Analytical solution of the geodesic equation in Kerr-(anti-) de Sitter space-times

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The complete analytical solutions of the geodesic equations in Kerr-de Sitter and Kerr-anti-de Sitter space-times are presented. They are expressed in terms of Weierstrass elliptic \( \wp \), \( \zeta \), and \( \sigma \) functions as well as hyperelliptic Kleinian \( \sigma \) functions restricted to the one-dimensional \( \theta \) divisor. We analyze the dependency of timelike geodesics on the parameters of the space-time metric and the test-particle and compare the results with the situation in Kerr space-time with vanishing cosmological constant. Furthermore, we systematically can find all last stable spherical and circular orbits and derive the expressions of the deflection angle of flyby orbits, the orbital frequencies of bound orbits, the periastron shift, and the Lense-Thirring effect.

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I. INTRODUCTION AND MOTIVATION

All observations in the gravitational domain can be explained by means of Einstein’s general relativity. While for small-scale gravitational effects (e.g. in the solar system) the standard Einstein field equations are sufficient, a consistent description of large-scale observations like the accelerated expansion of the Universe can be achieved by the introduction of a cosmological term into the Einstein field equation

\[
R_{\mu\nu} - \frac{1}{2}R g_{\mu\nu} + \Lambda g_{\mu\nu} = \kappa T_{\mu\nu},
\]

(1)

where \( \Lambda \) is the cosmological constant that has a value of \( |\Lambda| \lesssim 10^{-46} \text{ km}^{-2} \) [1].

Despite the smallness of the cosmological constant the question whether there might be measurable effects on solar system scales has attracted some attention. Within approximation schemes it was shown that all effects on this scale are too small to be detectable at present [2–4]. Nevertheless, there has been some discussion on whether the Pioneer anomaly, the unexplained acceleration of the Pioneer spacecraft toward the inner solar system of \( \Delta \text{Pioneer} = (8.47 \pm 1.33) \times 10^{-10} \text{ m/s}^2 \) [5], which is of the order of \( cH \) where \( H \) is the Hubble constant, may be related to the cosmological expansion and, thus, to the cosmological constant. The same order of acceleration is related to the cosmological expansion and, thus, to the Pioneer anomaly, the unexplained acceleration of the Pioneer spacecraft toward the inner solar system scales has attracted some attention. Within MOND which again is of the order of \( 10^{-10} \text{ m/s}^2 \) [6]. Because of this mysterious coincidence of characteristic accelerations appearing at different scales and due to the fact that all these phenomena appear in a weak gravity or weak acceleration regime, it might be not clear whether the linear approximation used in [2–4] really holds. Therefore, it is desirable to obtain analytical solutions of the equations of motion for a definite answer to these questions.

There has been also some discussion if the cosmological constant has a measurable effect on the physics of binary systems, which play an important role in testing general relativity. Although such an effect would be very small, it could influence the creation of gravitational waves [7,8]. In particular, the observation of gravitational waves originating from extreme mass ratio inspirals (EMRIs) is a main goal of the Laser Interferometer Space Antenna (LISA). The calculation of such gravitational waves benefits from analytical solutions of geodesic equations not only by improved accuracy, which is, in principle, arbitrary high, but also by the prospect of developing fast semi-analytically computation methods [9]. Also, analytical solutions offer a systematic approach to determine the last stable spherical and circular orbits, which are starting points for inspirals and, thus, important for the calculation of gravitational wave templates.

Finally, for a thorough understanding of the physical properties of solutions of the gravitational field equations it is essential to study the orbits of test particles and light rays in these space-times. On the one hand, this is important from an observational point of view, since only matter and light are observed and, thus, can give insight into the physics of a given gravitational field. On the other hand, this study is also important from a fundamental point of view, since the motion of matter and light can be used to classify a given space-time, to decode its structure and to highlight its characteristics. Furthermore, analytical solutions give a possibility to systematically study limiting cases like post-Newton, post-Schwarzschild, or post-Kerr...
expansions of geodesics and observables, which is also needed for a clear interpretation of the space-time.

Analytical solutions are especially useful for the analysis of the properties of a space-time not only from an academic point of view. In fact, they offer a frame for tests of the accuracy and reliability of numerical integrations due to their, in principle, unlimited accuracy. In addition, they can be used to systematically calculate all observables in the given space-time with the very high accuracy needed for the understanding of some observations. In 1931, Hagihara [10] first analytically integrated the geodesic equation of test-particle motion in a Schwarzschild gravitational field. This solution is given in terms of the elliptic Weierstrass $\wp$ function. The geodesic equations in Reissner-Nordström, Kerr, and Kerr-Newman space-times have the same mathematical structure [11] and can be solved analogously. For bound orbits in a Kerr space-time this has been elaborated recently [12,13]. The equations of geodesic motion in space-times with nonvanishing cosmological constant exhibit a more complicated structure. Recently two of us found the complete analytical solution of the geodesic equation in Schwarzschild-(anti-) de Sitter space-times based on the inversion problem of hyperelliptic integrals [14,15]. The equations of motion could be explicitly solved by restricting the problem to the divisor, an approach which was suggested by Enolskii, Pronine, and Richter who applied this method to the problem of the double pendulum [16]. The mathematical tool developed in these papers was also applied to geodesic motion in higher-dimensional spherically symmetric and static space-times [17] as well as to Taub-NUT-de Sitter and Plebański-Demiański space-times without acceleration [18].

In this paper we extend the approach developed in [14,15] to the case of the stationary and axially symmetric Kerr-(anti)-de Sitter space-times, thus generalizing the results of both [13,15]. We start with the derivation of the equation of motion for each coordinate dependent on proper time and decouple the equations for the $r$ and $\theta$ motion following an idea of Mino [19]. Then we discuss possible types of test-particle orbits with a focus on the motion following an idea of Mino [19]. Then we discuss the equation of motion for each coordinate dependent on results of both [13,15]. We start with the derivation of Kerr-(anti)-de Sitter space-times, thus generalizing the orbits, the periastron advance and the Lense-Thirring effect of observables of particle and light trajectories. For bound orbits, the periastron advance and the Lense-Thirring effect are given in terms of the fundamental orbital frequencies.

II. THE GEODESIC EQUATION

We consider the geodesic equation

$$0 = \frac{d^2x^\mu}{d\tau^2} + \left\{ \frac{\mu}{\rho \sigma} \right\} \frac{dx^\rho}{d\tau} \frac{dx^\sigma}{d\tau},$$

(2)

where $d\tau^2 = g_{\mu \nu} dx^\mu dx^\nu$ is the proper time along the geodesics and

$$\left\{ \frac{\mu}{\rho \sigma} \right\} = \frac{1}{2} g^{\mu \nu} \left( \partial_\rho g_{\sigma \nu} + \partial_\sigma g_{\rho \nu} - \partial_\nu g_{\rho \sigma} \right),$$

(3)

the Christoffel symbol, in a space-time given by the metric

$$d\tau^2 = \frac{\Delta_r}{\chi^2 \rho^2} (dt - \sin^2 \theta d\phi)^2 - \frac{\rho^2}{\Delta_r} dr^2 - \frac{\Delta_{\theta}}{\chi^2 \rho^2} \left( adt - (r^2 + a^2) d\phi \right)^2 - \frac{\rho^2}{\Delta_\theta} d\theta^2,$$

(4)

where

$$\Delta_r = \left( 1 - \frac{\Lambda}{3} r^2 \right) (r^2 + a^2) - 2Mr,$$

(5)

$$\Delta_\theta = 1 + \frac{a^2 \Lambda}{3} \cos^2 \theta,$$

(6)

$$\chi = 1 + \frac{a^2 \Lambda}{3},$$

(7)

$$\rho^2 = r^2 + a^2 \cos^2 \theta,$$

(8)

(in units where $c = G = 1$). This Boyer-Lindquist form of the Kerr-(anti)-de Sitter metric describes an axially symmetric and stationary vacuum solution of the Einstein equation and is characterized by $\tilde{M} = 2\tilde{M}$ related to the mass $M$ of the gravitating body, the angular momentum per mass $a = J/M$, and the cosmological constant $\Lambda$. Note that this metric has coordinate singularities on the axes $\theta = 0$, $\pi$ and on the horizons $\Delta_r = 0$. The only real singularity is located at $\rho^2 = 0$, i.e. at simultaneously $r = 0$ and $\theta = \frac{\pi}{2}$ assuming $a \neq 0$.

Analogously to the situation in Kerr space-time, we classify this form of the metric according to the number of (disconnected) regions where $\Delta_r > 0$, which depends on the parameters $\tilde{M}$, $a$, and $\Lambda$. We speak of “slow” Kerr-de Sitter if there are two regions and of “fast” Kerr-de Sitter if there is one region where $\Delta_r > 0$. The limiting case where two regions are connected by a zero $\Delta_r$ is called ‘extreme’ Kerr-de Sitter. Other cases are not possible, what can be seen by a comparison of coefficients in $\Delta_r = -\frac{1}{3} \rho^2 + (1 - \frac{\Lambda}{3} a^2) r^2 - 2Mr + a^2 = -\frac{1}{3} \prod_{i=1}^n \left( r - r_i \right)$

where $r_i$ denote the zeros of $\Delta_r$. Figure 1 shows the modification of regions of slow, fast, and extreme Kerr-de Sitter with varying $\Lambda$.

We can identify four constants of motion, two corresponding to the energy per unit mass $E$ and the angular momentum per unit mass in $z$ direction $L_z$ given by the generalized momenta $p_t$ and $p_\phi$

$$p_t = g_{tt} i + g_{t\phi} \phi =: E,$$

(9)

$$-p_\phi = -g_{\phi \phi} \phi - g_{\phi t} i =: L_z,$$

(10)
where the dot denotes a derivative with respect to the proper time \( \tau \). In addition, a third constant of motion is given by the normalization condition \( \delta = g_{\mu\nu}\dot{x}^{\mu}\dot{x}^{\nu} \) with \( \delta = 1 \) for timelike and \( \delta = 0 \) for lightlike geodesics. A fourth constant of motion can be obtained in the process of separation of the Hamilton-Jacobi equation

\[
2 \frac{\partial S}{\partial \tau} = g_{ij} \frac{\partial S}{\partial x^i} \frac{\partial S}{\partial x^j}
\]

(11)

using the ansatz

\[
S = \frac{1}{2} \delta \tau - Et + L_z \phi + S_r(r) + S_\theta(\theta).
\]

(12)

If we insert this into (11) we get

\[
\delta a^2 \cos^2 \theta + \Delta_\phi \left( \frac{\partial S_\phi}{\partial \phi} \right)^2 + \frac{\chi^2}{\Delta_\phi \sin^2 \theta} (aE \sin^2 \theta - L_z)^2 = -\delta r^2 + \Delta_r \left( \frac{\partial S_r}{\partial r} \right)^2 + \frac{\chi^2}{\Delta_r} (r^2 + a^2)E - aL_z^2,
\]

(13)

where each side depends on \( r \) or \( \theta \) only. This means that each side is equal to a constant \( K \), the famous Carter constant, [20].

