# The period function of reversible quadratic centers ${ }^{\tau}$ 

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#### Abstract

In this paper we investigate the bifurcation diagram of the period function associated to a family of reversible quadratic centers, namely the dehomogenized Loud's systems. The local bifurcation diagram of the period function at the center is fully understood using the results of Chicone and Jacobs [Bifurcation of critical periods for plane vector fields, Trans. Amer. Math. Soc. 312 (1989) 433-486]. Most of the present paper deals with the local bifurcation diagram at the polycycle that bounds the period annulus of the center. The techniques that we use here are different from the ones in [C. Chicone, M. Jacobs, Bifurcation of critical periods for plane vector fields, Trans. Amer. Math. Soc. 312 (1989) 433-486] because, while the period function extends analytically at the center, it has no smooth extension to the polycycle. At best one can hope that it has some asymptotic expansion. Another major difficulty is that the asymptotic development has to be uniform with respect to the parameters, in order to prove that a parameter is not a bifurcation value. We study also the bifurcations in the interior of the period annulus and we show that there exist three germs of curves in the parameter


[^0]space that correspond to this type of bifurcation. Moreover we determine some regions in the parameter space for which the corresponding period function has at least one or two critical periods. Finally we propose a complete conjectural bifurcation diagram of the period function of the dehomogenized Loud's systems. Our results can also be viewed as a contribution to the proof of Chicone's conjecture [C. Chicone, review in MathSciNet, ref. 94h:58072]. © 2005 Elsevier Inc. All rights reserved.

## 1. Introduction, main results and conjectures

In this work we study the bifurcation diagram of the period function associated to a family of quadratic centers. Chicone [2] has conjectured that if a quadratic system has a center with a period function which is not monotonic then, by an affine transformation and a constant rescaling of time, it can be brought to the Loud normal form

$$
\left\{\begin{array}{l}
\dot{x}=-y+B x y  \tag{1}\\
\dot{y}=x+D x^{2}+F y^{2}
\end{array}\right.
$$

and that the period function of these centers has at most two critical periods. In fact, there is much analytic evidence that the conjecture is true (see [5,15,20] for instance). On the other hand, it is proved in [7] that if $B=0$ then the period function of the center at the origin of system (1) is globally monotonous. So, from the point of view of the study of the period function, the most interesting stratum of quadratic centers is the family (1) with $B \neq 0$, which can be brought to $B=1$ by means of a rescaling, i.e., to

$$
\left\{\begin{array}{l}
\dot{x}=-y+x y  \tag{2}\\
\dot{y}=x+D x^{2}+F y^{2}
\end{array}\right.
$$

This is precisely the family of quadratic centers that we study in this paper and, following the terminology in [4], we call them dehomogenized Loud's systems. Note that the transformation $(x, y, t) \longmapsto(x,-y,-t)$ preserves the Loud normal form (1) and so these systems are reversible with respect to the $x$-axis. In fact any reversible quadratic center can be brought to Loud normal form by an affine transformation and a constant rescaling of time (see [22] for instance). Compactifying $\mathbb{R}^{2}$ to the Poincaré disc, the boundary of the period annulus of the center has two connected components, the center itself and a polycycle. We call them, respectively, the inner and outer boundary of the period annulus. It follows (see Lemma 2.7) that the bifurcation diagram of the period function consists of three parts:
(a) Bifurcations of the period function at the inner boundary (i.e., the center).
(b) Bifurcations of the period function at the outer boundary (i.e., the polycycle).
(c) Bifurcations of the period function in the interior of the period annulus.

For the precise definitions see Section 2. The local bifurcation diagram of the period function at the inner boundary is fully understood for the quadratic centers using the


Fig. 1. Numerical bifurcation diagram by Chicone and Jacobs.
results of Chicone and Jacobs [4]. The key point in their result (see Theorem 4.1) is that the period function extends analytically to the inner boundary because the center is nondegenerate. It thus admits a Taylor expansion whose coefficients are polynomials on the parameters of the initial system. For system (2), the zero level set of the first coefficient is an ellipse $\Gamma_{\mathrm{C}}$ given by $10 D^{2}+10 D F-D+4 F^{2}-5 F+1=0$ (see Fig. 11) and it corresponds to the local bifurcation diagram of the period function at the inner boundary. Around 1990, Chicone and Jacobs obtained a numerical computation of the complete bifurcation diagram, see Fig. 1, in which one can easily identify the ellipse $\Gamma_{\mathrm{C}}$. It also appears a strange kidney-shaped curve that would correspond to a numerical approximation of the local bifurcation diagram of the period function at the outer boundary. The main goal of this paper is to determine analytically this set.

One encounters two major difficulties in the study of the bifurcation diagram of the period function at the outer boundary of the period annulus. The first one is that, contrary to the situation in the inner boundary, the period function does not extend smoothly on the outer boundary. At best one can hope that it has some asymptotic development. The second one is that in order to prove that a parameter is not a bifurcation value one needs an asymptotic development which is uniform with respect to the parameters. This is not easily achieved because the shape of the polycycle in the outer boundary changes as the parameters vary.


Fig. 2. Bifurcation diagram of the period function at the outer boundary.

