Relation between solutions of the Schrödinger equation with transitioning resonance solutions of the gravitational three-body problem

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Abstract

It is shown that a class of approximate resonance solutions in the three-body problem under the Newtonian gravitational force are equivalent to quantized solutions of a modified Schrödinger equation for a wide range of masses that transition between energy states. In the macroscopic scale, the resonance solutions are shown to transition from one resonance type to another through weak capture at one of the bodies, while in the Schrödinger equation, one obtains quantized wave solutions transitioning between different energies. The resonance transition dynamics provides a classical model of a particle moving between different energy states in the Schrödinger equation. This methodology provides a connection between celestial and quantum mechanics.

1. Introduction

The purpose of this paper is to describe a mechanism to globally model the solutions of a modified Schrödinger equation and how they transition between energy states, with a special set of approximate resonance solutions to the classical gravitational Newtonian three-body problem, for a wide range of masses. These resonance solutions transition from one resonance to another through the process of weak capture.

We consider a special version of the three-body problem that has proved to be useful in understanding the complexities of three-body motion, in the macroscopic scale, going back to Poincaré [1]. This is the circular restricted three-body problem, where the motion of one body, \( P_0 \), is studied as it moves under the influence of the gravitational field of \( P_1 \), \( P_2 \), assumed to move in mutual circular orbits of constant frequency \( \omega \). It is also assumed that the mass of \( P_0 \), labeled \( m_0 \), is negligible with respect to the masses of \( P_1 \), \( P_2 \), labeled \( m_1 \), \( m_2 \), respectively. In this paper, we will also assume that \( m_3 \) is much smaller than \( m_1 \), \( m_2 \ll m_0 \). For example, in the case of planetary objects, one can take \( P_1 \), \( P_2 \) to be the Earth, Moon, respectively, and \( P_0 \) to be a rock. One can scale down \( m_0 \), \( m_1 \), \( m_2 \), as well as the relative distances between the particles, until the quantum scale is reached where the pure gravitational modeling is no longer sufficient. When using a rotating coordinate system that rotates about the center of mass of \( P_1 \), \( P_2 \) with constant frequency \( \omega \), it is well known that Hamiltonian function for the motion of \( P_0 \) is time independent defining a conservative system (see section 3).

For a general class of conservative systems, that includes, for example, the restricted three-body problem considered here, it is known that such systems can be associated to the Schrödinger equation. As is described in Lanczos [2], this can be done by computing the action function \( S(x) \) for the motion of \( P_0 \), where \( x \) is the position of \( P_0 \). \(^1\) \( S \) is a solution of the Hamilton-Jacobi partial differential equation associated to the restricted three-body problem. The \( x(t) \) are orthogonal to the surface \( S(x) = C \), for constant \( C \). In this sense, the iso-\( S \) surfaces locally determine the trajectories \( x(t) \), that is for \( t \) having sufficiently small variation. On the other hand, it can be shown that \( S \) is the phase of wave solutions to the Schrödinger equation. Thus, \( S(x) = C \) is a wave surface.

\(^1\) \( S \) is obtained as the integral of the Lagrangian function \( L(x, y) \), \( y = dx/dt \), over minimizing trajectories \( x(t) \).
Conversely, starting with the $S$ satisfying the Schrödinger equation, one obtains the Hamiltonian-Jacobi equation provided the Planck’s constant $\hbar \to 0$. This equivalence is local since in general $S$ can only be shown to locally exist as a solution to the Hamiltonian-Jacobi equation, which is the case for the three-body problem. For the full equivalence it is necessary that one restricts $\hbar \to 0$ that is not realistic. More significantly, this equivalence does not determine the global behavior of the solutions $x(t)$ and how they dynamically can model the transition of energy states for solutions of the Schrödinger equation.

This paper will use methods of dynamical systems to globally determine special resonance motions for $P_0$ that are shown to be equivalent to different energy states for a modified Schrödinger equation and where the transition between energy states is equivalent to the process of weak capture in the three-body problem described in this paper. This result does not require the restriction of Planck’s constant, and doesn’t use Planck’s constant in the modeling. The action function $S$ is not used. Due to the complexity of the motions described, it is seen that using the iso-$S$ surfaces to locally determine the global solutions would not seem feasible.

We describe a mechanism for the existence of a special family, $\mathcal{F}$, of approximate resonance motions of $P_0$ about $P_1$, that transition from one resonance to another by the process of weak capture by $P_2$. This is a temporary capture defined in section 2. These motions are approximately elliptical with frequency $\omega_1 \equiv \omega_1(m/n)$, $\omega_1(m/n)(t)$ or equivalently $\omega_1(t)$, are functions of the time $t$, where $\omega_1(m/n)$ is approximately equal to the constant values $(m/n)\omega_1$, $m$, $n$ are positive integers. That is, the period of motion of $P_0$ is approximately synchronized with the circular motion of $P_2$ about $P_1$, where in the time $P_1$ makes approximately $n$ revolutions about $P_1, P_2$ makes $m$ revolutions about $P_1$. The approximate resonance value of the frequency means that $|\omega_1(t) - (m/n)\omega_1| < \delta$, for a small tolerance $\delta$ as time varies for restricted time spans, described in section 2. When $P_0$ is moving on an approximate resonance orbit about $P_1$, it will eventually move away from this orbit and become captured temporarily about $P_2$, in weak capture. When $P_0$ escapes from this capture, it again moves about $P_1$ in another resonance elliptical orbit, with approximate resonance $m'/n'$. This process repeats either indefinitely, or ends when, for example, $P_0$ escapes the $P_1, P_2$-system. This also implies that the approximate two-body energy $E_1$ of $P_0$ about $P_1$ can only take on a discrete set of values, $E_1(m/n)(t)$ at each time $t$, which are approximately constant defined by the resonances. This is stated as result A in section 2 and as theorem A in section 3. The properties and dynamics of weak capture, and weak escape, are described section 3. Comets can perform such resonance transitions (see sections 2, 3).

A modified Schrödinger equation is defined for the motion of $P_0$ about $P_1$, under the gravitational perturbation of $P_2$. This is first considered in the case of macroscopic masses. It is given by,

$$-\frac{\sigma^2}{2\nu} \nabla^2 \Psi + \tilde{V} \Psi = E \Psi,$$

where $\nabla^2 \equiv \nabla \cdot \nabla$ is the Laplacian operator, $\tilde{V}$ is an averaged three-body gravitational potential, $E$ is the energy, and $\nu$ is the reduced mass for $P_0, P_1$, $\sigma$ is a function that depends on $m_0, m_1$ and $G$, the gravitational constant, $\sigma$ replaces $h = h/2\pi$, $h$ is Planck’s constant that is in the classical Schrödinger equation. In this case, for macroscopic values of the masses, since $P_0$ is not a wave, $\Psi$ is used to determine the probability distribution function, $|\Psi|^2$, of locating $P_0$ near $P_1$ as a macroscopic body. We show in section 4 that $E$ can only take on the following approximate quantized values,

$$E_{\tilde{n}} = \frac{-4\sigma}{\tilde{n}^2},$$

$\tilde{n} = 1, 2, 3, \ldots$. As seen in section 2, this implies that the frequency of $P_0$ takes a particularly simple form that is independent of any parameters. These frequencies have the approximate values, $8/\tilde{n}^2$. $\Psi$ is explicitly computed in section 4, $|\Psi|^2$ is shown to be exponentially decreasing as a function of the distance of $P_0$ from $P_1$. The general solution, $\Psi$, of the modified Schrödinger equation is described in result B in section 2.

A main result of this paper is that the quantized energy values $E_{\tilde{n}}$ correspond to a subset, $\mathbf{1}_E$, of the resonance orbit family, $\mathcal{F}$, of $P_0$ about $P_1$. This is listed as result C in section 2. This provides a global equivalence of the solutions of the modified Schrödinger equation with the transitioning resonance solutions of the three-body problem.

As a final result, we show that the solution, $\Psi$, for the location of $P_0$ for the modified Schrödinger equation for the macroscopic values of the masses, can be extended into the quantum-scale. This is summarized as result D in section 2. This gives a way to mathematically view the resonance motions in the quantum-scale, as an extension of the resonance solutions for macroscopic particles. Other models, such as the classical Schrödinger and Schrödinger-Newton equations are given in latter sections.

The results of this paper are described in detail and summarized in section 2. This section contains the main findings of this paper. Additional details, derivations, and proofs are contained in section 3 for weak capture and in section 4 for the modified Schrödinger equation.
2. Summary of results, definitions and assumptions

In this section we elaborate on the results described in the Introduction. The first set of results pertain to a family of resonance orbits about $P_1$ obtained from the three-body problem and the second set of results pertain to finding these orbits using a modified Schrödinger equation.

2.1. Resonance orbits in the three-body problem and weak capture

The motion of $P_0$ is defined for the circular restricted three-body problem described in the Introduction. It is sufficient to use the planar version of this model, without loss of generality for the purposes of this paper, where $P_0$ moves in the same plane of motion as that of the uniform circular motion of $P_0$, $P_2$ of constant frequency $\omega$ (see section 3). The macroscopic masses satisfy, $m_0/m_1 \ll 1$ and the mass of $m_0$ is negligibly small so that $P_0$ does not gravitationally perturb $P_1$, $P_2$, but $P_0$, $P_2$ perturb the motion of $P_0$. We consider an inertial coordinate system, $(X_1, X_2) \in \mathbb{R}^2$, whose origin is the center of mass of $P_1$, $P_2$.

The differential equations for $P_0$ are given by the classical system

$$\dot{X} = \Omega_X(X, t),$$

where $X = (X_1, X_2) \in \mathbb{R}^2$, $t \in \mathbb{R}$, $d\Omega_X = \Omega_X \Omega_X \Omega_X \Omega_X$, $\Omega_X \equiv \partial \Omega / \partial X$ and

$$\Omega = \frac{GM_1}{\eta(t)} + \frac{GM_2}{r_2(t)},$$

where $\eta(t) = |X - a(t)|$, $r_2 = |X - a_2(t)|$, $|\cdot|$ is the standard Euclidean norm. The mutual circular orbits of $P_1$, $P_2$ are given by $a_1(t) = \rho_1(\cos \omega_1 t, \sin \omega_1 t)$, $a_2(t) = -\rho_2(\cos \omega_2 t, \sin \omega_2 t)$, with constant circular frequencies $\omega_1$, $\omega_2$ of $P_0$, $P_2$, respectively. We have divided both sides of (3) by $\eta_0$ and then took the limit as $m_0 \to 0$. It is well known that the solutions of the circular restricted three-body problem for $P_0$ accurately model the motion of $P_0$ in the general three-body problem for circular initial conditions for $P_0$, $P_2$, with $m_0$ kept positive and negligibly small.

It is noted that all solutions $\xi(t) = (X(t), \dot{X}(t)) \in \mathbb{R}^4$ considered in this study will be $C^\infty$ in both initial conditions, $\xi(t_0) = (X(t_0), \dot{X}(t_0))$ at an initial time $t_0$. We refer to $C^\infty$ as smooth dependance. More exactly, this means that all derivatives of $\xi(t)$ with respect to $t$ of all orders are continuous and all partial derivatives of $\xi(t)$, $\xi(t_0)$ with respect to $X_1(t_0)$, $X_2(t_0)$, $\dot{X}_1(t_0)$, $\dot{X}_2(t_0)$, of all orders, are continuous.

Although $m_0$ is taken in the limit to be 0 in the definition of the differential equations for the motion of $P_0$ on, we will assume it is non-zero but still negligible in mass with respect $P_1$, $P_2$, $m_0 \gtrless 0$, in all equations that follow.

We transform to a $P_1$-centered coordinate system for the restricted three-body problem. In this system, $P_2$ moves about $P_1$ at a constant distance $\beta$, with constant circular frequency $\omega = \sqrt{G(m_1 + m_2) / \beta}$. Before stating our first result, two definitions are needed.

Assuming $m_1$ is much smaller than $m_2$, when $P_0$ moves about $P_1$ with elliptic initial conditions at an initial time $t = t_0$, this elliptic motion will be slightly perturbed by $P_0$'s. Let $a_1$ be the semi-major axis of $P_0$ with respect to $P_1$. As a function of time, $a_1$ will vary. If $m_2 = 0$, then $a_1$ is constant since $P_0$ will move on a pure ellipse. If $m_2$ is small, then $P_0$ moves in a nearly elliptic orbit about $P_1$, and $a_1(t)$ will be nearly constant for restricted time spans. This orbital element, along with the eccentricity, $e_1(t)$, true anomaly, $\theta_1(t)$, and other orbital elements, can be calculated for each $t$ using the variational differential equations obtained from (3). (See [3–6]). These are referred to as osculating elements. $e_1$ will likewise be nearly constant for a nearly elliptical orbit of $P_0$ about $P_1$.

The variation of the frequency $\omega_1(t)$ can be obtained from $a_1(t)$: The osculating two-body period, $T_1$, of $P_0$ is explicitly related to $a_1(t)$ by Kepler’s Third Law, $a_1^3 = (2\pi)^{-1} T_1^2 G(m_0 + m_1)$, and $\omega(t) = T_1^{-1}$.

**Definition 1.** An approximate resonance orbit, $\Phi_{m/n}(t)$, of $P_0$ moving about $P_1$ in a $P_1$-centered coordinate system, $Y = (Y_1, Y_2)$, as a function of $t$ in resonance with $P_0$, is an approximate elliptical orbit of frequency $\omega_1 = \omega_1(t)$, where $\omega_1 \approx (m/n) \omega$, $m$, $n$ are positive integers. Thus, $\omega_1$ is approximately constant as time varies. In phase space, $(Y, \dot{Y}) \in \mathbb{R}^4$, $\Phi_{m/n}(t) = (Y_1(t), Y_2(t), Y_1(t), Y_2(t))$. $\Phi_{m/n}(t)$ has a period $T_1 = \omega_1^{-1}$, approximately constant. $T_1 \approx (n/m) T$, $T$ is the constant circular period of $P_0$ about $P_1$, $T = \omega^{-1}$. For notational purposes, we refer to an approximate resonance orbit as a resonance orbit for short. A resonance orbit with $\omega_1 \approx (m/n) \omega$ is also referred to as a $n$: $m$ resonance orbit. *(Nearly resonant motion, related to approximate resonance motion, is described in [7].)*

The term ‘approximate’ in definition 1 means to within a small tolerance, $O(\delta)$, $\delta = m_2/m_1 \ll 1$. $O(\delta)$ is a function of time, $t$, and smooth in $t$. An approximate elliptic orbit means that the variation of the orbital parameters $(\omega_1, a_1, e_1)$ of $\Phi(t) = (Y(t), \dot{Y}(t))$, with respect to $P_1$ in a $P_1$-centered coordinate system, will slightly
vary by $O(\delta)$ due to the gravitational perturbation of $P_2$. The two-body energy of $\Phi_m/n(t)$ with respect to $P_1$ is labeled, $E_1(m/n) = E_1(m/n)(t)$, which is approximately constant. Note that the general Kepler energy $E_1(t)$ along a trajectory $(Y(t), \dot{Y}(t))$ is given by (7).