From the separation ansatz (12) we derive the equations of motion

\[
\rho^4 \dot{r}^2 = \chi^2 (E\tilde{M}^2\tilde{\pi})^2 - \Delta_r (\delta r^2 + K),
\]

(14)

\[
\rho^4 \dot{\theta}^2 = \Delta_\phi (K - a^2 \delta \cos^2 \theta) - \frac{\chi^2 (E\tilde{M}^2\tilde{\pi})^2}{\sin^2 \theta},
\]

(15)

\[
\frac{\rho^2}{\chi^2} \dot{\phi} = \frac{a}{\Delta_r} E\tilde{M} \tilde{\pi} - \frac{E\tilde{M}}{\Delta_\phi \sin^2 \theta} \tilde{\pi},
\]

(16)

\[
\frac{\rho^2}{\chi^2} \dot{i} = \frac{r^2 + a^2}{\Delta_r} E\tilde{M}^2 \tilde{\pi} - \frac{aE\tilde{M}}{\Delta_\phi} \tilde{\pi},
\]

(17)

where \( E\tilde{M}^2 \tilde{\pi} = (r^2 + a^2)E - aL_z \) and \( E\tilde{M} \tilde{\pi} = aE \sin^2 \theta - L_z \). The equations for \( \dot{r} \) and \( \dot{\theta} \) are coupled by \( \rho^2 = r^2 + a^2 \cos^2 \theta \). This difficulty can be overcome by introducing the Mino time \( \lambda \) [19], which is related to the proper time \( \tau \) by \( \frac{d\tau}{d\lambda} = \rho^2 \). For simplicity, we rescale the parameters appearing in Eqs. (14)–(17) such that they are dimensionless. Thus, we introduce

\[
\bar{r} = \frac{r}{\hat{M}}, \quad \bar{a} = \frac{a}{\hat{M}}, \quad \bar{\lambda} = \frac{1}{3} \lambda \hat{M}^2,
\]

(18)

and accordingly

\[
\Delta_\bar{r} = (1 - \bar{\lambda}\bar{r}^2) (\bar{r}^2 + \bar{a}^2) - \bar{r}, \quad (\Delta_\bar{r} = \bar{M}^2 \Delta_r)
\]

\[
\Delta_\bar{\phi} = 1 + \bar{a}^2 \bar{\lambda} \cos^2 \theta, \quad \chi = 1 + \bar{a}^2 \bar{\lambda}.
\]

(19)

In addition, we can absorb \((E\tilde{M})\) in the definition of \( \lambda \) by introducing

\[
D = \frac{\bar{L}_z}{E}, \quad \kappa = \frac{\bar{\Lambda}}{E}, \quad \delta_2 = \frac{\bar{\delta}}{E}, \quad \gamma = E\Lambda_\lambda.
\]

(20)

Then the Eqs. (14)–(17) decouple and read

\[
\left( \frac{d\bar{r}}{d\gamma} \right)^2 = R(\bar{r}) = \chi^2 \bar{\pi}^2 - \Delta_r (\bar{\delta} \bar{r}^2 + \kappa),
\]

(21)

\[
\left( \frac{d\bar{\theta}}{d\gamma} \right)^2 = \Theta(\theta) = \Delta_\phi (\kappa - \delta_2 \bar{a}^2 \cos^2 \theta) - \frac{\chi^2}{\sin^2 \theta} \bar{\pi}^2,
\]

(22)

\[
\frac{1}{\chi^2} \frac{d\phi}{d\gamma} = \bar{a} \bar{\pi} - \frac{1}{\Delta_\phi} \bar{\pi},
\]

(23)

\[
\frac{1}{\hat{M} \chi^2} \frac{dt}{d\gamma} = \frac{\bar{r}^2 + \bar{a}^2}{\Delta_r} \bar{\pi} - \frac{\bar{a}}{\Delta_\phi} \bar{\pi},
\]

(24)

where, as before,
\[ P = r^2 + \tilde{a}^2 - \tilde{a}\mathcal{D}, \quad (25) \]
\[ T = \tilde{a}\sin^2\theta - \mathcal{D}, \quad (26) \]

In Sec. IV, we will explicitly solve these equations.

### III. TYPES OF TIMELIKE GEODESICS

Before solving the equations of motion derived in the previous section, we analyze the structure of possible orbits dependent on the black hole parameters \( \tilde{a}, \Lambda \) and the particle parameters \( \delta, E, \mu, K \). The major point in this analysis is that \( (21) \) and \( (22) \) imply \( R(\tilde{r}) \geq 0 \) and \( \Theta(\theta) \geq 0 \) as a necessary condition for the existence of a geodesic. Although we concentrate here on timelike geodesics, light can be treated in the same manner and is in general easier to deal with.

First of all, we state two short theorems about possible values of the Carter constant. The corresponding theorems for vanishing cosmological constant can be found in [21].

**Theorem 1.** If a geodesic lies entirely in the equatorial plane \( \theta = \frac{\pi}{2} \) or if it hits the ring singularity \( \rho^2 = 0 \) then the modified Carter constant \( \tilde{Q} = K - \chi^2(\alpha E - L) \) is zero.

**Proof.** A geodesic lies entirely in the equatorial plane iff \( \theta(\gamma) = \frac{\pi}{2} \) for all \( \gamma \). This implies that \( \Theta(\theta) = (\frac{\delta}{\alpha})^2 = 0 \) and with

\[ \Theta(\theta = \frac{\pi}{2}) = \kappa - \chi^2(\tilde{a} - \mathcal{D})^2 = \tilde{Q}E^{-2}\tilde{M}^{-2} \]

it follows \( \tilde{Q} = 0 \). If a geodesic hits the ring singularity, then there is a \( \gamma \) such that \( \tilde{r}(\gamma) = 0 \) and \( \theta(\gamma) = \frac{\pi}{2} \). As \( R(\tilde{r}) \geq 0 \) and \( \Theta(\theta) \geq 0 \) for all \( \tilde{r} \) and \( \theta \), and, in particular, for \( \tilde{r} = 0, \theta = \frac{\pi}{2} \), it follows

\[ R(\tilde{r} = 0) = \chi^2\tilde{a}^2(\tilde{a} - \mathcal{D})^2 - \tilde{a}^2\kappa = \frac{-\tilde{a}^2Q}{E^2\tilde{M}^{-2}} \geq 0 \Rightarrow Q \leq 0 \]

and as above \( \Theta(\theta = \frac{\pi}{2}) = \frac{\tilde{Q}}{E\tilde{M}} \geq 0 \). \( \square \)

Note that this theorem implies that \( Q = 0 \) is a necessary condition for equatorial orbits, which is an important case found in many astrophysical objects like accretion discs and planetary systems. Note that the modified Carter constant \( Q \) depends on the cosmological constant, which also influences the next theorem.

**Theorem 2.** For \( \Lambda > -\frac{3}{\alpha^2} \) all timelike and null geodesics have \( K \geq 0 \). In this case \( K = 0 \) implies \( Q = 0 \) and the geodesic lies entirely in the equatorial plane.

**Proof.** A geodesic can only exist if there are values for \( \tilde{r}(\gamma) \) and \( \theta(\gamma) \) with \( R(\tilde{r}) \geq 0 \) and \( \Theta(\theta) \geq 0 \). From \( \Lambda > -\frac{3}{\alpha^2} \) it follows \( \Delta_\theta = 1 + \tilde{a}^2\Lambda \cos^2\theta > 1 - \cos^2\theta \geq 0 \). If \( K < 0 \) then \( (\kappa - \delta^2\tilde{a}^2\cos^2\theta) < 0 \) and

\[ \Theta(\theta) = \Delta_\theta(\kappa - \delta^2\tilde{a}^2\cos^2\theta) - \frac{\chi^2}{\sin^2\theta}T^2 < 0 \]

for all values of \( \theta \). Assume now \( K = 0 \). Consequently, \( \Theta(\theta) = \frac{-\delta^2\tilde{a}^2\cos^2\theta}{\sin^2\theta}T^2 \leq 0 \) and \( \Theta(\theta) = 0 \) only if \( \cos^2\theta = 0 \) and additionally \( T = \tilde{a}\sin^2\theta - \mathcal{D} = \tilde{a} - \mathcal{D} = 0 \). \( \square \)

Since from observation the cosmological constant has a small positive value, the condition \( \Lambda > -\frac{3}{\alpha^2} \) is always fulfilled.

From these two theorems it is obvious that, while \( K \) originates from the separation procedure, the modified Carter constant \( \tilde{Q} \) has a geometric interpretation since it is related to possible \( \theta \) values of the orbits. This relation will become more explicit in the following subsections.

In the remainder of the section we will study the consequences of the two conditions \( R(\tilde{r}) \geq 0 \) and \( \Theta(\theta) \geq 0 \).

#### A. Types of latitudinal motion

Geodesics can take an angle \( \theta \) if and only if \( \Theta(\theta) \geq 0 \). Thus, we want to determine which values of \( \tilde{a}, \Lambda, E, \mu, K \), and \( \theta \in [0, \pi] \) result in positive \( \Theta(\theta) \). For simplicity, we substitute \( \nu = \cos^2\theta \) giving

\[ \Theta(\nu) = (1 + \tilde{a}^2\Lambda \nu)(\kappa - \delta^2\tilde{a}^2\nu) - \chi^2(\tilde{a}^2(1 - \nu) - 2\tilde{a}\mathcal{D} + \frac{\mathcal{D}^2}{1 - \nu}). \quad (27) \]

Assume now that for a given set of parameters there exists a certain number of zeros of \( \Theta(\nu) \) in \([0, 1]\). If we vary the parameters, the position of zeros varies and the number of real zeros in \([0, 1]\) can change only if (i) a zero crosses 0 or 1 or (ii) two zeros merge. Let us consider case (i). Zero is a root of \( \Theta \).

\[ \Theta(\nu = 0) = \kappa - \chi^2(\tilde{a} - \mathcal{D})^2 = 0 \quad (28) \]

or

\[ \tilde{L}_\nu = \tilde{a}E \pm \frac{\sqrt{\tilde{K}}}{\chi}. \quad (29) \]

As \( \nu = 1 \) is in general a pole of \( \Theta(\nu) \) it is a necessary condition for 1 being a zero of \( \Theta(\nu) \) that this pole becomes a removable singularity. From (27) it follows that this is the case for \( \mathcal{D} = 0 \) or, equivalently, \( \tilde{L}_\nu = 0 \). Under this assumption we obtain

\[ \Theta(\nu = 1) = (1 + \tilde{a}^2\Lambda)(\kappa - \delta^2\tilde{a}^2) \quad \text{for} \quad \tilde{L}_\nu = 0. \quad (30) \]

If we additionally assume that \( \Lambda > \frac{3}{\alpha^2} \) (as in Thm. 2) we can conclude that \( \Theta(\nu = 1) = 0 \) iff \( \tilde{L}_\nu = 0 \) and \( \kappa = \delta^2\tilde{a}^2 \). Summarized, \( \tilde{L}_\nu = \tilde{a}E \pm \frac{\sqrt{\tilde{K}}}{\chi} \) and simultaneously \( \tilde{L}_\nu = 0 \) and \( \kappa = \delta^2\tilde{a}^2 \) (assuming \( \Lambda > \frac{3}{\alpha^2} \)) give us boundary cases of the \( \theta \) motion.

Now let us consider case (ii). If we exclude the coordinate singularities \( \theta = 0, \pi \) or \( \nu = 1 \) the zeros of \( \Theta(\nu) \) are given by the zeros of
\[ \Theta_r = (1 - \nu)(1 + \tilde{a}^2 \tilde{\Lambda} \nu)(\kappa - \delta_2 \tilde{a}^2 \nu) - \chi^2(\tilde{a} - D - \tilde{a} \nu)^2, \]  
(31)

which is in general a polynomial of degree 3. Then two zeros coincide at \( x \in [0, 1] \) if

\[ \Theta_r = (\nu - x)^2(a_1 \nu + a_0) \]  
(32)

for some real constants \( a_1, a_0 \). By a comparison of coefficients we can solve this equation for \( \tilde{L}_z(x) \) and \( E^2(x) \) dependent the remaining parameters \( \tilde{a}, \tilde{\Lambda}, \) and \( \tilde{K} \). This parametric representation of values of \( \tilde{L}_z \) and \( E^2 \) again corresponds to boundary cases of the \( \theta \) motion. Let us additionally consider the conditions for \( \nu = 1 \) being a double zero. With \( \tilde{L}_z = 0 \) and \( \kappa = \delta_2 \tilde{a}^2 \) (assuming \( \tilde{\Lambda} > \frac{1}{\tilde{a}^2} \)) it follows

\[ \frac{d\Theta_r}{d\nu} \bigg|_{\nu=1} = \tilde{a}^2 \chi(\chi - \delta_2), \]  
(33)

which is zero for \( \chi = \delta_2 \) or, equivalently, \( E^2 = \chi^{-1} \).

For given parameters of the black hole \( \tilde{a} \) and \( \tilde{\Lambda} \), we can use these informations to analyze the \( \theta \) motion of all possible geodesics in this space-time. As a typical example for timelike geodesics consider Fig. 2, where the curves divide the half plane into four regions (a)–(d) which correspond to different arrangement of zeros in \( [0, 1] \). A geodesic motion is only possible in regions (b) and (d) because in all other regions \( \Theta_r \) is negative for all \( \nu \in [0, 1] \). Note that for the special case of \( \kappa = \delta_2 \tilde{a}^2 \) (assuming \( \tilde{\Lambda} > \frac{1}{\tilde{a}^2} \)) strictly speaking regions (b) and (d) are divided by \( \tilde{L}_z = 0 \). However, in each region we have for \( \tilde{L}_z > 0 \) and \( \tilde{L}_z < 0 \) the same number of zeros in \( [0, 1] \) and, thus, the same type of motion. (More precisely, near \( \tilde{L}_z = 0 \) a zero \( \nu_0 < 1 \) of \( \Theta(\nu) \) approaches 1, but does not cross it.) Therefore, in each region we put the parts above and below \( \tilde{L}_z = 0 \) together. The arrangement of zeros in the two regions (b) and (d) correspond to the following different types of motion in \( \theta \) direction (cf. Figure 3)

(i) Region (b): \( \Theta_r \) has one real zero \( \nu_{\text{max}} \) in \( [0, 1] \) with \( \Theta_r \equiv 0 \) for \( \nu \in [0, \nu_{\text{max}}] \), i.e. \( \theta \) oscillates around the equatorial plane \( \theta = \frac{\pi}{2} \).