In order to formulate our main result let $\Gamma_{\mathrm{U}}$ be the union of dotted straight lines in Fig. 2. Consider also the bold curve $\Gamma_{\mathrm{B}}$. (Here the subscripts B and U stand for bifurcation and unspecified, respectively.) The curve $\Gamma_{\mathrm{U}}$ corresponds, except for the segment $(-1,-1 / 2) \times\{1 / 2\}$, to bifurcations of the phase portrait that affect the outer boundary of the period annulus (see Section 3.1). The curve $\Gamma_{\mathrm{B}}$ is the union of some explicit straight segments and a curve that joins the points $(-3 / 2,3 / 2)$ and $(-1 / 2,1)$. To be more precise, let us advance that this curve is defined as the zero level set of an explicit function that we introduce in Section 3.2.1. To draw it in Fig. 2 we have computed it numerically. Analytically, among other properties that are gathered in Proposition 3.11, we have proved that it is the graphic of an analytic function $D=\mathcal{G}(F)$. From Proposition 3.11 it follows in particular that $\Gamma_{\mathrm{B}}$ is a Jordan curve. We can consider therefore the bounded and unbounded components of $\mathbb{R}^{2} \backslash \Gamma_{\mathrm{B}}$, which we denote by $\mathcal{D}_{\mathrm{B}}$ and $\mathcal{I}_{\mathrm{B}}$ (for decreasing and increasing), respectively. With this notation we can now state our main result:

Theorem A. Denoting $\mu=(D, F)$, let $\left\{X_{\mu}, \mu \in \mathbb{R}^{2}\right\}$ be the family of vector fields in (2) and consider the period function of the center at the origin. Then the open set $\mathbb{R}^{2} \backslash\left\{\Gamma_{\mathrm{B}} \cup \Gamma_{\mathrm{U}}\right\}$ corresponds to local regular values of the period function at the outer boundary of the period annulus. In addition,
(a) If $\mu_{0} \in \mathcal{I}_{\mathrm{B}} \backslash \Gamma_{\mathrm{U}}$ then the period function of $X_{\mu_{0}}$ is monotonous increasing near the outer boundary.
(b) If $\mu_{0} \in \mathcal{D}_{\mathrm{B}} \backslash \Gamma_{\mathrm{U}}$ then the period function of $X_{\mu_{0}}$ is monotonous decreasing near the outer boundary.


Fig. 3. Conjectural bifurcation diagram of the period function.

Finally, the parameters in $\Gamma_{\mathrm{B}}$ are local bifurcation values of the period function at the outer boundary of the period annulus.

We have not determined the character of the parameters in $\Gamma_{\mathrm{U}}$. We conjecture that they are not bifurcation values at the outer boundary except for the segment $\{0\} \times$ $[0,1 / 2]$. The numerical picture in Fig. 1 fits relatively well with the bifurcation curve $\Gamma_{\mathrm{B}}$. There are, however, some striking differences. In particular, unlike the numerical picture, most parts of the bifurcation curve are straight segments.

The combination of the results of Chicone and Jacobs in [4] with the ones that we obtain in the present paper lead us to formulate the following conjecture about the complete bifurcation diagram (see Fig. 3) of the period function of the dehomogenized Loud's systems:

Conjecture. The bifurcation diagram of the period function of the dehomogenized Loud's family (2) consists in the union of the following curves:
(a) The ellipse $\Gamma_{\mathrm{C}}$, which corresponds to the local bifurcation values at the inner boundary.
(b) The Jordan curve $\Gamma_{\mathrm{B}}$ given in Theorem A together with the segment $\{0\} \times[0,1 / 2]$, which corresponds to the local bifurcation values at the outer boundary.
(c) Three simple curves $\delta_{1}, \delta_{2}$ and $\delta_{3}$ that connect $L_{1}$ with $(-2,2), L_{2}$ with $(0,0)$ and $L_{3}$ with $(-3 / 2,3 / 2)$, respectively, which correspond to the local bifurcation values in the interior.

The exact values of the parameters $L_{i}$ mentioned above are given in (36) and, following the terminology in [4], they correspond to the three weak centers of order two of system (2). In fact we prove in Theorem 4.3 the existence of the germs at $L_{i}$ of the conjectured curves $\delta_{i}$. We think of course that the period function of the dehomogenized Loud's systems has at most two critical periods. Fig. 3 shows the regions where we conjecture that there are 0,1 and 2 critical periods. We also sketch in balloons the changes in the monotonicity of the period function. Let us note finally that in the segment $\{0\} \times[0,1 / 2]$ occur two different types of bifurcation at the outer boundary. Indeed, crossing from left to right the segment $\{0\} \times[0,1 / 4]$ corresponds to the "disappearance" of two critical periods, while crossing $\{0\} \times[1 / 4,1 / 2]$ corresponds to a "rebound" of a critical period (see Fig. 14).

