

BACHELOR THESIS

Numerical propagation of trajectories of Earth-orbiting spacecraft

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We are just an advanced breed of monkeys on a minor planet of a very average star. But we can understand the Universe. That makes us something very special.

Stephen Hawking

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1 Introduction

2 Conics in a nutshell

In this section will revise some basic properties of conics, that will be used in the next sections for the study of the motion of two bodies under the influence of gravity.

2.1 General conics

Definition 1. A conic is the curve obtained as the intersection of a plane with the surface of a double cone (a cone with two *nappes*).

In Fig. 1 we show the 3 types of conics: the ellipse, the parabola, and the hyperbola, which differ on their eccentricity, as we will see later on. Note that the circle is a special case of the ellipse. The following

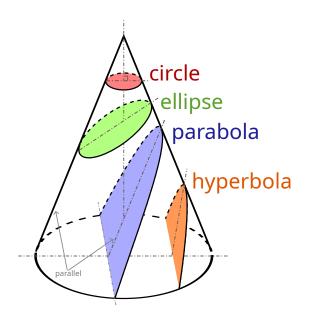


Figure 1: The black boundaries of the colored regions are conic sections. The other half of the hyperbola, which is not shown, is in the other nappe of the double cone. Source: [Mat12].

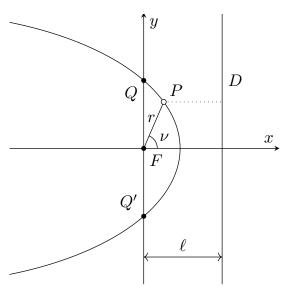


Figure 2: Reference frame centered at the focus of the conic and whose axes are such that the y-axis is parallel to the directrix and the x-axis is perpendicular to the directrix. The directions of the axes are chosen arbitrarily, subject to the constraint that a right-handed system (x,y) is obtained.

proposition gives a characterization of the conics.

Proposition 2. A conic is the set of all points P such that the distance from P to a fixed point F is a multiple of the distance from P to a fixed line D. Mathematically, this is expressed as:

$$d(P,F) = ed(P,D) \tag{1}$$

where d is the Euclidean distance. The point F is called *focus*; the line D, *directrix*, and the constant of proportionality e, eccentricity.

An important type of conic section is when e=0, but this definition is not valid in this case, as it reduces to a single point. In this case, the conic obtained, a *circle*, is defined as the set of all points P such that the distance from P to a fixed point F is constant. It can be thought as a limit case of Eq. (1) by letting the line D to infinity in a specific rate.

Note that using the polar coordinates (r, ν) centered at F (as in Fig. 2), we can rewrite Eq. (1) as:

$$r = e(\ell - r\cos\nu) \implies r = \frac{e\ell}{1 + e\cos\nu} =: \frac{p}{1 + e\cos\nu}$$
 (2)

where we have defined $p := e\ell$.

Definition 3. Let C be a conic and e be its eccentricity. We say that C is

- an ellipse if $0 \le e < 1$.
- a parabola if e = 1.

• a hyperbola if e > 1.

If e = 0, the conic is called *circle*.

2.2 Ellipse

From now on we will focus on the study of the ellipse. From Eq. (2), since e < 1, it follows that $r(\nu)$ is continuous, 2π -periodic and satisfies $r(0) = r(2\pi)$. Therefore, the ellipse is a bounded and closed curve, and it is the only conic section satisfying these two properties.

Let's now study the extrema of $r(\nu)$. An easy check shows that the minimum is attained at $\nu = 0$ and the maximum at $\nu = \pi$ and these values are given by:

$$r_{\min} = \frac{p}{1+e}$$
 and $r_{\max} = \frac{p}{1-e}$ (3)

When considering orbits of celestial bodies these points are called *periapsis* and *apoapsis*, respectively¹. The line connecting both points is called *line of apsides*. Let's seek now the extrema of $x = r \cos \nu$ and $y = r \sin \nu$. Differentiating with respect to ν yields:

$$x' = -\frac{p \sin \nu}{(1 + e \cos \nu)^2} \qquad y' = \frac{p(e + \cos \nu)}{(1 + e \cos \nu)^2}$$
(4)

On the one hand, the former expression vanishes at $\nu = 0, \pi$. Therefore, the extrema of x coincide with the periapsis and apoapsis points and at these points the y coordinate is equal to 0. This means that the line of apsides passes through the focus of the ellipse. On the other hand, y' vanishes at $\cos \nu = -e$. That is, at $\nu = \arccos(-e)$ and $\nu = 2\pi - \arccos(-e)$. Therefore, using that $\sin(\arccos x) = \sqrt{1-x^2}$, the values of y at these extrema are:

$$y_{\min} = \frac{p}{1 - e^2} \sin(2\pi - \arccos(-e)) = -\frac{p}{\sqrt{1 - e^2}}$$
 $y_{\max} = \frac{p}{1 - e^2} \sin(\arccos(-e)) = \frac{p}{\sqrt{1 - e^2}}$ (5)

Note that the x coordinate at these two points is the same: $-\frac{pe}{1-e^2}$.

Definition 4. Consider the reference frame of Fig. 3 centered at one focus. We define the *semi-major* $axis\ a$ as half the segment that connects the two extrema of the x coordinate. The *semi-minor* $axis\ b$ is defined as half the segment that connects the two extrema of the y coordinate. The length of those segments are also denoted by a and b, respectively. Thus, these are given by the following expressions:

$$a := \frac{x_{\text{max}} - x_{\text{min}}}{2} = \frac{r_{\text{max}} + r_{\text{min}}}{2} = \frac{p}{1 - e^2} \qquad b := \frac{y_{\text{max}} - y_{\text{min}}}{2} = \frac{p}{\sqrt{1 - e^2}}$$
 (6)

From here note that we can express b in terms of a and e as:

$$b = a\sqrt{1 - e^2} \tag{7}$$

Definition 5. We define the center of the ellipse O as the intersection of the semi-major axis and semi-minor axis. The *linear eccentricity* c is defined as the distance between the center O and the focus F.

We have just found a relation between a and b. Now, we would like to find a similar relation between a and c. To do so, let's calculate the distance from the focus F to one of the extrema of the y coordinate.

$$d\left(F, \left(-\frac{pe}{1 - e^2}, \pm \frac{p}{\sqrt{1 - e^2}}\right)\right) = \frac{p}{\sqrt{1 - e^2}} \sqrt{\frac{e^2}{1 - e^2} + 1} = \frac{p}{1 - e^2} = a \tag{8}$$

Hence, the value of c can be simplified to (see Fig. 3):

$$c^2 = a^2 - b^2 = a^2 - a^2(1 - e^2) = a^2 e^2 \implies c = ae$$
 (9)

Finally, one more property of the ellipse will be needed: its area.

¹Other names are used in the literature when the central body and the orbiter are particular ones. For example for the system Sun-Earth, the words *perihelion* and *aphelion* are used, whereas for the system Earth-Moon, the words *perigee* and *apogee* are used instead.

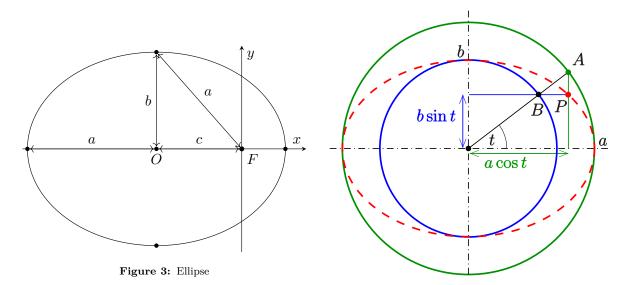


Figure 4: Reference frame centered at the center of the ellipse. Source: [Ag217].

Proposition 6. The area enclosed in an ellipse of semi-major axis a and semi-minor axis b is πab .

Proof. Consider the ellipse E centered at the origin and oriented as in Fig. 4. From Fig. 4 one can check that it can be parametrized by $(x, y) = (a \cos t, b \sin t)$ with $t \in [0, 2\pi)$. This parametrization satisfies:

$$\frac{x^2}{a^2} + \frac{y^2}{b^2} = 1\tag{10}$$

Hence, the area enclosed in the ellipse is can be parametrized by $(x,y)=(ar\cos t,br\sin t)$, with $r\in[0,1]$ and $t\in[0,2\pi)$. The Jacobian of this transformation is abr. Therefore, from the change of variable theorem we have that:

$$Area(E) = \iint_E dx dy = \int_0^{2\pi} \int_0^1 abr dr dt = \pi ab$$
 (11)

3 Introduction to astrodynamics and satellite tracking

3.1 The two body problem

3.1.1 Trajectory equation

We are interested in understanding the dynamics of a spacecraft in orbit around the Earth. These dynamics are governed by Newton's second law of motion, which assuming that both the Earth and the spacecraft are point masses (see Section 4 for a more realistic model), this can be written as

$$\ddot{\mathbf{r}} = -\frac{GM_{\oplus}}{r^2}\mathbf{e}_r \tag{12}$$

where \mathbf{r} is the position vector (also called radius vector) of the spacecraft with respect to the center of the Earth, $r := \|\mathbf{r}\|$, $\mathbf{e}_r = \frac{\mathbf{r}}{r}$ is the unit vector in the direction of \mathbf{r} , $M_{\oplus} \simeq 5.972 \times 10^{24}$ kg is the mass of the Earth, and $G \simeq 6.674 \times 10^{-11}$ m³·kg⁻¹·s⁻² is the universal gravitational constant. Note that the minus sign is due to the fact that the gravitational force is attractive, i.e. pointing towards the Earth. Here and throughout the document the dot notation $\ddot{\mathbf{r}}$ means that the derivatives are taken with respect to time. Cross-multiplying Eq. (12) by \mathbf{r} , we obtain

$$\frac{\mathrm{d}(\mathbf{r} \times \dot{\mathbf{r}})}{\mathrm{d}t} = \dot{\mathbf{r}} \times \dot{\mathbf{r}} + \mathbf{r} \times \ddot{\mathbf{r}} = -\frac{GM_{\oplus}}{r^3} (\mathbf{r} \times \mathbf{r}) = 0$$
(13)

Hence $\mathbf{r} \times \dot{\mathbf{r}} =: \mathbf{h}$ is constant. The physical intuition behind this is that the motion of the spacecraft around the Earth is confined to a plane, called the *orbital plane*, because the position \mathbf{r} and velocity $\dot{\mathbf{r}}$ are always perpendicular to \mathbf{h} .

We are interested now in what kind of curves may be described by a spacecraft orbiting the Earth, when considered both objects as point masses. That is, we want to somehow isolate \mathbf{r} (or r) from Eq. (12). In order to simplify the notation we will denote $\mu := GM_{\oplus}$.

Proposition 7 (Kepler's first law). Consider two point-mass bodies. The motion of one body orbiting the other can be described by a conic section. Hence, it can be expressed in the form:

$$r(t) = \frac{p}{1 + e\cos(\nu(t))} \tag{14}$$

for some parameters p and e.

Proof. Cross-multiplying Eq. (12) by \mathbf{h} we obtain

$$\frac{\mathrm{d}(\dot{\mathbf{r}} \times \mathbf{h})}{\mathrm{d}t} = \ddot{\mathbf{r}} \times \mathbf{h} = -\frac{\mu}{r^3} \mathbf{r} \times \mathbf{h} = -\frac{\mu}{r^3} \mathbf{r} \times (\mathbf{r} \times \dot{\mathbf{r}}) = \frac{\mu}{r^3} [(\mathbf{r} \cdot \mathbf{r})\dot{\mathbf{r}} - (\mathbf{r} \cdot \dot{\mathbf{r}})\mathbf{r}]$$
(15)

where in the last equality we have used the vector equality $\mathbf{u} \times (\mathbf{v} \times \mathbf{w}) = (\mathbf{u} \cdot \mathbf{w})\mathbf{v} - (\mathbf{u} \cdot \mathbf{v})\mathbf{w}$ for $\mathbf{u}, \mathbf{v}, \mathbf{w} \in \mathbb{R}^3$. Now note that:

$$\frac{\mathrm{d}}{\mathrm{d}t}\left(\frac{\mathbf{r}}{r}\right) = \frac{\dot{\mathbf{r}}}{r} - \frac{\dot{r}}{r^2}\mathbf{r} = \frac{1}{r^3}[(\mathbf{r}\cdot\mathbf{r})\dot{\mathbf{r}} - (\mathbf{r}\cdot\dot{\mathbf{r}})\mathbf{r}]$$
(16)

because $2r\dot{r} = \frac{d(r^2)}{dt} = \frac{d(\mathbf{r} \cdot \mathbf{r})}{dt} = 2\mathbf{r} \cdot \dot{\mathbf{r}}^2$. Thus:

$$\frac{\mathrm{d}(\dot{\mathbf{r}} \times \mathbf{h})}{\mathrm{d}t} = \mu \frac{\mathrm{d}}{\mathrm{d}t} \left(\frac{\mathbf{r}}{r}\right) \tag{17}$$

Integrating with respect to the time yields

$$\dot{\mathbf{r}} \times \mathbf{h} = \frac{\mu}{r} \mathbf{r} + \mathbf{B} \tag{18}$$

where $\mathbf{B} \in \mathbb{R}^3$ is the constant of integration. Observe that since $\dot{\mathbf{r}} \times \mathbf{h}$ is perpendicular to \mathbf{h} , $\dot{\mathbf{r}} \times \mathbf{h}$ lies on the orbital plane and so does \mathbf{r} . Hence, \mathbf{B} lies on the orbital plane too. Now, dot-multiplying this last equation by \mathbf{r} and using that $\mathbf{u} \cdot (\mathbf{v} \times \mathbf{w}) = (\mathbf{u} \times \mathbf{v}) \cdot \mathbf{w} \ \forall \mathbf{u}, \mathbf{v}, \mathbf{w} \in \mathbb{R}^3$ we obtain

$$h^{2} = \mathbf{h} \cdot \mathbf{h} = (\mathbf{r} \times \dot{\mathbf{r}}) \cdot \mathbf{h} = \mathbf{r} \cdot (\dot{\mathbf{r}} \times \mathbf{h}) = \frac{\mu}{r} \mathbf{r} \cdot \mathbf{r} + \mathbf{r} \cdot \mathbf{B} = \mu r + rB \cos \nu$$
 (19)

²Bear in mind that in general $\dot{r} \neq ||\dot{\mathbf{r}}||$. Indeed, if β denotes the angle between \mathbf{r} and $\dot{\mathbf{r}}$ we have that $\dot{r} = ||\dot{\mathbf{r}}|| \cos \beta$. In particular \dot{r} may be negative.

where $h := \|\mathbf{h}\|$, $B := \|\mathbf{B}\|$ and ν denotes the angle between \mathbf{r} and \mathbf{B} , called *true anomaly*. Rearranging the terms we finally obtain the equation of a conic section

$$r = \frac{h^2/\mu}{1 + (B/\mu)\cos(\nu)} \tag{20}$$

with
$$p := h^2/\mu$$
 and $e := B/\mu$.