(ii) Region (d): \( \Theta_r \) has two real zeros \( \nu_{\text{min}}, \nu_{\text{max}} \) in \( [0, 1] \) with \( \Theta_r \equiv 0 \) for \( \nu \in [\nu_{\text{min}}, \nu_{\text{max}}] \), i.e. \( \theta \) oscillates between \( \arccos(\pm \sqrt{\nu_{\text{min}}} \) and \( \arccos(\pm \sqrt{\nu_{\text{max}}} \).

The boundaries of region (b) are given by \( \tilde{L}_z = \tilde{a} E \pm \sqrt{K} \) and, therefore, the regions grows larger if \( \sqrt{K} \) grows, i.e. if \( \tilde{K} \) grows or \( \tilde{a}^2 \tilde{\Lambda} \) gets smaller. A change of \( \tilde{a} \) in addition causes region (b) to shift up or down. The dependence of region (d) on the parameters \( \tilde{K}, \tilde{a} \) and \( \tilde{\Lambda} \) is much more involved. The upper boundary of (d) is also the lower boundary of region (b). The lower boundary is given in a complicated parametric form which makes it apparently impossible to determine an explicit connection between the form of region (d) and the parameters. However, the point where the upper and lower boundaries of region (d) touch each other is where 0 is a double zero of \( \Theta_r \), which is given by \( x = 0 \) in \( E(x) \) and \( \tilde{L}_z(x) \) from (32),

\[ E(0) = \frac{1}{2} \tilde{a}^2 + \frac{\tilde{K} - \tilde{a}^2 \tilde{\Lambda} \tilde{K}}{\tilde{a} \sqrt{\tilde{K} \chi}}, \quad \tilde{L}_z(0) = \frac{1}{2} \tilde{a}^2 - \tilde{K} \chi. \]  
(34)

The regions (b) and (d) are characterized in a simple way in terms of the modified Carter constant \( Q \). As in region (d) \( \Theta_r(0) < 0 \) it follows that this region corresponds to \( Q < 0 \) because of \( \Theta_r(0) = \frac{Q}{E r^2} \). In the same way we can conclude that region (b), where \( \Theta_r(0) > 0 \), corresponds to \( Q > 0 \).

**B. Types of radial motion**

A geodesic can take a radial coordinate \( \tilde{r} \) if and only if \( \tilde{R}(\tilde{r}) \equiv 0 \). The zeros of \( \tilde{R} \) are extremal values of \( \tilde{r}(\gamma) \) and determine the type of geodesic. The polynomial \( \tilde{R} \) is in general of degree six in \( \tilde{r} \) and, therefore, has six possibly complex zeros of which the real zeros are of interest for the type of motion. As a Kerr-de Sitter space-time has no

![FIG. 2 (color online). Typical regions of different types of \( \theta \)-motion. Here, \( M = 2, \tilde{\Lambda} = 10^{-5}, \) and \( \tilde{K} = 3 \). The transition from slow Kerr-de Sitter (left, \( \tilde{a} = 0.4 \)) to fast Kerr-de Sitter (right, \( \tilde{a} = 0.8 \)) is continuous.](image-url)
singularity in $\bar{r} = 0$, $\theta \neq \frac{\pi}{2}$, we can also consider negative $\bar{r}$ as valid. However, $\bar{r} = 0$ is an allowed value of $\bar{r}(\gamma)$ iff

$$0 < R(0) = \chi^2(\bar{a}^2 - \bar{a}^2 \kappa)^2 - \bar{a}^2 \frac{Q}{E^2 r_0^2} = -\bar{a}^2 \Theta_r(0).$$

(35)

It follows that $\bar{r} = 0$ can only be crossed if $Q < 0$, which corresponds to region (d) of the $\theta$ motion. In region (b) of the $\theta$ motion where $Q > 0$ a transition from positive to negative $\bar{r}$ is not possible.

To clarify the discussion we introduce some types of orbits [21].

(i) Flyby orbit: $\bar{r}$ starts from $\pm \infty$, then approaches a periapsis $\bar{r} = r_0$ and back to $\pm \infty$.

(ii) Bound orbit: $\bar{r}$ oscillates between to extremal $r_1 \leq \bar{r} \leq r_2$ with $-\infty < r_1 < r_2 < \infty$.

(iii) Transit orbit: $\bar{r}$ starts from $\pm \infty$ and goes to $\mp \infty$ crossing $\bar{r} = 0$.

All other types of orbits are exceptional and treated separately. They are either connected with the ring singularity $\rho^2 = 0$ or with the appearance of multiple zeros in $R$, which simplifies the structure of the differential equation (21) considerably. Examples for the latter type are homoclinic orbits, cf. [22]. As large negative $\bar{r}$ correspond to negative mass of the black hole [23], we will assign the attribute “crossover” to flyby or bound orbits which pass from positive to negative $\bar{r}$ or vice versa. (By definition, a transit orbit is always a crossover orbit and, therefore, we will not explicitly state that.) Therefore, in region (d) of the $\theta$ motion exists a crossover orbit, whereas all orbits located in region (b) of the $\theta$ motion do not cross $\bar{r} = 0$.

For a given set of parameters we have a certain number of real zeros of $R$. If we vary the parameters this number can change only if two zeros merge to one. This happens at $\bar{r} = x$ iff

$$R = (\bar{r} - x)^2(a_1 \bar{r}^4 + a_2 \bar{r}^3 + a_2 \bar{r}^2 + a_1 \bar{r} + a_0)$$

(36)

for some real constants $a_i$. By a comparison of coefficients we can solve the resulting 7 equations for $E^2(x)$ and $L_z(x)$ dependent on the remaining parameters $\bar{a}$, $\Lambda$, and $\bar{K}$. A typical result in slow Kerr-de Sitter for small $\Lambda$ including the results of the foregoing subsection is shown in Figs. 5 and 6. An analysis of the influence of each of the parameters $\bar{a}$, $\Lambda$, and $\bar{K}$ is not done easily due to the complexity of the expressions for $E^2$ and $L_z$. However, some typical
Examples for varying $\tilde{K}$ are shown in Fig. 7. Examples of $R(\tilde{r})$ for different numbers of real zeros are given in Fig. 4.

We discuss now the resulting types of orbits. For simplicity, we restrict ourselves here to the case of slow Kerr-de Sitter although fast and extreme Kerr-de Sitter can be discussed analogously. This will be explicitly carried through in a future publication [24].

For comparison, let us first study the situation in slow Kerr with $\Lambda = 0$

Case $\Lambda = 0$. We recognize five regions of different types of $r$ motion. (Here, we always assume $r_i < r_{i+1}$.)

(i) Region (I): all zeros of $R$ are complex and $R(\tilde{r}) \geq 0$ for all $\tilde{r}$. Possible orbit types: transit orbit.

(ii) Region (II): $R$ has two real zeros $r_1, r_2$ and $R(\tilde{r}) \geq 0$ for $\tilde{r} \leq r_1$ and $r_2 \leq \tilde{r}$. Possible orbit types: two flyby orbits, one to $r_1$, one to $r_2$.

(iii) Region (III): all four zeros $r_i$, $1 \leq i \leq 4$, of $R$ are real and $R(\tilde{r}) \geq 0$ for $r_{2k-1} \leq \tilde{r} \leq r_{2k}$, $k = 1, 2$. Possible orbit types: two bound orbits.

(iv) Region (IV): again all four zeros of $R$ are real but $R(\tilde{r}) \geq 0$ for $\tilde{r} \leq r_1$, $r_2 \leq \tilde{r} \leq r_3$, and $r_4 \leq \tilde{r}$.

Possible orbit types: two flyby orbits, one to each of $\pm \infty$ and a bound orbit.

(v) Region (V): $R$ has two real zeros $r_1, r_2$ and $R(\tilde{r}) \geq 0$ for $r_1 \leq \tilde{r} \leq r_2$. Possible orbit types: a bound orbit.

Although there are the same number of real zeros the different orbit types in region (III)/(IV) and (II)/(V) are due to the different behavior of $R$ when $\tilde{r} \to \pm \infty$. For $\tilde{\Lambda} = 0$ the expression $E^2R = \sum_{i=1}^{4} a_i \tilde{r}^i$ is a polynomial of degree 4 with $\hat{a}_1 = E^2 - 1$ which for $\tilde{r} \to \pm \infty$ yields $R(\tilde{r}) \to \infty$ if $E^2 > 1$ and $R(\tilde{r}) \to -\infty$ if $E^2 < 1$.

Let us also analyze where we have crossover orbits. Region (II) is the only one which intersects region (b) as well as region (d) of $\theta$ motion. As region (I) can only contain a transit orbit, which is by definition a crossover orbit, it can only intersect region (d). All other regions contain only region (b) of $\theta$ motion and, therefore do not have any crossover orbits. The results of this paragraph together with the numbers of positive and negative zeros for each region are summarized in Table I.

Case $\Lambda > 0$. Let us analyze now which regions change compared to the case $\Lambda = 0$. At first, we recognize that
However, the possible orbit types in regions (III) and (IV) only contain region (b) of $\theta$ motion implying that there are no crossover orbits. Region (I) can only intersect region (d) because only transit orbits are possible. The remaining region (II) is the only one which intersects regions (b) and (d).

We conclude that for $E^2 > 1$ the types of orbits are not noticeably changed, whereas for $E^2 \leq 1$ there are significant changes. In the former region (V) [for $\Lambda = 0$], which is now in region (IV), and in region (III) we have two additional flyby orbits which are not present for $\Lambda = 0$. In a small vertical stripe left of $E^2 = 1$ there are even orbits which are bound for $\Lambda = 0$ but reaching infinity for $\Lambda > 0$. In particular, it is independent of the value of $E$ if a geodesic may reach infinity as expected from the repulsive cosmological force related to $\Lambda > 0$.

Note that for huge $\Lambda$ the separation in regions (I) to (IV) is no longer possible because the repulsive cosmological force becomes so strong that bound orbits are no longer possible. In this case we have only two regions, one with two real zeros corresponding to two flyby orbits and one with only complex zeros corresponding to a transit orbit.

All orbit types for small $\Lambda > 0$ are summarized in Table II.

Case $\Lambda < 0$. Here, region (V) from $\Lambda = 0$ merges with region (I) and region (III) becomes larger for $\Lambda < 0$ due to the shift of the separating $E^2 = 1$ line towards the left. A comparison of the possible orbit types for $\Lambda > 0$ with the one for $\Lambda = 0$ shows that in regions (I) and (II) there are no differences. However, these regions are slightly deformed (for small $\Lambda$) and a pair of parameters $(E^2, L_z)$ located in region (I) or (II) for $\Lambda > 0$ may be located in a different region for $\Lambda = 0$. However, the possible orbit types in regions (III) and (IV) changed compared to $\Lambda = 0$. Here, again we assume $r_i < r_{i+1}$.

(i) Region (III): all six zeros $r_i$ of $R$ are real and $R(\bar{r}) \geq 0$ for $\bar{r} = r_1$, $r_6 \leq \bar{r}$ and $r_{2k} \leq \bar{r} \leq r_{2k+1}$ for $k = 1, 2$. Possible orbit types: two flyby orbits, one to each of $\pm \infty$, and two different bound orbits.

(ii) Region (IV): $R$ has four real zeros and $R(\bar{r}) \geq 0$ for $\bar{r} \leq r_1$, $r_2 \leq \bar{r} \leq r_3$, $r_4 \leq \bar{r}$. Possible orbit types: two flyby orbits, one to each of $\pm \infty$ and a bound orbit.

Analogously to $\Lambda = 0$, regions (III) and (IV) only contain region (b) of $\theta$ motion implying that there are no crossover orbits. Region (I) can only intersect region (d) because only transit orbits are possible. The remaining region (II) is the only one which intersects regions (b) and (d).

Concerning crossover orbits regions (III) and (IV) again only contain region (b) of $\theta$ motion. Also region (II) inter-
TABLE I. Orbit types for $\Lambda = 0$. The $+$ and $-$ columns give the number of positive and negative real zeros of the polynomial $R$. Here, the thick lines represent the range of orbits. Turning points are shown by thick dots. The small vertical line denotes $\dot{r} = 0$.