The paper is organized in the following way. In Section 2 we introduce the precise definitions that we shall use. Section 3 is devoted to the proof of Theorem A. In all the cases that we study (see Fig. 4), the polycycle in the outer boundary of the period annulus has one or two singular points, which are saddles. In these cases, the symmetry of the Loud's systems allows to split up the period function and to consider only the time function associated to the passage around one saddle. The most complicated situations are those in which the period annulus is unbounded because then the saddle is at infinity and one has to consider meromorphic vector fields. In order to obtain the asymptotic development mentioned above we use a result proved in [13], which provides the first terms in the expansion of this type of time function (see Proposition 3.9). Theorem 3.3 deals with this situation and so it is the most difficult result to prove. In Section 4 we study the bifurcations of the period function in the interior of the period annulus and we show that there exist three germs of curves with this type of bifurcation values. Next, in Section 5, we determine some regions in the parameter space for which the corresponding period function has at least one or two critical periods. Finally in Section 6 we comment on the complete conjectural bifurcation diagram of the period function of the dehomogenized Loud's systems. We also pose some precise open questions remaining to prove its validity. In particular, for certain values of the parameters, the polycycle in the outer boundary of the period annulus has singular points at infinity that are resonant saddles or saddle-nodes. In these cases, tools analogous to Proposition 3.9 still have to be developed. The global study of the bifurcation values of the period function in the interior of the period annulus seems out of reach for the moment.

## 2. Basic definitions

We say that a critical point $p$ of a planar differential system is a center if it has a punctured neighborhood that consists entirely of periodic orbits surrounding $p$. The largest punctured neighborhood with this property is called the period annulus of the center and it will be denoted by $\mathcal{P}$. Compactifying $\mathbb{R}^{2}$ to the Poincaré disc, the boundary of $\mathcal{P}$ has two connected components, the center itself and a polycycle. We call them, respectively, the inner and outer boundary of $\mathcal{P}$.

Definition 2.1. Let $\Lambda$ be an open subset of $\mathbb{R}^{m}$ and consider a continuous family of analytic planar vector fields $\left\{X_{\mu}, \mu \in \Lambda\right\}$. Suppose that, for each $\mu \in \Lambda, X_{\mu}$ has a center at $p_{\mu} \in \mathbb{R}^{2}$. We say that the family of corresponding period annuli varies continuously if there exists a continuous family of analytic functions $\left\{\xi_{\mu}: \mu \in \Lambda\right\}$ such that, for each $\mu \in \Lambda, \xi_{\mu}:[0,1] \longrightarrow \mathbb{R P}^{2}$ verifies:
(a) $\xi_{\mu}(0)=p_{\mu}$ and $\xi_{\mu}(1)$ belongs to the outer boundary of $\mathcal{P}_{\mu}$,
(b) $\xi_{\mu}(s) \in \mathcal{P}_{\mu}$ for all $s \in(0,1)$,
(c) $\xi_{\mu}^{\prime}(s)$ is transverse to $X_{\mu}\left(\xi_{\mu}(s)\right)$ for all $s \in(0,1)$.

Note that $\xi_{\mu}$ is the parametrization of a transverse section for $X_{\mu}$ in $\mathcal{P}_{\mu}$. In general, for each fixed $\mu \in \Lambda$, it is always possible to take such a transverse section. Definition 2.1 requires the existence of one that varies continuously with the parameter. As we will see in Section 3.1, the period annuli of the family that we study vary continuously. Next remark shows, however, that this does not always occur.

Remark 2.2. The period annuli of the center at the origin of the 1-parameter family of potential systems

$$
\left\{\begin{array}{l}
\dot{x}=-y \\
\dot{y}=x+a x^{3}+x^{5}
\end{array}\right.
$$

do not vary continuously. Indeed, it is easy to show that, for $a<2$, the period annulus $\mathcal{P}_{a}$ is the whole plane, while for $a \geqslant 2$ there exists a positive constant $r$ (not depending on $a$ ) such that $\mathcal{P}_{a}$ is inside a disk of radius $r$.

Let $\left\{X_{\mu}, \mu \in \Lambda\right\}$ be a continuous family of analytic vector fields with a center $p_{\mu}$. Assume that the corresponding period annuli vary continuously and consider the family of transverse sections parametrized by $\left\{\xi_{\mu}, \mu \in \Lambda\right\}$. For each $(s ; \mu) \in(0,1) \times \Lambda$, we denote the period of the periodic orbit of $X_{\mu}$ passing through the point $\xi_{\mu}(s)$ by $P_{\mu}(s)$. We say then that $P_{\mu}$ is a parametrization of the period function of $X_{\mu}$. Note that $P_{\mu}$ is an analytic function on $(0,1)$. In order to study the qualitative properties of the period function we consider $Z_{\mu}(s)=P_{\mu}^{\prime}(s)$, which is a function defined on $(0,1)$ for all $\mu \in \Lambda$. The following definition deals with a slightly more general situation, but the convenience for this will be clear in a moment.

Definition 2.3. Let $\left\{I_{\mu}, \mu \in \Lambda\right\}$ be a continuous family of intervals in $\mathbb{R}$ and consider a continuous family of functions $\left\{Z_{\mu}: I_{\mu} \longrightarrow \mathbb{R}, \mu \in \Lambda\right\}$. We say that $\mu_{0} \in \Lambda$ is a regular value of the family $\left\{Z_{\mu}, \mu \in \Lambda\right\}$ if there exist a neighborhood $U$ of $\mu_{0}$ and an isotopy $\left\{h_{\mu}: I_{\mu} \longrightarrow I_{\mu_{0}}, \mu \in U\right\}$, with $h_{\mu_{0}}=i d$, such that

$$
\begin{equation*}
\operatorname{sgn}\left(Z_{\mu}(s)\right)=\operatorname{sgn}\left(Z_{\mu_{0}}\left(h_{\mu}(s)\right)\right) \tag{3}
\end{equation*}
$$

for all $s \in I_{\mu}$ and $\mu \in U$. A parameter $\mu_{0}$ which is not regular is called a bifurcation value.