Thus, $\omega_1 \approx (m/n)\omega$ is equivalent to $\omega_1 = (m/n)\omega + O(\delta)$. $a_1, \epsilon_1$ likewise vary within a variation $O(\delta)$. The variations $O(\delta)$ are all different functions for the different parameters, but the same notation is used. The tolerance on these orbital parameters is valid for finite times. We assume that $t$ varies over finite time spans. Thus, for a given variable, say $\omega_1(t)$, if $\epsilon_1 > 0$ is a given number, and $t \in [t_0, t_1]$, $t_1 > t_0, m_2$ can be taken small enough so that $|O(\delta)| < \epsilon_1$.

For a resonance orbit to be well defined, it is assumed that $m_2 > 0$. If $m_2 = 0$, then even though $\omega$ is defined, $P_2$ no longer exists. Thus, we assume $m_2 > 0$ throughout this paper, unless otherwise indicated. This assumption is also necessary for the definition of weak capture.

We define ‘weak capture’ of $P_0$ about $P_2$. In this case, we change to a $P_2$-centered coordinate system. This type of capture is discussed in section 3. Weak capture is where the two-body energy, $E_2$, of $P_0$ with respect to $P_2$ is temporarily negative. It is used to define an interesting region about $P_2$ described in the next section called the weak stability boundary. Chaotic motion occurs on and near this region.

**Definition 2.** $P_0$ has weak capture about $P_2$, in a $P_2$-centered coordinate system, $Z = (Z_1, Z_2)$, at a time $t_0$ if the two-body energy, $E_2$, of $P_0$ with respect to $P_2$, is negative at $t_0$, and for a finite time after where it becomes positive ($P_0$ escapes). More precisely, $E_2$ is given by

$$E_2 = \frac{1}{2} |Z|^2 - \frac{G(m_2 + m_2)}{r_2},$$

where $r_2 = |Z| > 0$. Let $Z(t) = (Z_1(t), Z_2(t), \dot{Z}_1(t), \dot{Z}_2(t))$ be a solution for the differential equations (3) in $P_2$-centered coordinates for $t \geq t_0$. $P_0$ is weakly captured at $t_0$ if $E_2(Z(t)) < 0$ for $t_0 < t < t_1$, $E_2(Z(t)) \geq 0$ for $t \geq t_1$. After $P_0$ leaves weak capture at $t_1$, we say that $P_0$ has weak escape from $P_2$ at $t = t_0$. Weak capture in backwards time from $t_0$ is similarly defined.

$P_0$ is captured at a point $(Z_1(t), Z_2(t))$ of a trajectory $Z(t)$ at $t^*$ if $E_2(Z(t^*)) < 0$. Capture at a point need not imply weak capture, in forward time, since $P_0$ could be captured for all time $t > t^*$.

This dynamics is summarized in result A and proven in section 3, where it is formulated more precisely as Theorem A.

**Result A.** Weak capture of $P_0$ about $P_2$ at a time $t_0$ yields resonance motion of $P_0$ about $P_1$, which repeats yielding a family, $\mathfrak{F}$, of resonance orbits. More precisely,

Assume $P_0$ is weakly captured by $P_2$ at time $t = t_0$, then

(i) As $t$ increases from $t_0$, $P_0$ first escapes from $P_2$ and then $P_0$ moves onto a resonance orbit, $\Phi_m/n(t)$, about $P_1$. $P_0$ performs a finite number of cycles about $P_1$ until it eventually moves again to weak capture by $P_2$, where the process continues and $P_0$ moves onto another resonance orbit. When $P_0$ moves from $\Phi_m/n(t)$ to another resonance orbit, $\Phi_m'/n'(t)$, $m', n'$ may or may not equal $m, n$. In general, a sequence of resonance orbits is obtained, $\{\Phi_m/n(t), \Phi_m'/n'(t), \ldots\}$. The process stops when $P_0$ escapes the $P_1P_2$-system, collides with $P_2$ or moves away from a resonance frequency. This set of resonance orbits forms a family, $\mathfrak{F}$, of orbits, that depend on the initial weak capture condition.

(ii) When $P_0$ moves onto a sequence of resonance orbits about $P_1$ as described in (i), then a discrete set of energies are obtained, $\{E_1(m, n), E_1(m', n'), \ldots\}$.

The proof of theorem A is given in detail in section 3.

Applications of theorem A to comet motions and numerical simulations is described in section 3.

A key result obtained in section 3 is,

**Lemma A.** The frequency $\omega_1(m/n)(t)$ of $\Phi_m/n(t)$ is given by,

$$\omega_1(m/n) = (m/n)\omega + O(\delta),$$

where $O(\delta)(t)$ is smooth in $t$. (i) Is proven in section 3 as theorem A. We prove (ii): consider the general two-body energy of $P_0$ with respect to $P_1$. It is given by,
where $Y = (Y_1, Y_2)$, $\eta = |Y|$, are inertial $P_1$-centered coordinates. $E_i$ can be written as

$$E_i = -G(m_0 + m_i)/(2a_i).$$

Using Kepler’s third law relating the period, $T_1$, to $a_1$, implies, $E_i = -A\omega_1^{2/3}$,

$$A = (1/2)(2\pi G(m_0 + m_i)^{2/3}.$$

(6) Implies $E_i$ can be written as,

$$E_i(\Phi_{m/n}(t)) \equiv E_i(m/n) = -[(m/n)\omega_1^{2/3}A + O(\delta)],$$

where the remainder is smooth in $t$. (8) Proves (ii). It is noted that in this case, $E_i \approx -[(m/n)\omega_1^{2/3}A$.

It is noted that $E_i$ can be written in an equivalent form to (7) by multiplying both sides of (7) by the reduced mass, $\nu = m_0m_i(m_0 + m_i)^{-1}$, yielding

$$\tilde{E}_i \equiv \nu E_i = \frac{\nu}{2}|Y|^2 - \frac{Gm_0m_i}{\eta},$$

which implies, $\tilde{E}_i = -Gm_0m_i/(2a_i).$ This scaled energy is more convenient to use when the modified Schrödinger equation is considered. To obtain the corresponding two-body differential equations with (9) as an integral, one multiplies the differential equations associated to (7) by $\nu$. The solutions are the same for both sets of differential equations. Thus, it is seen that $E_i = \nu E_i = -\nu A\omega_1^{2/3}$. Setting $\sigma = \nu A$ implies,

$$\tilde{E}_i = -\sigma\omega_1^{2/3},$$

$$\sigma = (1/2)(2\pi G)^{2/3}m_0m(m_0 + m_i)^{-1/3}. \tag{11}$$

$\omega_1$ remains the same since the solutions haven’t changed. This defines the function $\sigma$ that plays a key role in this paper.

2.2. A modified Schrödinger equation: macroscopic scale

Analogies A. It is noted that equation (10) has a form similar to the Planck-Einstein relation for quantum mechanics for the energy, $\mathcal{E}$, of a photon, $\mathcal{E} = h\lambda$, where $\lambda$ is the frequency of the photon. $\sigma$ is analogous to $h$ and $\omega_1^{2/3}$ is analogous to $\lambda$. There is another analogy for the case of an electron, $P_0$, moving about an atomic nucleus: when $P_0$ changes from one orbital to another, the energy of the photon absorbed or emitted is given by $\Delta E = h\lambda$, where $\Delta E = E_i - E_2$, where $E_i$ is the energy of $P_0$ in the $i$th orbital, $i = 1, 2$. This process is analogous to a macroscopic particle $P_0$ changing from one resonance orbit $\Phi_{m/n}(t)$ in $\mathcal{F}$ to another, $\Phi_{m'/n'}(t)$, through weak capture and escape, where $\Delta \tilde{E}_i = \tilde{E}_i(m/n) - \tilde{E}_i(m'/n')$. $\Delta \tilde{E}_i$ is analogous to $\Delta E$.

We consider a modified Schrödinger equation (1). This Schrödinger equation differs from the classical one by replacing $h$ by the function $\sigma = \sigma(m_0, m_i, G)$ and $V$ by a three-body potential $\tilde{V}$ derived from the circular restricted three-body problem. The choice of $\sigma$ is motivated by analogy A. In this modeling, the masses

\[\Phi_{m/n}\text{to another, } \Phi_{m'/n'}; \text{ about } P_1 \text{ through weak capture of } P_0 \text{ near } P_2.\]
m_0, m_1, m_2 and distance between them is assumed to not be in the quantum scale. Under this assumption, \( \Psi \) does not represent a wave motion and is used to measure the probability of locating \( P_0 \) at a distance \( r_1 \) from \( P_1 \).

The potential \( \Phi \) in (1), is derived from the restricted three-body problem modeling for the motion of \( P_0 \). In inertial coordinates, \( Y = (Y_0, Y_1, Y_2) \) centered at \( P_1 \), \( P_2 \) moves about \( P_1 \) on a circular orbit, \( \gamma(t) \) of constant radius \( \beta \) and circular frequency \( \omega = \sqrt{\gamma(m_1 + m_2)/\beta} \), \( \gamma(t) = -\beta(\cos \omega t, \sin \omega t) \). The potential for \( P_0 \) is given by

\[
\Phi = V_1 + V_2,
\]

\[
V_1 = -\frac{Gm_0m_1}{r}, \quad V_2 = -\frac{Gm_0m_2}{r_2},
\]

\[ r = |Y|, \quad r_2 = |Y - \gamma(t)|. \]

For simplicity of notation, we have replaced the symbol \( r_1 \) by \( r (n \equiv r) \).

We replace \( V_2 \) by an approximation given by an averaged value of \( V_2 \) over one cycle of \( \gamma(t), t \in [0, 2\pi/\omega] \),

\[
V_2 = -\frac{Gm_0m_2\omega}{2\pi} \int_0^{2\pi} \sqrt{r^2 + \beta^2 - 2\beta(Y_1\cos \omega t + Y_2\sin \omega t)} \, dt. \]

It is proven in section 4 that \( V_2 \) can be reduced to three cases in (48) depending on \( r < 0, r = 0, r > 0 \), respectively, where the first order term in \( V_2 \) has a form similar to \( V_1 \).

As an approximation to \( \Phi \) we use

\[
\Phi = V_1 + \tilde{V}_2.
\]

The modified Schrödinger equation that we consider is given by

\[
-\frac{\sigma^2}{2\nu} \nabla^2 \Psi + \tilde{V} \Psi = E\Psi,
\]

The solution of this modified Schrödinger equation is derived in section 4. The solution is summarized in result B.

Result B. The explicit solution of the modified Schrödinger equation, more generally in a three-dimensional \( P_1 \)-centered inertial coordinates, \((Y_0, Y_2, Y_3)\), (16), is given by,

\[
\Psi = R_{n\ell}(r) Y_{m\ell}(\phi, \theta) + \mathcal{O}(m_0m_2).
\]

\[ n = 0, 1, 2 \ldots; l = 0, 1, 2 \ldots; -l \leq m \leq l, \]

where \( \Phi \) is the distance from \( P_0 \) to \( P_1 \), \( r \geq 0, \phi \in [0, 2\pi] \) is the angle in the \( Y_0, Y_2 \)-plane relative to the \( Y_1 \)-axis, \( \theta \in [0, \pi] \) is the angle relative to the \( Y_3 \)-axis. \( R_{n\ell}(r) \) is given by (39) defined using Laguerre polynomials. \( Y_{m\ell} = \Phi_m(\phi) \Theta_{\ell m}(\theta) \) are spherical harmonic functions, where \( \Phi_m(\phi), \Theta_{\ell m}(\theta) \) are given by (35), (37), respectively. \( \Theta(\theta) \) is defined by Legendre polynomials.

\[ |\Psi|^2 \] is the probability distribution function of finding \( P_0 \) at a location \((r, \phi, \theta)\). In particular, the radial probability distribution function of finding \( P_0 \) at a radial distance \( r \) is given by

\[
P(r) = R^2(r) r^2 + \mathcal{O}(m_0m_2),
\]

\[ \Psi \text{ exists provided the energy, } E, \text{ is quantized as,} \]

\[
E \equiv \tilde{E}_\ell = -\frac{4\sigma}{n^2} + \mathcal{O}(m_0m_2),
\]

where the remainder term is smooth in \( Y_0, Y_2, Y_3 \).

It is noted that the solution of (16) is valid for \( \tilde{V} = 0 \). However, we assume \( \tilde{V} \approx 0 \) to compare with the resonance solutions of three-body problem, where \( Y_3 = 0 \). All the terms \( \mathcal{O}(m_0m_2) \) are smooth in \( Y_0, Y_2, Y_3 \).

It is assumed that \( m_0 \) can be taken sufficiently small, such that for any given small number, \( \varepsilon_2 > 0 \), and for \((Y_0, Y_2, Y_3) \in D, D \) compact, the term \( \mathcal{O}(m_0m_2) \) in (19) satisfies, \( |\mathcal{O}(m_0m_2)| < \varepsilon_2 \). In this sense, \( \tilde{E}_\ell \approx -\frac{4\sigma}{n^2} \).

The planar case is now assumed, \( Y_3 = 0 \), unless otherwise indicated.

Assumptions 1. The use of the approximate symbol for \( \tilde{E}_\ell \) using \( \mathcal{O}(m_0m_2) \), is different to the one given in definition 1, for \( \omega_1 \) using \( \mathcal{O}(\delta) \), \( \delta = m_2/m_1 \). In definition 1, \( \mathcal{O}(\delta) \) depends on \( t \), and to bound it by a given small number \( \varepsilon_1, t \) varies on a compact set and \( m_2 \) is taken sufficiently small. In the second case, \( \mathcal{O}(m_0m_2) \), depends on \((Y_0, Y_2)\), and to bound it by a small number \( \varepsilon_2 \), \( \tilde{Y}_0, \tilde{Y}_2 \) varies on a compact set \( D \) and \( m_0 \) is taken sufficiently small. To satisfy both cases, it is necessary to assume \( m_0, m_1 \) are sufficiently small. The use of \( \approx \) is taken from context.
2.3. Equivalence of quantized energies with resonance solutions

The quantized values of the energy, \( \tilde{E}_n \), (19), are for the modified Schrödinger equation, (16). These energies are not obtained for the three-body problem, but result from an entirely different modeling. When they are substituted for \( \tilde{E}_n \) in the two-body energy relation, (10), for \( P_0 \) moving about \( P_1 \), it is calculated in section 4, (53), that they yield rational values for the two-body frequency, \( \omega_1 \),
\[
\omega_1 \bigg|_{\tilde{E}_n=\tilde{E}_n} \equiv \omega_1(\tilde{n}) = \frac{8}{\tilde{n}^3} + O(m_0 m_2).
\]
(20)

It is significant that the leading dominant term of \( \omega_1(\tilde{n}) \) is independent of masses and distances. This implies that to first order the frequencies do not depend on the masses or any other parameters.

