213 WRITE OUTPUT TAPE 6,214,V
214 FORMAT (5H XMAX 1PE20.12,2(5H 1PE20.12))
215 GO TO 3
216 CALL INTEGR
2161 NCOUNT = NCOUNT + 1
2162 NTRACE = NTRACE + 1
225 IF (NTRACE - NTOTAL) 226, 300, 3
226 NSTEPS = T/VINC

```

```

227 IF (NCOUNT - NSTEPS) 210, 230, 3
230 IF (T - 2.0) 231, 3, 240
231 T = 3.0
240 S = -S

241 GO TO 207
300 DELX = V(2) - XPERM
301 DELXDT = V(3) - XDPERM
302 RAD = DELX**2 + DELXDT**2
303 RAD = SQRT(RAD)
304 PHI = ATAN(DELXDT/DELX)/CONV
3041 IF (DELX) 3045, 3042, 3042
3042 IF (DELXDT) 3043, 305, 305
3043 PHI = PHI + 360.0
3044 GO TO 305
3045 PHI = PHI + 180.0
305 BOTPHI = 1.0 - (ECCEN * COS(PHI * CONV))**2
306 RADPHI = SQRT(TOP/BOTPHI)
307 RUERIN = RAD/RADPHI
3071 FMOSER = PHI - PHIIN
3072 GMOSER = RAD - RADIN
308 NCYCLE = NCYCLE + 1
310 WRITE OUTPUT TAPE 6, 75, NCYCLE, V, RAD, PHI, RUERIN, FMOSER, GMOSER
316 IF (NCYCLE - MLTPLS) 100, 317, 350
317 IF (PHIIN - PHIFIN) 318, 3, 3
318 PHIIN = PHIIN + PHIINC
319 GO TO 28
350 CALL SYSTEM
351 END

SYSTEM
ATAN
INTEGR
ABSF
SIN
SQRT
COS
PRINT
READ

```

```

1000  SUBROUTINE INTEGR
      DIMENSION D(4),E(4),F(4),VIP(3),VD(3),V(3)
      COMMON V,VD,VINC,VIP,D,E,F,PZERO,S
1010  DO 1018 J = 1, 4
1011  VD(1) = 1.0
1012  VD(2) = V(3)
1013  VD(3) = -PZERO * S * V(2)**3
1014  DO 1018 K = 1, 3
1015  VD(K) = VINC * VD(K)
1016  R = D(J) * (VD(K) - E(J) * VIP(K))
1017  VIP(K) = VIP(K) + 3.0 * R + F(J) * VD(K)
1018  V(K) = V(K) + R
1019  RETURN
1020  END
      SYSTEM

```

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## VITA

Curtis Arthur Wagner was born on [REDACTED] [REDACTED] in [REDACTED] [REDACTED], where he attended public schools. He entered the University of Wisconsin in 1956, and received the degree of Bachelor of Arts in Physics, with High Honors in the Major, in June, 1960. The same year he enrolled in the Graduate College of the University of Illinois, majoring in physics and minoring in mathematics and astronomy. He was a Woodrow Wilson Honorary Fellow, and held National Science Foundation National and Cooperative Fellowships for four years. During this time he spent a year at the University of Copenhagen studying celestial mechanics. In the succeeding years he has held both teaching and research assistantships in the Physics Department while working with Professor J. H. Bartlett in the field of nonlinear mechanics. He has published one paper, on the restricted problem of three bodies. He married the former Margret Anne Schoenbohm on November 11, 1967.