Note that the domain of definition of $Z_{\mu}$ depends on $\mu$. To be more precise, by a continuous family of functions we mean with respect to the induced topology on $\cup_{\mu \in \Lambda} I_{\mu} \times\{\mu\}$ as a subset of $\mathbb{R} \times \Lambda$.

Definition 2.4. Let $\left\{X_{\mu}, \mu \in \Lambda\right\}$ be a continuous family of analytic vector fields with a center $p_{\mu}$ and assume that the corresponding period annuli vary continuously.
(a) We say that $\mu_{0} \in \Lambda$ is a regular (respectively, bifurcation) value of the period function if for some parametrization of the period function $P_{\mu}$ we have that $\mu_{0}$ is a regular (respectively, bifurcation) value of the family $\left\{P_{\mu}^{\prime}:(0,1) \longrightarrow \mathbb{R}, \mu \in \Lambda\right\}$.
(b) We say that $\mu_{0} \in \Lambda$ is a local regular value of the period function in the interior if there exists some parametrization of the period function $P_{\mu}$ such that for any $c \in$ $(0,1)$ there exists a continuously varying neighborhood $I_{\mu}(c)$ of $c$ in $(0,1)$ such that $\mu_{0}$ is a regular value of the family $\left\{P_{\mu}^{\prime}: I_{\mu}(c) \longrightarrow \mathbb{R}, \mu \in \Lambda\right\}$. A parameter which is not a local regular value in the interior is called a local bifurcation value in the interior.
(c) We say that $\mu_{0} \in \Lambda$ is a local regular value of the period function at the inner (respectively, outer) boundary if for some parametrization of the period function $P_{\mu}$ there exists a continuously varying neighborhood $I_{\mu}(c)$ of $c=0$ (respectively, $c=1)$ such that $\mu_{0}$ is a regular value of the family $\left\{P_{\mu}^{\prime}: I_{\mu}(c) \cap(0,1) \longrightarrow \mathbb{R}, \mu \in\right.$ $\Lambda\}$. A parameter which is not a local regular value at the inner (respectively, outer) boundary is called a local bifurcation value at the inner (respectively, outer) boundary.
(d) We say that the period function of $X_{\mu_{0}}$ is monotonous increasing (respectively, decreasing) at the inner boundary if for some parametrization of the period function $P_{\mu}$ there exists $\varepsilon>0$ such that $P_{\mu_{0}}^{\prime}(s)>0$ (respectively, $P_{\mu_{0}}^{\prime}(s)<0$ ) for all $s \in$ $(0, \varepsilon)$. The monotonicity in the outer boundary is defined exactly the same way using $(1-\varepsilon, 1)$ instead of $(0, \varepsilon)$.

Remark 2.5. In the above definitions one can replace "some parametrization" by "any parametrization". Indeed, assume for instance that $\mu_{0} \in \Lambda$ is a regular value using $P_{\mu}$ and consider another parametrization, say $\widetilde{P}_{\mu}$. Then, following the notation of Definition 2.3, take $\widetilde{h}_{\mu}:=\tau_{\mu_{0}} \circ h_{\mu} \circ \tau_{\mu}^{-1}$ where $\tau_{\mu}$ is the Poincaré mapping from the
transverse section given by $\xi_{\mu}$ to the one given by $\widetilde{\xi}_{\mu}$. Now, taking $P_{\mu}(s)=\widetilde{P}_{\mu}\left(\tau_{\mu}(s)\right)$ and $\tau_{\mu}^{\prime}(s)>0$ into account, it is easy to verify that $\mu_{0}$ is a regular value using $\widetilde{P}_{\mu}$.

Remark 2.6. There are two situations in which it is very easy to decide whether a parameter $\mu_{0}$ is a local regular value or not:
(a) If any neighborhood of $\mu_{0}$ contains two parameters $\mu_{+}$and $\mu_{-}$such that $X_{\mu_{+}}$and $X_{\mu_{-}}$have different monotonicity at the inner (respectively, outer) boundary, then $\mu_{0}$ is a local bifurcation value at the inner (respectively, outer) boundary.
(b) If for some parametrization of the period function $P_{\mu}$ there exists a neighborhood $U$ of $\mu_{0}$ and $\varepsilon>0$ such that $P_{\mu}^{\prime}(s) \neq 0$ for all $\mu \in U$ and $s \in(0, \varepsilon)$ (respectively, $s \in(1-\varepsilon, 1)$ ), then $\mu_{0}$ is a local regular value in the inner (respectively, outer) boundary.

Lemma 2.7. Let $\left\{X_{\mu}, \mu \in \Lambda\right\}$ be a continuous family of analytic vector fields with a center $p_{\mu}$ and assume that the corresponding period annuli vary continuously. Then the bifurcation diagram of the period function is the union of the local bifurcation diagrams at the inner and outer boundary and in the interior.