<table>
<thead>
<tr>
<th>Region</th>
<th>$-$</th>
<th>$+$</th>
<th>Range of $\dot{r}$</th>
<th>Types of orbits</th>
</tr>
</thead>
<tbody>
<tr>
<td>Id</td>
<td>0</td>
<td>0</td>
<td></td>
<td>Transit</td>
</tr>
<tr>
<td>IIb</td>
<td>1</td>
<td>1</td>
<td></td>
<td>2x flyby</td>
</tr>
<tr>
<td>IIId</td>
<td>2</td>
<td>0</td>
<td></td>
<td>Flyby, crossover flyby</td>
</tr>
<tr>
<td>IIIb</td>
<td>0</td>
<td>4</td>
<td></td>
<td>2x bound</td>
</tr>
<tr>
<td>IVb</td>
<td>1</td>
<td>3</td>
<td></td>
<td>2x flyby, bound</td>
</tr>
<tr>
<td>Vb</td>
<td>0</td>
<td>2</td>
<td></td>
<td>Bound</td>
</tr>
</tbody>
</table>

TABLE II. Orbit types for small $\Lambda > 0$. For the description of the $+$, $-$ and range of $\dot{r}$ columns see Table I.

<table>
<thead>
<tr>
<th>Region</th>
<th>$-$</th>
<th>$+$</th>
<th>Range of $\dot{r}$</th>
<th>Types of orbits</th>
</tr>
</thead>
<tbody>
<tr>
<td>Id</td>
<td>0</td>
<td>0</td>
<td></td>
<td>Transit</td>
</tr>
<tr>
<td>IIb</td>
<td>1</td>
<td>1</td>
<td></td>
<td>2x flyby</td>
</tr>
<tr>
<td>IIId</td>
<td>2</td>
<td>0</td>
<td></td>
<td>Flyby, crossover flyby</td>
</tr>
<tr>
<td>IIIb</td>
<td>1</td>
<td>5</td>
<td></td>
<td>2x flyby, 2x bound</td>
</tr>
<tr>
<td>IVb</td>
<td>1</td>
<td>3</td>
<td></td>
<td>2x flyby, bound</td>
</tr>
</tbody>
</table>

TABLE III. Orbit types for $\Lambda < 0$. For the description of the $+$, $-$ and range of $\dot{r}$ columns see Table I.

<table>
<thead>
<tr>
<th>Region</th>
<th>$-$</th>
<th>$+$</th>
<th>Range of $\dot{r}$</th>
<th>Types of orbits</th>
</tr>
</thead>
<tbody>
<tr>
<td>Id</td>
<td>1</td>
<td>1</td>
<td></td>
<td>Crossover bound</td>
</tr>
<tr>
<td>IIb</td>
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<td></td>
<td>2x bound</td>
</tr>
<tr>
<td>IIId</td>
<td>3</td>
<td>1</td>
<td></td>
<td>Bound crossover bound</td>
</tr>
<tr>
<td>IIIb</td>
<td>0</td>
<td>4</td>
<td></td>
<td>2x bound</td>
</tr>
<tr>
<td>IVb</td>
<td>2</td>
<td>4</td>
<td></td>
<td>3x bound</td>
</tr>
</tbody>
</table>
sects both (b) and (d), whereas region (I) can only contain region (d) of $\theta$ motion.

Summarizing, the types of orbits significantly change if $E^2 > 1$. The transit orbit in region (I) for $\Lambda = 0$ is transformed to a bound orbit for $\Lambda < 0$ as well as the flyby orbits in regions (II) and (IV). Although region (V) for $\Lambda = 0$ merges with region (I) for $\Lambda < 0$, the types of orbits do not change there. In general, because of $R \to -\infty$ if $\tilde{r} \to \pm \infty$ we cannot have orbits reaching $\tilde{r} = \pm \infty$ at all as expected due to the attractive cosmological force related to $\Lambda < 0$.

All orbit types for $\Lambda < 0$ are summarized in Table III.

IV. ANALYTIC SOLUTIONS OF THE EQUATIONS OF MOTION

We will now analytically solve the geodesic equation in Kerr-de Sitter space-time (21)–(24). Each equation will be treated separately.

A. $\theta$ motion

We begin with the differential equation (22)

$$\left(\frac{d\vartheta}{d\gamma}\right)^2 = \Theta = \Delta \vartheta (\kappa - \delta \tilde{a}^2 \cos^2 \vartheta) - \frac{\lambda^2 \mp 2}{\sin^2 \vartheta},$$

which can be simplified by the substitution $\nu = \cos^2 \vartheta$ yielding

$$\frac{1}{4} \left(\frac{d\nu}{d\gamma}\right)^2 = \nu \Theta_{\nu},$$

(37)

where $\Theta_{\nu}$ is defined as in (31). This differential equation can be solved easily if $\nu \Theta_{\nu}$ has a zero with multiplicity 2 or more. In this case (37) can be rewritten as

$$4(\gamma - \gamma_0) = \int_{\nu_0}^{\nu} \frac{d\nu'}{\nu' - \nu_i} \sqrt[4]{P_2(\nu')},$$

(38)

where $\gamma_0$ and $\nu_0$ are initial values, $P_2$ is a polynomial with maximum degree 2, and $\nu_i$ is a zero of $\nu \Theta_{\nu}$ with multiplicity $2j$ or $2j + 1, j = 1, 2$. The integral on the right-hand side can then be solved by elementary functions [25]. As in this case the explicit expression provides no further insight and some case distinctions would be necessary we skip the solution procedure.

If $\nu \Theta_{\nu}$ has only simple zeros the differential equation (37) is of elliptic type and first kind and can be solved in terms of the Weierstrass elliptic function $\rho$. In contrast to the $r$ motion considered in the next subsection, this structure does not simplify if we consider only light with $\delta = 0$. To obtain a solution we transform $\nu \Theta_{\nu}$ to the Weierstrass form $(4\xi^3 - g_2\xi - g_3)$ for some constants $g_2$ and $g_3$: First, we substitute $\nu = \xi^{-1}$ giving

$$\frac{1}{4} \left(\frac{d\xi}{d\gamma}\right)^2 = \Theta_{\xi},$$

(39)

where

$$\Theta_{\xi} := \xi^3 (\kappa - \chi^2(\tilde{a} - D))^2 + \xi^2 (\tilde{a}^2 (\kappa \Lambda - \delta_2) - \kappa + 2 \tilde{a}^2 (\tilde{a} - D)) + \tilde{a}^2 (\delta_2 (1 - \tilde{a}^2) - \chi^2 - \Lambda \kappa) \xi + \delta_2 \tilde{a}^4 \tilde{\Lambda}$$

(40)

Second, we substitute $\xi = \frac{1}{a_4} (4y - \frac{a_2}{y})$ yielding

$$\frac{1}{4} \left(\frac{dy}{d\gamma}\right)^2 = 4y^3 - g_2y - g_3,$$

(41)

where

$$g_2 = \frac{a_2^2}{12} - \frac{a_1 a_3}{4},$$

$$g_3 = \frac{1}{48} a_1 a_2 a_3 - \frac{1}{16} a_0 a_2^2 - \frac{1}{216} a_2^2.$$

The differential equation (41) is elliptic of first kind, which can be solved by

$$y(\gamma) = \rho(2\gamma - \gamma_{\theta, in}; g_2, g_3).$$

(43)

Accordingly, the solution of (22) is given by

$$\theta(\gamma) = \arccos \left(\pm \frac{a_3}{4\rho(2\gamma - \gamma_{\theta, in}; g_2, g_3) - \frac{a_2}{\gamma}}\right),$$

(44)

where $\gamma_{\theta, in} = 2 \gamma_0 + \int_{\gamma_0}^{\gamma} \frac{dv'}{\sqrt{4v'^2 - g_2v' - g_3}}$ with $\gamma_0 = \frac{a_0}{4\cos^2(\theta_0)} + \frac{a_1}{12}$ depends on the initial values $\gamma_0$ and $\theta_0$ only. The sign of the square root depends on whether $\theta(\gamma)$ should be in $(0, \pi)$ (positive sign) or in $(\frac{\pi}{2}, \pi)$ (negative sign) and reflects the symmetry of the $\theta$ motion with respect to the equatorial plane $\theta = \frac{\pi}{2}$. If the motion is located in region (b) from the previous section this implies that the two solutions have to be glued together along $\theta(\gamma) = \frac{\pi}{2}$ if the whole $\theta$ motion should be considered.

B. $r$ motion

The differential equation that describes the dynamics of $r$ (21)

$$\left(\frac{d\tilde{r}}{d\gamma}\right)^2 = R = \chi^2 \varrho^2 - \Delta \gamma (\delta_2 \tilde{r}^2 + \kappa)$$

is more complicated because $R$ is a polynomial of a degree up to 6. If $R$ has a zero of multiplicity 4 or more, or if $R$ has two zeros of multiplicity 2 or more, the differential equation (21) can be written as

$$\gamma - \gamma_0 = \int_{\tilde{r}_0}^{\tilde{r}} \frac{d\tilde{r}'}{\prod_{i=1}^{k} (\tilde{r}' - \tilde{r}_i) \sqrt[4]{P_2(\tilde{r}')}},$$

(45)

where $\gamma_0$ and $\tilde{r}_0$ are initial values, $P_2$ is a polynomial with maximum degree 2, $\tilde{r}_i$ are zeros of $R$ with multiplicity $2j_i$ or $2j_i + 1$ where $j_i = 1, 2$, and $k = 2$ if there are two zeros.
of multiplicity 2 or more and \( k = 1 \) else. The integral on the right-hand side can then be solved by elementary functions [25]. As the explicit expression provides no further insight and some case distinctions would be necessary we skip the solution procedure.

If we consider null geodesics, i.e. \( \delta = 0 \), \( R \) is in general of degree 4 and the differential equation (21) is of elliptic type and first kind. Then it can be handled using the method presented in the foregoing subsection: With the substitutions \( \tilde{r} = \frac{1}{b_1} + \tilde{r}_R \), where \( \tilde{r}_R \) is a zero of \( R \), and 
\[
\tilde{\xi} = \frac{1}{b_1}(4y - \frac{b_2}{y}),
\]
where \( b_1 = \frac{1}{4\beta}(b_2^2 - \beta y^2) \), we arrive at a form (41). This can then again be solved in terms of Weierstrass elliptic functions. The result is
\[
\tilde{r}(\gamma) = \frac{b_3}{4\beta(y - \gamma_{r,in})} + \tilde{r}_R,
\]
where \( \gamma_{r,in} = y_0 + \int_{y_0}^{y} \frac{d\gamma'}{R_0(r - r_1)^{\frac{1}{2}}} \), with \( y_0 = \frac{b_1}{4\beta_0 - \beta a} + \frac{b_2}{4\beta_0} \) depends only on the initial values \( y_0 \) and \( \tilde{r}_0 \) and \( g_2 \), \( g_3 \) are defined as in (42) with \( a_i = b_i \).

The differential equation (21) is also of elliptic type but of third kind if \( R \) has a double or triple zero \( \tilde{r}_1 \). In this case (21) reads
\[
\gamma - y_0 = \int_{r_0}^{r} \frac{d\tilde{r}'}{P_3(\tilde{r})},
\]
where \( P_3 \) is a polynomial of degree 4. This equation can be solved for \( \tilde{r}(\gamma) \) analogous to the method which will be presented in Sec. IV C.