Thus, for \( \tilde{n} = 1, 2, \ldots \), infinitely many frequencies are obtained, \( \omega_1(\tilde{n}) \). We would like to show that these frequencies correspond to \( \tilde{n} \) resonance orbits of \( \mathcal{F} \) for the three-body problem. Hence, we need to compare the frequencies \( \omega_1(\tilde{n}) \), given by (20), with the frequencies \( \omega_1(m/n) \), defined by (6). The following result is obtained,

**Result C.** The quantized energy values, (19), of the modified Schrödinger equation can be put into a one to one correspondence with a subset, \( \mathcal{U} \), of the resonance solutions \( \Phi_{m/n}(t) \) \( \in \mathcal{F} \) of the circular restricted three-body problem, where
\[
\mathcal{U} = \{ \Phi_{m/n}(t) | m = 8, n = \tilde{n}^3, \tilde{n} = 1, 2, \ldots \}.
\]
(21)

This is proven in section 4 by scaling the restricted three-body problem and using the fact that this scaling does not effect the leading order term \( \frac{8}{n^3} \) of \( \omega_1(\tilde{n}) \).

It is noted that \( m = 8, n = \tilde{n}^3 \) implies that in the time it takes \( P_0 \) to make \( \tilde{n}^3 \) cycles about \( P_1 \), \( P_2 \) makes 8 cycles about \( P_1 \).

As previously noted, the limiting case of \( m_0 = 0 \) has been excluded in this paper since it is degenerate in the sense that the resonance families of solutions no longer exist. One can make a comparison with quantized two-body elliptic orbits of \( P_0 \) about \( P_1 \) with the classical Schrödinger equation for \( m_2 = 0 \) (see [9], page 263), but this case does not yield the transitioning resonance solutions described in this paper.

2.4. Quantum scale

The results presented thus far are for mass values that are not in the quantum-scale. Consider the family, \( \mathcal{U} \subset \mathcal{F} \), of resonance periodic orbits for \( P_0 \) in the three-body problem, whose frequencies, \( \omega_1(m/n) \), given by (6), where \( m/n \approx 8/\tilde{n}^3 \). These frequencies correspond to the quantized energy values, \( \tilde{E}_n \), of the modified Schrödinger equation. When the masses, \( m_0, k = 0, 1, 2 \), get smaller and smaller, along with the relative distances between the particles, as they approach the quantum-scale, \( \omega_1, \omega_2 \), increase in value as \( \rho_1^{-1/2}, \beta^{-1/2} \) as \( \rho_1, \beta \rightarrow 0 \), respectively. The particles remain gravitationally bound to each other. The mass of \( P_0 \) is negligible with respect to that of \( P_1, P_2 \). As the distances decrease, the motions of the particles produces a gravitational field by the circular motion of \( P_1, P_2 \) and the resonance motion of \( P_0 \). We refer to this gravitational field as a *resonance gravitational field*.

When the system of three particles reaches the quantum scale they take on a wave-particle duality. The differential equations for the three-body problem are no longer defined. The previous resonance motion of the particles takes on a wave character.

The three-body problem is no longer defined in the quantum scale and therefore result A is no longer valid. However, the modified Schrödinger equation is still well defined. We can now assume the three-dimensional wave solutions. The quantized energy values are still defined, for \( \tilde{n} = 1, 2, \ldots \). Now, they are identified with pure wave solutions \( \Psi(r, \theta) \) given in result B. The values of \( \tilde{E}_n \), can be viewed as taking on wave resonance values. This is summarized in,

**Result D.** The resonance solutions \( \Psi_{m/n}(t) \in \mathcal{U} \) for \( P_0 \) for the three-body problem, which are given by the solutions \( \Psi, (17) \), of the modified Schrödinger equation are also given by \( \Psi \) when the masses are reduced to the quantum scale. This provides a quantization of the gravitational dynamics of \( P_1 \) for the motion of \( m_0 \) corresponding to the energies \( \tilde{E}_n \), given by (19).

In the quantum scale, where \( \sigma \rightarrow 0 \) as \( m_0, m_1 \rightarrow 0 \), shown in section 4, there is a transition of the resonance solutions into wave solutions, as summarized in result D, using \( \Psi \). However, to make these wave solutions more physically relevant, we would like to have \( \sigma(m_0, m_1, G) = \hbar \).

It is shown in section (4), proposition 4.1, that as \( m_0, m_1 \rightarrow 0 \), there exist mass values where \( \sigma(m_0, m_1, G) = \hbar \). These mass values lie on an algebraic curve in \((m_0, m_1)\)-space. For these values of \( m_0, m_1 \), the term \( -\frac{\sigma}{2\hbar} \nabla^2 \Psi \) of the modified Schrödinger equation matches the same term of classical Schrödinger
equation. In this case, only the gravitational potential is present. To make this accurate for atomic interaction, for example, for the motion of an electron about a nucleus of the Hydrogen atom, the gravitational potential needs to be replaced by the Coulomb potential.

If we consider the modified Schrödinger equation, it can be further altered by adding, for example, a Coulomb potential. If the masses are chosen so that \( \sigma = \hbar \), then one obtains a classical Newton–Schrödinger equation model [10, 11]. This could also be studied with \( \sigma \neq \hbar \).

The wave solutions of the modified Schrödinger equation could be considered in the quantum scale where \( \sigma \neq \hbar \). This is not studied in this paper.

3. Weak capture and resonant motions in the three-body problem

In this section we show how to prove result A. The idea of the proof of result A is to utilize the geometry of the phase space about \( P_1, P_2 \), where the motion of \( P_0 \) is constrained by Hill’s regions. Within the Hill’s regions, the dynamics associated to weak capture from near \( P_1 \) together with the global properties of the invariant hyperbolic manifolds around \( P_1 \) will yield the proof.

The planar circular restricted three-body problem in inertial coordinates is defined in section 2 by (3) for the motion of \( P_0 \). If \( P_0 \) moves about \( P_1 \) with elliptic initial conditions, and a rotating coordinate system is assumed that rotates with the same constant circular frequency \( \omega \) between \( P_1 \) and \( P_2 \), then the motion is understood by the KAM theorem. It says that nearly all initial elliptic initial conditions of \( P_0 \) with respect to \( P_1 \), give rise to quasi-periodic motion, of the two frequencies, \( \omega_1, \omega_2 \) where \( \omega_1 \) is the frequency of the elliptic motion of \( P_0 \) about \( P_1 \), provided they satisfy the condition that \( \omega_1/\omega_2 \) is sufficiently non-rational. For the relatively small set of motions of \( P_0 \) where \( \omega_1/\omega_2 \) is sufficiently close to a rational number, the motion is chaotic. It is also necessary to assume that \( \mu = m_2/(m_0 + m_2) \) is sufficiently small.

The planar modeling is assumed without loss of generality. This follows since the resonance orbits we will be considering for \( P_0 \) moving about \( P_1 \) are approximately two-body in nature. This implies approximate planar motion. These same orbits result from weak capture conditions and escape, which imply that the plane of motion of \( P_0 \) about \( P_1 \) will approximately be the same plane of motion as that of \( P_2 \) about \( P_1 \). Thus, co-planar modeling assumed in the restricted three-body problem is a reasonable assumption.

Whereas the motion of \( P_0 \) about \( P_1 \) is well understood by the KAM theorem for small \( m_2 \), the general motion of \( P_0 \) about \( P_2 \) is not well understood since it’s considerably more unstable. The instability arises due to the fact that \( m_2 \) is much smaller than \( m_1 \), and the KAM theorem cannot be easily applied unless \( P_0 \) moves infinitely close to \( P_2 \) [12]. This implies that if \( P_0 \) starts with an initial two-body elliptic state with respect to \( P_2 \), its trajectory is substantially perturbed by the gravitational effect of \( P_1 \). The resulting motion of \( P_0 \) about \( P_2 \) is unstable and generally rapidly deviates from the initial elliptic state. Numerical simulations show the motion to be chaotic in nature. Results described in this section provide a way to better understand motion about \( P_2 \).

The notion of weak capture (defined in section 2) of \( P_0 \) about \( P_2 \) is useful in trying to understand the motion of \( P_0 \) about \( P_2 \) with initial elliptic conditions. The idea is to numerically propagate trajectories of the three-body problem with initial conditions that have negative energy, \( E_2 < 0 \), with respect to \( P_2 \), and measure how they cycle about \( P_2 \) described in more detail later in this section. Generally, if \( P_0 \) performs \( k \) complete cycles about \( P_2 \), relative to a reference line emanating from \( P_2 \), without cycling about \( P_1 \), then the motion of \( P_2 \) is called ‘stable’, provided it returns to the line with \( E_2 < 0 \), while if does not return to the line after \( (k - 1) \) complete cycles, and cycles about \( P_1 \), the motion is called ‘unstable’. It is also called unstable if \( P_0 \) does return to the line, but where \( E_2 > 0 \).(see figure 4) The line represents a two-dimensional surface of section in the four-dimensional phase space, \( \mathbb{R}^4 \). The set of all stable points about \( P_2 \) defines the ‘\( k \)th stable set’, \( W_{k}^s \), and the set of all unstable points is called the ‘\( k \)th unstable set’, \( W_{k}^u \). Points that lie on the boundary between \( W_{k}^s \) and \( W_{k}^u \) define a set, \( W_{k} \), called the ‘\( k \)th weak stability boundary’. The boundary points are determined algorithmically, by iterating between stable and unstable points [13].

Points that belong to \( W_{k}^u \) are in weak capture with respect to \( P_2 \) since they start with \( E_2 < 0 \), which lead to escape with \( E_2 > 0 \) (proposition 3.1). However, this may not be the case for points in \( W_{k}^s \) since after they cycle about \( P_2 \), \( k \) times, it is possible they can remain moving about \( P_2 \) for all future time and \( E_2 \) will be negative each time \( P_0 \) intersects the line.

\( W_{k} \) was first defined in [14], for the case \( k = 1 \). This set has proved to have important applications in astrodynamics to enable spacecraft to transfer to the Moon and automatically go into weak capture about the Moon, that requires no fuel for capture. This was a substantial improvement to the Hohmann transfer, which requires substantial fuel for capture [8, 15]. It also has applications in astrophysics on the Lithopanspermia

\[4\] It was first used operationally in 1991 to rescue a Japanese lunar mission by providing a new type of transfer from the Earth to the Moon used by its spacecraft, Hiten.
Hypothesis [16]. The weak stability boundary was generalized to $k$-cycles, $k > 1$, with new details about its geometric structure in [17], [13] makes an equivalence of $W_k$ with the stable manifolds of the Lyapunov orbits associated to collinear Lagrange points.

$W_1^s$, $W_1^u$, $W_2^s$, $W_2^u$ are defined more precisely: we transform from $X = (X_1, X_2)$ defined in (3) to a rotating coordinate system, $x = (\xi, \eta)$, that rotates with frequency $\omega$, so that in this system, $P_1$, $P_2$ are fixed on the $x_1$-axis. Scaling $m_1 = 1 - \mu$, $m_2 = \mu$, $\mu > 0$, $G = 1$, $\beta = 1$, $\omega = 1$, as mentioned in section 2, we place $P_1$ at $x = (\mu, 0)$ and $P_2$ at $(-1 + \mu, 0)$. Then becomes,

\[ \dot{x} + 2(-\dot{\xi}, \dot{\eta}) = \tilde{\Omega} x, \]

\[ \tilde{\Omega} = (1/2)|x|^2 + (1 - \mu)\tau_1^2 + \mu_2 \tau_1 + (1/2)\mu(1 - \mu), \tau_1 = |x - (\mu, 0)|, \tau_2 = |x - (-1 + \mu, 0)|. \]

The Jacobi integral function, $f(x, \dot{x})$ for this system is given by

\[ J = 2\tilde{\Omega} - |x|^2. \]

The differential equations have 5 well known equilibrium points, $L_0$, $i = 1, 2, 3, 4, 5$, where $L_1$, $L_2$, $L_3$ are the collinear Lagrange points, and $L_4$, $L_5$ are equilateral points. We assume the convention that $L_2$ lies between $P_1$, $P_2$, $L_4 = C_3$, where $C_3 = C_4 = C_5 < C_2 < C_1$. The collinear points are all local saddle-center points with eigenvalues, $\pm \alpha$ and $\pm i\beta$, $\alpha > 0$, $\beta > 0$, $i^2 = -1$. The equilateral points $L_4, L_5$ are locally elliptic points.

We will focus on $L_1$, $L_2$ in our analysis. As is described in [8, 18], the distance of $L_1$, $L_2$ to $P_2$, $\tau_{1,j}$, $j = 1, 2$, is $\tau_{1,j} = O(\mu^{1/3})$, $C_4 = 3 + O(\mu^{1/3}) \gg 3$ for $\mu \geq 0$.

Projecting the three-dimensional Jacobi surface $J^{-1}(C)$ into physical $(x_1, x_2)$-space, yields the Hill’s regions, where $P_0$ is constrained to move. (see [8], figure 3.6) For $C$ slightly greater than $C_2$, $C \gtrsim C_2$, the Hill’s regions about $P_0$, $P_2$, labeled $H_0$, $H_2$, respectively, are not connected, so that $P_0$ cannot pass from one region to another. There is also a third Hill’s region, $H_3$, that surrounds both $P_0$, $P_2$ disconnected from $H_1$, $H_2$, where $P_0$ can move about both primaries. When $C = C_2$, $H_0$, $H_2$ are connected at the single Lagrange point $L_2$ and $P_0$ still cannot pass between the primaries. When $C \lesssim C_1$, a small opening occurs between $P_0$, $P_2$ near the $L_2$ location, we refer to as a neck region, $N_2$, first discussed in [19]. When $C$ decreases further, $C \gtrsim C_1$, another opening occurs near $L_1$ and forms another neck region, $N_1$, that connects $H_2$ with the outer Hill’s region, $H_3$.

A retrograde unstable hyperbolic periodic orbit is contained in $N_2$, we label $\gamma_2$. $\gamma_2$ has local stable and unstable two-dimensional manifolds $M_1^s(\gamma_2)$, $M_1^u(\gamma_2)$, $j = 1, 2$, which extend from $N_2$ into $H_2$. These manifolds are topologically equivalent to two-dimensional cylinders. It is shown in [19] that orbits can only pass from $H_2$ to $H_1$, or from $H_1$ to $H_2$, by passing within the three-dimensional region contained inside $M_1^s(\gamma_2)$, $M_1^u(\gamma_2)$, which are called transit orbits. For example, to pass from $H_2$ to $H_1$, $P_0$ must pass into the three-dimensional region inside $M_1^s(\gamma_2) \subset H_2$, and out from the region inside $M_1^u(\gamma_2) \subset H_1$. When a transit orbit passes between $H_2$ and $H_1$, then on the Jacobi surface, $P_0$ passes from $S_1^2$ to $S_2^2$. The bounding spheres separate $H_0$, $H_2$ from $N_2$.