From here on, let's assume that B is small enough to satisfy e < 1, as this is the primary case of interest. We've seen in Section 2.2 the range of values that can r take in that case, and we deduced an equation for the semi-major axis. Note that this latter quantity can also be expressed as:

$$a = \frac{r_{\text{max}} + r_{\text{min}}}{2} = \frac{p}{1 - e^2} = \frac{h^2}{\mu(1 - e^2)}$$
 (21)

Note that at \mathbf{r}_{\min} , we have $\nu = 0$ and so $\mathbf{r} \parallel \mathbf{B}$. Hence, \mathbf{B} points towards the periapsis of the orbit.

Definition 8. Let $\mathbf{r}(t)$ be the position of the spacecraft at time t and A(t) be the area swept by the radius vector $\mathbf{r}(t)$ in the time interval [0,t]. We define the areal velocity as $\frac{\mathrm{d}A(t)}{\mathrm{d}t}$.

Proposition 9 (Kepler's second law). The areal velocity remains constant.

Proof. Recall that the area of a parallelogram generated by two vectors $\mathbf{u}, \mathbf{v} \in \mathbb{R}^3$ is given by $\|\mathbf{u} \times \mathbf{v}\|$. Thus, approximating the difference A(t+k) - A(t) by half of the area of the parallelogram generated by $\mathbf{r}(t)$ and $\mathbf{r}(t+k)$ (see Fig. 5) we obtain

$$\frac{\mathrm{d}A(t)}{\mathrm{d}t} = \lim_{k \to 0} \frac{A(t+k) - A(t)}{k} = \lim_{k \to 0} \frac{\|\mathbf{r}(t) \times \mathbf{r}(t+k)\|}{2k} = \lim_{k \to 0} \frac{\|\mathbf{r}(t) \times (\mathbf{r}(t+k) - \mathbf{r}(t))\|}{2k} = \frac{\|\mathbf{r}(t) \times \dot{\mathbf{r}}(t)\|}{2} = \frac{h}{2} \quad (22)$$

where the penultimate equality is due to the continuity and linearity of the cross product. \Box

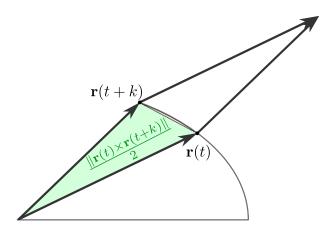


Figure 5: Graphical representation of the error made (red region) when approximating the area swept by the radius vector by half the area of the parallelogram generated by $\mathbf{r}(t)$ and $\mathbf{r}(t+k)$ (green region).

Definition 10. Let T be the orbital period of the satellite. We define the mean motion as $n := 2\pi/T$.

Proposition 11 (Kepler's third law). The mean motion is related to the semi-major axis by:

$$n = \sqrt{\frac{\mu}{a^3}} \tag{23}$$

Proof. Integrating Eq. (22) with respect to time between 0 and T (the period) yields:

$$\pi ab = A(T) = \int_{0}^{T} A'(t) dt = \int_{0}^{T} \frac{h}{2} dt = \frac{hT}{2} \implies n = \frac{2\pi}{T} = \frac{h}{ab} = \frac{h}{a^2 \sqrt{1 - e^2}} = \sqrt{\frac{\mu}{a^3}}$$
 (24)

where we have used Eqs. (7) and (21).

3.1.2 Kepler's equation

So far we have been able to describe the geometry of motion of a body orbiting another one. However, we have not concerned about the specific position of the body as a function of time. That is how to obtain $\nu(t)$ at each instant of time. In order to do this, we may think the area A as a function of ν , that measures the area swept by the radio vector from an initial instant ν_0 . Thus, from differential calculus we know that:

$$A(\nu) = \int_{\nu_0}^{\nu} \int_{0}^{r(\theta)} r \, \mathrm{d}r \, \mathrm{d}\theta = \int_{\nu_0}^{\nu} \frac{r(\theta)^2}{2} \, \mathrm{d}\theta \implies \frac{\mathrm{d}A}{\mathrm{d}\nu} = \frac{r^2}{2}$$
 (25)

Using the chain rule and Eq. (22) we obtain that:

$$\frac{h}{2} = \frac{\mathrm{d}A}{\mathrm{d}t} = \frac{\mathrm{d}A}{\mathrm{d}\nu} \frac{\mathrm{d}\nu}{\mathrm{d}t} = \frac{r^2}{2}\dot{\nu} \tag{26}$$

So from Eqs. (20) and (26), we get the following differential equation that must satisfy ν :

$$\dot{\nu} = \frac{h}{r^2} = \frac{h}{p^2} (1 + e \cos \nu)^2 \tag{27}$$

which, when integrated with respect to the time, lead us to an elliptic integral which can be very computationally expensive. Our goal in this section is to find an easier way to compute the exact position of the satellite at each instant of time [MG05]. This will lead us to the so-called *Kepler's equation*. For this purpose we are forced to introduce a new parameter, E, called *eccentric anomaly*. It is defined as the angle between the segment from the origin of the ellipse to the periapsis, and the line passing through the center of the ellipse and the point in a circle (of radius a and same center as the ellipse) which is just above the position of the satellite (see Fig. 6 for a better understanding). Clearly, using the reference

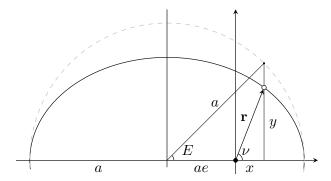


Figure 6: Ellipse orbit of the satellite together with an auxiliary circle of radius a needed to define the eccentric anomaly.

frame of Fig. 6, the position of the satellite is determined by $x = r \cos \nu$, $y = r \sin \nu$. But we would like to find an expression of x and y in terms of E rather than ν . To do this, note that $a \cos E = ae + x$, so:

$$x = a(\cos E - e) \tag{28}$$

We can also get an expression of r in terms of E by solving the equation:

$$r = \frac{p}{1 + e\cos\nu} = \frac{a(1 - e^2)}{1 + e\frac{x}{\pi}} = \frac{ra(1 - e^2)}{r + ae(\cos E - e)} \implies r = a(1 - e\cos E)$$
 (29)

Finally from Eqs. (28) and (29) we get:

$$y^{2} = r^{2} - x^{2} = a^{2}(1 - e^{2})(\sin E)^{2} \implies y = a\sqrt{1 - e^{2}}\sin E$$
(30)

Expressing now the areal velocity h as a function of E we have:

$$h = x\dot{y} - y\dot{x} \tag{31}$$

$$= a^{2}(\cos E - e)\sqrt{1 - e^{2}(\cos E)\dot{E} + a^{2}(\sin E)^{2}\dot{E}\sqrt{1 - e^{2}}}$$
(32)

$$= a^2 \sqrt{1 - e^2} \dot{E} (1 - e \cos E) \tag{33}$$

From Eq. (21) we know that $h = \sqrt{\mu a(1 - e^2)}$. Thus, substituting this in the latter equation we deduce that E must satisfy the following differential equation:

$$\dot{E}(1 - e\cos E) = \sqrt{\frac{\mu}{a^3}} = n \tag{34}$$

where the last equality follows from Proposition 11. Integrating this equation with respect to time yields the *Kepler's equation*:

$$E(t) - e\sin E(t) = n(t - t_0)$$
(35)

where t_0 is the time at which E vanishes. Using the reference frame of Fig. 6 (also known as the perifocal frame, see Definition 24) this corresponds to the time at which the satellite is at the perigee. The value $M := n(t - t_0)$ is called *mean anomaly*. Note that contrarily to E and ν , the mean anomaly increases linearly with time.

Kepler's equation is the key to solve the problem of finding the position of the satellite at each instant of time. Later on we will discuss techniques to effectively solve this equation for E, given e and M.

3.2 Time and reference systems

3.2.1 Julian day

The measurement of time has undergone significant changes throughout the centuries, and its interpretation continues to vary across different cultures. Astronomy, however, has always needed a universal time measurement system and easy-to-work with. In view of this, the so-called *Julian date* is used as the standard day-counting system in astronomy.

Definition 12 (Julian date). The *Julian date* (JD) is the number of days, of length $24 \cdot 3600 = 86400$ seconds, elapsed since the beginning of the *Julian period*, that is, since January 1st, 4713 BC at 12:00 (noon) in the Julian calendar³. A *Julian year* is defined as 365.25 days and therefore a *Julian century* as 36525 days.

As the unit of second is not constant among the time systems that we will use throughout the document (UT1, TT, UTC, etc.) (see Section 3.2.2), as a consequence, the JD will vary depending on the time system used. We will distinguish them by adding a subscript to the JD, for instance JD_{UT1} , JD_{TT} or JD_{UTC} . Later on, we will give formulas that relate these time systems, and usually they will be expressed as the number of Julian centuries elapsed since January 1st, 2000 at 12:00 TT (noon) in the Julian calendar. This number is given by:

$$\frac{\text{JD} - 2451545}{36525} \tag{36}$$

Here 2451545 corresponds to the Julian date of January 1st, 2000 at 12:00 TT (noon).

3.2.2 Time measurement

As human beings, we are naturally interested in how time passes and therefore the correct measure of it becomes an essential necessity for us. As it is the Sun that governs our daily activity, it is natural to define time from it. But first we need some preliminary definitions:

Definition 13. We define the *equatorial plane* as the plane in \mathbb{R}^3 that contains the Earth equator. We define the *ecliptic plane* as the orbital plane in \mathbb{R}^3 of the Earth around the Sun.

Definition 14. We define the *celestial sphere* as an abstract sphere of infinite radius centered at the center of mass of the Earth. All the celestial objects are thus, projected naturally on the celestial sphere, identifying them with two spherical coordinates, known as *right ascension* and *declination*, which we define below. The intersection of the equatorial plane with the celestial sphere is called *celestial equator*. The intersection of the ecliptic plane with the celestial sphere is called *ecliptic* (see Fig. 7 for a better understanding).

A first important thing to note is that, since the celestial sphere is centered at the Earth, the Sun moves along the ecliptic. Moreover, note that both the celestial equator and the ecliptic are two different great circles on the celestial sphere. Hence, they intersect at exactly two points.

³According to [Val13], the convention to start the JD at noon each day benefits astronomers (who often work at night) because they can make all their observations on a single day.

Definition 15. Consider the two points of intersection between the celestial equator and the ecliptic. We define the *vernal equinox* as the point Υ , from these two, where the Sun crosses the celestial equator from south to north.

The angle measured along the equator of any object on the celestial sphere from the vernal equinox is called *right ascension*, whereas the angle measured along the meridian of the object from the position of the object to the equator is called *declination* (see Fig. 7).

An apparent solar day is defined to be the time between two successive transits of the Sun across our local meridian. One should note that the Earth has to rotate on itself slightly more than one revolution in order to complete one solar day. In addition, due to the non-circular orbit of the Earth around the Sun, the length of an apparent solar day is not constant, as the Earth has to rotate on itself slightly more in the perihelion that in the aphelion. The apparent sidereal day is defined as the time it takes for the Earth to complete a rotation (see Fig. 8 for a better understanding). From the point of view of the celestial sphere, the apparent solar time is the angle (measured along the celestial equator) between the local meridian⁴ and the meridian of the Sun at that epoch, which is not uniform because the angular velocity of the Earth is not constant [MG05].

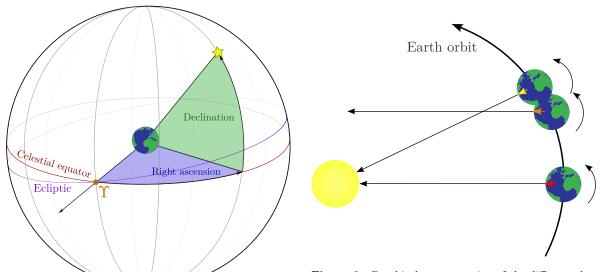


Figure 7: Right ascension and declination of a star in the celestial sphere

Figure 8: Graphical representation of the difference between a solar day (yellow) and a sidereal day (orange). (not to scale)

Due to Kepler's second law, the Earth's non-circular orbit around the Sun results in some days being shorter than others. As a result, the actual position of the Sun is not ideal for precise time measurement. The introduction of a $mean\ Sun$ is, thus, a necessity.