If we consider particles, i.e. \( \delta = 1 \), and assume that \( R \) has only simple zeros the differential equation (21) is of hyperelliptic type. It can be solved in terms of derivatives of the Kleinian \( \sigma \) function with the method developed in [14]. For this, we have to cast (21) in the standard form by a substitution \( \tilde{r} = \pm \frac{1}{a} + \tilde{r}_R \) with a zero \( \tilde{r}_R \) of \( R \). This yields
\[
\left( \frac{u}{d\gamma} \right)^2 = c_5 R_u,
\]
where
\[
R_u = \sum_{i=0}^{5} \frac{c_i}{c_5} u^i, \quad c_i = \frac{(\pm 1)^i}{(6 - i)!} d^{(6-i)} \langle 0 \rangle (\tilde{r}_R),
\]
The sign in the substitution has been chosen such that the constant \( c_5 \) is positive and, therefore, depends on the choice of \( \tilde{r}_R \) and the sign of \( \tilde{A} \). The differential equation (48) is of first kind and can be solved by
\[
\gamma - y_0 = \int_{r_0}^{r} \frac{d\tilde{r}'}{\sqrt{c_5} - \gamma_{r,in}},
\]
where \( \gamma_{r,in} = \sqrt{c_5} y_0 + \int_{y_0}^{y} \frac{d\gamma'}{\sqrt{c_5} - \gamma_{r,in}} \) and \( u_0 = \pm (\tilde{r}_0 - \tilde{r}_R)^{-1} \) depends only on the initial values \( y_0 \) and \( \tilde{r}_0 \). Here, \( f \) is the function that describes the \( \theta \) divisor, i.e. \( \sigma_1(f(x),y)^2 = 0 \), cf. [14]. The radial distance \( \tilde{r} \) is then given by
\[
\tilde{r}(\gamma) = \frac{\sigma_2}{\sigma_1} \left( \frac{f(\sqrt{c_5} \gamma - \gamma_{r,in})}{\sqrt{c_5} \gamma - \gamma_{r,in}} \right) + \tilde{r}_R.
\]
where the sign depends on the sign chosen in the substitution \( \tilde{r} = \pm \frac{1}{a} + \tilde{r}_R \), i.e. is such that \( c_5 \) in (49) is positive.

C. \( \phi \) motion

We treat now the most complicated equation of motion in Kerr-de Sitter space-time, namely, the equation for the azimuthal angle (23)
\[
\frac{d\phi}{d\gamma} = \chi^2 \left[ \frac{a}{\Delta x} P^\parallel - \frac{1}{\Delta x \sin^2 \theta} \right].
\]
This equation can be split in a part dependent only on \( \tilde{r} \) and in a part only dependent on \( \theta \). Integration yields
\[
\phi - \phi_0 = \chi^2 \left[ \int_{\gamma_0}^{\gamma} \frac{a}{\Delta x} \frac{d\gamma}{\Delta x \sin^2 \theta(\gamma)} - \int_{\theta_0}^{\theta} \frac{1}{\Delta x \sin^2 \theta(\Theta)} d\theta' \right],
\]
where we substituted \( \tilde{r} = \tilde{r}(\gamma) \), i.e. \( \frac{d\tilde{r}}{d\gamma} = \sqrt{\tilde{R}} \), in the first and \( \theta = \theta(\gamma) \), i.e. \( \frac{d\theta}{d\gamma} = \frac{\tilde{R}}{\sqrt{\tilde{R}}} \), in the second integral.

We will solve now the two integrals in (52) separately.

1. The \( \theta \)-dependent integral

Let us consider the integral
\[
I_\theta := \int_{\theta_0}^{\theta} \frac{a - D - a \nu'}{\Delta x \nu^2 \nu' \tilde{R}} d\nu',
\]
which can be transformed to the simpler form
\[
I_\theta = \frac{1}{2} \int_{\gamma_0}^{\gamma} \frac{\tilde{r} - \tilde{D} - \tilde{a} \nu'}{\Delta x \nu^2 \nu' \tilde{R}} d\nu',
\]
by the substitution \( \nu = \cos^2 \theta \), where \( \Theta_{\nu} \) is defined in (31) and \( \Delta_x = 1 + \tilde{a} \tilde{R} \tilde{\nu} \). Here, we have to pay special attention to the integration path. If \( \theta \in (0, \frac{\pi}{2}) \) we have \( \cos \theta = +\sqrt{\tilde{\nu}} \) but for \( \theta \in \left[ \frac{\pi}{2}, \pi \right] \) it is \( \cos \theta = -\sqrt{\tilde{\nu}} \). Accordingly, we first have to split the integration path from \( \theta_0 \) to \( \theta \) such that every piece is fully contained in the interval \( (0, \frac{\pi}{2}) \) or \( \left[ \frac{\pi}{2}, \pi \right] \) and then to choose the appropriate sign of the square root of \( \nu \). In the following we assume for simplicity that \( \cos \theta = +\sqrt{\tilde{\nu}} \).

Analogous to Sec. IV A the integral \( I_\theta \) can be solved by elementary functions if \( \nu \Theta_{\nu} \) has at least a double zero [25]. If \( \nu \Theta_{\nu} \) has only simple zeros \( I_\theta \) is of elliptic type and of third kind. If this is the case, the solution to \( I_\theta \) is given by
\[ I_\theta = \frac{|a_3|}{2a_3} \left\{ (\bar{a} - D)(v - v_0) - \sum_{i=1}^{4} \frac{a_3}{4\chi \rho^i(v_i)} \left\{ \xi(v_i)(v - v_0) + \log \frac{\sigma(v - v_i)}{\sigma(v_0 - v_i)} + 2\pi ik_i \right\} \cdot (\bar{a}^2 \Lambda (\chi - \bar{a} \bar{\Lambda} D) \right. \]
\[ \times (\delta_{i1} + \delta_{i2}) + \mathcal{D}(\delta_{i3} + \delta_{i4}) \right\}, \tag{55} \]

where the constants \( a_i \) are defined as in Sec. IV A, \( \rho(v_i) = a_i - \frac{1}{2} \bar{a}^2 \Lambda \), \( \rho(v_0) = a_i - \frac{1}{2} \bar{a}^2 \Lambda_0 \), \( \rho(v) = a_i - \frac{1}{2} \bar{a}^2 \Lambda \) with \( \gamma_{r,\text{in}} \) as in (43) and \( v_0 = v(\gamma_0) \). The integers \( k_i \) correspond to different branches of log. The details of the computation can be found in Appendix A.

2. The \( \bar{r} \)-dependent integral

We solve now the first, \( \bar{r} \)-dependent integral in (52)

\[ I_\bar{r} := \int_{\bar{r}_0}^{\bar{r}} \frac{\bar{a}(\bar{r}^2 + \bar{a}^2 - \bar{a} \bar{D})d\bar{r}}{\Delta \bar{r} \sqrt{R}}. \tag{56} \]

Analogous to Sec. IV B this integral can be solved by elementary functions if \( R \) has a zero with multiplicity 4 or more or two zeros with multiplicity 2 or more [25].

If we consider light, i.e. \( \delta = 0 \), \( R \) is in general of degree 4 and \( I_\bar{r} \) is of elliptic type and third kind. In this case it can be solved analogously to \( I_\theta \). The same substitutions \( \bar{r} = \frac{1}{\bar{r}} + \bar{r}_R \) and \( \bar{\xi} = \frac{1}{\bar{b}_1}(4\bar{y}_r - \bar{a}_r \bar{y}) \) as in Sec. IV B for the case \( \delta = 0 \), a subsequent partial fraction decomposition, and the final substitution \( y = \rho(v) \) result in

\[ \frac{b_3}{|b_3|} I_\bar{r} = \sum_{i=1}^{4} C_i \int_{v_0}^{v} \frac{dv}{\rho(v) - y_i} \]
\[ - \frac{\bar{a}(\bar{r}_R^2 + \bar{a}^2 - \bar{a} \bar{D})}{\Delta \bar{r} \bar{r}_R} \int_{v_0}^{v} d\nu, \tag{57} \]

where \( y_i \) are the four zeros of \( \Delta(y,\bar{r}) \), \( b_3 \) defined as in (46), and \( C_i \) are the coefficients of the partial fractions dependent on the parameters and \( \bar{r}_R \). The four functions \( f_i(v) = (\rho(v) - y_i)^{-1} \) have simple poles in \( v_{i1}, v_{i2} \) with \( \rho(v_{i1}) = y_i = \rho(v_{i2}) \) and have to be integrated with the method presented in Appendix A. Then \( I_\bar{r} \) is given by

\[ \frac{b_3}{|b_3|} I_\bar{r} = \sum_{i=1}^{4} \sum_{j=1}^{2} C_i (\nu_{ij})(v - v_0) + \log \sigma(v - v_i) \]
\[ - \frac{\bar{a}(\bar{r}_R^2 + \bar{a}^2 - \bar{a} \bar{D})(v - v_0)}{\Delta \bar{r} \bar{r}_R}, \tag{58} \]

where \( v = v(\gamma) = \gamma - \gamma_{r,\text{in}}, v_0 = v(\gamma_0) \) with \( \gamma_{r,\text{in}} \) as in (46). In the same way \( I_\bar{r} \) can be solved if \( R \) has a double or triple zero.

If we consider particles, i.e. \( \delta = 1 \), and assume that \( R \) has only simple zeros, \( I_\bar{r} \) is of hyperelliptic type and third kind. The details of the solution method can be found in Appendix B but we give an outline here: First, we transform analogously to Sec. IV B to the standard form by \( \bar{r} = \pm 1/u + \bar{r}_R \) with a zero \( \bar{r}_R \) of \( R \). Afterward we simplify the integrand by a partial fraction decomposition which allows us to express \( I_\bar{r} \) in terms of the canonical holomorphic differentials \( d\bar{\xi} (B10) \) and the canonical differential of third kind \( dP(x_1, x_2) (B7) \). These differentials can then be expressed in dependence of the normalized Mino time \( \gamma \).

We define \( w = w(\gamma) = \sqrt{c_5} \gamma - \gamma_{r,\text{in}} \) and \( w_0 = w(\gamma_0) \) the result is

\[ I_\bar{r} = -\frac{\bar{a}u_0}{\sqrt{c_5} u_0} \left\{ C_1 (w - w_0) + C_0 (f(w) - f(w_0)) \right. \]
\[ + \sum_{i=1}^{4} C_{2,i} \left[ \frac{1}{2} \log \frac{\sigma(W^+(w))}{\rho(W^-(w))} \right. \]
\[ \left. - (f(w) - f(w_0), w - w_0) \left( \int_{u_i}^{u_i} \frac{d\bar{r}}{\rho} \right) \right] \tag{59} \]

where the constants \( C_i \) are the coefficients of the partial fractions, \( u_i \) are the four zeros of \( \Delta_{\bar{r},\pm 1/u + \bar{r}_R}, u_0 = \pm (\bar{r} - \bar{r}_R)^{-1}, \) and \( c_5 \) and \( R_u \) as in (48). The functions \( W^\pm \) are defined by \( W^\pm (w) = (f(w), w)^t - 2 \int_{\infty}^{a_i} d\bar{\xi} \), where the points \( u_i = u_i, \pm \sqrt{R_u} \) on the Riemann surface of \( \gamma^2 = \bar{R}_u \) are the pole \( u_i \) located on the positive and negative branch of the square root.

D. \( t \) motion

The equation for \( t \) (24)

\[ \frac{dt}{d\gamma} = \chi^2 M \left[ \frac{\bar{r}^2 + \bar{a}^2}{\bar{r}^2} \right] \bar{\rho} - \frac{\bar{a}}{\Delta \bar{r} \bar{r}_R} \tag{57} \]

is of the same structure as the equation for \( \phi \) motion. An integration yields

\[ t - t_0 = \chi^2 M \left[ \int_{\gamma_0}^{\gamma} \frac{\bar{r}^2 + \bar{a}^2 \bar{\rho}}{\Delta \bar{r} \bar{r}_R} \right. \]
\[ \left. \int_{\gamma_0}^{\gamma} \bar{\rho} d\gamma \right] = \chi^2 M \left[ \int_{\bar{r}_0}^{\bar{r}} \frac{\bar{r}^2 + \bar{a}^2 \bar{\rho}}{\Delta \bar{r} \bar{r}_R} d\bar{r} - \int_{\theta_0}^{\theta} \frac{\bar{a}}{\Delta \bar{r} \bar{r}_R} d\theta \right] \]
\[ = \chi^2 M [I_\bar{r} - I_\theta]. \tag{58} \]

Because we already demonstrated the solution procedure, we only give here the results for the most general cases.

If \( \nu \Theta_t \) in (35) has only simple zeros the solution of the \( \theta \)-dependent part is given by

\[ \tilde{I}_\theta = a_3 (v - v_0) - \sum_{i=1}^{2} a_3 \bar{a}^2 \bar{\Lambda} \]
\[ + \log \sigma(v - v_i) - \log \sigma(v_0 - v_i) \tag{58} \]

where \( a_3 \) is defined as in (40), \( \rho(v_i) = a_i - \frac{1}{2} \bar{a}^2 \Lambda \), \( \rho(v_2) = 2\gamma - \gamma_{r,\text{in}}, v_0 = v(\gamma_0) \) as in (55).

If we consider light, i.e. \( \delta = 0 \), the solution for the \( \bar{r} \)-dependent part is given by
where \( b_1 \) is defined as in (46), \( \hat{C}_i \) are the coefficients of the partial fractions, \( \varphi(v_{ij}) = y_j = \varphi(v_{ij}) \) with the four zeros \( y_i \) of \( \Delta_{ij} = 0 \) and \( v = v(\gamma) = \gamma - \gamma_{\text{in}}, v_0 = v(\gamma_0) \) as in (46).