For $C \lesssim C_1, N_1$ contains the Lyapunov orbit $\gamma_1$. Manifolds, $M_1^s(\gamma_1), M_1^u(\gamma_1), j = 2, 3$, are similarly obtained where transit orbits can pass between $H_2$ and $H_3$, passing through the respective bounding spheres. The geometry in this case is shown in figure 3.

Figure 2. Hill’s regions. $H_0$, $H_2$ connected in neck region about unstable Lyapunov orbit $\gamma_2$, $C \lesssim C_2$. Cylindrical stable and unstable manifolds, $M_1^s(\gamma_2), k = 1, 2$, in respective regions $H_0$, shown projected from four-dimensional position–velocity space into position space. The dot inside $\gamma_2$ is the location of $L_2$ when $C = C_2$. $P_0$ can only move from $H_2$ to $H_1$ through $M_1^s$ and then $M_1^u$ from an unstable point on $\mathcal{L}$. This illustrates the separatrix property of the manifolds. (This projection is not to scale and shows $M_1^s$, $M_1^u$, that exist in 4-dimensional phase space, projected into physical space. It is meant to give an idea of the geometry).
It is noted that in center of mass, rotating coordinates, \((x_1, x_2)\), (5) becomes,

\[
E_{2R} = \frac{1}{2}|x|^2 + \frac{1}{2}|x|^2 - \omega(x_2x_1 - x_1x_2) = \frac{G(m_0 + m_2)}{r_2},
\]

with \(\omega = 1, m_0 = 0, m_2 = \mu, G = 1\).

Finally, a translation, \(z_1 = x_1 - (1 + \mu), z_2 = x_2\), is made to a \(P_2\)-centered coordinate system, \((z_1, z_2)\), where \(P_2\) is at \((1, 0)\), and \(r_2 = |z|\). For notation, we set \(r \equiv r_2\). We refer to \(E_{2R}\) in center of mass rotating coordinates \((x_1, x_2)\) and also in \(P_2\)-centered rotating coordinates \((z_1, z_2)\), where \(E_{2R}\) is a different expression from (24).

The line \(L\) emanates from \(P_2\) and makes an angle \(\theta_2 \in [0, 2\pi]\) with respect to the \(z_1\)-axis. Trajectories of \(P_0\) are propagated from \(L\) such that at each point on \(L\) at a distance \(r > 0\), the eccentricity, \(e_2\), is kept fixed to a value, \(e_2 \in [0, 1)\), by adjusting the velocity magnitude, whose initial direction is perpendicular to \(L\). As the velocity direction is assumed to be clockwise (similar results are obtained for counter clockwise propogation). The initial points of propagation on \(L\) are periapsis points of an osculating ellipse of velocity \(v_p = \sqrt{G(m_0 + m_2)(1 + e_2)/r} - \omega r, \omega = 1, m_0 = 0, m_2 = \mu\). It is noted that \(L\) makes a two-dimensional surface of section, defined in polar coordinates, \(S_p = \{r_2, \theta_2, \theta_3, \theta_3, \theta_4 = \theta_2, \theta_3 > 0\}\). It is also noted that as \(r\) changes on \(L\), the Jacobi energy also changes. This implies that \(W_k^u, W_k^s, W_k^v\) do not lie on a fixed Jacobi surface. Also the Hill’s regions vary within these sets.

As is described in [13, 17], a sequence of consecutive open intervals, \(I^k_j\), are obtained along \(L\), for a fixed \(\theta_2, e_2\), that alternate between stable and unstable points, for \(k\) cycles. That is, \(I^k_0 = \{r^k_0 < r < r^k_1\}, r^k_0 = 0\), are stable points, \(I^k_1 = \{r^k_1 < r < r^k_2\}\) are unstable points, etc (see figure 4) There are \(N_k(\theta_2, e_2)\) stable sets, and unstable sets, for an integer \(N_k \geq 1\). The boundary points \(r^k_j, j = 1, 2, \ldots, N_k(\theta_2, e_2)\), represent the transition between stable and unstable points relative to \(k\) cycles, where the \(k\)th unstable points lead to stable motion for \(k - 1\) cycles and are unstable on the \(k\)th cycle. The \(kth\) stable set for a given value of \(\theta_2, e_2\) is given by,

\[
W_k^s(\theta_2, e_2) = \bigcup_{j=0}^{N_k(\theta_2, e_2)} (r^k_{2j}, r^k_{2j+1}).
\]

This is a slice of the entire stable set, \(W_k^s\), by varying \(\theta_2, e_2\), given by

\[
W_k^s = \bigcup_{\theta_2 \in [0, 2\pi], e_2 \in [0, 1)} W_k^s(\theta_2, e_2).
\]

We define \(W_k = \partial W_k^s\). \(W_k\) has a Cantor-like structure as is described in [13]. The numerical estimation of \(W_k^s, W_k\) is given in [13, 17, 20], for different values of \(k, \mu, \theta_2, e_2\). The motion of \(P_0\) is seen to be unstable and sensitive for initial conditions near \(W_k\). It is remarked that due to limitations of computer processing time, \(k\) is not taken too large.

A main result of [13], is that \(W_k\) about \(P_2\) is equivalent to the set of global stable manifolds, \(M^s(\gamma) \cup M^s(\gamma_2)\), to the Lyapunov orbits, \(\gamma_1, \gamma_2\), respectively, about the collinear Lagrange points, \(L_1, L_2, L_3\), on
either side of $P_2$ for $C \leq C_1$ and $\mu$ sufficiently small, in $H_2$. Similarly, one could restrict $C \leq C_2$ and have equivalence to only the global manifold $M^s_1(\gamma_2)$ in $H_2$. This is demonstrated numerically by examining the intersections of $M^s_2(\gamma_2)$ on surfaces of section $S_0$ satisfying $t = 0$, $E_2 < 0$ for $\mu$ sufficiently small, and varying $0 \leq \Omega_2 \leq 2\pi$. It is shown in [13] that a very small set of points exist on $W_k$ that do not satisfy this equivalence. These points are not considered.

The reason this equivalence is true is due to the separatrix property of the manifolds (see [13, 19]). Assume $C \leq C_2$. The separatrix property means that if a trajectory point is inside of the region bounded by $M^s_1(\gamma_2)$ on $L$, it will wind about $P_2$ staying inside the region contained by $M^s_1(\gamma_2)$ as $M^s_2(\gamma_2)$ winds about $P_2$. $P_0$ can’t go outside this manifold region. Eventually, $M^s_2(\gamma_2)$ will go to $\gamma_2$, and $P_0$ will pass through $N_2$ into $H_1$ as a transit orbit, after it makes $k - 1$ complete cycles, before completing the $k$th cycle. This corresponds to an unstable point on the set $W_k^u$. If a trajectory point is outside $M^s_2(\gamma_2)$ on $L$, $P_0$ will remain in $H_2$, making $k$ complete cycles about $P_2$ near the outside of $M^s_2(\gamma_2)$, but it can’t escape to $H_1$. Thus, $M^s_2(\gamma_2)$ itself is equivalent to $W_k$. That is, $M^s_2(\gamma_2)$ separates between stable and unstable motion.

The intersections of $M^s_2(\gamma_2)$ on $L$ in physical space as it cycles around $P_2$ give rise to the alternating intervals between stable and unstable motion, $[t_1^k, t_2^k, \ldots]$, where there are $N_k(\theta_2, \epsilon)$ such intervals. Points inside $M^s_2(\gamma_2)$ on $L$ correspond to points in the set $W_k^u$, and points outside of $M^s_2(\gamma_2)$ on $L$, and close to it, correspond to points of the set $W_k^s$.

The relationship between the manifolds and $W_k^s$, $W_k^u$ is shown in figure 5.

It can be shown that if $M^s_2(\gamma_2)$ has transverse intersections, which is numerically demonstrated, where the manifold tube breaks, the separatrix property is still satisfied even though a section through the tube no longer gives a circle, but rather parts of broken up circles.

Numerical simulations in [13, 18], indicate for $C \leq C_1$, that $M^s(\gamma_2)$ can intersect $M^u_1(\gamma_2)$ transversally, giving rise to a complex network of invariant manifolds about $P_2$, and for $C \leq C_2$, $M^s(\gamma_2)$ can intersect $M^u_2(\gamma_2)$ transversally, for a set of $\mu$ and $C$. This supports the fact that the motion near $W_k$ is sensitive.

Let $z(t) = (z_1(t), z_2(t), z_3(t)) \in \mathbb{R}^3$ be the trajectory of $P_0$ in rotating $P_2$-centered coordinates, and $\zeta(t) = (\zeta_1(t), \zeta_2(t)) \in \mathbb{R}^2$ the trajectory of $P_0$ in physical coordinates. Similarly, in inertial $P_2$ centered coordinates, we define, $\tilde{z}(t) = (\tilde{z}_1(t), \tilde{z}_2(t), \tilde{z}_3(t))$, and $\tilde{\zeta}(t) = (\tilde{\zeta}_1(t), \tilde{\zeta}_2(t))$. We will use inertial and rotating coordinates to describe the motion of $P_0$.

The following result, referenced previously, is proven.

**Proposition 3.1** ($P_0 \in W_k^u$ implies $P_0$ is weakly captured by $P_2$). Assume $P_0 \in W_k^u$ at time $t = t_0$, which implies $E_2(t_0) < 0$. There are two possibilities: (i.) $P_0$ cycles about $P_2$ $k - 1$ times, then moves to cycle about $P_1$ without cycling about $P_2$. This implies that $P_0$ weakly escapes $P_2$. That is, there exists a time $t^* > t_0$ after the $(k - 1)$st cycle where $E_2(t^*) = 0$, $E_2 \geq 0$ for $t \geq t^*$, and $E_2(t) < 0$ for $t_0 \leq t < t^*$, (ii.) $P_0$ does $k$ complete cycles about $P_2$, where on the $k$th cycle $P_0$ returns to $L$ with $E_2 > 0$. (It is assumed the set of collision orbits to $P_1$ and $P_2$, $I$, are excluded which are a set of measure 0.)

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5 It is noted that the proof of equivalence of $W_k$ with the global stable manifolds to $L_k$, $L_2$ in [13] is numerically supported and based on rigorous analytical estimates. Thus, the proof is rigorous in that sense. This is also true of the structure of $W_k$ obtained in [17]. A purely analytic proof for the global manifold structure about $P_2$ and $W_k$ is not available at this time. However, in the case of motion about $P_0$, the analogous structure of $W_k$ is analytically proven [21].

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Figure 4. Alternating stable and unstable intervals $I^j$ on $L$ for $k = 1$ cycles illustrating the algorithm. The boundary points $s^j$ belong to the weak stability boundary. Unstable points in the intervals leading to unstable motion (labeled un). These points are weakly captured. Stable motion, labeled s, for points in stable intervals. (Sketch, not to scale).
Proof of proposition 3.1-(i). This is shown to be true by noting that when $P_0$ does a cycle about $P_1$, it will cross the $x_1$-axis, where $r_2 > 1 + c$, $c > 0$. \[ \text{(24)} \] implies that $E_{2R} = (1/2) \dot{x}_1^2 + (1/2) \mu^2 + \dot{x}_0 (1 + c) - \mu (1 + c)^{-1}$, where $\dot{x}_0 > 0$, $\mu \ll 1$. This implies there exists a time $t^*$ where $E_{2R} > 0$. Since $E_{2R} < 0$ at $t = t_0$, then there exists a time $t_0 < t^* < t^{**}$ where $E_{2R} = 0$ and $E_2 \geq 0 \text{ for } t \geq t^*$. (ii.) This yields weak capture since $E_2 < 0$ for $t = t_0$ and $E_2$ becomes positive. Thus, there exists a time $t^*$ where $E_2(t^*) = 0$, then becomes slightly positive.

Global trajectory after weak capture

We rigorously prove result A, that after weak capture with respect to $P_2$, $P_0$ can move onto a resonance orbit about $P_1$ in resonance with $P_2$, and then return to weak capture. This is done by a series of Propositions.

The following sets are defined for trajectories for $P_0$ starting in weak capture at $t = t_0$ that go to weak escape at a time $t_1 > t_0$.

Assumptions A.

Type I = \{ $\zeta = (z, \dot{z})$ at $t_0$ is on or near $W_{E_s}^u$ (\( E_2 < 0 \), $\dot{r} = 0$ or $|\dot{r}| \gtrless 0$, resp.) $\}$

Type II = \{ $\zeta(t_0) = (z(t_0), \dot{z}(t_0))$ is not near $W_{E_s}^u$, where $|\dot{r}(t_0)|$ is not near 0 and $E_2(t_0) < 0$ $\}$

Type IIA = \{ $\zeta(t_0)$ is a Type II point where $\zeta(t)$ goes to weak escape at $t_1 > t_0$ with $|\dot{r}(t_0)|$ not near 0 (i.e. there is no cycling about $P_2$.) $\}$

Type IIB = \{ $\zeta(t_0)$ is a Type II point where there exists a time $\dot{r} > t_0$, such that $\zeta(\dot{r})$ on or near $W_{E_s}^u$, for some integer $k' \geq 1$ $\}$

Case A = \{ $P_0$ cycles about $P_1 (k-1)$ times, then moves to cycle about $P_1$ $\}$

Case B = \{ $P_0$ does not cycle about $P_1$ after $(k-1)$ st cycle. Instead, on $k$-th cycle about $P_2$, $P_0$ returns to $L$ with $E_2 > 0$ $\}$

$\Gamma = \{ P_0$ goes to collision with $P_1$ or $P_2$ for $t > t_0$ $\}$

Result A is stated more precisely as,

**Theorem A.** Assume $P_0$ is weakly captured at a distance $r$ from $P_2$ at $t = t_0$. Assume the weak capture point, $\zeta(t_0)$ is of Type I, Type IIB, Case A, which are numerically observed to be generic \([13]\), and assume the following sets are ruled out: Type IIA, Case B, Gamma (numerically observed to be small \([13]\)). Assume also that $C \lesssim C_2$ is sufficiently small.
Then $P_0$ will escape $P_2$ through $N_2$ by passing within the region contained within $M_i^k(\gamma_j) \subset H_2$, and moving into the $H_1$ through the region within $M_i^n(\gamma_j)$. This escape is approximately parabolic since $E_2 \approx 0$ on $S_1^{E2} \subset H_2$. (Parabolic escape is when there exists a time $t > t_0$ where $E_2 = 0$. $\zeta(t)$ evolves into an approximate resonance orbit about $P_2$ with an apoapsis near $S_1^{E2}$ of $N_2$, performing several cycles about $P_2$; then returns to $S_1^{E2}$ passing through $N_2$ within the region contained within $M_i^k(\gamma_j)$ and exiting $N_2$ through the interior of $M_i^n(\gamma_j)$ and moving onto weak capture about $P_2$. The process repeats unless $P_0$ moves on any of the sets: Type Iia, Case B, $\Gamma$, or $P_0$ escapes the $P_1$, $P_2$ system.