Definition 16. The mean Sun is a fictitious Sun that moves along the celestial equator at a constant rate [Val13]. This rate is determined in order to make the real Sun and the mean Sun coincide at the vernal equinox. We define the mean solar time as the hour angle (along the celestial equator) between the local meridian and the meridian of the mean Sun.

It is worth noting that unlike the real Sun, the mean Sun moves along the celestial equator, rather than the ecliptic.

Definition 17. We define the *prime meridian* or *zero meridian* as the meridian on the celestial sphere that passes through the Royal Observatory in Greenwich, England (when projected on the Earth).

Definition 18. The *Greenwich Mean Time* (GMT) or *Universal Time* (UT) is the hour angle of the mean Sun measured from the prime meridian and counted from midnight. That is, when the prime meridian and mean Sun meridians coincide, the GMT is 12:00.

The use of two distinct names, namely GMT and UT, to refer to the same time can be attributed to historical reasons. Initially, GMT was defined as the mean solar time at the prime meridian with 00:00

⁴We consider the local meridian as the meridian on the celestial sphere that is projected onto the Earth and moves in synchronization with it.

GMT coinciding with the moment when the mean Sun was at the that meridian. On the other hand, UT was introduced as a 12-hour translation of GMT, intended for civilian purposes. Eventually, GMT was redefined to align with UT.

In the middle of the 20th century, Ephemeris Time (ET) was introduced to cope with the irregularities of the Earth's rotation (see Section 3.2.3). This time was defined from historical observations of planets in a Newtonian physics framework, isolating the time from the equations, and the origin of time was chosen coherently with GMT, at January 0, 1900, 12:00 GMT (JD 2415020.0). This timescale provided a uniform time, although it was more difficult to measure than the mean solar time. In the meantime, atomic clocks were invented and at 1967 the atomic time (TAI, from French Temps Athomique International) was adopted as the SI unit of second. The origin was chosen such that the TAI matched UT at the 00:00:00 UT of January 1st, 1958, and at that time the ET was displaced from UT by 32.184 seconds. At the end of the 20th century, the Terrestrial Time (TT) was introduced within a relativistic framework in order to replace ET and provide a smooth and more accurate continuation of it yielding the relation [MG05]:

$$TT = ET = TAI + 32.184 s$$
 (37)

For astronomical calculations, it is convenient to consider a timescale defined directly from Earth's rotation, known as *sidereal time*. Namely, *Greenwich Mean Sidereal Time* (GMST) is defined as the angle between the prime meridian and the mean vernal equinox of date (see Section 3.2.3). Due to unpredictable irregular changes on the rotation of the Earth, the GMST cannot be computed directly with a formula in terms of the TAI or TT.

The *Universal Time 1* (UT1) is the presently used form of Universal time, and it is defined in terms of Earth's rotation with the following deterministic formula given in [Aok+81]. For each day, the 00:00 UT1 is defined as the time instant in which GMST has the value:

$$GMST(0h UT1) = 24110.54841 + 8640184.812866T_{UT1,0} + 0.093104T_{UT1,0}^{2} - 6.2 \cdot 10^{-6}T_{UT1,0}^{3}$$
 (38)

where $T_{\rm UT1,0} = \frac{\rm JD(0h~UT1)-2451545}{36525}$ denotes the number of Julian centuries that have passed since January 2000, 1.5 UT1 at the beginning of the day. The units of the coefficients are seconds. For any instant of time during the day, the following formula is used:

$$GMST(UT1) = 24110.54841 + 8640184.812866T_{UT1} + 1.002737909350795 \cdot 240UT1 + 0.093104T_{UT1}^{2} - 6.2 \cdot 10^{-6}T_{UT1}^{3}$$
 (39)

where $T_{\rm UT1} = \frac{\rm JD(UT1)-2451545}{36525}$ and UT1 are measured in seconds. The coefficient $\omega := 1.002737909350795$ is the Earth's mean angular velocity in degrees per second and the coefficient $240 = \frac{3600}{15}$ is the number of seconds in one degree⁵. Similarly to the GMST, there is no simple conversion between the UT1 and the TT or TAI. Instead, IERS⁶ (International Earth Rotation and Reference Systems Service) provides regularly a bulletin with the difference $\Delta T := {\rm TT} - {\rm UT1}$ at several dates. Interpolating these values we can obtain UT1 from TT at any epoch.

Finally, our everyday clock is based on the *Coordinated Universal Time* (UTC). It is defined to be as uniform as the TAI but always kept closer than 0.9 seconds to the UT1 in order to resemble a mean solar time (see Fig. 9). Scientists achieve this by introducing a *leap second* (see Fig. 10), which is an extra second added to UTC at irregular intervals. Fig. 9 summarizes all the time systems introduced in the document.

We summarize here useful conversions between time systems:

$$\begin{aligned} \text{GMST(UT1)} &= 24110.54841 + 8640184.812866T_{\text{UT1}} + 1.002737909350795 \cdot 240\text{UT1} + \\ &\quad + 0.093104T_{\text{UT1}}{}^2 - 6.2 \cdot 10^{-6}T_{\text{UT1}}{}^3 \\ \text{UT1} &= \text{TT} - \Delta T \\ \text{TT} &= \text{TAI} - 32.184 \\ \text{TAI} &= \text{UTC} + \delta \end{aligned}$$

⁵Recall that $15^{\circ} = 1$ h, $60' = 1^{\circ}$, 60'' = 1', where " are arc seconds and ' are arc minutes.

⁶More information on https://www.iers.org (Accessed on June 8, 2023).

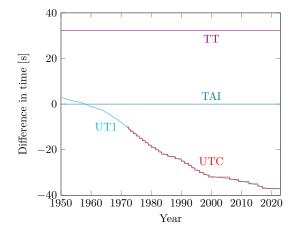


Figure 9: Evolution of times TT, UT1 and UTC in comparison with TAI. [Obsa]

1998	December	31,	23h	59m	59s
1998	December	31,	23h	59m	60s
1999	January	1,	00h	00m	00s
1999	January	1,	00h	00m	01s

Figure 10: Leap second introduced to the UTC time at the end of the December 1998. [RS98b]

where ΔT is the difference between the TT and UT1, and δ is a piecewise constant function that counts the number of leap seconds introduced since 1972, when they were introduced for the first time. All the numbers have units of second. Note that since the rotation of the Earth cannot be predicted accurately, ΔT can only be determined retrospectively, and is given by the *International Earth Rotation and Reference Systems Service* (IERS).

3.2.3 Reference systems

It is well-known that Newton's second law is only valid when applied to an *inertial reference frame*, that is, a frame of reference that is not undergoing any acceleration. In practice, however, almost any frame of reference is non-inertial. So in this chapter we will describe a quasi-inertial frame of reference which will be used to integrate Newton's second law. On the other hand, since the Earth is not a body with a homogeneous density of mass, there are zones with higher mass density than others, and therefore with higher gravitational pull (see Section 4.2.3). Therefore, we will need the longitude and latitude of the satellite with respect to the Earth at each time of integration in order to place it in the correct position with respect to the Earth.

Based on our study of satellite motion around the Earth, it is natural to place all the origins of the reference frames considered throughout the document at the center of mass of the Earth.

The first reference frame we must consider is the *celestial* one. In the celestial frame, the x-axis is defined as the line of intersection between the equatorial plane and the ecliptic plane. The positive direction is chosen to point towards the vernal equinox. The z-axis is chosen to be perpendicular to the equatorial plane and the y-axis is such that the triplet (x, y, z) is a right-handed system.

However, due to the presence of other solar system planets (and other smaller perturbations), the orbital plane of the Earth is not fixed in space, but is subjected to a small variation called *planetary precession*. Moreover, the gravitational attraction of the Sun and Moon on the Earth's equator cause Earth's axis of rotation to precess similarly to the action of a spinning top with a period of about 26000 years [MG05]. This motion is called *lunisolar precession*. On the other hand, smaller perturbations in amplitude and shorter period (around < 18.6 years [Wikb]) superposed with the precessional motion creates a motion called *nutation*. When this latter oscillations are averaged out, the vernal equinox and the equator are referred to *mean* values, rather than *true* values.

In addition, the Earth's axis of rotation undergoes a slight periodic motion around a reference axis of rotation, that passes through the IRP (IERS Reference Pole). This motion is called polar motion and is caused by the redistribution of the Earth's mass due to the seasonal variations of the atmosphere and the oceans. The polar motion is usually described by the CEP (Celestial Ephemeris Pole), which is the intersection of the Earth's axis of rotation with the celestial sphere. In Fig. 11 we summarize the different types of perturbations on the Earth's axis of rotation.

In view of this time-dependent orientation of both the ecliptic and the equator, the standard-reference frame chosen is based on the mean equator, ecliptic and mean equinox of a fixed time, the beginning of

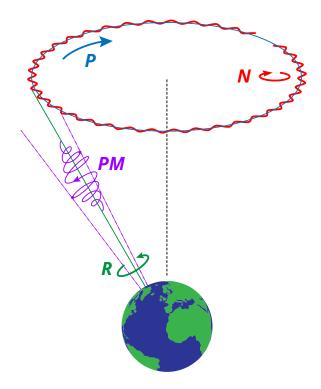


Figure 11: Graphical explanation of the perturbation by precession (blue), nutation (red) and polar motion (violet) of the Earth's axis of rotation (green).

the year 2000, namely at 12:00 TT on 1st January 2000, the so-called J2000 epoch.

Definition 19 (Earth-centered inertial frame). We define the J2000 reference frame as the reference frame whose x-axis is the intersection of the mean celestial equator and the ecliptic of the J2000 epoch, pointing at the mean vernal equinox of the same epoch; the z-axis is perpendicular to the mean equator, and the y-axis is chosen such that the triplet (x, y, z) is a right-handed system. This frame of reference is also called Earth-centered inertial (ECI) frame.

Let's introduce now an Earth-fixed reference frame.

Definition 20 (Earth-centered, Earth-fixed frame). We define the *Earth-centered, Earth-fixed frame of reference* (ECEF) as the frame of reference whose x-axis is pointing to the prime meridian, the z-axis is perpendicular to the Earth's equator and the y-axis is chosen such that the triplet (x, y, z) is a right-handed system.

These two coordinate systems have, as mentioned earlier, the origin at the center of mass of the Earth. Note that, the ECEF frame is not inertial, since it is rotating with the Earth.

3.2.4 Conversion between reference systems

As we noted in the previous section the angle ε between the celestial equator and ecliptic planes is not constant due to the planetary precession.

Our goal in this section is to transform the position of the satellite from the ECI system to the ECEF system and vice versa. This rotation transformation is given by a product of 4 rotations matrices:

- the precession matrix **P**,
- the nutation matrix N,
- the Earth rotation matrix Θ , and
- the polar motion matrix Π .

These matrices are such that:

$$\mathbf{r}_{\text{ECEF}}(t) = \mathbf{\Pi}(t)\mathbf{\Theta}(t)\mathbf{N}(t)\mathbf{P}(t)\mathbf{r}_{\text{ECI}}(t)$$
(40)

where $\mathbf{r}_{\text{ECEF}}(t)$ is the position vector of the satellite in the ECEF frame at time t and $\mathbf{r}_{\text{ECI}}(t)$ is the position vector of the satellite in the ECI frame at time t. From here on, we will omit the evaluation on the time t for the sake of simplicity. Let's now argue why the transformation has this particular form.

The precession matrix is responsible for *eliminating* all the movement due to the planetary and lunisolar precession. Thus, \mathbf{P} transforms the mean equator and mean equinox at time J2000 to the respective values at time t. With the help of Fig. 12, one can check that this transformation is given by:

$$\mathbf{P} = \mathbf{R}_z(-90 - z)\mathbf{R}_x(\theta)\mathbf{R}_z(90 - \zeta) \tag{41}$$

And with a bit of algebra it can be simplified to:

$$\mathbf{P} = \mathbf{R}_z(-z)\mathbf{R}_y(\theta)\mathbf{R}_z(-\zeta) \tag{42}$$

Recall that the fundamental rotation matrices $\mathbf{R}_x(\theta)$, $\mathbf{R}_y(\theta)$ and $\mathbf{R}_z(\theta)$ are with respect to the axis of the J2000 frame, and they are given by:

$$\mathbf{R}_{x}(\varphi) = \begin{pmatrix} 1 & 0 & 0 \\ 0 & \cos \varphi & \sin \varphi \\ 0 & -\sin \varphi & \cos \varphi \end{pmatrix} \quad \mathbf{R}_{y}(\varphi) = \begin{pmatrix} \cos \varphi & 0 & -\sin \varphi \\ 0 & 1 & 0 \\ \sin \varphi & 0 & \cos \varphi \end{pmatrix} \quad \mathbf{R}_{z}(\varphi) = \begin{pmatrix} \cos \varphi & \sin \varphi & 0 \\ -\sin \varphi & \cos \varphi & 0 \\ 0 & 0 & 1 \end{pmatrix}$$
(43)

where we have used the convention of signs given by [GPS02]. The reader may wonder why we have used

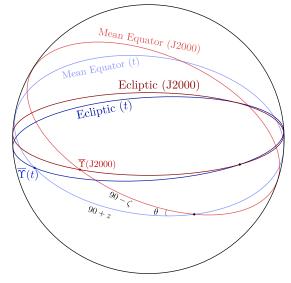


Figure 12: Celestial sphere showing the ecliptic and the equator of both the epoch J2000 and the current epoch t. Dark colors represent the ecliptic while light colors represent the equator. On the other hand, red colors represents the J2000 epoch and blue colors represents the current epoch t. Source: based on [MG05].