If \( R \) has only simple zeros and we consider timelike geodesics \( \delta = 1 \) the solution of the \( \hat{r} \)-dependent part is given by

\[
\frac{b_3}{b_1} \hat{I}_r = -\frac{u_0}{|u_0| \sqrt{c_5}} \left\{ \hat{C}_1 (w - w_0) + \hat{C}_0 (f(w) - f(w_0)) + \frac{4}{\sqrt{R_{u_j}}} \left[ \frac{1}{2} \log \frac{\sigma(W^+(w))}{\sigma(W^-(w))} - \frac{1}{2} \log \frac{\sigma(W^+(w_0))}{\sigma(W^-(w_0))} \right] \right\} - (f(w) - f(w_0), w - w_0) \left( \int_{u_j} \frac{d \hat{r}}{u_j} \right) \right\},
\]

where the notation is as in (59) and \( \hat{C}_0, \hat{C}_1, \hat{C}_{2,i} \) are the coefficients of the partial fractions.

**V. DISCUSSION OF SOME GEODESICS**

In Sec. III, we discussed the general features of the different types of timelike geodesic motion in Kerr-de Sitter and Kerr-anti-de Sitter space-time. With the analytical solution derived in Sec. IV at hand we want to discuss now some chosen geodesics.

We start with orbits which highlight the influence of \( \Lambda \) on the geodesics. From the results of Sec. III we conclude that for \( \Lambda > 0 \) there are four parameter regions where the changes compared to \( \Lambda = 0 \) are most obvious. The first two are the regions (III) and (IV) with \( E^2 < 1 \), where we have additional flyby orbits not present for \( \Lambda = 0 \). Third and fourth, the shift from region (V) of \( \Lambda = 0 \) to region (II) of \( \Lambda > 0 \) for \( E^2 = 1 - \epsilon, \epsilon > 0 \) small, and the shift from region (III) of \( \Lambda = 0 \) to region (IV) of \( \Lambda > 0 \), again for \( E^2 = 1 - \epsilon \) are most interesting as the (outer) bound orbit becomes a flyby orbit. A plot of the corresponding orbits can be found in Fig. 8.

An important feature of geodesics in stationary axisymmetric space-times is motion of the nodes where the orbit of a test particle or light intercepts the equatorial plane. This motion is caused by the \( g_{00} \) components of the space-time metric and known as the Lense-Thirring effect. In the weak field regime it becomes visible by a precession of the orbital plane, cf. Fig. 9 for an obvious example. This orbital precession has been confirmed within an accuracy of about 10% by the LAGEOS (Laser Geodynamics Satellite) mission [26,27].

Let us also discuss some exceptional orbits related to multiple zeros of \( R \), i.e. spherical orbits with constant \( r \) and orbits asymptotically approaching a constant \( r \). There are two types of spherical orbits: stable and unstable. Stable spherical orbits with \( \hat{r}(\gamma) = \hat{r}_0 \) occur if radial coordinates adjacent to \( \hat{r}_0 \) are not allowed due to \( R(\hat{r}) < 0 \), which happens if \( \hat{r}_0 \) is a maximum of \( R \). Unstable spherical orbits with \( \hat{r}(\gamma) = \hat{r}_0 \) are trajectories where radial coordinates \( \hat{r} \) in the neighborhood of \( \hat{r}_0 \), with \( \hat{r} < \hat{r}_0 \) or \( \hat{r} > \hat{r}_0 \) are allowed. Therefore, these orbits are related to a minimum or to an inflection point of \( R \). If \( \hat{r}_0 \) is an inflection point, an asymptotic approach to \( \hat{r}_0 \) is only possible from one side of \( \hat{r}_0 \) whereas this is possible from both sides if \( \hat{r}_0 \) is a minimum of \( R \). Asymptotic orbits can also be divided into two types: unbound and bound. The latter case corresponds to orbits which approach for both \( t \to \infty \) and \( t \to -\infty \) a spherical orbit. Bound asymptotic orbits are also known as homoclinic orbits. If the asymptotic orbit is unbound it reaches \( \hat{r} = \infty \) for either \( t \to \infty \) or \( t \to -\infty \).

FIG. 8 (color online). Flyby orbits for \( \Lambda = 10^{-5}, \hat{a} = 0.4 \), and \( \hat{K} = 3 \). For \( \Lambda = 0 \) the parameters of (c) would have been located in region (Vb) and the parameters of (d) in region (IIIb). Light grey cones correspond to extremal \( \theta \) and dark grey spheres to horizons.

(a) \( E^2 = 0.9, \bar{L}_z = -1 \) (region (IVb))
(b) \( E^2 = 0.94, \bar{L}_z = 0.6 \) (region (IIIb))
(c) \( E^2 = 0.97, \bar{L}_z = -1 \) (region (IIb))
(d) \( E^2 = 0.97, \bar{L}_z = 1.8 \) (region (IVb))
For asymptotic bound or unbound trajectories corresponding to an unstable spherical orbit the equations of motion simplify considerably. In this case the equation for $\dot{r}(\gamma)$ as well as the $\dot{r}$-dependent integrals in the $\phi$ and $t$ equations are of elliptic type and can be solved in terms of Weierstrass elliptic functions, see (47) and (58). Note that these solutions are not limited to the case of equatorial circular orbits but are valid for all types of asymptotic orbits and, thus, generalize the analytical solutions for homoclinic orbits in [22] not only to Kerr-de Sitter space-time but also to arbitrary inclinations.

From all spherical orbits the Last Stable Spherical Orbit (LSSO), in particular, the Innermost Stable Circular Orbit (ISCO) in the equatorial plane are of importance as they represent the transition from stable orbits to those which fall through the event horizon. The corresponding multiple zero of $R$ appears at the boundaries of the different regions of $r$ motion, cf. Fig. 5. Because necessarily $\tilde{K} = \chi^2 E^2 (\tilde{a} - D)^2$ for equatorial orbits, from this we can determine the LSSO for given $\tilde{a}$, $\tilde{\Lambda}$, $\tilde{K}$ and the ISCO for given $\tilde{a}$, $\tilde{\Lambda}$ by solving first

$$ R(\tilde{r}) = 0, \quad \frac{dR}{d\tilde{r}}(\tilde{r}) = 0, \quad \text{and} \quad \frac{d^2R}{d\tilde{r}^2}(\tilde{r}) = 0 \quad (64) $$

for $\tilde{r} \geq \tilde{r}_h$ with the event horizon $\tilde{r}_h$. The solutions are limiting cases of the LSSO or ISCO and are given by the corner points on the borders of region (III) of the $\tilde{r}$ motion (as corners on other boundaries correspond to $\tilde{r} < \tilde{r}_h$). From the results of (64) we search for the smallest possible double zero $\tilde{r}$ which is a maximum. In the case of the ISCO in the equatorial plane we are now done. For the LSSO, we have to check in addition whether the corresponding values of $E^2(\tilde{r})$ and $L_z(\tilde{r})$ (given by (36)) are located in an allowed region of the $\theta$ motion. If this is the case, we found the LSSO. If not, we can determine the LSSO as the intersection point of the boundary of region (III) with a boundary of an allowed $\theta$ region. Note that it is not possible to determine an LSSO (for given $\tilde{a}$, $\tilde{\Lambda}$, and $\tilde{K}$) if there is no spherical orbit at all outside the event horizon which happens if no boundary of the $\tilde{r}$ motion is located in an allowed region of the $\theta$ motion. As an example, this is the case for $\tilde{\Lambda} = 10^{-5}$, $\tilde{a} = 0.1$, and $\tilde{K} = 0.1$. Also, the LSSO is identical with the ISCO if it is given as an intersection point with the boundary of region (b) of the $\theta$ motion. For examples of spherical orbits see Figs. 10 and 11. Note that within the event horizon there may be additional stable spherical orbits.

VI. ANALYTIC EXPRESSIONS FOR OBSERVABLES

For the understanding of characteristic features of space-times by measurements of geodesics in that space-time it is crucial to identify certain theoretical quantities of observables. For flyby orbits, this can be the deflection angle of the geodesic whereas for bound orbits it is of interest to

\[ \frac{\partial}{\partial \tilde{r}} \approx \frac{\partial}{\partial \tilde{\theta}} \approx \frac{\partial}{\partial \tilde{\phi}} \]

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\[ \frac{\partial}{\partial \tilde{r}} \approx \frac{\partial}{\partial \tilde{\theta}} \approx \frac{\partial}{\partial \tilde{\phi}} \]
The orbital frequency of the motion with respect to $\gamma$ is then given by \( \frac{2\pi}{\omega_r} \). For the calculation of \( \Omega_r \), which represents the orbital frequency with respect to $t$, we need in addition the average rate $\Gamma$ at which $r$ accumulates with $\gamma$. This will be determined below. For the calculation of the orbital frequency $\Omega_\theta$ we have to determine the orbital period $\omega_\theta$ such that $\theta(\gamma + \omega_\theta) = \theta(\gamma)$. The $\theta$ motion is likewise bounded by $\theta_{\text{min}} \leq \theta \leq \theta_{\text{max}}$ for two real zeros $\theta_{\text{min}}, \theta_{\text{max}} \in (0, \pi)$ of $\theta$ and, therefore, \[
\omega_\theta = \frac{2\pi}{\theta_{\text{max}}}, \]
Again, the orbital frequency of the $\theta$ motion with respect to $\gamma$ is given by $\frac{2\pi}{\omega_\theta}$.

The orbital periods of the remaining coordinates $t$ and $\phi$ have to be treated somewhat differently because they depend on both $r$ and $\theta$. The solutions $\tilde{r}(\gamma)$ and $\tilde{\phi}(\gamma)$ consist of two different parts, one which represents the average rates $\Gamma$ and $Y_\phi$ at which $r$ and $\phi$ accumulate with $\gamma$ and one which represents oscillations around it with periods $\omega_r$ and $\omega_\theta [13, 29]$. The periods $\Gamma$ and $Y_\phi$ can be calculated by [13] 

\[
\Gamma = \frac{2}{\omega_r} \int_{r_p}^{r_a} \frac{\sqrt{\nu}}{\Delta_r \sqrt{R}} \frac{d\tilde{r}}{\sqrt{\tilde{r}}} - \frac{2}{\omega_\theta} \int_{\theta_{\text{min}}}^{\theta_{\text{max}}} \frac{d\tilde{\theta}}{\Delta_\theta \sqrt{\Theta}}, \quad (68)
\]

\[
Y_\phi = \frac{2}{\omega_r} \int_{r_p}^{r_a} \frac{\sqrt{\nu}}{\Delta_r \sqrt{R}} \frac{\tilde{\phi} d\tilde{r}}{\sqrt{\tilde{r}}} - \frac{2}{\omega_\theta} \int_{\theta_{\text{min}}}^{\theta_{\text{max}}} \frac{d\tilde{\theta}}{\Delta_\theta \sqrt{\Theta}}, \quad (69)
\]

The orbital frequencies $\Omega_r$, $\Omega_\theta$, and $\Omega_\phi$ are then given by 

\[
\Omega_r = \frac{2\pi}{\omega_r} \Gamma, \quad \Omega_\theta = \frac{2\pi}{\omega_\theta} \Gamma, \quad \Omega_\phi = \frac{Y_\phi}{\Gamma}. \quad (70)
\]

If the integral expressions for $\omega_r$, and the $\tilde{r}$-dependent parts of $\Gamma$ and $Y_\phi$ degenerate to elliptic or elementary type, i.e. if we consider light or $R$ possesses multiple zeros, we can find analytical expressions for (66), (68), and (69) with the techniques presented in [13]. If $R$ has only simple zeros and $\delta = 1$, the integral $\omega_r$ is an entry of the fundamental period matrix $\omega$ which enters in the definition of the period lattice $\Gamma$ of the holomorphic differentials $d\zeta$, $\Gamma = \{\omega v + \omega' v' \mid v, v' \in \mathbb{Z}^2\}$. The more complicated integrals involving $R$ in (68) and (69) can be rewritten in terms of periods of the differentials of second kind $d\tilde{\phi}$ and of third kind $dP(x_1, x_2)$ by a decomposition in partial fractions. In this way, expressions for $\omega_r$, $\Gamma$ and $Y_\phi$, which are totally analogous to the elliptic case can be obtained.