If $C \leq C_0$, then $P_0$ can parabolically escape $P_2$ through $S_1^{E2}$ of $N_2$, as previously described obtaining a sequence of resonance orbits in $H_1$, or it can parabolically escape $P_2$ through $S_1^{E2} = 0N_1$ into $H_0$ by passing through the region within $M_i^k(\gamma_j)$ and exiting from the region within $M_i^n(\gamma_j)$, and form a larger resonance orbit about $P_2$ with a periapsis near $S_1^{E2}$ in $H_0$ which eventually returns to weak capture about $P_2$ through $N_1$, reversing the previous pathway. This process terminates if $P_0$ moves on any of the sets Type Iia, Case B, $\Gamma$ or escapes the $P_1$, $P_2$ system.

This yields a sequence of approximate resonance orbits depending on the choice of the weak capture initial condition. The set of all such resonance orbits form the family, $\mathcal{F}$. The frequencies of these orbits satisfy (6) of lemma A.

**Proposition 3.2 (Capture by $P_2$ implies weak capture).** Let $P_0$ be captured with respect to $P_2$ at a distance $r$ from $P_2$ at a time $t_0$, where $E_2(Z(t_0)) = E_2(\zeta(t_0)) < 0$. Then, $P_0$ is weakly captured at $t_0$. That is, $P_0$ moves to weak escape at a time $t = t^* > t_0$, where $E_2(Z(t^*)) = 0$, $E_2(\zeta(t)) \gtrless 0$, $t \gtrless t^*$. (It is assumed Type Iia, Case B, $\Gamma$ points are excluded.)

**Proof of proposition 3.2.** We distinguish several types of weak capture points.

Type I is where $P_0$ is on or near $W_0^u$ at $t = t_0$. In this case, $P_0$ is at a distance $r$ from $P_2$ where $r = 0$ or $|r| \gtrsim 0$, where we have made use of the fact $W_0^u$ is open, so that $e_1 < 1$ and $E_2(t_0) < 0$. Thus, $P_0$ is captured at $t = t_0$.

For $t > t_0$, the proof follows by proposition 3.1.

Type II is where $P_0$ is not near $W_0^u$ since $|r|$ is not near 0 at $t = t_0$. There are two types. Type Iia is where $P_0$ starts at $t = t_0$ with $E_2$ and then to weak escape, with no cycling, by definition. If $P_0$ starts on a Type Ib point for $t = t_0$, then for $t > t_0$ there will be a time $\tilde{t} > t_0$ where $r = 0$ or $|r| \gtrsim 0$. In that case, $P_0$ is on or near $W_0^u$ at $t = \tilde{t}$, for some $k' > 1$. This yields a Type I point, that implies weak capture.

In all these cases, $P_0$ moves to weak escape at a time $t = t^* > t_0$, where $E_2(\zeta(t^*)) = 0$, $E_2(\zeta(t)) \gtrsim 0$, $t \gtrsim t^*$. This proves proposition 3.2.

As in [13], we exclude Type Iia points as they are not generic. Points on $\Gamma$ are a set of measure 0 and can be omitted. Case B points are non-generic and excluded.

We now determine what kind of motion $P_0$ has about $P_2$ for times up to weak escape at $t^*$. Consider the trajectory of $P_0$ as it undergoes counterclockwise cycling about $P_2$ after leaving points on or near $W_0^u$ on a line $L$ in both Types I, Ib. (similar results are obtained for clockwise cycling) As $P_0$ performs $k - 1$ cycles, it either has weak escape prior to completing the $k$th cycle, where $E_2 = 0$, and then when it intersects $L$, $E_2 > 0$, we call Case B, or it moves to cycle $P_1$ after the ($k - 1$)st cycle where it was shown in proposition 3.1 that $P_0$ achieves weak capture, we refer to as Case A.

**Proposition 3.3 ($P_0$ escapes from $P_2$ through the $N_0$, $N_2$ regions).** Assume $P_0 \in W_0^u$ at $t = 0$, $C \lesssim C_0$, assuming Case A, and excluding Case B. Then after $(k - 1)$-cycles about $P_2$, $P_0$ moves away from $P_2$, passing through the interior region of $M_i^k(\gamma_j)$ into $N_2$, between $H_2$ and $H_1$, through the interior of $M_i^n(\gamma_j)$, into $H_1$ where it starts to cycle $P_1$. When $P_0$ is within $N_2$, $E_2 \lesssim 0$. If $C \lesssim C_0$, then after $(k - 1)$-cycles about $P_2$, $P_0$ moves away from $P_2$, passing through the interior of $M_i^k(\gamma_j)$ into $N_2$ between $H_2$ and $H_1$ and out into $H_1$ as before, or $P_0$ passes through the interior of $M_i^n(\gamma_j)$ into $N_2$, between $H_2$ and $H_3$, and out into $H_3$ through the interior of $M_i^n(\gamma_j)$. When $P_0$ is within $N_2$, $E_2 \lesssim 0$.

**Proof of proposition 3.3.** Case A is considered with $C \lesssim C_2$. $P_0$ starts on $L$ at $t = t_0$ with $E_2 < 0$, $r = 0$. It cycles about $P_2$, completes the $(k - 1)$st cycle, then moves to $H_1$ where it starts to cycle about $P_1$, where $\theta_1 > 0$, for $0 \leq t_0 \leq 2\pi$ (see [8]). By the separatrix property $P_0$ is within the interior region contained by $M_i^k(\theta_0)$ on $L$ at $t = t_0$, and it must pass from $P_2$, through $N_2$, where it is a transit orbit [13]. When $P_0$ passes through $N_2$, it must pass inside the region bounded by $M_i^k(\gamma_j)$, and emerge from $N_2$ into the region bounded by $M_i^n(\gamma_j)$ at $S_1^{E2}$, where it will begin to cycle about $P_2$. For $\mu$ sufficiently small, the width of $N_2$ is near 0, and geometrically this implies the velocity of $P_0$ with respect to $P_2$ is near 0 since it passes close to $L_2$ in phase space.

When $P_0 \in S_1^{E2}$ in $H_1$, the distance from $P_0$ to $P_2$ can be estimated. The value of $C \lesssim C_2$, and $C_2 = 3 + (\mu/3)^{2/3} + O(\mu/3)$. $P_2$ is near $L_2$. It directly follows that $r \equiv r = (\mu/3)^{2/3} + O((\mu/3)^{2/3}))$. (This implies, $n = 1 - |O((\mu/3)^{2/3})| \lesssim 1$ since $P_0$ is slightly to the right of $L_2$ at $S_1^{E2}$.)
The estimate of $r_2$ implies that for $\mu$ sufficiently small $r_2 \approx \mu/3$. Also, at $L_2$, $\dot{x} = 0$. Thus, equation (24) implies
\begin{equation}
E_{2B} = (-3^{1/3} + (1/2)3^{-2/3})\mu^{-2/3} + O(\mu^0) \lesssim 0,
\end{equation}

$b > 2/3$.

Case A is considered with $C \lesssim C_1$. (As $r$ increases along $L$, keeping a constant eccentricity, $C$ will decrease and move slightly below the other value, $C_1$ for $L_{11}$, 180 degrees away from $L_2$ on the anti-$P_1$ side of $P_2$, $C \lesssim C_1$.)

As $t$ increases from $t_0$, by the separatrix property, $P_0$ has two possibilities: (i) $P_0$ can pass through the region bounded by $M_2^r(\gamma_1)$, through $N_2$ and exiting within the region bounded by $M_1^r(\gamma_1)$ into $H_2$ intersecting $S_{21}^r = \partial N_2$ at a time $t$. It can then start to cycle about $P_1$ in $H_3$ for $t > t_1$ where $\theta_{t_1} > 0$. This implies unstable motion occurs, where after $(k - 1)$-cycles about $P_2$, $P_2$ starts to cycle about $P_1$ in the $H_2$ region. It is similarly verified that (27) is satisfied on $S_{21}^r$ at $t = t_1$. This is different from when $P_0$ cycles about $P_1$ after $(k - 1)$-cycles as it emerges from $N_2$ into the $H_1$ region. However, in both cases, as seen, $E_2 \lesssim 0$ in the neck region bounding spheres. (ii) $P_0$ passes through $N_2$ into $H_1$. This yields the same results as in Case A. (It is verified that $C \lesssim C_1$ is sufficient to yield the same estimates in this proof for $E_{2B}$ as obtained for $C \lesssim C_2$.)

In summary, given $P_0 \in W^u_2$ at time $t = t_0$ there exists a time $t_1 > t_0$ where $E_2 \lesssim 0$ which occurs at $S_{21}^r \subset H_1$ for $C \lesssim C_2$, and for $C \lesssim C_1$, $E_2 \leq 0$ on $S_{21}^r \subset H_3$, or on $S_{21}^r \subset H_1$. Thus, in both cases, approximate parabolic escape occurs.

In the next step, we see what happens as $P_0$ starts to move about $P_1$ after leaving $S_{21}^r = \partial N_2$ in $H_1$, or $S_{21}^r = \partial N_1$ in $H_3$, through $M_2^r(\gamma_1)$, or $M_1^r(\gamma_1)$, respectively.

**Proposition 3.4** ($P_0$ leaves $S_{21}^r (S_{21}^l)$, moves in approximate resonance orbit about $P_1$, returns to $S_{21}^r (S_{21}^l)$ and then to weak capture by $P_2$). Assume $P_0 \in S_{21}^r \subset H_1$ at $t = t_1$. $P_0$ moves from $S_{21}^r$ for $t > t_1$ into an approximate resonance orbit about $P_1$. After $j$ cycles, $j \geq 1$, $P_0$ returns to $S_{21}^r$ where $E_2 \lesssim 0$. It then moves through $N_2$ to weak capture by $P_2$.

(Similarly, assuming $P_0 \in S_{21}^l \subset H_1$ at $t = t_1$, $P_0$ moves from $S_{21}^l$ for $t > t_1$ into an approximate resonance orbit about $P_1$ in the outer Hills region $H_1$. After $j$ cycles, $j \geq 1$, about $P_1$, $P_0$ returns $S_{21}^l$ where it then moves through $N_1$ to weak capture by $P_2$.)

**Proof of proposition 3.4.** The case of $P_0 \in S_{21}^r \subset H_1$ is considered first, where $t = t_1$, $C \lesssim C_2$. [21] is referenced since it determines the set $W_2$ about the larger primary $P_1$ analytically.

When $P_0 \in S_{21}^r$ for $t = t_1$, this implies it lies in the three-dimensional region bounded by $M_2^u$. Moreover, for $t > t_1$, due to the separatrix property, $P_0$ stays within this region inside $M_2^u$ for all time moving forward [21]. This manifold stays within a bounded region, $\mathfrak{M}_1$, bounded by the following: $S_{21}^r$, the boundary of $H_1$ (a zero velocity curve), and a two-dimensional McGehee torus, $T_\mathfrak{M}$, about $P_1$ [18,21]. The width of $\mathfrak{M}_1$ is $O(\mu^{1/2})$.

There are two cases. The first is where $M_2^u$ is a homoclinic two-dimensional tube which transitions from $M_2^u$ to $M_1^u$ which goes to $S_{21}^r$. This implies $P_0$ returns to $S_{21}^r$ at a later time. Now, if $M_2^u$ intersects $M_1^u$ transversely, then these manifolds intersect in a complex manner, where the image of $M_2^u$ on two-dimensional sections, $S_i^\delta$, are not circles, but parts of circles after several cycles of $P_0$ about $P_1$. However, the separatrix property is still preserved, and $P_0$ still returns to $S_{21}^r$ [21].

Let’s assume it returns to $S_{21}^r$ after a time $T$, $(t_2 = t_1 + T)$. $P_0$ is a transit orbit and must pass through $N_2$ for $t > t_2$ into $H_2$ through the interior region bounded by $M_2^u$, where it is again weakly captured by $P_2$. This follows since when $P_0$ passes through $N_2$ into $H_2$, within the interior region bounded by $M_2^u$, it will intersect $S_{21}^r \subset N_2$ in $H_2$. The estimate obtained in (27) is also obtained at $S_{21}^r$. This implies $P_0$ is captured by $P_2$ at $S_{21}^r \subset N_2$ at a time $t_1 + \delta$, $\delta > 0$. Under the previous assumptions on capture points in theorem A, $P_0$ is weakly captured and weakly escapes $P_2$.

The motion of $P_0$ as it leaves weak capture near $P_2$, passing into the $H_1$ region and moving back to the $H_2$ to weak capture is illustrated in figure 6.

It is noted that there exists a time $\delta_2 < t_2$ where $E_2 > 0$, which follows from the proof of proposition 3.1. Thus, $P_0$ is weakly captured in backwards time at $t = t_2 + \delta$.

A similar argument holds for $C \lesssim C_1$. Within $H_2$ there are openings at $N_1$ to the left of $P_2$ and $N_2$ to the right. $P_0$ can now move into $H_2$ through $N_1$, in addition to moving into $H_1$ through $N_2$, from weak capture points on $W^u_2 \cap \mathfrak{M}_1$ in $H_2$ after cycles. If $P_0$ moves into $H_1$, it does so from the region bounded by $M_2^u(\gamma_2)$ and the same argument follows from the case $C \lesssim C_2$.

If $P_0$ moves about $P_1$ in $H_1$, it moves in a bounded region $\mathfrak{M}_3$. This region is bounded by $S_{21}^l = \partial N_1$ in $H_3$, a McGehee torus $T_\mathfrak{M}$ about $P_1$ in $H_3$ and the boundary of $H_3$, $P_0$ moves inside the region enclosed by $M_1^u(\gamma_1)$ and $T_\mathfrak{M}$ exists due to the fact that KAM tori on $|f| = C$ cannot exist too close to $P_2$. 

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E Belbruno

**J. Phys. Commun.** 4 (2020) 015012
stays within it as this manifold either transitions into $M^{-1}(\gamma_0)$ as a homoclinic tube or if the manifolds have transverse intersection. The separatrix property is satisfied, and $P_0$ will cycle about $P_1$ in $H_1$ times until transits into $H_2$ and intersects $S^1_2 = \partial N_1$ in $H_2$, where (27) is satisfied which implies $P_0$ is captured by $P_2$ at a time $t_2 + \delta$. Under the assumptions of theorem A, $P_0$ is weakly captured at $P_2$.

It is noted $P_0$ is weakly captured in backwards time, since there is a time $t_2 < t_1$ where $E_2 > 0$, when $P_0$ was moving in $H_2$, using the same argument as in the proof of proposition 3.1.

The final part of the proof of proposition 3.4 is to prove that $P_0$ moves in resonance orbits about $P_1$.