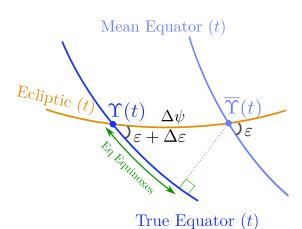


Figure 13: True equator and mean equator, and true equinox (Υ) and mean equinox $(\overline{\Upsilon})$ at a given epoch t together with the ecliptic at that time. Source: based on [MG05].

the notation 90 - z and $90 - \zeta$ instead of z and ζ (for example) for the angles in question. The reason is related to the precise definition of these angles from the pole of the celestial sphere rather than from where we have defined them, but we will not elaborate on this point here. Nonetheless, we have chosen this notation to maintain consistency with related articles [Lie+77].

The nutation perturbations are removed by the nutation matrix \mathbf{N} . This matrix transforms the coordinates of the mean equator and equinox at epoch t to those of the true equator and equinox at the same epoch, respectively. Hence, from figure Fig. 13 we can see that the nutation matrix is given by:

$$\mathbf{N} = \mathbf{R}_x(-\varepsilon - \Delta\varepsilon)\mathbf{R}_z(-\Delta\psi)\mathbf{R}_x(\varepsilon) \tag{44}$$

In Section 3.2.2 we defined the GMST as the hour angle between the mean vernal equinox and the prime meridian, measured on the true equator. Similarly, we define the *Greenwich Apparent Sidereal Time* (GAST) as the hour angle between the true vernal equinox and the prime meridian, measured on

the true equator too. The difference between these two times is given by the so-called *equation of the* equinoxes, which up to first order in the nutation angles is given by (see Fig. 13):

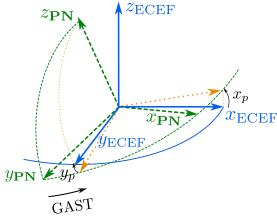
$$GAST - GMST \simeq \Delta \psi \cos(\varepsilon + \Delta \varepsilon)$$
 (45)

The Earth rotation matrix Θ is responsible for aligning satellite with his actual meridian at the time of the observation. This matrix is given by:

$$\mathbf{\Theta} = \mathbf{R}_z(\text{GAST}) \tag{46}$$

Finally the polar motion movement around the CEP axis is modelled with the matrix Π given by:

$$\mathbf{\Pi} = \mathbf{R}_y(-x_p)\mathbf{R}_x(-y_p) \tag{47}$$



GAST

Figure 14: Three reference frames used to transform from the ECI frame to the ECEF frame. In green, the ECI frame once applied the precession and nutation transformations. In orange, the transformation of the green system once applied the rotation matrix Θ . Finally, in blue, the transformation of the orange system once applied the polar motion matrix Π .

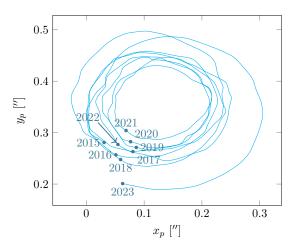


Figure 15: Graphical representation of evolution of the CEP with respect to IRP (placed at the origin of the graphic) due to polar motion. The blue dots indicate the position of the CEP at the beginning of the respective years. Data from [Obsb].

In Fig. 14 we can see a graphical representation of these two latter transformation matrices. The parameters x_p and y_p are constantly updated in the IERS bulletins [RS98a]. All the other parameters are given in [Lie+77], and we provide here a summary of them:

$$\varepsilon = 23.4392911^{\circ} - 46.8150''T - 0.00059''T^{2} + 0.001813''T^{3}$$
(48)

$$\zeta = 2306.2181'' + 0.30188''T + 0.017998''T^2 \tag{49}$$

$$\theta = 2004.3109''T - 0.42665''T^2 - 0.041833''T^3 \tag{50}$$

$$z = 2306.2181'' + 1.09468''T^2 + 0.018203''T^2$$
(51)

where $T = \frac{\text{JD}_{TT} - 2451545}{36525}$ are the Julian centuries that have elapsed since the J2000 epoch. The values for nutations $\Delta \psi$ and $\Delta \varepsilon$ can be expanded into [MG05]:

$$\Delta \psi = \sum_{i=1}^{106} (\Delta \psi)_i \sin \phi_i \qquad \Delta \varepsilon = \sum_{i=1}^{106} (\Delta \varepsilon)_i \cos \phi_i$$
 (52)

where $\phi_i = p_{\ell,i}\ell + p_{\ell',i}\ell' + p_{F,i}F + p_{D,i}D + p_{\Omega,i}\Omega$, $p_{\ell,i}$, $p_{\ell',i}$, $p_{F,i}$, $p_{D,i}$, $p_{\Omega,i}$ are integer coefficients and the other variables are the Moon's mean anomaly (ℓ) , the Sun's mean anomaly (ℓ') , the mean distance of the Moon from the ascending node (F), the difference between the mean longitudes of the Sun and the Moon (D), and the mean longitude of the ascending node of the Moon (Ω) . We will not go into detail about these variables, there are explicit expressions for them in [MG05] as a function of $T = \frac{\text{JD}_{TT} - 2451545}{36525}$:

$$\ell = 134^{\circ}57'46.733'' + 477198^{\circ}52'02.633''T + 31.310''T^2 + 0.064''T^3$$
(53)

$$\ell' = 357^{\circ}31'39.804'' + 35999^{\circ}03'01.224''T - 0.577''T^2 - 0.012''T^3$$
(54)

$$F = 93^{\circ}16'18.877'' + 483202^{\circ}01'03.137''T - 13.257''T^2 - 0.011''T^3$$
(55)

$$D = 297^{\circ}51'01.307'' + 445267^{\circ}06'41.328''T - 6.891''T^{2} + 0.019''T^{3}$$
(56)

$$\Omega = 125^{\circ}02'40.280'' - 1934^{\circ}08'10.539''T + 7.455''T^{2} + 0.008''T^{3}$$
(57)

Some of these integer coefficients together with the values of the amplitudes $(\Delta \psi)_i$ and $(\Delta \varepsilon)_i$ are given in Table 1.

i	$p_{\ell,i}$	$p_{\ell',i}$	$p_{F,i}$	$p_{D,i}$	$p_{\Omega,i}$	$(\Delta \psi)_i \ [0.0001'']$	$(\Delta \varepsilon)_i [0.0001'']$
1	0	0	0	0	1	-171996 - 174.2T	92025 + 8.9T
2	0	0	0	0	2	2062 + 0.2T	-895 + 0.5T
3	-2	0	2	0	1	46	-24
4	2	0	-2	0	0	11	0
5	-2	0	2	0	2	-3	1
6	1	-1	0	-1	0	-3	0
7	0	-2	2	-2	1	-2	1
8	2	0	-2	0	1	1	0
9	0	0	2	-2	2	-13187 - 1.6T	5736 - 3.1T
:	:	:	:	:	:	:	:
•	•	•	•	•	•	•	

Table 1: First nutation coefficients [MG05]

3.3 Keplerian orbital elements

In this section we will introduce the Keplerian orbital elements, which are a set of variables that completely determine the orbit of a satellite and, therefore, they are very useful in the storage of the orbital information of a satellite.

3.3.1 Keplerian orbital elements from position and velocity

We first give some preliminary definitions.

Definition 21. Consider a satellite orbiting the Earth. The *orbital plane* is the plane that contains the orbit of the satellite. The *line of nodes* is the line of intersection between the orbital plane and the equatorial plane. Finally, the *ascending node* is the point on the line of nodes and the orbit of the satellite where the satellite crosses the equatorial plane from south to north.

Definition 22 (Orbital elements). The *Keplerian orbital elements* of a satellite are five independent quantities that completely determine its orbit. If moreover the exact position of the satellite on the orbit is wanted, a sixth quantity is needed. The first five orbital elements are:

- 1. The *semi-major axis a* of the orbit.
- 2. The eccentricity e of the orbit.
- 3. The inclination i, which is the angle between the equatorial plane and the orbital plane.
- 4. The longitude of the ascending node Ω , which is the angle between the vernal equinox and the ascending node.
- 5. The argument of perigee ω , which is the angle between the ascending node and the periapsis, measured along the orbit.

The sixth quantity is the true anomaly ν , which is the angle between the periapsis and the position of the satellite on the orbit.

Fig. 16 shows a schematic representation of these elements. The elements a, e and i are always well-defined. However, the elements Ω , ω and ν are not always well-defined, namely for e=0 or i=0. We discuss this in more detail below. For the moment assume $e \neq 0$ and $i \neq 0$. In order to express these elements in terms of the position and velocity of the satellite, we need to introduce the following vectors:

Definition 23. Let $\mathbf{e}_z = (0,0,1)$ be the unit vector perpendicular to the equatorial plane. We define the vector $\mathbf{n} := \mathbf{e}_z \times \mathbf{h}$ and the eccentricity vector \mathbf{e} as $\mathbf{e} := \mathbf{B}/\mu$, whose norm is the eccentricity e.

Note that $\mathbf{n} \perp \mathbf{e}_z$ and $\mathbf{n} \perp \mathbf{h}$ which imply that \mathbf{n} must lie on the orbital plane and equatorial plane, and therefore, in the line of nodes pointing towards the ascending node, by the right-hand rule. On the other hand, since \mathbf{B} points towards the periapsis, so does \mathbf{e} . From here, let's obtain the orbital elements in terms of \mathbf{r} and $\dot{\mathbf{r}}$.

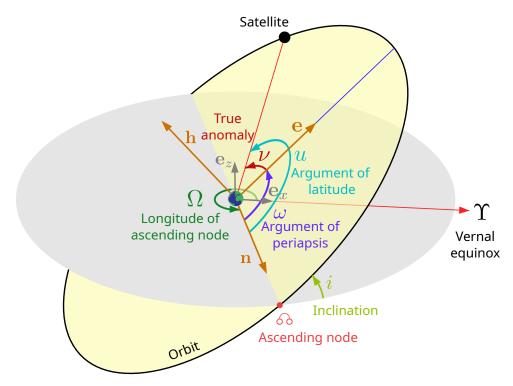


Figure 16: Orbital elements of a satellite. Source: based on [Las07].

First, we calculate $\mathbf{h} = \mathbf{r} \times \dot{\mathbf{r}}$ and from here note that the vectors \mathbf{n} and \mathbf{B} can be computed directly from \mathbf{r} , $\dot{\mathbf{r}}$ and \mathbf{h} . Thus, we can get e by taking the norm of \mathbf{e} . Then, looking at Fig. 16 it's not hard to see that the angles i, Ω , ω and ν are given by:

$$i = \arccos\left(\frac{\mathbf{h} \cdot \mathbf{e}_z}{h}\right) \quad \Omega = \arccos\left(\frac{\mathbf{n} \cdot \mathbf{e}_x}{n}\right) \quad \omega = \arccos\left(\frac{\mathbf{n} \cdot \mathbf{e}}{ne}\right) \quad \nu = \arccos\left(\frac{\mathbf{e} \cdot \mathbf{r}}{er}\right)$$
 (58)

Here $\mathbf{e}_x = (1,0,0)$ denotes the basis unit vector that points towards the vernal equinox. By convention the angles i, Ω , ω are always given in the interval $[0,2\pi)$, so in these formulas a small correction is needed when the angles are negative. The angle $\nu \in \mathbb{R}$ does not need any correction, because we will use it as a counter of the number of revolutions of the satellite. Finally, we can obtain a from Eq. (21):

$$a = \frac{h^2}{\mu(1 - e^2)} \tag{59}$$

Now we study the singular cases. If e=0 (and therefore ${\bf e}=0$) and $i\neq 0$, the orbit is called *circular inclined* [Val13]. In this case the elements ω , ν are not well-defined because there is no periapsis, or in other words, all the points lie at the same distance from the center of the Earth. To correct this we replace these variables by the *argument of latitude u*, which measures the angle between the ascending node and the position of the satellite on the orbit. The argument of latitude can be computed with the formula:

$$u = \arccos\left(\frac{\mathbf{n} \cdot \mathbf{r}}{nr}\right) \tag{60}$$

Note that in the case $e \neq 0$ and $i \neq 0$, $u = \omega + \nu$ (see Fig. 16). If i = 0 and $e \neq 0$, the orbit lies in the equatorial plane and it is called *elliptical equatorial*. Note that in this situation we have $\mathbf{n} = 0$ and the variables Ω and ω are undefined. By convention, we set $\Omega = 0$ and:

$$\omega = \arccos\left(\frac{\mathbf{e} \cdot \mathbf{e}_x}{e}\right) \tag{61}$$

If e=0 and i=0, the orbit is called *circular equatorial* and all these three variables $(\Omega, \omega \text{ and } \nu)$ are undefined. In this case Ω is set to 0 and the other two variables are replaced with the *true longitude* λ , which is the angle between the vernal equinox and the position of the satellite on the orbit:

$$\lambda = \arccos\left(\frac{\mathbf{r} \cdot \mathbf{e}_x}{r}\right) \tag{62}$$

3.3.2 TLE sets

The positions of satellites are recorded and stored in a particular way, called *Two Line Element* sets (TLEs). It is base on a two-line text with different data that facilitate the computation of position and velocity of the satellite at that specific instant of time. The following table summarizes all the information on it:

Line			tell			lass		In			ati na									000		`										[re	'n ev/	/2 da		1					[re		/6 'da	ıy ³	1			(dra	B ag		rm	1)		Model	Tonor		ΓLE	E er	heck
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1	2	3 4	3	6	1	5 9	10	11	12	13	14	15	16	1/ 1	8 1	20	21	22	23	24	25	26	21	28	29	30 .	31 .	32 3	5	4	36) 37	38	39	40	41	42	43 4	4 4	S 4	9 47	48	49	50	51	52	53 5	3 5	5 56	37	58	59	S 1	F. 6.	2 63	3 64	65 6	00 0	08	69
1		5 5	1	2	4 l	J	9	8	0	6	7	U	R		2	2 3	0	8	6		1	6	7	7	8	9	9	4	ľ		0	1	5	4	0	4	0	7	ľ	1	4	1	0	2	-	2		2	6	6	3	4	- 7	2	0			9 9	9	1
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2		3 4 5 5				8 9		11					T	7 1	8 1	9 20			0		7					0							38						T	7 1	T				П		53 5								9 9		1 :			П

Table 2: TLE data set from the NUTSAT satellite. The white (empty) cells designate space characters and the green and yellow ones are used to distinguish consecutive data blocks. The cells labeled with S or E respresent cells reserved for the negative sign and the exponent of a number respectively, while the red dots denote that an implicit decimal point is assumed.