It follows that the periastron shift is given by 

\[
\Delta_{\text{periastron}} = \Omega_\phi - \Omega_r = \left(\frac{Y_\phi}{\omega_r} - \frac{2\pi}{\omega_r} \right)^2 \Gamma, \quad (71)
\]

and the Lense-Thirring effect by 

\[
\Delta_{\text{Lense-Thirring}} = \Omega_\phi - \Omega_\theta = \left(\frac{Y_\phi}{\omega_\theta} - \frac{2\pi}{\omega_\theta} \right)^2 \Gamma. \quad (72)
\]

Another way to access information encoded in the orbits is through a frequency decomposition of the whole orbit [29]. This will be analyzed elsewhere.

\section*{VII. SUMMARY AND OUTLOOK}

In this paper we derived the analytical solution for both timelike and lightlike geodesic motion in Kerr-(anti-) de Sitter space-time. The analytical expressions for the orbits $(r, \theta, \phi, t)$ are given by elliptic Weierstrass and hyperelliptic Kleinian functions in terms of the normalized Mino time. We also presented a method for solving differential equations of hyperelliptic type and third kind, and applied it to the equations of motion for $\phi$ and $t$. We classified possible types of geodesic motion by an analysis of the zeros of the polynomials underlying the $\theta$ and $\tilde{r}$ motion and discussed the influence of a nonvanishing cosmological constant on the orbit types. Some particular interesting orbits not present in Kerr space-time were shown and a systematic approach for determining the last stable spherical and circular orbits was presented. In addition, we derived the analytic expressions for observables connected with geodesic motion in Kerr-de Sitter space-time, namely, the deflection angle for escape orbits as well.

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as the orbital frequencies, the periastron shift, and Lense-Thirring effect of bound orbits.

The results of this paper can be viewed as the starting point for the analysis of several features of geodesics in Kerr-de Sitter space-time not treated in this publication. Although mathematical analogues to the case of slow Kerr-de Sitter a complete discussion of orbits in fast and extreme Kerr-de Sitter space-times may lead to special features and should be carried through. Also, it would be interesting to study bound geodesics crossing $\tilde{r} = 0$ (and maybe also the Cauchy horizon for positive $\tilde{r}$) in general and, in particular, their causal structure. In this context the analysis of closed timelike trajectories is also of interest. In addition, we not yet considered geodesics lying entirely on the axis $\theta = 0$, $\pi$ or even crossing it. Until now, we only considered the Boyer-Lindquist form of the Kerr-de Sitter metric which is not a good choice for considering geodesics which fall through a horizon. Therefore, for a future publication it would be interesting to use a coordinate-singularity free version of the metric.

The methods for obtaining analytical solutions of geodesic equations presented in this paper are not only limited to Kerr-de Sitter space-time. Indeed, they had already been used to solve the geodesic equation in Schwarzschild-de Sitter [14,15] and Reissner-Nordström-de Sitter [17] space-times. Also the geodesic equations in higher-dimensional static spherically symmetric space-times [17] were solved by these methods. The same type of differential equations is also present in the Plebański-Demiański space-time without acceleration, which is the most general space-time with separable Hamilton-Jacobi equation. The analytical solution of the geodesic equation in this space-time is given in [18] but will be elaborated in an upcoming publication. It will also be interesting to apply the presented methods to higher-dimensional stationary axially symmetric space-times like the Myers-Perry solutions.

The same structure of equations we solved in this paper is also present in the geodesic equation of the effective one-body formalism of the relativistic two-body problem. The effective metric in this formalism can be described as a perturbed Schwarzschild or Kerr metric, where the perturbation is given in powers of the radial coordinate $r$ [30–32]. Therefore, we expect that the polynomial appearing in the resulting equations of motion will have a higher degree than the corresponding polynomial in the Schwarzschild or Kerr case and, thus, that it is necessary to generalize the elliptic functions used in these cases to hyperelliptic functions used in this paper. A similar situation can be found in the expressions of axisymmetric gravitational multipole space-times. For example, some types of geodesics in Erez-Rosen space-time, which reduces to the Schwarzschild case if the quadrupole moment is neglected, were already solved analytically [33,34]. We expect that the methods presented in this paper will be helpful to solve geodesics in space-times with multipoles.

Analytic solutions are the starting point for approximation methods for the description of real stellar, planetary, comet, asteroid, or satellite trajectories (see e.g. [35]). In particular, it is possible to derive post-Kerr, post-Schwarzschild, or post-Newton series expansions of analytical solutions. Because of the, in principle, arbitrary high accuracy of analytic solutions of the geodesic equation they can also serve as test beds for numerical codes for the dynamics of binary systems in the extreme stellar mass ratio case (extreme mass ratio inspirals, EMRIs) and also for the calculation of corresponding gravitational wave templates. For the case of Kerr space-time with vanishing cosmological constant it has already been shown that gravitational waves from EMRIs can be computed more accurately by using analytical solutions than by numerical integration [13].

Because of the high precision, the analytical expressions for observables in Kerr-de Sitter space-time may be used for comparisons with observations where the influence of the cosmological constant might play a role. This could be the case for stars moving around the galactic center black hole or binary systems with extreme mass ratios where one body serves as test-particle. For example, quasar QJ287 (cf. [36]) could be a candidate for observing the effects of a nonvanishing cosmological constant. In this context it would also be interesting for a future publication to derive post-Kerr, post-Schwarzschild, or post-Newton expressions for observables.

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APPENDIX A: INTEGRATION OF ELLIPTIC INTEGRALS OF THE THIRD KIND

In this appendix we will demonstrate the details of the integration method of the elliptic integrals of third kind which appear in the $\theta$-dependent part of the $\phi$ and $t$ motion (52) and (60) and in the special case of the $r$ motion where $R$ has a double or triple zero. We will explain the procedure for the example of the integral $I_\phi$ in (52). As this integral is only elliptic if $\Theta$ has only simple zeros we assume in the following that this is the case.

Before we demonstrate the calculational steps, we will summarize them for convenience:

(1) Cast the expression under the square root in the standard Weierstrass form $4y^3 - g_2y - g_3$ for some constants $g_2, g_3$, the so-called Weierstrass invariants.

(2) Decompose the integrand (without the square root) in partial fractions.

(3) Substitute $y = \varphi(\nu)$. 

APPENDIX B: INTEGRATION OF ELLIPTIC INTEGRALS OF THE THIRD KIND

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(2) Decompose the integrand (without the square root) in partial fractions.

(3) Substitute $y = \varphi(\nu)$. 

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(4) For every partial fraction, rewrite the integrand in terms of the Weierstrass \( \wp \) (double pole) or \( \zeta \) function (simple pole).

(5) Integrate the Weierstrass \( \wp \) and \( \zeta \) functions and assemble all parts.

Let us start from Eq. (54)

\[
I_\theta = \frac{1}{2} \int_{\nu_0}^{\nu} \frac{\bar{\alpha} - \bar{D} - \bar{\alpha} v'}{\Delta_{\nu'}(1 - \nu' \sqrt{\nu' \Theta_{\nu'}})} \, dv',
\]

where we assumed that the original integration path was fully contained in \((0, \pi/2)\). With the substitutions \( \nu = \xi^{-1} \) and \( \xi = \frac{1}{a_3}(4y - \frac{a_0}{4}) \), as in Sec. IV.A, we obtain

\[
I_\theta = \frac{1}{2} \int_{\xi_0}^{\xi} \frac{(\bar{\alpha} - \bar{D}) \xi' - \bar{\alpha} \xi' d\xi'}{\Delta_{\xi'}(\xi' - 1) \sqrt{\Theta_{\xi'}}}
\]

\[
= \frac{|a_3|}{2a_3} \int_{y_0}^{y} \frac{(\bar{\alpha} - \bar{D})(4y' - \frac{a_0}{4}) - \bar{\alpha}a_3(4y' - \frac{a_0}{4})dy'}{a_3 \Delta_y (4y' - \frac{a_0}{4} - a_3) \sqrt{4y'^2 - g_2y' - g_3}},
\]

(\text{A1})

where \( \Theta_{\xi} \) is defined in (40), \( g_2, g_3 \) are defined as in (42),

\( \Delta_\xi = \xi + \frac{a_2}{a_3} \bar{\Lambda} \), and \( \Delta_y = d_3 y + d_4 = \frac{a_2}{a_3} y - \frac{a_0}{3a_3} + \bar{\Lambda} \).

Now we simplify the integrand in (A1) (without the square root) by a partial fraction decomposition

\[
\frac{(\bar{\alpha} - \bar{D})(4y - \frac{a_0}{4} - \bar{\alpha}a_3)(4y - \frac{a_0}{4})}{a_3 \Delta_y (4y - \frac{a_0}{4} - a_3)}
\]

\[
= (\bar{\alpha} - \bar{D}) - a_3 \frac{\bar{\Lambda}(\bar{\alpha} - \bar{\Lambda})}{4y + \frac{a_0}{4}} + \frac{\bar{D}}{y - d_3}.
\]

(\text{A2})

where \( d_3 = \frac{a_0}{12} + \frac{a_2}{4} \). With the substitution \( y = \wp(\nu) \) we can get rid of the square root in (A1) as \( \wp'(\nu) = \pm \sqrt{4\wp^3(\nu) - g_2\wp(\nu) - g_3} \) where the sign has to be chosen according to the sign of \( \wp' \) and the branch of the square root. The function \( \wp'(\nu) \) is negative for \( \nu \in [0, \omega_1] \) and positive for \( \nu \in [\omega_2, \omega_2 + \omega_1] \) where \( 2\omega_1 \in \mathbb{R} \) and \( 2\omega_2 \in \mathbb{C} \) are the fundamental periods of \( \wp \). As \( \theta = \frac{\pi}{2} \) corresponds to \( y = \infty \) we will have \( \nu \in [0, \omega_1] \) in most cases.

Altogether, the integral \( I_\theta \) now reads

\[
I_\theta = \frac{|a_3|}{2a_3} \left\{ (\bar{\alpha} - \bar{D}) \int_{\nu_0}^{\nu} dv' - a_3 \frac{\bar{\Lambda}(\bar{\alpha} - \bar{\Lambda})}{4y + \frac{a_0}{4}} \right\}
\]

\[
\times \left[ a_3 \bar{\Lambda} \int_{\nu_0}^{\nu} \frac{\bar{D}dv'}{\wp'(\nu')} + \int_{\nu_0}^{\nu} \frac{\bar{D}dv'}{\wp'(\nu')} + d_3 \right].
\]

(\text{A3})

The second and third integral are of the third kind because \( f_1(\nu) = (\wp(\nu) + \frac{d_3}{d_1})^{-1} \) and \( f_2(\nu) = (\wp(\nu) - d_3)^{-1} \) have simple poles. We will rewrite now \( f_1 \) and \( f_2 \) in terms of the Weierstrass \( \sigma \) function, which has a simple zero in 0. The reason is, that \( \zeta \) can easily be integrated in terms of the Weierstrass \( \sigma \) function.

\[
\int_{\nu_0}^{\nu} \zeta(\nu')dv' = \log(\sigma(\nu) - \log(\sigma(\nu_0)).
\]

(A4)

We only demonstrate the procedure for \( f_1 \) which is totally analogous to the procedure for \( f_2 \). The function \( f_1 \) has two simple poles \( v_1 \) and \( v_2 \) lying in the fundamental domain \( \{2a_0\omega_1 + 2b\omega_2 | a, b \in [0, 1]\} \), where \( \omega_1 \) and \( \omega_2 \) as above, with \( \sigma(v_1) = \frac{-d_2}{d_4} = \sigma(v_2) \). An expansion of \( f_1 \) and \( \sigma(v) + \frac{d_2}{d_4} \) in neighbourhoods of \( v_i \) yields

\[
f_1(\nu) = a_{-1,i}(\nu - v_i)^{-1} + \text{holomorphic part},
\]

(A5)

\[
\wp(v) + \frac{d_2}{d_1} = \wp'(v_i)(\nu - v_i) + \text{higher order terms}
\]

(A6)

for some constants \( a_{-1,i} \). Now a comparison of coefficients gives

\[
1 = f_1(\nu) \left( \wp(v) + \frac{d_2}{d_1} \right) = a_{-1,i} \wp'(v_i) + \text{higher order terms},
\]

that is \( a_{-1,i} = (\wp'(v_i))^{-1} \).