We consider $C \subseteq C_2$, where $P_0$ is moving in $H_1$, $P_0$ is on an approximate resonance orbit in $H_1$ about $P_1$ for $t_1 \leq t \leq T$. This is proven as follows: $P_0$ moves in $\mathcal{W}_0$. The orbit for $P_0$ will not deviate too much for $\mu$ sufficiently small, by the amount $O(\mu^{1/3})$ [18]. It is an approximate elliptic Keplerian orbit about $P_1$, since its energy $E_1 < 0$ (proven in the following text, see proposition A). It has a uniform approximate Keplerian period, $T_1$, for $\mu$ sufficiently small, with approximate frequency $\omega_1 = T_1^{-1}$. Once $P_0$ moves away from $S^1_2$ for $t > t_1$, and returns to $S^1_2$ for $t = t_1 = T + t_1$. When $P_0$ returns to $S^1_2$ it returns to near $P_1$, to approximately the distance $O((\mu/3)^{1/3})$, as follows from the proof of proposition 3.3.

Since $P_0$ returns to $S^1_2$, near to $P_2$, $T$ must approximately be an integer multiple, $n$, of the period, $T_2$, of $P_2$ about $P_1$. That is, $T \approx nT_2$. Also, since $P_0$ returns to near where it started, $T \approx mT_2$. Thus, $nT_2 \approx mT_2$.

Since the approximate elliptic orbits of $P_0$ have an apoapsis distance from $P_1$ that is approximately the distance of $S^1_2$ to $P_1$.

This can also be visualized in inertial coordinates, $(Y_0, Y_2)$, centered at $P_1$. When $P_0$ has started its motion on a near ellipse, for $t > t_2$, it has just left weak capture from near $P_2$ at the location, $Y^*$, $P_0$ then cycles about $P_1$ and keeps returning to near $Y^*$ each approximate period $T_1$. When it arrives near $Y^*$, $P_2$ needs to be nearby as when $P_0$ started its motion. Otherwise, $P_0$ won’t become weakly captured by $P_2$ and leave the ellipse to move to the $H_2$ to weak capture by $P_2$. In that case it will continue cycling about $P_1$. If it does return to near $Y^*$ and $P_0$ has also returned near to where it started also near $Y^*$, then this means $P_2$ has gone around $P_1$ approximately $n$ times and $P_0$ has gone around $P_1$ approximately $m$ times.

In the case where $P_0$ moves in the $H_2$ region after leaving $S^1_2$ for $t > t_1$, one also obtains a resonance orbit by an analogous argument. In this case, $P_0$ has a periapsis near $S^1_2 = \partial N_1$ with respect to $P_1$, where $S^1_2 = \partial N_1$ is near $P_2$ at a distance of approximately, $O((\mu/3)^{1/3})$. These resonance orbits in $H_3$ are much larger than the resonance orbits in $H_1$, since they move about both $P_1$, $P_2$.

**Proposition A.** $E_1 < 0$ when $P_0$ moves in $H_1$, about $P_1$ in a resonance orbit.

**Proof of proposition A.** When $P_0$ moves for $t > t_1$ it moves in an approximate two-body manner for finite time spans, where the osculating eccentricity $e_{10}$ and semi-major axis $a_{10}$ vary only for a small amount amount by $O(\mu^{1/3})$ since $P_0$ moves within $\mathcal{W}_0$. The energy $E_1$ is estimated (in an inertial frame). Since at $t = t_1$, $V_1(t_1) \approx 1 - |z(t_1)|$. We can estimate $|z(t_1)|$ as roughly the distance of $L_2$ to $P_2$. This implies,
Thus, \( E_1 \approx (1/2)(1 - \alpha^{1/3} + \mathcal{O}(\alpha^{2/3}))^2 - (1 - \mu)(1 - \alpha^{1/3} + \mathcal{O}(\alpha^2/3))^{-1}. \) (28)

Thus \( E_1 \approx -(1/2) + \mathcal{O}(\alpha^{1/3}). \) This implies that \( E_1 < 0 \) for \( \mu \) sufficiently small. \( P_0 \) will then be moving on an approximate ellipse about \( P_1 \) of an eccentricity, \( e_1 < 1. \) The apoapsis of this ellipse will approximately be \( r_a \approx \eta(t_1) \approx 1 - (\mu/3)^{1/3}. \) The semi-major axis of the ellipse at \( t = t_1 \) is approximately, \( a_1 \approx -(1 - \mu)/(2E_1) \approx (1 - \mu) + \mathcal{O}((\mu/3)^{1/3}). \) \( e_1 \approx 1 - (r_a/a_1) < 1. \)

When \( P_0 \) moves in resonance orbits in \( H_3 \) for \( C \lesssim C_1, \) similar estimates are made where \( E_1 < 0. \)

End of proof of proposition A.

The resonance orbits of \( P_0 \) about \( P_1 \) move in an approximate two-body fashion where the perturbation due to \( P_2 \) is negligible for finite time spans. Thus, \( \omega_1(t) \approx (m/n)\omega \) until it enters either neck to move to weak capture near \( P_2. \)

We assume \( C \lesssim C_2 \) and examine what happens to the motion of \( P_0 \) near \( Y^* \) when in \( n: m \) resonance with \( P_2. \) \( P_0 \) is at a minimal distance to \( P_2 \) when near \( Y^*. \) In the rotating system, it is near \( S_{12}^2 \) and lies in the three-dimensional region contained within \( M_1^3 \) since it has been moving about \( P_1 \) within this region by the separatrix property. At this minimal distance, \( M_1^3 \) is close enough to \( \gamma_2, \) so that it can connect with it, and \( P_0 \) can move as a transit orbit and move through \( N_2 \) and exit into \( H_1 \) through the three-dimensional interior region bounded by \( M_1^n. \) It is then captured by \( P_2, \) with \( E_2 \lesssim 0 \) at \( S_{12}^2. \) An analogous argument holds when \( C \lesssim C_1. \)

Assuming the generic assumptions are satisfied for theorem A, \( P_0 \) will weakly escape \( P_2 \) and again move into \( H_1, \) or \( H_3, \) obtaining resonance orbits, satisfying, \( \omega_1 \approx (m'/n')\omega \) for integers \( n' \geq 1, \) \( m' \geq 1. \) The set of all such resonance orbits forms the family \( \tilde{S}. \) This concludes the proof of theorem A.

An example of the geometry of resonance transitions for an observed comet, Oterma, from a 2: 3 to 3: 2, in 1936, and then back from a 3: 2 to a 2: 3, in 1962, (§22, 23) is illustrated in figure 7. (See section 3.1 at the end of this section.)

It is noted that the estimates of \( E_0, E_2 \) in the proof of theorem 2 while \( P_0 \) moves in \( H_0, i = 1, 2, 3, \) are observed in the motions of the resonating comets studied in [22]. It can be seen in [22] that when the comet Gehrels 3 was weakly captured by Jupiter(\( P_2 \)) from a 2: 3 resonance orbit into an approximate 3: 2 resonance orbit, \( E_2 \lesssim 0. \) Also, when the comet moved about the Sun(\( P_1 \)) in an approximate 2: 3 resonance orbit, \( E_2 > 0 \) and \( E_1 < 0. \)

For each resonance orbit obtained from the choice of the weak capture initial condition, (6) is satisfied, proving lemma A.
3.1. Examples of resonance orbits in \( \mathfrak{F} \) and \( \Omega \).

Result A describes a dynamical mechanism of resonance orbits about \( P_1 \). The resonance motion described in this paper is observed both in nature and numerically.

It was originally inspired by the fact that comets are observed to perform it. More exactly, there exists a special set of comets that move about the Sun that transition between approximate resonance orbits about the Sun due to weak capture at Jupiter. This is studied in [22–24]. Several comets are described in [22, 24], that perform this motion. For example, the comet \( \text{Oterma} \) transitions between a 3:2-resonance with respect to the Sun, where \( \omega_1 = (2/3) \omega_2 \). \( \omega_1 \) is the frequency of Jupiter, to a 2:3-resonance. When passing between these resonances, the comet, \( P_0 \), is weakly captured by Jupiter. There are many others, listed in [22] (table 1), and in [24]. These comets include Helin-Roman-Crockett (3: 2 \( \rightarrow \) 3: 2), Harrington-Abell (5: 3 \( \rightarrow \) 8: 5). It is important to note that the modeling used to describe the resonance orbits of these comets is not exactly the model used in this paper. It models the true orbit of Jupiter about the Sun using the planetary ephemeris and the observed orbits of the comets for initial conditions which are not exactly planar. This model is very close to the planar restricted three-body problem. The definition of approximate resonance orbits in this paper for the restricted three-body problem is well suited to the resonances comets perform.

The existence of orbits performing resonance transitions as in \( \mathfrak{F} \) can also be found in the planar circular restricted three-body problem used in this paper. A special case where the resonance orbit precisely returns to its initial condition after preforming a transition was shown to exist in [25]. This yields an exact periodic orbit that repeats the same transition over and over. Other simulations of approximate resonance orbits as in \( \mathfrak{F} \) for the planar circular restricted three-body problem and models very close to that model are done in [26, 27].

An interesting example of orbits that occur in nature can be obtained for \( \Omega \) given in Result C. These are the subset of resonance transition orbits that have frequencies, \( \omega_1(m/n) \approx (m/n) \omega_2 \), \( m = 8, n = \bar{n}^3 \). \( \bar{n} = 1, 2, \ldots \) That is, the orbits have \( \bar{n}^3 \): 8 resonances. A special case of these resonances is an 8: 8 resonance for \( \bar{n} = 2 \). On the other hand, a 1: 1 resonance orbit is a special case of an 8: 8 resonance orbit. An example of this is for the Trojan asteroids, where \( P_1 \) is the Sun, \( P_2 \) is Jupiter, and \( P_3 \) is a Trojan asteroid. Many other examples can be found by asteroids located near the equilateral Lagrange points with respect to a body \( P_2 \), orbiting \( P_1 \).

4. Modeling resonance motions with the modified Schrödinger equation

In this section some of the results are expanded upon in section 2.

The family \( \mathfrak{F} \) of resonance periodic orbits, \( \Phi_{m/n} \), are modeled in the plane by the restricted three-body problem. The planar model is justified in section 3. To try and model \( \mathfrak{F} \) with quantum mechanical ideas, we therefore use planar modeling. Thus, we consider the planar, time independent, modified Schrödinger equation given in the Introduction by (1), obtained from the classical Schrödinger equation, (2), by replacing \( h \) by \( \sigma \), and the potential is given by the three-body potential \( \tilde{V} \) derived from the planar restricted three-body problem. This partial differential equation is time independent.

The motivation of replacing \( h \) by \( \sigma \) is given by analogy A in section 2, where \( \sigma \) given by (11). The potential \( \tilde{V} \) is given by (15), obtained from \( \tilde{V} = V_1 + V_2 \). It is recalled that \( V_1 \) is the potential due to \( P_1 \) and \( V_2 \) is the potential due to \( P_2 \) in an inertial \( P_1 \)-centered coordinate system, \( (Y_1, Y_2) \). It is also recalled that \( P_2 \) moves about \( P_1 \) on the circular orbit, \( \gamma(t), m_2 > 0 \), and as described in section 2 for Results A, it is necessary that \( m_0, m_2 \) are sufficiently small for the approximations in assumptions 1.

As described in section 2 the modified Schrödinger equation is given by (16), where \( \tilde{V} \) is the average of \( V \), obtained by averaging \( V_2 \) over a cycle of \( P_2 \) about \( P_1 \) on \( \gamma \), given by (14). For reference, we recall (16),

\[
\frac{-\sigma^2}{2\nu} \nabla^2 \Psi + \tilde{V} \Psi = E \Psi, \tag{29}
\]

In the macroscopic scale for the masses, and relative distances, two dimensions is required, where, in the inertial \( P_1 \)-centered coordinates, \( (Y_1, Y_2) \), \( \nabla^2 \equiv \frac{\partial^2}{\partial Y_1^2} + \frac{\partial^2}{\partial Y_2^2} \). When solving (29), the three-dimensional problem is solved for generality. The three-dimensional problem is discussed when considering the quantum scale.

It is noted that the units of \( \sigma \) are \( m^2 \text{kg}^{-4/3} \). This needs to match the units of \( h \) which are \( m^2 \text{kg} \text{s}^{-1} \). (\( h = 6.626 \times 10^{-34} \text{m}^2 \text{kg} \text{s}^{-1} \)) Thus, we need to multiply \( \sigma \) by \( \rho = 1^{4/3} \), \( \sigma^* \equiv \sigma \rho \). \( \sigma^* \) has the same units as \( h \). Keeping the same notation, \( \sigma^* \equiv \sigma \). At the end of this section it is seen that when the masses approach the quantum scale, \( \sigma \) also gets small, and \( m_0, m_2 \) can be adjusted so that \( \sigma = h \).

We recall (15), \( \tilde{V} = V_1 + V_2 \), where \( V_1 = -Gm_0m_2/r_1, r_1 = |Y| \) and \( V_2 = -Gm_0m_2/r_2, r_2 = |Y - \gamma(t)| \). It is shown in this section, \( \tilde{V} \) is given by (48) which has a form similar to \( V_1 \). This enables solving (16).

The solution to (29) is done in two steps. In the first step, we solve (29) in the absence of gravitational perturbations due to \( P_2 \) for \( m_2 = 0 \), where \( \tilde{V} = 0 \). Thus, in this case \( \tilde{V} = V_1 \). It is then solved for \( m_2 > 0 \) where
$V_2$ is non-zero. (Since the form of $V_2$ is shown to have a form analogous to $V_1$, we can modify the solution obtained for $V_1$ for the addition of $V_2$.)

We explicitly solve (29), with $m_2 = 0$, by separation of variables. This is for the two-body motion of $P_0$ about $P_1$. It is first solved for three-dimensions, then restricted to the planar case studied in this paper for the macroscopic scale. The three-dimensional solution is referred to when discussing solutions in the quantum scale. (29) is transformed to spherical coordinates, $r$, $\phi$, $\theta$. It is assumed the solution is of the form,

$$\Psi = R(r)Y(\phi, \theta).$$

(30)

$r \geq 0$, and $\phi$ is the angle relative to the $Y_1$-axis, $0 \leq \phi \leq 2\pi$. $\theta$ is the angle relative to the $Y_3$-axis, $0 \leq \theta \leq \pi$. $|\Psi|^2$ is the probability of finding $P_0$ at distance $r$ from $P_1$.

The solution of (29) for $V_2 = 0$ follows the method described in [28] for the case of an electron in the Hydrogen atom moving about the nucleus. This is a standard approach used in solving the classical Schrödinger equation in quantum mechanics found in many references. There are some minor modifications. The Coulomb potential is used in [28] for $V_1$ and here we are using the gravitational potential between two particles $P_0$, $P_1$; however, they are of the same form, both proportional to $r^{-1}$, $r^2 = Y_1^2 + Y_2^2 + Y_3^2$. Instead of the proportionality term of $Gm_0m_1$, the Coulomb potential has the term, $Ze^2/4\pi\varepsilon_0$, where $Z$ is the atomic number, $Z = 1$ for Hydrogen, $e$ is the charge of the electron and the charge of the nucleus, $P_0$ and an atomic nucleus, $P_1$, $\varepsilon_0$ is the permittivity of vacuum, $\hbar$ is replaced by $\sigma$. The reduced mass $\nu = m_0m_1/(m_0 + m_1)$ is defined for either the gravitational or Coulomb modeling. When referring to [28], one replaces $e^2/(4\pi\varepsilon_0)$ by $Gm_0m_1$.