Proof. It is obvious that a regular value is a local regular at the inner and outer boundary and in the interior. Let us prove the converse. Let $\mu_{0} \in \Lambda$ be a local regular value at the inner boundary, the outer boundary and the interior. Note that by Remark 2.5 we can assume that we use the same parametrization, say $P_{\mu}$, of the period function. By the local regularity at the inner and outer boundary, there is a neighborhood $U$ of $\mu_{0}$ and continuously varying neighborhoods $I_{\mu}(0)$ and $I_{\mu}(1)$, of the inner and outer boundary, respectively, on which an isotopy $h_{\mu}$ as in the Definition 2.3 exists for $Z_{\mu}=P_{\mu}^{\prime}$. By analyticity, $P_{\mu_{0}}^{\prime}(s)$ has at most a finite number of zeros, say $c_{1}, \ldots, c_{k}$, in an open neighborhood $J$ of $(0,1) \backslash\left(I_{\mu}(0) \cup I_{\mu}(1)\right)$. Using the local regularity of $\mu_{0}$ in the interior, for each $i=1,2, \ldots, k$, there exists a continuously varying closed interval $I_{\mu}\left(c_{i}\right)$ containing $c_{i}$ and an isotopy $h_{\mu}$ such that equality (3) holds for all $s \in I_{\mu}\left(c_{i}\right)$ and $\mu \in U$. Reducing $U$ and each $I_{\mu}\left(c_{i}\right)$ if necessary, we can assume in addition that $I_{\mu}\left(c_{1}\right), \ldots, I_{\mu}\left(c_{k}\right)$ are pairwise disjoint and that

$$
P_{\mu}^{\prime}(s) \neq 0 \quad \text { for } s \in J \backslash\left(\bigcup_{i=1}^{k} I_{\mu}\left(c_{i}\right)\right) \text { and } \mu \in U
$$

On the other hand, reducing also $I_{\mu}(0)$ and $I_{\mu}(1)$ if necessary, we can assume that $I_{\mu}\left(c_{1}\right), \ldots, I_{\mu}\left(c_{k}\right)$ do not intersect $I_{\mu}(0)$ and $I_{\mu}(1)$ neither. It remains therefore to define the isotopy in a finite disjoint union of open intervals. In each of these intervals we define it as an affine map whose values at the endpoints are already defined.

The above result shows that if $\mathcal{P}_{\mu}$ varies continuously, then in order to obtain the bifurcation diagram it is enough to study the three possible types of local bifurcations given in (b) and (c) of Definition 2.4. However, dealing with a family of centers such
that the period annuli do not vary continuously, it may occur that some bifurcation does not correspond to any of these three types. In fact this is the case of the period function of the centers in Remark 2.2 (see [11] for details).

As we already mentioned, the local bifurcation diagram at the inner boundary is fully understood for the quadratic centers (see Section 4) thanks to the results of Chicone and Jacobs [4]. Let us point out that their definition of bifurcation value at the inner boundary is not equivalent to ours. Their definition allows to describe better the bifurcation, but its usefulness is strongly based on the fact that the period function of a nondegenerate center can be extended analytically to the inner boundary. In general this is not possible in the outer boundary, which is the case that we study. We want, on the other hand, a unified definition for both boundaries because otherwise a result as Lemma 2.7 is very difficult to obtain. This is the reason why we use here a different definition. We point out, however, that, for the quadratic centers, the bifurcation values at the inner boundary are the same with both definitions (see Remark 4.2).

## 3. Bifurcation at the outer boundary

This section is devoted to the proof of Theorem A and it is divided into four subsections. In the first one we study the phase portrait of the dehomogenized Loud's systems and we focus on the shape of the period annulus of the center at the origin. In brief, we show that, apart from a parameter subset which consists of some straight lines, there are four different types of period annuli. We turn then to the study of the period function in each situation. We consider the two cases in which the period annulus is unbounded in Section 3.2, and the two cases in which it is bounded in Section 3.3. Finally in Section 3.4 we prove Theorem A.

### 3.1. Study of the phase portrait

In the sequel, setting $\mu=(D, F)$, we shall denote by $\left\{X_{\mu}, \mu \in \mathbb{R}^{2}\right\}$ the family of vector fields corresponding to the dehomogenized Loud's systems, i.e.,

$$
X_{\mu}=y(x-1) \partial_{x}+\left(x+D x^{2}+F y^{2}\right) \partial_{y} .
$$

For each value of $\mu$, the vector field $X_{\mu}$ has a center at the origin, whose period function is our object of study. In order to do this, we need to determine the period annulus $\mathcal{P}_{\mu}$ of $X_{\mu}$ as well as its outer boundary, which is a polycycle in some compactification of $\mathbb{R}^{2}$. Usually, one takes the Poincaré disk but, for the sake of simplicity in the computations, we will use instead the real projective plane $\mathbb{R} \mathbb{P}^{2}$. We consider $\mathbb{R}^{2}$ covered by the charts $(x, y),(u, v)=\left(\frac{1}{1-x}, \frac{y}{1-x}\right)$ and $(\zeta, \omega)=\left(\frac{1-x}{y}, \frac{1}{y}\right)$. The expressions of $X_{\mu}$ in
$(u, v)$ and $(\zeta, \omega)$ coordinates are given, respectively, by