Let's clarify the meaning of some of the data blocks. The satellite number is a unique identifier assigned by NORAD (*North American Aerospace Defense Command*) for each earth-orbiting artificial satellite [Kel]. The classification of the satellite (Class) is divided into three categories: unclassified (U), classified (C), and secret (S). The international designator is comprised of three parts: the launch year (Year), the launch number of the year (Launch) and the piece of the launch (Piece). In the epoch year (Epoch), the first two numbers indicate the last two digits of the year, and the day of the year is represented as a fraction starting from 1. The model category refers to the orbital model used to generate the data, as specified in [Kel; Wikc]. The element set number (TLE number) is incremented by one when a new TLE (Two-Line Element) is generated for this satellite. On the second line, the number of revolutions indicates the number of times the satellite has orbited the Earth since its launch. Finally, the checksum (modulo 10) is used to verify the integrity of the data⁷.

3.3.3 Position and velocity in terms of the TLEs' orbital elements

We are now interested in proceed the other way around, that is, given the orbital elements or more precisely the TLE data, we want to compute the position and velocity of the satellite. In order to do that, we need to introduce the basis $(\mathbf{P}, \mathbf{Q}, \mathbf{W})$ linked to the orbit.

Definition 24 (Perifocal coordinate system). Consider the orbit of a satellite. We define its associated *perifocal coordinate system* ($\mathbf{P}, \mathbf{Q}, \mathbf{W}$) as follows. The center is on the Earth's center of mass. The unit vectors \mathbf{P} and \mathbf{Q} lie on the orbital plane and are such that \mathbf{P} points towards the periapsis, that is $\mathbf{P} := \mathbf{B}/B$. The unit vector \mathbf{W} is defined as $\mathbf{W} := \mathbf{h}/h$, that is perpendicular to the orbital plane, and $\mathbf{Q} := \mathbf{W} \times \mathbf{P}$.

Recall that in Section 3.1.2 we have seen that the position of the satellite in the perifocal frame is given by:

$$\mathbf{r}_{\text{Peri}} = a(\cos E - e)\mathbf{P} + a\sqrt{1 - e^2}\sin E\mathbf{Q}$$
(63)

Differentiating yields:

$$\dot{\mathbf{r}}_{\text{Peri}} = -a(\sin E)\dot{E}\mathbf{P} + a\sqrt{1 - e^2}(\cos E)\dot{E}\mathbf{Q}$$
(64)

where $\dot{E} = \frac{n}{1 - e \cos E}$ is given by Eq. (34). These two quantities depend on the two unknown variables a and E, because n and e are given in the TLE data set. The semi-major axis can be easily obtained from n due to the Kepler's third law (Proposition 11). For the eccentric anomaly, we need to solve the Kepler equation:

$$E - e\sin E = M \tag{65}$$

Lemma 25. Let $e \in [0,1)$, $M \in \mathbb{R}$ and $\bar{M} = M \mod 2\pi$ be such that $\bar{M} \in [0,2\pi)$. Then, the function

$$f(E) = E - e\sin E - M \tag{66}$$

 $^{^{7}}$ Taking into account that the negative sign is counted as 1, and all the other cells without a number as 0.

has a unique solution in the interval [M, M+e] if $\bar{M} \in [0, \pi)$ and in the interval [M-e, M] if $\bar{M} \in [\pi, 2\pi)$.

Proof. We first prove the uniqueness. Clearly $f \in C^1(\mathbb{R})$ and $f'(E) = 1 - e \cos E > 0$ for all $E \in [0, 2\pi)$ because e < 1. Thus, f is strictly increasing and so it has at most one zero. Now, if $0 \le \bar{M} < \pi$, then:

$$f(M) = -e\sin M \le 0$$
 and $f(M+e) = e(1-\sin(M+e)) \ge 0$

So by Bolzano's theorem, f has a solution in [M, M + e]. If $\pi \leq \bar{M} < 2\pi$, then:

$$f(M) = -e \sin M \ge 0$$
 and $f(M - e) = -e(1 + \sin(M - e)) \le 0$

So again by Bolzano's theorem, f has a solution in [M-e, M].

We will use the Newton's method to find the zero of this non-linear equation. For small eccentricities e, the natural choice for the initial guess is $E_0 = M$. For large eccentricities (e > 0.8) the initial guess $E_0 = \pi$ should be used in order to avoid convergency problems [MG05]. Alternatively, we can do one or two steps of the bisection method to get a better initial guess, and then apply the Newton's method. Once obtained E, the position and velocity of the satellite in the J2000 frame are given by:

$$\mathbf{r}_{\mathrm{ECI}} = \mathbf{T}\mathbf{r}_{\mathrm{Peri}} \qquad \dot{\mathbf{r}}_{\mathrm{ECI}} = \mathbf{T}\dot{\mathbf{r}}_{\mathrm{Peri}}$$
 (67)

where T is the rotation matrix that transform one frame to another and is given by (look at Fig. 16):

$$\mathbf{T} = \mathbf{R}_z(-\Omega)\mathbf{R}_x(-i)\mathbf{R}_z(-\omega) \tag{68}$$

4 Earth's gravitational field

So far we have only considered the gravitational force acting between point masses. In reality, the Earth is not a point mass, neither a spherically symmetric mass distribution. In this section we will delve into the details of a more realistic model of the Earth's gravitational field.

4.1 Geopotential model

4.1.1 Continuous distribution of mass

In Section 3.1 we saw that the motion of a body orbiting another one can be described by a conic section. However, we have not been concerned about the mass distribution of the large body, in our case the Earth. In this section we will see that the motion of the smaller body, the satellite, is slightly perturbed by the mass distribution of the Earth as well as the presence of other forces such as atmospheric drag, solar radiation pressure, or the gravitational pull of the Moon and Sun, which we will discuss later on. Nevertheless, the perturbations are relatively small and the orbits of the satellites are still approximating ellipses, but as we will corroborate experimentally in Section 5 they are essential to obtain accurate results.

Consider a body confined in a compact region $\Omega \subseteq \mathbb{R}$ with a continuous density of mass $\rho: \Omega \to \mathbb{R}$. We would like to know the gravitational pull on a point mass m located at position \mathbf{r} from the center of mass of the body. To do this, consider a covering of Ω in a set of disjoint cubes Q_i , $i=1,\ldots,N$, small enough to be considered as point masses and let $R_i := Q_i \cap \Omega$. Then, $\bigsqcup_{i=1}^N R_i = \Omega$. If each R_i has mass m_i , volume V_i , density ρ_i , and its center is located at $\mathbf{s}_i \in \mathbb{R}^3$, then the total gravitational acceleration \mathbf{g} exerted on m is the sum of the contributions of all the forces exerted by the N point masses, and it is given by:

$$\mathbf{g} = -\sum_{i=0}^{N} \frac{m_i}{\|\mathbf{r} - \mathbf{s}_i\|^3} (\mathbf{r} - \mathbf{s}_i) = -\sum_{i=0}^{N} \frac{\rho_i}{\|\mathbf{r} - \mathbf{s}_i\|^3} (\mathbf{r} - \mathbf{s}_i) V_i$$
 (69)

Note that Eq. (69) is a Riemann sum and so letting $N \to \infty$ we get:

$$\mathbf{g} = -\int_{\Omega} \frac{\rho(\mathbf{s})}{\|\mathbf{r} - \mathbf{s}\|^3} (\mathbf{r} - \mathbf{s}) d^3 \mathbf{s}$$
 (70)

where $d^3\mathbf{s} := dx' dy' dz'$, if $\mathbf{s} = (x', y', z')$.

Theorem 26. Let Ω be a compact region in \mathbb{R}^3 with a continuous density of mass $\rho: \Omega \to \mathbb{R}$. Then, the gravitational acceleration field \mathbf{g} is conservative. That is, there exists a function $f: \mathbb{R}^3 \to \mathbb{R}$ such that $\mathbf{g} = \nabla f$.

Proof. An easy computation shows that fixed $\mathbf{s} \in \mathbb{R}^3$ we have:

$$\nabla \left(\frac{1}{\|\mathbf{r} - \mathbf{s}\|} \right) = -\frac{1}{\|\mathbf{r} - \mathbf{s}\|^3} (\mathbf{r} - \mathbf{s})$$
(71)

So we need to justify whether the following exchange between the gradient and the integral is correct:

$$\mathbf{g} = -\int_{\Omega} \frac{\rho(\mathbf{s})}{\|\mathbf{r} - \mathbf{s}\|^3} (\mathbf{r} - \mathbf{s}) d^3 \mathbf{s} = \int_{\Omega} \rho(\mathbf{s}) \mathbf{\nabla} \left(\frac{1}{\|\mathbf{r} - \mathbf{s}\|} \right) d^3 \mathbf{s} = \mathbf{\nabla} \int_{\Omega} \frac{\rho(\mathbf{s})}{\|\mathbf{r} - \mathbf{s}\|} d^3 \mathbf{s}$$
(72)

Without loss of generality it suffices to justify that

$$\frac{\partial}{\partial x} \int_{\Omega} \frac{\rho(\mathbf{s})}{\|\mathbf{r} - \mathbf{s}\|} d^3 \mathbf{s} = \int_{\Omega} \frac{\partial}{\partial x} \left(\frac{\rho(\mathbf{s})}{\|\mathbf{r} - \mathbf{s}\|} \right) d^3 \mathbf{s}$$
 (73)

assuming $\mathbf{r} = (x, y, z)$ and $\mathbf{s} = (x', y', z')$. In order to apply the theorem of derivation under the integral sign we need to control $\frac{\partial}{\partial x} \left(\frac{\rho(\mathbf{s})}{\|\mathbf{r} - \mathbf{s}\|} \right) = -\rho(\mathbf{s}) \frac{x - x'}{\|\mathbf{r} - \mathbf{s}\|^3}$ by an integrable function $h(\mathbf{s})$. Using spherical coordinates centered at \mathbf{r} and writing $(\mathbf{r} - \mathbf{s})_{\mathrm{sph}} = (\rho_{\mathbf{r}}, \theta, \phi)$, the integrand to bound becomes (in spherical coordinates):

$$\left| -\rho(\mathbf{s}) \frac{x - x'}{\|\mathbf{r} - \mathbf{s}\|^3} \rho_{\mathbf{r}}^2 \sin \phi \right| = |\rho(\mathbf{s})| \left| \frac{\rho_{\mathbf{r}} \cos \theta \sin \phi}{\rho_{\mathbf{r}}^3} \rho_{\mathbf{r}}^2 \sin \phi \right| \le |\rho(\mathbf{s})| \le K$$
 (74)

where the last inequality follows for certain $K \in \mathbb{R}$ by Weierstrass theorem (ρ is continuous and Ω is compact). Thus, since $h(\mathbf{s}) = K$ is integrable, because Ω is bounded, the equality of Eq. (73) is correct.

Physically speaking, the gravitational force ${\bf F}$ being conservative means that the work W done by the force along a path C

$$W = \int_{C} \mathbf{F} \cdot d\mathbf{s} \tag{75}$$

depends only on the initial and final positions of it. Moreover, due to historical reasons, we will write $\mathbf{g} = -\nabla V$ (with the minus sign) and call V the gravitational potential. The minus sign is chosen according the convention that work done by gravitational forces decreases the potential.

4.1.2 Laplace's equation for V

Theorem 27. Consider a distribution of matter of density ρ in a compact region Ω . Then, the gravitational potential V satisfies the Laplace equation

$$\Delta V = 0 \tag{76}$$

for all points outside Ω^8 .