When solving (29) by separation of variables, (30) is substituted into (29), obtaining differential equations for $R(r)$ and $Y(\phi, \theta)$,

$$\frac{d^2R}{dr^2} + 2r^{-1}\frac{dR}{dr} + \left[2\nu\sigma^{-2}(E + Gm_0m_1) - \alpha r^{-2}\right]R = 0,$$

(31)

$$(\sin \theta)^{-1}\frac{\partial}{\partial \theta}\left(\sin \theta \frac{\partial Y}{\partial \theta}\right) + (\sin \theta)^{-2}\frac{\partial^2 Y}{\partial \phi^2} + \alpha Y = 0.$$

(32)

$\alpha$ is a separation constant.

(32) is solved first, yielding spherical harmonics. Using separation of variables, solutions are obtained in the form, $Y(\phi, \theta) = \Phi(\phi)\Theta(\theta)$. This gives the differential equations,

$$\Theta^{-1}\sin \theta \frac{d}{d\theta}\left(\sin \theta \frac{d\Theta}{d\theta}\right) + \alpha \sin^2 \theta - \beta = 0,$$

(33)

$$\frac{d^2\Phi}{d\phi^2} + \beta \Phi = 0,$$

(34)

where $\beta$ is a separation constant [28].

(34) gives the solution,

$$\Phi(\phi) \equiv \Phi_{m}(\phi) = (1/2\pi)^{1/2}e^{im\phi},$$

(35)

$\beta = m_1^2$, $m_0 = 0$, $\pm 1, \pm 2, \ldots$, $l^2 = -1$. $\Phi(\phi)$ varies in the $Y_1$, $Y_2$-plane.

The solution of (33) follows by setting $x = \cos \theta$, $G(x) \equiv \Theta(\cos \theta)$ transforming (33) into an associated Legendre type differential equation,

$$\left(1 - x^2\right)\frac{d^2G}{dx^2} - 2x\frac{dG}{dx} + \left(\alpha - m^2(1 - x^2)^{-1}\right)G = 0,$$

(36)

$\alpha = l(l + 1)$, $l = |m|$, $|m| + 1$, $|m| + 2, \ldots$ [28], $l$ varies between $\pm |m|$. The solutions of (36) are given by associated Legendre polynomials $P_{l|m|}(x)$. (see [29] for tables of these polynomials) The solutions of (33) are given by ([28], page 527),

$$\Theta_{l,m}(\theta) = \left[\begin{array}{c} 2l + 1 \\
\end{array}\right]^{1/2}P_{l|m|}(\cos \theta).$$

(37)

It is remarked that in the two-dimensional problem, with coordinates $(Y_1, Y_2)$, $\theta = \pi/2$. In this case, there is no variation with respect to $\theta$ and $\Theta$ is only defined at $\theta = \pi/2$.

The solution $Y(\phi, \theta)$ is given by the spherical harmonics $Y_{l,m} = \Phi_{m}(\phi)\Theta_{l,m}(\theta)$. To solve (31), set $u = Rr$.

$$\frac{d^2u}{dr^2} + (\alpha r^{-1} - \beta r^{-2})u = \lambda^2u,$$

(38)

where $\lambda^2 = 2\nu|E|\sigma^{-2}$, $\alpha = 2\nu\sigma^{-2}Gm_0m_1$, $b = l(l + 1)$. It is verified that solving this differential equation yields the solution of (31),
\[ R \equiv R_{a,l}(r) = -\left[ \frac{2}{\tilde{n}} \left( \frac{n - l - 1}{2n[(n + l)!]} \right) \right] \rho^{\frac{1}{2}} L_{n+l+1}^{2l+1}(\rho) e^{-\rho/2}, \quad (39) \]

where, \( \tilde{n} = 1, 2 \ldots, l = 1, \ldots, \tilde{n} - 1, \rho = (2/(\tilde{n}a))r, a = (\nu Gm_0 m_l)^{-1}\sigma^2 \), and \( L_j \) are the associated Laguerre polynomials ([28], page 8, table 3.2), and

\[ E \equiv \hat{E}_{a} = -\frac{2\nu^2(Gm_0 m_l)^2}{\sigma^2 n^2}. \quad (40) \]

\( E \) can be written as,

\[ \hat{E}_{a} = -\frac{4\sigma}{\tilde{n}^2}. \quad (41) \]

This follows by the identity,

\[ 2\nu^2\tilde{n}^2\sigma^2 = 4, \]

where \( \tilde{p} = Gm_0 m_l \).

The solution for \( E \) yields quantized values of the energy, which quantizes the gravitational field between \( P_2 \), \( P_1 \).

The general solution to (29) is given by multiplying (39) with \( Y_{l,m_l} \),

\[ \Psi(r, \phi, \theta) \equiv \Psi_{h,m,l} = R_{a,l}(r) \Phi_{h}(\phi) \Theta_{l,m_l}(\theta). \quad (43) \]

\( \tilde{n}, l, m_l \) are the quantum numbers associated with \( \Psi \), \( \tilde{n} \) is called the principal quantum number and is independent of \( l, m_l \). It specifies the energy value and limits the value of \( l \). The quantum numbers \( l, m_l \) occur by consideration of the spherical harmonics.

We compute the probability distribution function, \( F \), of locating \( P_0 \) relative to \( P_0 \) at a given point \((r, \phi, \theta)\). Since we are currently considering the masses and the relative distances to be macroscopic, this probability is used to measure the location of a macroscopic particle, and not as a wave as is done in the quantum scale, that is considered later in this section. By definition, \( F(r, \phi, \theta) = |\Psi_{h,m,l}(r, \phi, \theta)|^2 \), depending on the quantum numbers.

It is more convenient to compute the probability at a given radial distance \( r \) independent of \( \phi, \theta \). It is labeled \( P(r) \).

It is verified that \( P(r) = R^2r^2 \). This is valid in the two-dimensional case as well for \((Y_0, Y_2)\).

As an example, calculate \( P(r) \) at the lowest energy value corresponding to \( \tilde{n} = 1, l = 0 \).

\[ P(r) = 4a^{-2}r^2 e^{-2r/a}. \quad (44) \]

(see [28], table 3.2, where \( a \) is given in the Hydrogen atom case). It is verified that \( P(r) \) has a maximum at \( r = a \), with \( P(0) = 0 \) and where \( P(r) \to 0 \) as \( r \to \infty \). It yields a curve \( \{(r, P(r)) | r \in [0, \infty]\} \) analogous to the Hydrogen atom case (see [28], figure 3.20). The numerical values of \( P(r) \) will differ from the Hydrogen atom case, since in the gravitational case \( a = (\nu Gm_0 m_l)^{-1}\sigma^2 \).

The maximum of the distribution function at \( r = a \) says that \( P_0 \), as a macroscopic body, has the highest probability of being located at this distance. In the case of the Hydrogen atom, where \( P_0 \) has wave-particle duality, this distance corresponds to the Bohr radius, which is the most probable location to find an electron in general, referred to as the 1s-orbital (1s denotes \( l = 0 \)).

In the same way, \( P(r) \) can be computed for \( \tilde{n} = 1, 2 \ldots, l = 0, 1, 2, \ldots, \tilde{n} - 1 \), which determines most probable radial locations for \( P_0 \) to be located.

If one considers computing \( F(r, \phi, \theta) = |\Psi_{h,m,l}(r, \phi, \theta)|^2 \) over \( \tilde{n}, l, m_l \), where \(-|m_l| \leq l \leq |m_l|\), then one obtains precise regions about \( P_1 \), where \( P_0 \) is most probable to be located. These are well known in the case of the Hydrogen atom [28]. It is remarkable these are observed to occur. In the gravitational case considered in this paper, the regions will have a similar geometry, but with different scaling.

The solution of (29) has been obtained in three dimensions for \( m_2 = 0 \). We solve this for \( m_2 > 0, m_2 \ll m_1 \), and then reduce to the planar case in order to compare to the planar restricted three-body problem in the macroscopic scale.

When the gravitational perturbation due to \( P_2 \) is included, the previous results are obtained with a small perturbation. It is also seen that the frequencies \( \omega_2(\tilde{n}) \) correspond to the subset, \( \Omega_1 \), of the resonant family \( \tilde{\Omega} \).

**Three-body potential**

The previous analysis can be done for a more general three-body potential by taking into account the gravitational perturbation due to \( P_2 \). We do this by using an averaged potential, \( \bar{V}_2 \), obtained from the potential \( V_2 \), due to the gravitational interaction of \( P_0, P_2 \). This yields the three-body potential, \( \bar{V} = V_1 + \bar{V}_2 \) that approximates \( \bar{V} = V_1 + V_2 \). This is done as follows,
We do this analysis in three-dimensional intertial coordinates, \( Y = (Y_1, Y_2, Y_3) \), centered at \( P_1 \), \( P_2 \) moves about \( P_1 \) on a circular orbit of radius \( \beta \) and angular frequency \( \omega \), in the \((Y_1, Y_2)\)-plane, \( \gamma(t) = \beta(\cos \omega t, \sin \omega t, 0) \), \( \beta \) is a constant. The potential for \( P_0 \) due to the perturbation of \( P_2 \) is given by

\[
\bar{V} = V_1 + V_2 = -\frac{Gm_0m_2}{r} - \frac{Gm_0m_2}{r_2},
\]

(45)

where \( r = |Y| \), \( r_2 = |Y - \gamma(t)| \). We can write \( r_2 \) as

\[
r_2 = \sqrt{r^2 + \beta^2 - 2\beta(Y_1 \cos \omega t + Y_2 \sin \omega t)},
\]

We consider \( V_2 \) and take the average of it over one cycle of \( P_2 \), where \( t \in [0, 2\pi/\omega] \),

\[
\bar{V}_2 = -\frac{Gm_0m_2}{2\pi} \int_0^{2\pi} \frac{dt}{\sqrt{r^2 + \beta^2 - 2\beta(Y_1 \cos \omega t + Y_2 \sin \omega t)}}.
\]

(46)

This averaged potential term is an approximation to \( V_2 \) felt by \( P_0 \) at a point \((Y_1, Y_2, Y_3)\) over the circular orbit of \( P_1 \) in the \((Y_1, Y_2)\)-plane. It is advantageous to use since it eliminates the time dependence in \( V_2 \), and as we’ll show, can be written so that it approximately takes the form of \( V_1 \). This implies we can solve (29) as before, with minor modifications. Approximating \( V_2 \) in this way yields \( \bar{V}_2 \).

Expressing \( Y_1, Y_2 \) in polar coordinates, \( Y_1 = r \cos \theta, Y_2 = r \sin \theta \), and making a change of the independent variable, \( t, \phi = \omega t \), we obtain

\[
\bar{V}_2 = -\frac{Gm_0m_2}{2\pi} \int_0^{2\pi} \frac{d\phi}{\sqrt{r^2 + \beta^2 - 2\beta r \cos(\phi - \theta)}}.
\]

(47)

\( \bar{V}_2 \) is simplified by considering three cases, \( r < \beta \), and \( r > \beta, r = \beta \), and expanding \( \bar{V}_2 \) as a binomial series. We prove,

**Summary A.** The general three-body potential \( \bar{V} = V_1 + V_2 \), (12), for \( P_0 \) can be approximated by replacing \( V_2 \), due to the perturbation of \( P_2 \), with the averaged potential \( \bar{V}_2 \), (47). \( \bar{V}_2 \) can be written as,

\[
\bar{V}_2 = \begin{cases} 
-\frac{Gm_0m_2}{2}\beta^{-1} + O(m_0m_2), & r > \beta \\
-\frac{Gm_0m_2}{2}\beta^{-1} + O(m_0m_2), & r < \beta \\
-\frac{Gm_0m_2}{2\sqrt{r}} + O(m_0m_2), & r = \beta.
\end{cases}
\]

(48)

The quantized energy, \( \hat{E}_\hbar \), for the approximated three-body potential \( \bar{V} = V_1 + \bar{V}_2 \) is given by,

\[
\hat{E}_\hbar = \begin{cases} 
-4\sigma\hbar^{-2}(1 + \mu + \mu^2) + O(m_0m_2), & r > \beta \\
-4\sigma\hbar^{-2} - Gm_0m_2\beta^{-1} + O(m_0m_2), & r < \beta \\
-4\sigma\hbar^{-2}(1 + \frac{1}{2}\mu + \frac{1}{2}\mu^2) + O(m_0m_2), & r = \beta
\end{cases}
\]

(49)

which reduces to (41) for \( m_2 = 0 \), and where \( \mu = m_2/m_1 \).

From the form of \( \sigma \), (49) implies,

\[
\hat{E}_\hbar = -4\sigma\hbar^{-2} + O(m_0m_2).
\]

(50)

Similarly,

**Summary B.**

\[
R_{\hbar, l}(r) = \left[ \frac{2}{\hbar} \left( \frac{\hbar - l - 1}{2\hbar} \right)! \right]^{1/2} \left[ \frac{1}{2\hbar l + 1} \right]^{1/2} \mu l^{1/4} \mu e^{-\mu r/2} + O(m_0m_2),
\]

(51)

for \( r \geq \beta, r < \beta \). The probability distribution function is generalized to,

\[
P(r) = R^2(r) r^2 + O(m_0m_2).
\]

(52)

When adding the gravitational perturbation due to \( P_2 \), represented by \( \bar{V}_2 \), one obtains smooth dependence on this term in all the calculations. This proves summary B.

**Equivalence of solutions of the modified Schrödinger equation with the family \( \tilde{\beta} \) of the three-body problem**

The two-dimensional case is now considered to compare the solutions of the modified Schrödinger equation, (29), to the family of solutions \( \tilde{\beta} \) of the planar restricted three-body problem.
We set \( Y_0 = 0, \theta = \pi/2 \). The quantized energy \( \tilde{E}_m \) \((41)\) for \((29)\) of the two-body motion of \( P_2 \) about \( P_1 \), with \( m_2 = 0 \), is not defined in the same way as the two-body Kepler energy \( \tilde{E}_1 \) \((10)\), \( \tilde{E}_m \) is computed from the modified Schrödinger equation and \( \tilde{E}_1 \) is computed for the Kepler problem for general elliptic motion of \( P_0 \) about \( P_1 \). These are different expressions. However, they both represent the energy for \( P_0 \) for the two-body gravitational potential. Also, \( \tilde{E}_m \) is quantized and \( \tilde{E}_1 \) is not quantized.