$$
\begin{aligned}
X_{\mu}(u, v)= & \frac{1}{u}\left(-u v \partial_{u}+\left(-u+u^{2}+D(u-1)^{2}+(F-1) v^{2}\right) \partial_{v}\right) \\
X_{\mu}(\zeta, \omega)= & \frac{1}{\omega}\left(\left((1-F) \zeta+D \zeta^{3}+(D+1) \zeta \omega^{2}-(2 D+1) \omega \zeta^{2}\right) \partial_{\zeta}\right. \\
& \left.+\omega(-F-(\omega-\zeta)((D+1) \omega-D \zeta)) \partial_{\omega}\right) .
\end{aligned}
$$

It is easy to check (see [18] for instance) that if $F \notin\left\{0,1, \frac{1}{2}\right\}$ then the vector field $X_{\mu}$ has a Darboux type first integral given by

$$
\begin{equation*}
H_{\mu}(x, y)=(1-x)^{-2 F}\left(\frac{1}{2} y^{2}-q_{\mu}(x)\right), \tag{4}
\end{equation*}
$$

where $q_{\mu}(x)=a(\mu) x^{2}+b(\mu) x+c(\mu)$ with

$$
\begin{equation*}
a(\mu):=\frac{D}{2(1-F)}, \quad b(\mu):=\frac{D-F+1}{(1-F)(1-2 F)} \quad \text { and } \quad c(\mu):=\frac{F-D-1}{2 F(1-F)(1-2 F)} . \tag{5}
\end{equation*}
$$

The line at infinity $L_{\infty}$ (with respect to the $(x, y)$-coordinates), the conic $\mathcal{C}_{\mu}:=\left\{\frac{1}{2} y^{2}-\right.$ $\left.q_{\mu}(x)=0\right\}$ and the line $L_{1}:=\{x=1\}$ are invariant curves of $X_{\mu}$. The determinant associated to the conic $\mathcal{C}=\mathcal{C}_{\mu}$, which coincides with the discriminant of $q_{\mu}(x)$, is given by

$$
\Delta(\mu):=\frac{(D+F)(D+1-F)}{(1-2 F)^{2}(1-F) F}
$$

Thus, we can see that $\mathcal{C}$ degenerates into two lines when $(D+F)(D+1-F)=0$. Indeed, it is easy to check that the conic $\mathcal{C}$ splits into two real lines when $F=-D$ and $D \notin[-1,0]$. On the other hand, if $F=D+1$ (respectively, $F=-D \in(0,1)$ ), then the conic $\mathcal{C}$ becomes two complex conjugated lines having the center $(x, y)=(0,0)$ (respectively, $(x, y)=(-1 / D, 0))$ as the unique real common point. In the other cases the affine type of $\mathcal{C}$ can be determined by the sign of $\Delta$ and $a$ in the following way:

- If $a<0$ and $\Delta<0$ then the conic $\mathcal{C}$ has no real points.
- If $a<0$ and $\Delta>0$ then the conic $\mathcal{C}$ is an ellipse.
- If $a>0$ then the conic $\mathcal{C}$ is a hyperbola and we have two subcases depending on the sign of $\Delta$. If $\Delta>0$ then $\mathcal{C}$ cuts the $x$-axis in two points which will be denoted in the sequel by $p_{1}$ and $p_{2}$ with $p_{1}<p_{2}$. If $\Delta<0$ then the hyperbola $\mathcal{C}$ has no common point with $\{y=0\}$.
- If $a=0$ then the conic $\mathcal{C}$ is a parabola (this only occurs when $D=0$ ).

It is well known that every quadratic system has seven singularities (in the projective complex domain and counting multiplicities). Taking the pairwise intersections of the
invariant curves $L_{1}, L_{\infty}$ and $\mathcal{C}$ we obtain five singular points. Moreover, apart from the center at the origin $(x, y)=(0,0)$, we have the singular point $[-1,0, D]$, which, in general, does not live on none of the invariant curves $L_{1}, L_{\infty}$ or $\mathcal{C}$. Now we proceed to study in some detail each singular point of $X_{\mu}$ :