Proof. Recall that $\Delta V = \operatorname{\mathbf{div}}(\nabla V)$. So since $\mathbf{g} = -\nabla V$ it suffices to prove that $\operatorname{\mathbf{div}}(\mathbf{g}) = 0$. Note that if $\mathbf{r} \in \Omega^c$, then $\exists \delta > 0$ such that $\|\mathbf{r} - \mathbf{s}\| \ge \delta > 0 \ \forall \mathbf{s} \in \Omega$ because Ω is closed. As a result, $\frac{\mathbf{r} - \mathbf{s}}{\|\mathbf{r} - \mathbf{s}\|^3}$ is differentiable and:

$$\mathbf{div}\left(\frac{\mathbf{r} - \mathbf{s}}{\|\mathbf{r} - \mathbf{s}\|^3}\right) = \frac{\partial}{\partial x} \left(\frac{x - x'}{\|\mathbf{r} - \mathbf{s}\|^3}\right) + \frac{\partial}{\partial y} \left(\frac{y - y'}{\|\mathbf{r} - \mathbf{s}\|^3}\right) + \frac{\partial}{\partial z} \left(\frac{z - z'}{\|\mathbf{r} - \mathbf{s}\|^3}\right) =$$

$$= \frac{\|\mathbf{r} - \mathbf{s}\|^2 - 3(x - x')^2}{\|\mathbf{r} - \mathbf{s}\|^5} + \frac{\|\mathbf{r} - \mathbf{s}\|^2 - 3(y - y')^2}{\|\mathbf{r} - \mathbf{s}\|^5} + \frac{\|\mathbf{r} - \mathbf{s}\|^2 - 3(z - z')^2}{\|\mathbf{r} - \mathbf{s}\|^5} = 0$$

Hence, as in Theorem 26, we have that for each $\mathbf{r} \in \Omega^c \exists \varepsilon, \delta > 0$ such that $\forall \tilde{\mathbf{r}} \in B(\mathbf{r}, \varepsilon) \subset \Omega^c$ we have:

$$\left| \rho(\mathbf{s}) \frac{\|\tilde{\mathbf{r}} - \mathbf{s}\|^2 - 3(\tilde{x} - x')^2}{\|\tilde{\mathbf{r}} - \mathbf{s}\|^5} \right| \le \frac{4|\rho(\mathbf{s})|}{\|\tilde{\mathbf{r}} - \mathbf{s}\|^3} \le \frac{4|\rho(\mathbf{s})|}{\delta^3}$$

which is integrable by Weierstrass theorem. Thus, by the theorem of derivation under the integral sign:

$$\mathbf{div}(\mathbf{g}) = -\mathbf{div} \int_{\Omega} \frac{\rho(\mathbf{s})}{\|\mathbf{r} - \mathbf{s}\|^3} (\mathbf{r} - \mathbf{s}) d^3 \mathbf{s} = -\int_{\Omega} \rho(\mathbf{s}) \, \mathbf{div} \left(\frac{\mathbf{r} - \mathbf{s}}{\|\mathbf{r} - \mathbf{s}\|^3} \right) d^3 \mathbf{s} = 0$$
 (77)

So far we have seen that the gravitational potential V satisfies the Laplace equation. If, moreover, we choose the origin of potential to be at the infinity, that is, if we impose $\lim_{\|\mathbf{r}\|\to\infty}V=0$, then the gravitational potential created by a distribution of mass in a compact region Ω is a solution of the following exterior

Dirichlet problem:

$$\begin{cases} \Delta V = 0 & \text{in } \Omega^c \\ V = f & \text{on } \partial \Omega \\ \lim_{\|\mathbf{r}\| \to \infty} V = 0 \end{cases}$$
 (78)

If Ω represents the Earth, then $f = f(\theta, \phi)$ is the boundary condition concerning the gravitational potential at the surface of the Earth as a function of the longitude θ and colatitude ϕ .

We will see now that Eq. (78) has a unique solution. To do that we invoke the maximum principle, which we will not prove here (see [Eva10] for more details).

⁸It can be seen that V satisfies in fact the Poisson equation $\Delta V = 4\pi G\rho$ for any point $\mathbf{r} \in \mathbb{R}^3$, which reduces to Laplace equation when $\mathbf{r} \in \Omega^c$, because there we have $\rho(\mathbf{r}) = 0$.

Theorem 28 (Maximum principle). Let $U \subset \mathbb{R}^n$ be open and bounded and $u \in \mathcal{C}^2(U) \cap \mathcal{C}(\overline{U})$. Suppose that u is harmonic within U, that is, $\Delta u = 0$ in U. Then, $\max_{\overline{U}} u = \max_{\partial U} u$.

Corollary 29. The Dirichlet problem of Eq. (78) has a unique solution.

Proof. Suppose we have two solutions V_1 , V_2 of Eq. (78). Then, $W := V_1 - V_2$ is harmonic in Ω^c , W = 0 on $\partial \Omega$ and $\lim_{\|\mathbf{r}\| \to \infty} W = 0$. So $\forall \varepsilon > 0$, $\exists n \in \mathbb{N}$ large enough such that $\Omega \subseteq B(0,n)$ and $|W| \le \varepsilon$ on $\mathbb{R}^3 \setminus \overline{B(0,n)}$. Thus, by the maximum principle, $|W| \le \varepsilon$ on $\overline{B(0,n)} \cap \Omega^c$. Since the ε is arbitrary, we must have W = 0 on Ω^c , that is, $V_1 = V_2$.

Now that we know that the Dirichlet problem has a unique solution, we can proceed to find it. In the next section we will construct an explicit solution for the gravitational potential of created by the Earth.

4.2 Spherical harmonics

4.2.1 Legendre polynomials, regularity and orthonormality

In this section we aim to introduce a class of functions that will appear later on in the general solution of the Laplace equation (see Section 4.2.2). To accomplish this, we need first to introduce the Legendre polynomials. There are several ways to define them, but the most convenient one for our purposes is from the following differential equation. Consider the following second-order differential equation called Legendre differential equation:

$$(1 - x^2)y'' - 2xy' + \lambda y = 0 (79)$$

for $\lambda \in \mathbb{R}$. This equation can be rewritten as:

$$((1 - x^2)y')' + \lambda y = 0 \tag{80}$$

Seeking for analytic solutions of this equation using the power series method [Mez], i.e. looking for solutions of the form $y(x) = \sum_{j=0}^{\infty} a_j x^j$, we see that:

$$0 = (1 - x^{2}) \sum_{j=0}^{\infty} a_{j+2}(j+1)(j+2)x^{j} - 2x \sum_{j=0}^{\infty} a_{j+1}(j+1)x^{j} + \lambda \sum_{j=0}^{\infty} a_{j}x^{j} = \sum_{j=0}^{\infty} a_{j+2}(j+1)(j+2)x^{j} - \sum_{j=0}^{\infty} a_{j}(j-1)jx^{j} - \sum_{j=0}^{\infty} 2a_{j}jx^{j} + \sum_{j=0}^{\infty} \lambda a_{j}x^{j} = \sum_{j=0}^{\infty} [a_{j+2}(j+1)(j+2) - a_{j}(j(j+1) - \lambda)]x^{j}$$
(81)

Equating the general term of the series to 0 we obtain this recursion:

$$a_{j+2} = \frac{j(j+1) - \lambda}{(j+1)(j+2)} a_j \quad j = 0, 1, 2, \dots$$
 (82)

From here we can obtain two independent solutions by setting the initial conditions a_0 and a_1 of the iteration. For example, setting $a_1 = 0$ we obtain a series that has only even powers of x. On the other hand, setting $a_0 = 0$ we obtain a series that has only odd powers of x. These two series converge on the interval (-1,1) by the ratio test (by looking at Eq. (82)) and can be expressed compactly as [Mez]:

$$y_{e}(x) = a_{0} \sum_{j=0}^{\infty} \left[\prod_{k=0}^{j-1} (2k(2k+1) - \lambda) \right] \frac{x^{2j}}{(2j)!} \quad y_{o}(x) = a_{1} \sum_{j=0}^{\infty} \left[\prod_{k=0}^{j-1} ((2k+1)(2k+2) - \lambda) \right] \frac{x^{2j+1}}{(2j+1)!}$$
(83)

Here, the empty product (that is, when k ranges from 0 to -1) is defined to be 1. However, for each $\lambda \in \mathbb{R}$ either one of these series or both diverge at $x = \pm 1$, as they behave as the harmonic series in a neighborhood of $x = \pm 1$. We are interested, though, in the solutions that remain bounded on the whole interval [-1,1]. Looking at the expressions of Eq. (83) one can check that the only possibility to make the series converge in [-1,1] is when $\lambda = n(n+1)$, $n \in \mathbb{N} \cup \{0\}$. In this case, for each $n \in \mathbb{N} \cup \{0\}$ exactly one of the series is in fact a polynomial of degree n. If, furthermore, we choose a_0 or a_1 be such that the polynomial evaluates to 1 at x = 1, these polynomials are called *Legendre polynomials*, and they are denoted by $P_n(x)$. The other (divergent) series is usually denoted in the literature by $Q_n(x)$ (check [RHB99; Mez]) and it is independent of $P_n(x)$. Thus, the general solution of Eq. (80) for $\lambda = n(n+1)$ can be expressed as a linear combination of P_n and Q_n , because the space of solutions form a vector space of dimension 2. The following proposition will be of our interest in the next section [RHB99].

n	$P_n(x)$
0	1
1	x
2	$\frac{\frac{1}{2}(3x^2 - 1)}{\frac{1}{2}(5x^3 - 3x)}$
3	$\frac{1}{2}(5x^3-3x)$
4	$\frac{1}{8}(35x^4-30x^2+3)$
5	$\frac{1}{9}(63x^5 - 70x^3 + 15x)$
6	$\frac{1}{16}(231x^6 - 315x^4 + 105x^2 - 5)$
7	$\frac{\frac{1}{16}(231x^6 - 315x^4 + 105x^2 - 5)}{\frac{1}{16}(429x^7 - 693x^5 + 315x^3 - 35x)}$

Table 3: First eight Legendre polynomials

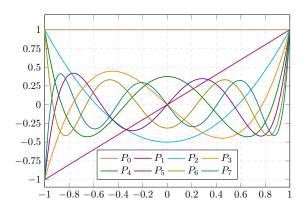


Figure 17: Graphic representation of the first eight Legendre polynomials.

Proposition 30. Let y(x) be a solution to the Legendre differential equation. Then, $\forall m \in \mathbb{N} \cup \{0\}$ the function

$$w_m(x) = (1 - x^2)^{m/2} \frac{\mathrm{d}^m y(x)}{\mathrm{d}x^m}$$
(84)

solves the general Legendre differential equation:

$$(1 - x^2)y'' - 2xy' + \left(\lambda - \frac{m^2}{1 - x^2}\right)y = 0$$
(85)

In particular if $\lambda = n(n+1)$ for $n \in \mathbb{N} \cup \{0\}$, then $w_m(x)$ is denoted as

$$P_{n,m}(x) := (1 - x^2)^{m/2} \frac{\mathrm{d}^m P_n}{\mathrm{d} x^m}$$
(86)

and it is called the associated Legendre polynomial of degree n and order m.

Note that although we opted to call the functions $P_{n,m}$ as polynomials, they are only true polynomials when m is even. But we have opted to call them in that manner as it is the common practice in the literature (see [Wei; RHB99; Mez]).

Moreover, from the definition of $P_{n,m}$, we can see that $P_{n,0} = P_n$ and that $P_{n,m} = 0$ if m > n. So we can restrict the domain of m to the set $\{0, 1, \ldots, n\}$.

n	$P_{n,1}(x)$	n	$P_{n,2}(x)$
1	$\sqrt{1-x^2}$	2	$3(1-x^2)$
2	$3x\sqrt{1-x^2}$	3	$15x(1-x^2)$
3	$\frac{3}{2}(5x^2-1)\sqrt{1-x^2}$	4	$\frac{15}{2}(7x^2-1)(1-x^2)$
4	$\frac{5}{2}x(7x^2-3)\sqrt{1-x^2}$	5	$\frac{105}{2}x(3x^2-1)(1-x^2)$
5	$\frac{15}{8}(21x^4 - 14x^2 + 1)\sqrt{1 - x^2}$	6	$\frac{105}{8}(33x^4 - 18x^2 + 1)(1 - x^2)$

Table 4: First associated Legendre polynomials for m = 1 and m = 2.

Definition 31. Let $n \in \mathbb{N} \cup \{0\}$ and $m \in \{0, 1, ..., n\}$. We define the real spherical harmonics $Y_{n,m}^{c}$ and $Y_{n,m}^{s}$ as:

$$Y_{n,m}^{c}(\theta,\phi) = \sqrt{(2 - \delta_{0,m})(2n+1)\frac{(n-m)!}{(n+m)!}} P_{n,m}(\cos\phi)\cos m\theta$$
 (87)

$$Y_{n,m}^{s}(\theta,\phi) = \sqrt{(2 - \delta_{0,m})(2n+1)\frac{(n-m)!}{(n+m)!}} P_{n,m}(\cos\phi)\sin m\theta$$
 (88)

Here the coordinates $\theta \in [0, 2\pi)$ and $\phi \in [0, \pi]$ are the spherical coordinates such that any point on the sphere can be written uniquely as $(x, y, z) = (\sin \phi \cos \theta, \sin \phi \sin \theta, \cos \phi)$. The factor $N_{n,m} := \sqrt{(2 - \delta_{0,m})(2n+1)\frac{(n-m)!}{(n+m)!}}$ is called the *normalization factor* of the spherical harmonics and $\delta_{0,m}$ is the Kronecker delta and will become clear in the next section.