(A7)

It follows that the function \( f_1(\nu) - \sum_{i=1}^{2} \frac{\zeta(\nu - v_i)}{\wp'(v_i)} \) is an elliptic function without poles and, therefore, equal to a constant \( A \) which can be determined by \( 0 = f(0) \). This yields

\[
A = -\sum_{i=1}^{2} \frac{\zeta(\nu - v_i)}{\wp'(v_i)}
\]

(A8)

Note that \( \wp'(v_2) = \wp'(2\omega_j - v_1) = -\wp'(v_1) \). The expression for \( \wp'(v_1) \) can be determined by the differential equation \( \wp'(x) = \pm \sqrt{4\wp^3(x) - g_2\wp(x) - g_3} \), where again the sign of the square root has to be chosen according to the sign of \( \wp' \). An integration of \( f_1 \) yields

\[
\int_{\nu_0}^{\nu} f_1(\nu')dv' = \sum_{i=1}^{2} \frac{1}{\wp'(v_i)} \left\{ \zeta(v_i)(\nu - v_0) + \log(\sigma(\nu - v_i)) - \log(\sigma(v_0 - v_i)) \right\}.
\]

(A9)

In the same way we can integrate \( f_2 \), where we assume that \( v_3 \) and \( v_4 \) are the simple poles of \( f_2 \) in the fundamental domain with \( \wp(v_3) = d_3 = \wp(v_4) \). Summarized, \( I_\theta \) is given by [cf. (55)]

\[
I_\theta = \frac{|a_3|}{2a_3} \left\{ (\bar{\alpha} - \bar{D})(\nu - v_0) - \sum_{i=1}^{4} \frac{a_3}{4x}\wp'(v_i) \right\}
\]

\[
\times \left( \zeta(v_i)(\nu - v_0) + \log(\sigma(\nu - v_i)) - \log(\sigma(v_0 - v_i)) + 2\pi ik \right)
\]

\[
\cdot (\bar{\alpha} \bar{\Lambda}(\bar{\alpha} - \bar{\Lambda})D)(\delta_{i1} + \delta_{i2}) + D(\delta_{i3} + \delta_{i4}) \right\}.
\]
where \( \nu = \nu(\gamma) = 2\gamma - \gamma_{\theta\text{in}} \) by (43) and \( \nu_0 = \nu(\gamma_0) \).
The integers \( k_i \) correspond to different branches of log.

**APPENDIX B: INTEGRATION OF HYPERELLPTIC INTEGRALS OF THE THIRD KIND**

In this appendix we will demonstrate the details of the integration method of the hyperelliptic integrals of third kind which appear in the \( \bar{r} \)-dependent part of the \( \phi \) and \( t \) motion (52) and (60). We will explain the procedure for the example of the integral \( I_r \) in (52). As this integral is only hyperelliptic if we consider timelike geodesics, i.e. \( \delta = 1 \) and \( R \) has only simple zeros we assume in the following that this is the case.

Before we demonstrate the solution steps, we will summarize them for convenience:
1. Cast the expression under the square root in the standard form \( u^2 + \sum_{i=0}^4 c_i u^i \) for some constants \( c_i \).
2. Decompose the integrand (without the square root) in partial fractions.
3. If integrals of first or second kind are present, rewrite them as functions of \( \gamma \).
4. Rewrite the integrals of third kind in terms of the canonical integral of third kind \( dP(x_1, x_2) \).
5. Rewrite the canonical integrals of third kind in terms of the Kleinian sigma functions, such reducing them to functions of integrals of first kind.
6. Express the integrals of first kind in terms of \( \gamma \) and assemble all parts.

Let us start with Eq. (56)

\[
I_r = \int_{\bar{r}_0}^{\bar{r}} \frac{\bar{a}(\bar{r}^2 + \bar{a}^2 - \bar{D}) d\bar{r}}{\Delta \bar{r} \sqrt{\bar{R}}},
\]

which can analogously to Sec. IV B be transformed to the standard form by \( \bar{r} = \pm 1/u + \bar{r}_R \) with a zero \( \bar{r}_R \) of \( R \) and get

\[
I_r = -\bar{a} \int_{u_0}^{u} \frac{(\pm 1/\bar{r} + \bar{r}_R)^2 + \bar{a}(\bar{a} - \bar{D})}{\Delta \bar{r} \sqrt{\bar{R}}} \frac{d\bar{r}}{u^2}
= -\bar{a} \int_{u_0}^{u} \frac{[\bar{r}^2_\bar{R} + \bar{a}(\bar{a} - \bar{D})] u^2 + 2\bar{r}_R u + 1}{\sqrt{\bar{c}_5 \Delta_u \sqrt{\bar{R}}} u^3} \frac{d\bar{r}}{u},
\]

(B1)

where \( R_u \) and \( c_5 \) are defined as in Sec. IV B and \( \frac{1}{\Delta} \Delta_u = \Delta_{u/\sqrt{\bar{r}}}, \) i.e.

\[
\Delta_u = (u^2(1 - \bar{A}^2) + 2\bar{r}_R \bar{A} u - \bar{A})(\bar{r}^2_\bar{R} + \bar{a}^2)
\]

\[
= 2\bar{r}_R u + 1) + u^2 - \bar{r}_R u^4,
\]

(B2)

which is a polynomial of degree 4 in \( u \). Note that for geodesic motion, the coordinate \( \bar{r} \) is always contained in an interval bounded by two adjacent real zeros of the polynomial \( R \) or by a real zero and infinity. This implies that \( u = (\bar{r} - \bar{r}_R)^{-1} \) for a real zero \( \bar{r}_R \) of \( R \) does not change sign on the integration path and, therefore, we can neglect the absolute value of \( u \) appearing in the integrand if we multiply the whole integral with \( \text{sign}(u_0) = u_0/|u_0| \).

Consequently,

\[
I_r = \frac{-|u_0|}{u_0} \int_{u_0}^{u} \frac{[\bar{r}_R^2 + \bar{a}(\bar{a} - \bar{D})] u^2 + 2\bar{r}_R u + 1}{\sqrt{\bar{c}_5 \Delta_u \sqrt{\bar{R}}} u^3} du.
\]

(B3)

The expression for the integrand in (B3) can be simplified by a partial fraction decomposition

\[
\frac{1}{\Delta u} I_r = C_1 \int_{u_0}^{u} \frac{u du}{\sqrt{\bar{R}}} + C_0 \int_{u_0}^{u} \frac{du}{\sqrt{\bar{R}}}
+ \frac{4}{\sqrt{\bar{c}_5}} \int_{u_0}^{u} \frac{du}{(u - u_1) \sqrt{\bar{R}}},
\]

(B4)

where \( u_i, 1 \leq i \leq 4 \) denote the zeros of \( \Delta_u \) and \( C_0, C_1, C_{2,i} \) are quite complicated expressions dependent on the parameters and the zero \( \bar{r}_R \) of \( R \), which may be calculated by a computer algebra system.

The first two integrals in this expression are of first kind and can be solved analogous to Sec. IV B (48), i.e.

\[
\int_{u_0}^{u} \frac{u du}{\sqrt{\bar{R}}} = \sqrt{\bar{c}_5} (\gamma - \gamma_0),
\]

(B5)

\[
\int_{u_0}^{u} \frac{du}{\sqrt{\bar{R}}} = \int_{u_0}^{u} \frac{du}{\sqrt{\bar{R}}} + \int_{u_0}^{u} \frac{du}{\sqrt{\bar{R}}}
= -f(\sqrt{\bar{c}_5} \gamma_0 - \gamma_{\text{in}}) + f(\sqrt{\bar{c}_5} \gamma - \gamma_{\text{in}}),
\]

(B6)

where again \( \gamma_{\text{in}} = \sqrt{\bar{c}_5} \gamma_0 + \int_{u_0}^{u} \frac{u du}{\sqrt{\bar{R}}} \) with \( u_0 = (\bar{r}_0 - \bar{r}_R)^{-1} \) only depends on the initial values \( \gamma_0 \) and \( u_0 \), and \( f \) describes the \( \theta \) divisor, i.e. \( f((\gamma, z), (z), y) = 0 \).

The four integrals in (B4) containing \( (u - u_1)^{-1} \) are in general of third kind and can be expressed in terms of the canonical integral of third kind \( \int dP(x_1, x_2) \). The most simple construction of a differential of third kind

\[
dP(x_1, x_2) = \left( \frac{y + y_1}{x - x_1} - \frac{y + y_2}{x - x_2} \right) dx
\]

(B7)

has simple poles in the points \( (x_1, y_1) \) and \( (x_2, y_2) \) of the Riemann surface of \( y^2 = g(x), g \) a polynomial, with residual +1 and −1, respectively, (cf. [37,38]). In particular, we get

\[
\int_{u_0}^{u} \frac{du}{(u - u_1) \sqrt{\bar{R}}} = \frac{1}{\sqrt{\bar{R}_{u_1}}} \sqrt{\bar{R}_u} \int_{u_0}^{u} dP(u_1^+, u_i),
\]

(B8)

where \( u_i^+ = (u_i, \sqrt{\bar{R}_{u_i}}) \) is the pole \( u_i \) located on the positive branch of the square root and \( u_i^- = (u_i, -\sqrt{\bar{R}_{u_i}}) \) is the pole \( u_i \) located on the negative branch of the square root. Based
on Riemann’s vanishing theorem (see e.g. [37]) the canonical differential of third kind can be expressed in terms of Kleinian $\sigma$ functions by

$$
\int_{u_0}^{u} dP(u_i^+, u_i^-) = \frac{1}{2} \log \left( \frac{\int_{u_0}^{u_i} \frac{d\bar{z}}{d\bar{r}_{\sigma}} - 2 \int_{u_0}^{u_i} \frac{d\bar{z}}{d\bar{r}_{\sigma}}}{\sigma \int_{u_0}^{u_i} \frac{d\bar{z}}{d\bar{r}_{\sigma}} - 2 \int_{u_0}^{u_i} \frac{d\bar{z}}{d\bar{r}_{\sigma}}} \right) - \frac{1}{2} \log \left( \frac{\int_{u_0}^{u_i} \frac{d\bar{z}}{d\bar{r}_{\sigma}} - 2 \int_{u_0}^{u_i} \frac{d\bar{z}}{d\bar{r}_{\sigma}}}{\sigma \int_{u_0}^{u_i} \frac{d\bar{z}}{d\bar{r}_{\sigma}} - 2 \int_{u_0}^{u_i} \frac{d\bar{z}}{d\bar{r}_{\sigma}}} \right) - \left( \int_{u_0}^{u} \frac{d\bar{z}}{d\bar{r}_{\sigma}} \right) \left( \int_{u_0}^{u} \frac{d\bar{r}}{d\bar{r}_{\sigma}} \right),
$$

where $d\bar{z}$ is the vector of canonical differentials of the first kind and $d\bar{r}$ the vector of canonical differentials of the second kind

$$
d\bar{z}_i = \frac{u^{(i-1)} du}{\sqrt{R_{\sigma}}}, \quad i = 1, 2.
$$

Finally, we rewrite (B9) in terms of the affine parameter $\gamma$. By (B5) and (B6) we can express $\int_{u_0}^{u} d\bar{z}$ as well as the arguments of the $\sigma$ functions $\int_{u_0}^{u} d\bar{z} = \int_{u_0}^{u} \frac{d\bar{z}}{d\bar{r}_{\sigma}}$ as functions of $\gamma$. If we define $w = w(\gamma) = \sqrt{\epsilon_{\sigma}} \gamma - \gamma_{i,\text{in}}$ and $w_0 = w(\gamma_0)$ the integral $I_r$ is given by [cf. (59)]

$$
I_r = -\frac{\tilde{a} u_0}{\sqrt{\epsilon_{\sigma} u_0}} \left[ C_1 (w - w_0) + C_0 (f(w) - f(w_0)) \right]
+ \sum_{i=1}^{4} \frac{C_{2i}}{\sqrt{R_{\sigma}}} \left[ \frac{1}{2} \log \left( \frac{\sigma(W^+(w))}{\sigma(W^-((w_0)))} \right) \right]
- (f(w) - f(w_0), w - w_0) \left( \int_{u_0}^{u} \frac{d\bar{z}}{d\bar{r}_{\sigma}} \right).
$$

where $W^+(w) := (f(w), w) - 2 \int_{u_0}^{u} \frac{d\bar{z}}{d\bar{r}_{\sigma}}$ and $W^-(w) = f(w), w) - 2 \int_{u_0}^{u} \frac{d\bar{z}}{d\bar{r}_{\sigma}}$.

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[27] Another method to observe the influence of the gravitomagnetic components $g_0$ is through the precession of gyroscopes known also as the Schiff effect. Such a measurement has been carried through by gravity probe B [28]. While the Lense-Thirring effect is an orbital effect involving the motion of the whole orbit thus constituting a global measurement, the Schiff effect is a local effect showing the dragging of local inertial frames due to the existence of the $g_0$ components.