- The two points $L_{1} \cap \mathcal{C}=\left\{\left[1, \pm \sqrt{\frac{-(D+1)}{F}}, 1\right]\right\}$ are real when $(D+1) F<0$. The linear part of $X_{\mu}$ at these points has eigenvalues $\lambda_{1}= \pm 2 F \sqrt{\frac{-(D+1)}{F}}$ and $\lambda_{2}= \pm \sqrt{\frac{-(D+1)}{F}}$.
- The two points $L_{\infty} \cap \mathcal{C}=\left\{\left[1, \pm \sqrt{\frac{D}{1-F}}, 0\right]\right\}$ are real when $(1-F) D>0$. Working in $(u, v)$-coordinates, the linear part of $u X_{\mu}$ at these points has eigenvalues $\lambda_{1}=$ $\mp 2 D \sqrt{\frac{1-F}{D}}$ and $\lambda_{2}=\mp \sqrt{\frac{D}{1-F}}$.
- At the point $L_{1} \cap L_{\infty}=[0,1,0]$ the linear part of $\omega X_{\mu}$ has eigenvalues $\lambda_{1}=-F$ and $\lambda_{2}=1-F$.
- At the point $[-1,0, D]$, with $D \neq 0$, the linear part of $X_{\mu}$ has eigenvalues $\lambda_{i}=$ $\pm \sqrt{1+\frac{1}{D}}$.
Fig. 4 shows the bifurcation diagram of the phase portrait of the dehomogenized Loud's systems. It is important to note that in each phase portrait we place the center $(0,0)$ on the left of the centered invariant line $L_{1}=\{x=1\}$. In addition the conic $\mathcal{C}$ appears in boldface type when it is relevant. We will describe next in brief all the bifurcations occurring in this diagram:
- Along $D=-1$ there is a collapse of the three singularities $L_{1} \cap \mathcal{C}$ and $[-1,0, D]$. We point out that if $F>1$ then this bifurcation does not affect the period annulus.
- The bifurcation at $D=0$ occurs when the three singularities $L_{\infty} \cap \mathcal{C}$ and $[-1,0, D]$ collapse. This bifurcation always affects the period annulus.
- The bifurcations at $F=0$ and $F=1$ can also be easily described. Indeed, the three singular points $L_{1} \cap L_{\infty}$ and $L_{1} \cap \mathcal{C}$ collapse giving raise to a saddle-node at infinity, whose strong separatrix is on $L_{\infty}$ when $F=0$ and on $L_{1}$ when $F=1$. Note that these bifurcations only affect the period annulus when $D \in[-1,0]$.
- Along $F+D=0$ the conic $\mathcal{C}$ degenerates, but this does not affect the outer boundary of the period annulus if $F \in(0,1)$.
- Along $F=D+1$ the conic $\mathcal{C}$ also degenerates, but this never affects the period annulus.
- Finally the bifurcation at $F=1 / 2$ is more subtle because there is no confluence of singularities. The position of the conic depends on $F>1 / 2$ or $F<1 / 2$ and it "explodes" to the limit set $L_{1} \cup L_{\infty}$ as $F$ tends to $1 / 2$. Note in addition that the singular point $L_{1} \cap L_{\infty}$ is a saddle for $0<F<1$. In fact it can be shown that this saddle is orbitally linearizable for $F \neq 1 / 2$. In contrast, if $F=1 / 2$ then the singular point $L_{1} \cap L_{\infty}$ (which belongs to the outer boundary of $\mathcal{P}_{\mu}$ when $D \in[-1,0]$ ) becomes a resonant saddle with hyperbolicity ratio equal to one. Notice, however, that this bifurcation never affects the structure of the outer boundary of the period annulus.


Fig. 4. Phase portraits of the dehomogenized Loud's systems.
Remark 3.1. The discussion above shows that the bifurcations (in the structure) of the outer boundary of the period annulus of the center at the origin occur only on the dotted curve in Fig. 4. Let us point out that we shall not study the period function corresponding to these parameters.

Lemma 3.2. The family of period annuli of the center at the origin of the dehomogenized Loud's systems $\left\{X_{\mu}, \mu \in \mathbb{R}^{2}\right\}$ varies continuously.
Proof. Let $(b(\mu), 0)$ be the intersection point of the outer boundary of $\mathcal{P}_{\mu}$ with the positive $x$-axis (see Fig. 4). This provides us a natural parametrization for $\mathcal{P}_{\mu}$. Indeed,
notice that $b(\mu)$ depends continuously on $\mu$ and that $X_{\mu}$ is transverse to the segment $\{(x, 0): 0<x<b(\mu)\}$ for all $\mu \in \mathbb{R}^{2}$. It suffices therefore to consider $\xi_{\mu}:[0,1] \longrightarrow$ $\mathbb{R}^{2}$ defined by means of $\xi_{\mu}(s):=(b(\mu) s, 0)$.

As we already mentioned, in order to study the behavior of the period function near the outer boundary of the period annulus we must treat separately the four different types of polycycle that bound it. We gather them in two sections according to whether the period annulus (considered as a subset of $\mathbb{R}^{2}$ ) is bounded or not. In Section 3.2 we deal with the unbounded case, which is divided in two subcases: Section 3.2.1 corresponds to period annuli with the outer boundary contained in $\mathcal{C} \cup L_{\infty}$ and Section 3.2.2 to period annuli with the outer boundary contained in $L_{1} \cup L_{\infty}$. As a matter of fact this last case was already treated in [13] and here, for the sake of completeness, we only recall the result that we obtained. Finally, Section 3.3 deals with the cases in which the period annulus is bounded. From Fig. 4 it follows that there are two possibilities for the polycycle in the outer boundary, namely, a saddle loop or a bicycle. Let us conclude this section pointing out that Fig. 4 does not constitute a new result (see [19,22] for instance).