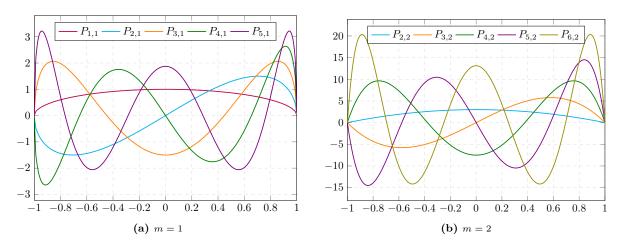


Figure 18: Graphic representation of the first associated Legendre polynomials for m=1 and m=2.

n	m	$Y_{n,m}^{\rm c}(\theta,\phi)$	n	m	$Y_{n,m}^{\mathrm{c}}(\theta,\phi)$
0	0	1	2	2	$\frac{\sqrt{15}}{2}(\sin\phi)^2\cos 2\theta$
1	0	$\sqrt{3}\cos\phi$	3	0	$\frac{\sqrt{7}}{2}\cos\phi(5(\cos\phi)^2-3)$
1	1	$\sqrt{3}\sin\phi\cos\theta$	3	1	$\frac{\sqrt{42}}{4} (5(\cos\phi)^2 - 1)\sin\phi\cos\theta$
2	0	$\frac{\sqrt{5}}{2}(3(\cos\phi)^2-1)$	3	2	$\frac{\sqrt{105}}{2}(\sin\phi)^2\cos\phi\cos2\theta$
2	1	$\sqrt{15}\sin\phi\cos\phi\cos\theta$	3	3	$\frac{\sqrt{70}}{4}(\sin\phi)^3\cos3\theta$

Table 5: First cosine spherical harmonics.

The associated Legendre polynomials satisfy an orthogonality relation:

Lemma 32. Let $n_1, n_2 \in \mathbb{N} \cup \{0\}$ and $m \leq \min\{n_1, n_2\}$. Then:

$$\int_{0}^{1} P_{n_{1},m}(x) P_{n_{2},m}(x) dx = \frac{2}{2n_{1}+1} \frac{(n_{1}+m)!}{(n_{1}-m)!} \delta_{n_{1},n_{2}}$$
(89)

where δ_{n_1,n_2} denotes the Kronecker delta.

Similarly, it can be shown that the spherical harmonics from an orthonormal family of functions:

Proposition 33. The family of spherical harmonics $\{Y_{n,m}^{c}(\theta,\phi),Y_{n,m}^{s}(\theta,\phi):n\in\mathbb{N}\cup\{0\},m\leq n\}$ is orthonormal in the following sense:

$$\frac{1}{4\pi} \int_{0}^{2\pi} \int_{0}^{\pi} Y_{n_1,m_1}^i(\theta,\phi) Y_{n_2,m_2}^j(\theta,\phi) d\Omega = \delta_{n_1,n_2} \delta_{m_1,m_2} \delta_{i,j}$$
(90)

where $d\Omega = \sin \phi \, d\phi \, d\theta$ is the solid angle element, which measures the element of area on a sphere of radius 1.

Proof. Let N_{n_1,m_1} , N_{n_2,m_2} be the normalization factors of the spherical harmonics Y_{n_1,m_1} , Y_{n_2,m_2} respectively. Note that we can separate the variables in the integral of Eq. (90). So if $i \neq j$, the integral over θ becomes $\int_0^{2\pi} \cos(m_1\theta) \sin(m_2\theta) \, d\theta$ which is equal to 0 regardless of the values of m_1 and m_2 . So from now on assume that i = j. Due to the symmetry between the cosine and the sine we can suppose that i = c. Thus:

$$\int_{0}^{2\pi} \int_{0}^{\pi} Y_{n_{1},m_{1}}^{i}(\theta,\phi) Y_{n_{2},m_{2}}^{j}(\theta,\phi) d\Omega =
= N_{n_{1},m_{1}} N_{n_{2},m_{2}} \int_{0}^{\pi} P_{n_{1},m_{1}}(\cos\phi) P_{n_{2},m_{2}}(\cos\phi) \sin\phi d\phi \int_{0}^{2\pi} \cos(m_{1}\theta) \cos(m_{2}\theta) d\theta \quad (91)$$

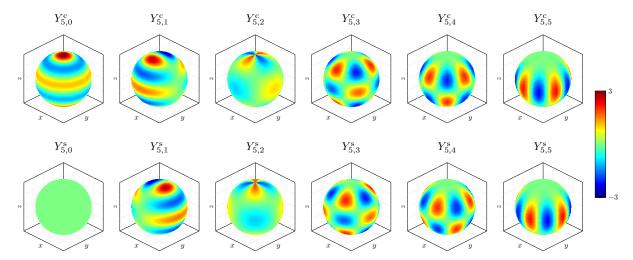


Figure 19: 3D heat map of the spherical harmonics of degree n = 5. The first row correspond to the cosine spherical harmonics and the second row correspond to the sine spherical harmonics.

An easy check shows that if $m_1 \neq m_2$ then the integral over θ is zero (and the same applies with sines). So suppose $m_1 = m_2 = m$. In that case, if $m \neq 0$ we have $\int_0^{2\pi} (\cos m\theta)^2 d\theta = \int_0^{2\pi} (\sin m\theta)^2 d\theta = \pi$ and if m = 0, the cosine integral evaluates to 2π whereas the sine integral is 0. We can omit this latter case because $Y_{n,0}^s$ is identically zero. Thus:

$$\frac{2\pi}{2 - \delta_{0,m}} N_{n_1,m} N_{n_2,m} \int_{0}^{\pi} P_{n_1,m}(\cos\phi) P_{n_2,m}(\cos\phi) \sin\phi d\phi = \frac{2\pi}{2 - \delta_{0,m}} N_{n_1,m} N_{n_2,m} \int_{-1}^{1} P_{n_1,m}(x) P_{n_2,m}(x) dx$$
(92)

By Lemma 32 this latter integral is $\frac{2}{2n_1+1}\frac{(n_1+m)!}{(n_1-m)!}\delta_{n_1,n_2}$. Finally, if $n_1=n_2=n$, putting all normalization factors together we get:

$$\frac{2\pi}{2 - \delta_{0,m}} N_n^m N_n^m \frac{2}{2n+1} \frac{(n+m)!}{(n-m)!} = 4\pi$$
(93)

Moreover, an important result in the Sturm-Liouville Theory of second order differential equations ([Wikd; Wan+09]) says that the family of spherical harmonics $\{Y_{n,m}^c(\theta,\phi),Y_{n,m}^s(\theta,\phi):n\in\mathbb{N}\cup\{0\},m\leq n\}$ form a complete set in the sense that any smooth function defined on the sphere $f:\mathbf{S}^2\to\mathbb{R}$ can be expanded in a series of spherical harmonics:

$$f(\theta,\phi) = \sum_{n=0}^{\infty} \sum_{m=0}^{n} (c_{n,m} Y_{n,m}^{c}(\theta,\phi) + s_{n,m} Y_{n,m}^{s}(\theta,\phi))$$
(94)

This will be useful in Section 4.2.3 when expanding the gravitational potential created by the Earth at some arbitrary point in spherical harmonics.

4.2.2 Laplace's equation in spherical coordinates

Definition 34. Let $f: \mathbb{R}^3 \to \mathbb{R}$ be a twice-differentiable function. The Laplace equation is the equation

$$\Delta f = \frac{\partial^2 f}{\partial x^2} + \frac{\partial^2 f}{\partial y^2} + \frac{\partial^2 f}{\partial z^2} = 0 \tag{95}$$

where Δ is the Laplace operator.

The next proposition gives the Laplace equation in spherical coordinates.

Proposition 35. Let $f: \mathbb{R}^3 \to \mathbb{R}$ be a twice-differentiable function. Then:

$$\Delta f = \frac{1}{r^2} \frac{\partial}{\partial r} \left(r^2 \frac{\partial f}{\partial r} \right) + \frac{1}{r^2 \sin \phi} \frac{\partial}{\partial \phi} \left(\sin \phi \frac{\partial f}{\partial \phi} \right) + \frac{1}{r^2 (\sin \phi)^2} \frac{\partial^2 f}{\partial \theta^2}$$
 (96)

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where $r \in [0, \infty)$ denotes the radial distance, $\theta \in [-\pi, \pi)$ denotes the longitude, and $\phi \in [0, \pi]$, the

$$x = r \sin \phi \cos \theta$$

$$y = r \sin \phi \sin \theta$$

$$z = r \cos \phi$$
(97)

We are now interested in solving the Laplace equation. Theorem 36 gives the solution of it as a function of the spherical harmonics.

Theorem 36. The regular solutions in a bounded region $\Omega \subseteq \mathbb{R}^3$ such that $0 \notin \overline{\Omega}$ of the Laplace equation in spherical coordinates are of the form

$$f(r,\theta,\phi) = \sum_{n=0}^{\infty} \sum_{m=0}^{n} (a_n r^n + b_n r^{-n-1}) P_{n,m}(\cos\phi) (c_{n,m}\cos(m\theta) + s_{n,m}\sin(m\theta))$$
(98)

$$= \sum_{n=0}^{\infty} \sum_{m=0}^{n} (a_n r^n + b_n r^{-n-1}) (\tilde{c}_{n,m} Y_{n,m}^{c}(\theta, \phi) + \tilde{s}_{n,m} Y_{n,m}^{s}(\theta, \phi))$$
(99)

where $a_n, b_n, c_{n,m}, s_{n,m}, \tilde{c}_{n,m}, \tilde{s}_{n,m} \in \mathbb{R}$.

Proof. Let $f(r,\theta,\phi)$ be a solution of Eq. (96). Using separation variables $f(r,\theta,\phi) = R(r)\Theta(\theta)\Phi(\phi)$ we can write:

$$\frac{\Theta\Phi}{r^2}(r^2R')' + \frac{R\Theta}{r^2\sin\phi}(\sin\phi\Phi')' + \frac{R\Phi}{r^2(\sin\phi)^2}\Theta'' = 0$$
(100)

Here, we are making and abuse of notation denoting all the derivatives with a prime, but the reader should have no confusion with it. Isolating R from Θ and Φ yields:

$$\frac{\left(r^2R'\right)'}{R} = -\frac{1}{\sin\phi\Phi}\left(\sin\phi\Phi'\right)' - \frac{1}{\left(\sin\phi\right)^2\Theta}\Theta''$$
(101)

Since the left-hand side depends entirely on r and the right-hand side does not, it follows that both sides must be constant. Therefore:

$$\frac{\left(r^2R'\right)'}{R} = \lambda \tag{102}$$

$$\frac{(r^2R')'}{R} = \lambda$$

$$\frac{1}{\sin\phi\Phi}(\sin\phi\Phi')' + \frac{1}{(\sin\phi)^2\Theta}\Theta'' = -\lambda$$
(102)

with $\lambda \in \mathbb{R}$. Similarly, separating variables from Eq. (103) we obtain that the equations

$$\frac{1}{\Theta}\Theta'' = -m^2 \tag{104}$$

$$\frac{\sin\phi}{\Phi}(\sin\phi\Phi')' + \lambda(\sin\phi)^2 = m^2 \tag{105}$$

must be constant with $m \in \mathbb{C}$ (a priori). The solution to the well-known Eq. (104) is a linear combination of the $\cos(m\theta)$ and $\sin(m\theta)$. Note, though, that since Θ must be a 2π -periodic function, that is satisfying $\Theta(\theta+2\pi)=\Theta(\theta)\ \forall \theta\in\mathbb{R},\ m$ must be an integer. On the other hand making the change of variables $x = \cos \phi$ and $y = \Phi(\phi)$ in Eq. (105) and using the chain rule, that equation becomes:

$$(1 - x^2) \frac{\mathrm{d}^2 y}{\mathrm{d}x^2} - 2x \frac{\mathrm{d}^2 y}{\mathrm{d}x^2} + \left(\lambda - \frac{m^2}{1 - x^2}\right) y = 0$$
 (106)

which is the associate Legendre equation. We have argued in Proposition 30 that we need $\lambda = n(n+1)$ and $m \leq n$ in order to obtain regular solutions at $x = \cos \phi = \pm 1$. Moreover, these solutions are $P_{n,m}(\cos\phi).$

Finally, note that equation Eq. (102) is a Cauchy-Euler equation (check [Wika]) and so the general solution of it is given by

$$R(r) = c_1 r^n + c_2 r^{-n-1} (107)$$

because $\lambda = n(n+1)$ (the reader may check that r^n and r^{-n-1} are indeed two independent solutions of Eq. (102)). So the general solution becomes a linear combination of each solution found varying $n \in \mathbb{N} \cup \{0\}$ and $m \in \{0, 1, \dots, n\}$:

$$f(r,\theta,\phi) = \sum_{n=0}^{\infty} \sum_{m=0}^{n} (a_n r^n + b_n r^{-n-1}) P_{n,m}(\cos\phi) (c_{n,m}\cos(m\theta) + s_{n,m}\sin(m\theta))$$
(108)

We ignore the singularity at r = 0 of Eq. (99) from now. See Section 4.2.3 for more details.