A key observation of this paper is that when one solves for the Kepler frequency \( \omega_1 \) in \((10)\) as a function of \( \tilde{E}_1 \) and substitutes \( \tilde{E}_m \) in place of \( \tilde{E}_1 \), a simple equation is obtained for \( \omega_1 \),

\[
\omega_{1|\tilde{E}_m=\tilde{E}_a} \equiv \omega_1(\tilde{n}) = 8\tilde{n}^{-3}. \tag{53}
\]

It is remarked that this equation is also valid in three-dimensions since \((10)\) is also valid for the three-dimensional Kepler two-body problem.

This follows by noting that \((10)\) implies,

\[
\omega_1 = [-\sigma^{-1}\tilde{E}_1]^{1/2}. \tag{54}
\]

This yields \((53)\).

Inclusion of gravitational perturbation of \( P_2 \) \((m_2 > 0)\), implies more generally,

\[
\omega_{1|\tilde{E}_m=\tilde{E}_a} \equiv \omega_1(\tilde{n}) = 8\tilde{n}^{-3} + \mathcal{O}(m_0 m_2). \tag{55}
\]

\((53)\) does not depend on the masses or any other physical parameter. This says substituting the quantized energy from the modified Schrödinger equation, with \( m_2 = 0 \), into the Kepler energy for the frequency, yields a frequency of motion for \( P_0 \) moving about \( P_1 \) in elliptical orbits that is same for all masses only depending on the wave number \( \tilde{n} \). This discretizes the Kepler frequencies.

The discretization of the Kepler frequencies restricts the elliptical two-body motion of \( P_0 \). This result becomes relevant in the three-body problem for \( \tilde{g} \) when the gravitational perturbation \( P_2 \) is included since it selects a the set of resonance orbits, \( \mathfrak{U} \subset \tilde{g} \) (see result C, section 2).

We prove result C from section 2, which we state as \( \omega_1(\tilde{n}) \), given by \((55)\), which are the frequencies in the restricted three-body problem corresponding to the modified Schrödinger equation energies, form a subset \( \mathfrak{U} \subset \tilde{g} \) of resonance orbits where \( m = 8, n = \tilde{n} \).

**Proof of result C.** This is proven by first noting that the circular restricted three-body problem can rescaled so that \( m_1 = 1 - \mu, m_2 = \mu, G = 1, \beta = 1 \) where \( \mu = m_2/(m_1 + m_2) \) \([8]\). This scaling does not reduce the generality of the mass values nor \( \beta \). This scaling implies \( \omega = 1 \). Thus, \((6)\) becomes,

\[
\omega_1(m/n) = (m/n) + \mathcal{O}(\delta). \tag{56}
\]

A key observation is that this scaling does not effect the leading term \( 8/\tilde{n}^3 \) of \( \omega_1(\tilde{n}) \) given by \((55)\). Thus, after the scaling, subtracting \((55)\) from \((56)\) yields,

\[
\omega_1(m/n) - \omega_1(\tilde{n}) = \frac{m}{n} - \frac{8}{\tilde{n}^3} + \mathcal{O}(m_0 m_2). \tag{57}
\]

Thus, taking \( m = 8, n = \tilde{n}^3 \) and assuming \( m_2 \) is sufficiently small, implies,

\[
\omega_1(m/n) \approx \omega_1(\tilde{n}). \tag{58}
\]

This condition is preserved by rescaling to general \( m_1, m_2, \beta \), yielding \( \omega_1(m/n) \approx (8/\tilde{n}^3)\omega \) for \( \mathfrak{U} \).

**From the macro to quantum scale**

When \( m_0, m_2, \beta \) are in the macroscopic scale then as described in section 2, in results B and C, the family \( \mathfrak{U} \subset \tilde{g} \) of near resonance orbits can be described by the solution \((17)\) of \((29)\).

For the masses in the quantum scale, the family \( \mathfrak{U} \) are no longer valid. \( \Psi \) by \((17)\) is still valid but now as pure wave solutions. This is summarized in result D, section 2. Thus, \( \Psi \) is defined for both macroscopic and quantum scales. In the macroscopic scale, \( \Psi \) is interpreted as a probability, whereas in the quantum scale, \( \Psi \) is a pure wave solution. The quantized energies \( \tilde{E}_m \) are still well defined. \( \mathcal{P} \) is still defined and yields the single domain \( \mathcal{D} \), for \( \tilde{n} = 1, 2, 3 \ldots \) This is discussed in section 2.

This section is concluded with an analysis of \( \sigma \), referred to in section 2.

**Proposition 4.1** \( \sigma = h \) is satisfied for a one-dimensional algebraic curve \((59)\) in \((m_0, m_3)\)-space. This is proven by noting \( \sigma = (1/2)(2\pi G)^{3/2}m_0 m_1(m_0 + m_3)^{-1/3} < (1/2)(2\pi G)^{3/2}m_0 m_1 m_0^{-1/3} \). Hence, \( \sigma < (1/2)(2\pi G)^{3/2}m_0^{-1/3}m_3 \). Thus, \( \sigma \to 0 \) as \( m_0, m_3 \to 0. \) This implies there exists values of \( m_0, m_3 \) such that \( \sigma = h. \) This is equivalent to the equation,

\[
m_3^2 m_1^3 - h a^{-3}(m_0 + m_1) = 0, \tag{59}
\]
\[ a = (1/2)(2\pi G)^{2/3}. \quad (59) \] yields a one-dimensional algebraic curve, \( \Gamma \), in the coordinates \( m_0, m_1 \). This proves proposition 4.1.

The relevancy of possible wave motions for \( (m_0, m_1) \) not on \( \Gamma \), with \( m_0, m_1, m_2 \) in the quantum scale, is not considered in this paper. This is discussed in section 2. Different models are also discussed in section 2.

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Appendix. Proof of summary A

Summary A is proven as follows: the integrand, \( \tilde{V}_2 \), of \( V_2 \) is given by

\[
I = \frac{1}{\sqrt{r^2 + \beta^2 - 2\beta r \cos(\phi - \theta)}} = \frac{1}{\sqrt{r^2 + \beta^2}} \frac{1}{\sqrt{1 - \frac{2\beta r}{r^2 + \beta^2} \cos(\phi - \theta)}}. \quad (60)
\]

Case 1: \( r > \beta \)

This implies,

\[
\frac{1}{\sqrt{r^2 + \beta^2}} = \frac{1}{r \sqrt{1 + \frac{\beta^2}{r}}} = \frac{1}{r}(1 + O(x)), \quad (61)
\]

where \( x = \beta/r < 1 \). This results from expanding the fraction containing the square root into a binomial series. Likewise, we can also expand the first term on the right in \( (60) \) in a binomial expansion since \( |\cos(\phi - \theta)| \leq 1 \) and

\[
\frac{2\beta r}{r^2 + \beta^2} < 1, \quad (62)
\]

yielding

\[
\frac{1}{\sqrt{1 - \frac{2\beta r}{r^2 + \beta^2} \cos(\phi - \theta)}} = 1 + O(y), \quad (63)
\]

where

\[
y = \frac{2\beta r}{r^2 + \beta^2}\cos(\phi - \theta), \quad (64)
\]

\( |y| < 1 \). Thus, \( (60) \) becomes,

\[
I = \frac{1}{r}(1 + O(w)), \quad (65)
\]

\(|w| < 1, w = \max\{x, y\}\). Thus, in this case,

\[
\tilde{V}_2 = -\frac{Gm_0m_2}{r}(1 + O(w)) = -\frac{Gm_0m_2}{r} + O(m_0m_2). \quad (66)
\]

This implies in the derivation of \( \tilde{E}_\beta \), we proceed as before and replace the numerator \( Gm_0m_1 \) of \( V_1 \) in \( (31) \) with \( Gm_0(m_1 + m_2) \) and adding \( O(m_0m_2) \) to this term. This yields,

\[
\tilde{E}_\beta = -\frac{2\nu \pi^2 (Gm_0(m_1 + m_2))^2}{a^2 \hbar^2} + O(m_0m_2). \quad (67)
\]

This can be written as,

\[
\tilde{E}_\beta = -\frac{4\sigma}{\hbar^2}(1 + \mu + \mu^2) + O(m_0m_2), \quad (68)
\]

where \( \mu = m_2/m_1 \), and where we have used \( (42) \).

Case 2: \( r < \beta \)
This case is done in a similar way as in Case 1. Instead of factoring out \( r^{-1} \) from (61), we factor out \( b-1 \), which yields,

\[
\frac{1}{\sqrt{r^2 + \beta^2}} = \frac{1}{\beta \sqrt{1 + r^2 / \beta^2}} = \frac{1}{\beta} (1 + \mathcal{O}(\tilde{\epsilon})),
\]

where \( \tilde{\epsilon} = r / \beta < 1 \). Proceeding as in Case 1, we obtain,

\[
I = \frac{1}{\beta} (1 + \mathcal{O}(\tilde{\omega}))
\]

(70)

\(|\tilde{\omega}| < 1, \tilde{\omega} = \max\{\tilde{\epsilon}, \gamma\}\). This implies,

\[
\tilde{V}_2 = -\frac{G m_0 m_2}{\beta} (1 + \mathcal{O}(\tilde{\omega})) = -\frac{G m_0 m_2}{\beta} + \mathcal{O}(m_0 m_2).
\]

(71)

Since \( \beta \) is a constant, then in the derivation of \( \tilde{E}_\alpha \), we replace \( E \) with \( E + (G m_0 m_2 / \beta) + \mathcal{O}(m_0 m_2) \) in (31) keeping \( V_1 \) as was used in the case of \( m_2 = 0 \) in (31). This yields,

\[
\tilde{E}_\alpha = -\frac{2\pi}{\sigma^2 H^2} m_0 m_2 - \frac{G m_0 m_2}{\beta} + \mathcal{O}(m_0 m_2)
\]

(72)

This can be reduced to,

\[
\tilde{E}_\alpha = -\frac{4\pi}{\tilde{n}^2} m_0 m_2 - \frac{G m_0 m_2}{\beta} + \mathcal{O}(m_0 m_2).
\]

(73)

Case 3: \( r = \beta \)

In this final case,

\[
\frac{1}{\sqrt{r^2 + \beta^2}} = \frac{1}{r \sqrt{1 + r^2 / \beta^2}} = \frac{1}{\sqrt{2} \tilde{r}}
\]

(74)

Thus,

\[
I = \frac{1}{\sqrt{r^2 + \beta^2}} \sqrt{1 - \frac{2\tilde{r}}{r^2 + \beta^2} \cos(\phi - \theta)} = \frac{1}{\sqrt{2} \tilde{r}} \sqrt{1 - h(\psi)}
\]

(75)

where \( h = \cos(\psi), \psi = \phi - \theta \). We assume that \( P_0 \) does not collide with \( P_2 \), implying \( \psi \approx 0, \pm 2\pi j, \) \( j = 1, 2, 3, \ldots \) Thus, \(|h| < 1\). We can write \( I \) as,

\[
I = \frac{1}{\sqrt{2} \tilde{r}} (1 + \mathcal{O}(h)).
\]

(76)

Hence,

\[
\tilde{V}_2 = -\frac{G m_0 m_2}{\sqrt{2} \tilde{r}} + \mathcal{O}(m_0 m_2).
\]

(77)

Proceeding as in Case 1,

\[
\tilde{E}_\alpha = -\frac{2\pi}{\sigma^2 H^2} \left( m_0 \left( m_1 + \frac{1}{\sqrt{2}} m_2 \right) \right)^2 + \mathcal{O}(m_0 m_2)
\]

(78)

This can be written as,

\[
\tilde{E}_\alpha = -\frac{4\pi}{\tilde{n}^2} \left( 1 + \frac{1}{\sqrt{2}} \mu + \frac{1}{2} \mu^2 \right) + \mathcal{O}(m_0 m_2).
\]

(79)

\[\text{References}\]

[1] Siegel C L and Moser J K 1971 Lectures on Celestial Mechanics (Grundlehren Series) (Heidelberg-Berlin: Springer)
[2] Lanczos C 1949 The Variational Principals of Physics (Toronto: University of Toronto Press)
[3] Pollard H 1976 Celestial Mechanics, The Carus Mathematical Monographs (Washington, DC: Mathematical Association of America) no. 13

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[4] Széchenyi V 1967 *Theory of Orbits* (New York: Academic)
[5] Moulton F R 1970 *An Introduction to Celestial Mechanics* (London: Macmillan) Reprinted by Dover, New York, 1970
[6] Stiefel E L and Scheifele G 1971 *Linear and Regular Celestial Mechanics* vol 174 (New York: Springer)
[7] Naidu P and Endo S 2017 *Rep. Prog. Phys.* 80 056601
[8] Belbruno E 2004 *Capture Dynamics and Chaotic Motions in Celestial Mechanics* (Princeton, NJ: Princeton University Press)
[9] Sommerfeld A 1921 *Atombau und Spektrallinien* (Braunschweig: Friedr. Vieweg & Sohn)
[10] Diósi L 1984 *Phys. Lett.* A 105 199–202
[11] Penrose R 1996 *Gen. Relativity and Gravitation* 28 581–600
[12] Kummer M 1979 *Am. J. Math.* 101 1333–54
[13] Belbruno E, Gidea M and Topputo F 2010 *SIAM J. Appl. Dyn. Sys.* 9 1061–89
[14] Belbruno E 1987 *Proceedings of the AIAA/DGLR/JSASS Inter. Elec. Propl. Conf.* no. 87–1054
[15] Belbruno E, Miller J and Guid J 1993 *Control Dyn. Astr.* 16 770–5
[16] Belbruno E, Moro-Martín A, Malhotra R and Savransky D 2012 *Astrobiology placed after 2012* 12 1–21
[17] García G and Gomez G 2007 *Cel. Mech. Dyn. Astr.* 97 87–100
[18] Llibre J, Martínez R and Simo C 1985 *J. Diff. Equ.* 58 104–56
[19] Conley C 1968 *SIAM J. Appl. Math.* 16 732–46
[20] Topputo F and Belbruno E 2009 *Cel. Mech. Dyn. Astr.* 105 3–17
[21] Belbruno E, Gidea M and Topputo F 2013 *Qual. Theory Dyn.. Sys.* 12 53–66
[22] Belbruno E and Marsden B 1997 *Astron. J.* 113 1433–44
[23] Koon W S, Lo M W, Marsden J E and Ross S D 2001 *Celest. Mech. Dyn. Astron.* 81 27–38
[24] Ohtsuka K, Ito T, Yoshikawa M, Asher D J and Arakiida H 2008 *Astron. Astrophys.* 489 1335–62
[25] Koon W S, Lo M W, Marsden J E and Ross S D 2000 *Chaos* 10 427–69
[26] Belbruno E, Topputo F and Gidea M 2008 *Adv. Space Res.* 42 1330–52
[27] Belbruno E 1997 *Annals New York Academy of Sciences (Near Earth Objects)* 822 195–225 New York: Wiley
[28] Atkins P and Friedman R 2005 *Molecular Quantum Mechanics* 4th edn (Oxford: Oxford University Press)
[29] Abramowitz M and Stegun I 1974 *Handbook of Mathematical Functions* (New York: Dover)