4.2.3 Expansion in spherical harmonics

We have just seen that if V satisfies the exterior Dirichlet problem for the Laplace equation, then, by uniqueness of solutions, it can be expressed as:

$$V(r,\theta,\phi) = \sum_{n=0}^{\infty} \sum_{m=0}^{n} (a_n r^n + b_n r^{-n-1}) (\tilde{c}_{n,m} Y_{n,m}^{c}(\theta,\phi) + \tilde{s}_{n,m} Y_{n,m}^{s}(\theta,\phi))$$
(109)

where $a_n, b_n, \tilde{c}_{n,m}, \tilde{s}_{n,m} \in \mathbb{R}$. If we impose V to satisfy the third condition of Eq. (78), we must have $a_n = 0$. Finally, if we choose R_{\oplus} as a reference radius for a spherical model of the Earth, using the boundary condition on $\partial \Omega$

$$f(\theta,\phi) = \sum_{n=0}^{\infty} \sum_{m=-n}^{n} \frac{b_n}{R_{\oplus}^{n+1}} (\tilde{c}_{n,m} Y_{n,m}^{c}(\theta,\phi) + \tilde{s}_{n,m} Y_{n,m}^{s}(\theta,\phi))$$
(110)

and the orthogonality of the spherical harmonics, we can deduce that the coefficients $b_n \tilde{c}_{n,m}$ and $b_n \tilde{s}_{n,m}$ are given by:

$$b_n \tilde{c}_{n,m} = \frac{R_{\oplus}^{n+1}}{4\pi} \int_0^{2\pi} \int_0^{\pi} f(\theta, \phi) Y_{n,m}^{c}(\theta, \phi) \sin \phi \, d\phi \, d\theta$$
 (111)

$$b_n \tilde{s}_{n,m} = \frac{R_{\oplus}^{n+1}}{4\pi} \int_0^{2\pi} \int_0^{\pi} f(\theta, \phi) Y_{n,m}^{s}(\theta, \phi) \sin \phi \, d\phi \, d\theta$$
 (112)

Hence, introducing the gravitational constant G and the Earth's mass M_{\oplus} into the equation, our final expression for the gravitational potential is

$$V(r,\theta,\phi) = \frac{GM_{\oplus}}{R_{\oplus}} \sum_{n=0}^{\infty} \sum_{m=0}^{n} \left(\frac{R_{\oplus}}{r}\right)^{n+1} (\bar{C}_{n,m} Y_{n,m}^{c}(\theta,\phi) + \bar{S}_{n,m} Y_{n,m}^{s}(\theta,\phi))$$
(113)

where the coefficients $\bar{C}_{n,m}, \bar{S}_{n,m} \in \mathbb{R}$ are given by the formulas:

$$\bar{C}_{n,m} = \frac{1}{4\pi} \frac{R_{\oplus}}{GM_{\oplus}} \int_{0}^{2\pi} \int_{0}^{\pi} f(\theta, \phi) Y_{n,m}^{c}(\theta, \phi) \sin \phi \, d\phi \, d\theta$$
 (114)

$$\bar{S}_{n,m} = \frac{1}{4\pi} \frac{R_{\oplus}}{GM_{\oplus}} \int_{0}^{2\pi} \int_{0}^{\pi} f(\theta, \phi) Y_{n,m}^{s}(\theta, \phi) \sin \phi \, d\phi \, d\theta$$
 (115)

The coefficients $\bar{C}_{n,m}$, $\bar{S}_{n,m}$ are called geopotential coefficients, and they describe the dependence on the Earth's internal structure. They are obtained from observation of the perturbations seen in the orbits of other satellites [MG05], because it is not possible to measure the Earth's density directly. Other methods for obtaining such data are through surface gravimetry, which provides precise local and regional information about the gravity field, or through altimeter data, which can be used to provide a model for the geoid of the Earth (which is the shape that the ocean surface takes under the influence of the gravity of Earth) which in turn can be used to obtain the geopotential coefficients.

4.3Numerical computation

Gravitational acceleration

Up to this point, we have only studied the gravitational potential exerted by the non-homogeneous Earth on a satellite. But, in order to integrate the equations of motion of the satellite, we need to compute the gravitational acceleration $\mathbf{g} = -\nabla V$ instead. In order to do this efficiently, we will make use of the following formulas given in [MG05; Cun70]. First, let

$$V_{n,m}(\theta,\phi) = \left(\frac{R_{\oplus}}{r}\right)^{n+1} P_{n,m}(\cos\phi)\cos(m\theta) \qquad W_{n,m}(\theta,\phi) = \left(\frac{R_{\oplus}}{r}\right)^{n+1} P_{n,m}(\cos\phi)\sin(m\theta) \qquad (116)$$

Thus, we can write:

$$V = \frac{GM_{\oplus}}{R_{\oplus}} \sum_{n=0}^{\infty} \sum_{m=0}^{n} (\bar{C}_{n,m} N_{n,m} V_{n,m} + \bar{S}_{n,m} N_{n,m} W_{n,m})$$
(117)

Let $C_{n,m} := \overline{C}_{n,m} N_{n,m}$ and $S_{n,m} := \overline{S}_{n,m} N_{n,m}$. If $\mathbf{g} = (\ddot{x}, \ddot{y}, \ddot{z})$, then:

$$\ddot{x} = \sum_{n=0}^{\infty} \sum_{m=0}^{n} \ddot{x}_{n,m} \qquad \ddot{y} = \sum_{n=0}^{\infty} \sum_{m=0}^{n} \ddot{y}_{n,m} \qquad \ddot{z} = \sum_{n=0}^{\infty} \sum_{m=0}^{n} \ddot{z}_{n,m}$$
(118)

where the partial accelerations $\ddot{x}_{n,m}$, $\ddot{y}_{n,m}$, $\ddot{z}_{n,m}$ are given by:

$$\ddot{x}_{n,m} = \begin{cases}
-\frac{GM_{\oplus}}{R_{\oplus}^{2}} C_{n,0} V_{n+1,1} & \text{if } m = 0 \\
-\frac{GM_{\oplus}}{R_{\oplus}^{2}} \cdot \frac{1}{2} \left[C_{n,m} V_{n+1,m+1} + S_{n,m} W_{n+1,m+1} - \frac{(n-m+2)!}{(n-m)!} \left(C_{n,m} V_{n+1,m-1} + S_{n,m} W_{n+1,m-1} \right) \right] & \text{if } m > 0 \end{cases}$$
(119)

$$\ddot{y}_{n,m} = \begin{cases}
-\frac{GM_{\oplus}}{R_{\oplus}^{2}} C_{n,0} W_{n+1,1} & \text{if } m = 0 \\
-\frac{GM_{\oplus}}{R_{\oplus}^{2}} \cdot \frac{1}{2} \left[C_{n,m} W_{n+1,m+1} - S_{n,m} V_{n+1,m+1} - \frac{(n-m+2)!}{(n-m)!} (C_{n,m} W_{n+1,m-1} - S_{n,m} V_{n+1,m-1}) \right] & \text{if } m > 0 \\
\ddot{z}_{n,m} = -\frac{GM_{\oplus}}{R_{\oplus}^{2}} (n-m+1) (C_{n,m} V_{n+1,m} + S_{n,m} W_{n+1,m})
\end{cases} (120)$$

$$\ddot{z}_{n,m} = -\frac{GM_{\oplus}}{R_{\oplus}^{2}}(n-m+1)\left(C_{n,m}V_{n+1,m} + S_{n,m}W_{n+1,m}\right)$$
(121)

and the functions $V_{n,m}$, $W_{n,m}$ are calculated using the following recurrence relations:

$$\begin{cases} V_{n,m} = \frac{2n-1}{n-m} \frac{R_{\oplus}}{r} \cos \phi V_{n-1,m} - \frac{n+m-1}{n-m} \frac{R_{\oplus}^2}{r^2} V_{n-2,m} \\ W_{n,m} = \frac{2n-1}{n-m} \frac{R_{\oplus}}{r} \cos \phi W_{n-1,m} - \frac{n+m-1}{n-m} \frac{R_{\oplus}^2}{r^2} W_{n-2,m} \end{cases} & \text{if } 0 \le m \le n-2 \\ \begin{cases} V_{n,n-1} = (2n-1) \frac{R_{\oplus}}{r} \cos \phi V_{n-1,n-1} \\ W_{n,n-1} = (2n-1) \frac{R_{\oplus}}{r} \cos \phi W_{n-1,n-1} \end{cases} & \text{if } m = n-1 \\ \begin{cases} V_{n,n} = (2m-1) \frac{R_{\oplus}}{r} \sin \phi [\cos \theta V_{n-1,n-1} - \sin \theta W_{n-1,n-1}] \\ W_{n,n} = (2m-1) \frac{R_{\oplus}}{r} \sin \phi [\cos \theta W_{n-1,n-1} + \sin \theta V_{n-1,n-1}] \end{cases} & \text{if } m = n \end{cases}$$

starting from the initial quantities $V_{00} = \frac{R_{\oplus}}{r}$ and $W_{00} = 0$ and using the following scheme [MG05]:

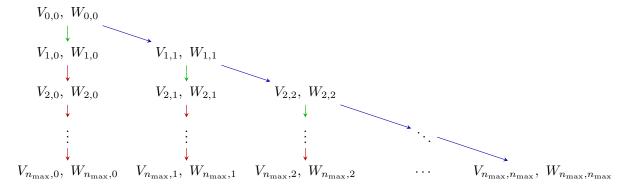


Figure 20: Scheme for the calculation of the coefficients $V_{n,m}$ and $W_{n,m}$ for $0 \le m \le n \le n_{\text{max}}$. The red arrows indicate that the coefficients are calculated using the first of the above recursions; the green arrows indicate that they are calculated using the second recursion; and the blue arrows indicate that they are calculated using the third recursion.

4.3.2 Position and velocity

With all the ingredients in place, the position and velocity of the satellite can be calculated from the following system of differential equations:

$$\ddot{\mathbf{r}} = \mathbf{g} \implies \begin{cases} x' = v_x \\ y' = v_y \\ z' = v_z \\ v_x' = \ddot{x} \\ v_y' = \ddot{y} \\ v_z' = \ddot{z} \end{cases}$$

where \ddot{x} , \ddot{y} , \ddot{z} are given in Eq. (118). The initial conditions will be the position and velocity obtained from the TLE.

In order to solve this system of differential equations, we have opted to use the Runge-Kutta-Fehlberg of order 78.

4.4 Other perturbations

atmospheric drag, solar radiation pressure, and the gravitational pull of the Moon and Sun, POSAR-HO SI AL FINAL FAIG SIMULACIO AMB AIXO, SI NO, NOOOO.

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References

- [Ag217] Ag2gaeh. Elliko-sk. Accessed: June 04, 2023. Mar. 2017. URL:
- [Aok+81] S. Aoki et al. "The new definition of Universal Time." In: Astronomy and Astrophysics 105 (Dec. 1981), pp. 359–361.
- [Cun70] Leland E. Cunningham. "On the computation of the spherical harmonic terms needed during the numerical integration of the orbital motion of an artificial satellite." In: Celestial mechanics 2 (Aug. 1970), pp. 207–216.
- [Eva10] Lawrence C. Evans. *Partial Differential Equations*. Second Edition. Americal Mathematical Society, 2010. ISBN: 978-0-8218-4974-3.
- [GPS02] H. Goldstein, C.P. Poole, and J. Safko. Classical Mechanics. Pearson, 2002. ISBN: 978-81-317-5891-5.
- [Kel] T.S. Kelso. Frequently Asked Questions: Two-Line Element Set Format. Accessed: April 29, 2023. URL:
- [Las07] Lasunncty. Orbit1. Accessed: June 04, 2023. Oct. 2007. URL:
- [Lie+77] J. H. Lieske et al. "Expressions for the Precession Quantities Based upon the IAU (1976) System of Astronomical Constants." In: Astronomy and Astrophysics 58 (June 1977), pp. 1–16.
- [Mat12] Magister Mathematicae. Conic Sections. Accessed: April 29, 2023. Mar. 2012. URL: .
- [Mez] Hamid Meziani. Legendre Polynomials and Applications. Accessed: April 7, 2023. URL: ...
- [MG05] Oliver Montenbruck and Eberhard Gill. Satellite Orbits: Models, Methods, and Applications. Springer, 2005. ISBN: 978-3-540-67280-7.
- [Obsa] U.S. Naval Observatory. *IERS Rapid Service / Prediction Center*. Accessed: April 17, 2023. URL: ☑.
- [Obsb] U.S. Naval Observatory. *IERS Rapid Service / Prediction Center Daily Earth Orientation Parameters Solutions*. Accessed: June 9, 2023. URL: 🗹.
- [RHB99] K.F. Riley, M.P. Hobson, and S.J. Bence. *Mathematical Methods for Physics and Engineering:* A Comprehensive Guide. Cambridge University Press, 1999. ISBN: 978-0-521-86153-3.
- [RS98a] International Earth Rotation and Reference Systems Service. *IERS Bulletin A*. Accessed: April 17, 2023. July 1998. URL: □.
- [RS98b] International Earth Rotation and Reference Systems Service. *IERS Bulletin C.* Accessed: April 17, 2023. July 1998. URL: □.
- [Val13] David A. Vallado. Fundamentals of Astrodynamics and Applications. Ed. by Wayne D. McClain. The Space Technology Library, 2013. ISBN: 978-1-881-88318-0.
- [Wan+09] Aiping Wang et al. "Completeness of Eigenfunctions of Sturm-Liouville Problems with Transmission Conditions." In: *Methods Appl. Anal.* 16 (Sept. 2009). DOI: □.
- [Wei] Eric W. Weisstein. Associated Legendre Polynomial. From MathWorld—A Wolfram Web Resource. Accessed: April 7, 2023. URL:
- [Wika] The Free Encyclopedia Wikipedia. Cauchy-Euler equation. Accessed: April 7, 2023. URL:
- [Wikb] The Free Encyclopedia Wikipedia. Earth-centered inertial. Accessed: April 14, 2023. URL: 🗹.
- [Wikc] The Free Encyclopedia Wikipedia. Simplified perturbations models. Accessed: April 29, 2023. URL: ☑.
- [Wikd] The Free Encyclopedia Wikipedia. Sturm-Liouville theory. Accessed: April 7, 2023. URL: 🗹.