### 3.2. Unbounded period annulus

### 3.2.1. The case $F>1, F+D>0$ and $D<0$.

In this subsection we study the period function of the center at the origin of $X_{\mu}$ in case that the parameter $\mu$ belongs to

$$
U:=\left\{(D, F) \in \mathbb{R}^{2}: F>1, F+D>0 \text { and } D<0\right\}
$$

We shall prove the existence of a curve $\Gamma_{1}$ such that setting $\Gamma_{2}=U \cap\{F=2\}$ then the following holds:

Theorem 3.3. Let $\left\{X_{\mu}, \mu \in \mathbb{R}^{2}\right\}$ be the family of vector fields in (2) and consider the period function of the center at the origin. Then the open set $U \backslash\left\{\Gamma_{1} \cup \Gamma_{2}\right\}$ corresponds to local regular values of the period function at the outer boundary of the period annulus. Moreover, for these parameters, the period function is monotonous near the outer boundary and the corresponding character is shown in Fig. 5.

To be more precise, $\Gamma_{1}$ is the zero level set of an explicit function which is given in (28). In order to draw $\Gamma_{1}$ in Fig. 5 we have computed it numerically. Analytically, among other properties that are gathered in Proposition 3.11, we have proved that $\Gamma_{1}$ is the graphic of an analytic function $D=\mathcal{G}(F)$.

In order to prove Theorem 3.3 we shall study the asymptotic development of the period function near the outer boundary of $\mathcal{P}_{\mu}$. For the parameter values under consideration (see Fig. 4), recall that the outer boundary of the period annulus of the center is made up of the line at infinity and a branch of the conic $\mathcal{C}_{\mu}=\left\{\frac{1}{2} y^{2}-q_{\mu}(x)=0\right\}$,


Fig. 5. Monotonicity of the period function at the outer boundary of $\mathcal{P}_{\mu}$.
where $q_{\mu}(x)=a x^{2}+b x+c$ with the coefficients $a, b, c$ defined in (5). For $\mu \in U$, the conic has two different intersection points with $y=0$, namely

$$
p_{1}:=\frac{-b-\sqrt{b^{2}-4 a c}}{2 a} \quad \text { and } \quad p_{2}:=\frac{-b+\sqrt{b^{2}-4 a c}}{2 a}
$$

which one can verify that $0<p_{1}<p_{2}$ and $p_{1}<1$. Notice in particular that ( $p_{1}, 0$ ) belongs to the outer boundary of the period annulus. Since one can check that $X_{\mu}$ is transverse to $\left\{(x, 0): 0<x<p_{1}\right\}$, we have a global parametrization of the set of periodic orbits in $\mathcal{P}_{\mu}$. Thus, for $(s, \mu) \in\left(0, p_{1}\right) \times U$, we denote by $P(s ; \mu)$ the period of the periodic orbit of $X_{\mu}$ passing through the point $\left(p_{1}-s, 0\right)$.

Notice that one can easily normalize $P(s ; \mu)$ to obtain a parametrization of the period function defined for $s \in(0,1)$ and so that the inner and outer boundary correspond to $s \approx 0$ and $s \approx 1$, respectively. However, for convenience in the computations we prefer to use the previous one instead, for which we stress that the outer boundary corresponds to $s \approx 0$.

Theorem 3.3 follows almost directly from Theorem 3.6, which gives the first nontrivial term of the asymptotic development of $P_{s}(s ; \mu)$ at $s=0$. In its statement we use the following definitions:

Definition 3.4. Let $W$ be an open subset of $\mathbb{R}^{m}$. We denote by $\mathcal{I}(W)$ the set of germs of analytic functions $h(s ; \mu)$ defined on $(0, \varepsilon) \times W$ for some $\varepsilon>0$ such that

$$
\lim _{s \rightarrow 0} h(s ; \mu)=0 \quad \text { and } \quad \lim _{s \rightarrow 0} s \frac{\partial h(s ; \mu)}{\partial s}=0
$$

uniformly (on $\mu$ ) on every compact subset of $W$.

Let us also denote by $\mathcal{I}_{0}(W)$ the set of germs of analytic functions $h(s ; \mu)$ defined on $(-\varepsilon, \varepsilon) \times W$ for some $\varepsilon>0$ such that $h(0 ; \mu) \equiv 0$. Note therefore that $\mathcal{I}_{0}(W) \subset \mathcal{I}(W)$.

Definition 3.5. The function defined for $s>0$ and $\alpha \in \mathbb{R}$ by means of

$$
\omega(s ; \alpha)= \begin{cases}\frac{s^{\alpha-1}-1}{\alpha-1} & \text { if } \alpha \neq 1, \\ \log s & \text { if } \alpha=1\end{cases}
$$

is called the Roussarie-Ecalle compensator.
Let us define in addition

$$
\lambda(\mu):=\frac{1}{2(F-1)}
$$

and introduce the covering of the parameter space $U$ given by the open subsets

$$
\begin{align*}
& U_{1}:=\{\mu \in U: F<3 / 2\}, \quad U_{2}:=\{\mu \in U: F>3 / 2\} \quad \text { and } \\
& U_{3}:=\{\mu \in U: 5 / 4<F<2\}, \tag{6}
\end{align*}
$$

which one can verify that correspond to $\lambda(\mu)>1, \lambda(\mu)<1$ and $1 / 2<\lambda(\mu)<2$, respectively.

Now, with the definitions and notation introduced above, we prove the following:
Theorem 3.6. Denote

$$
\Delta_{0}(\mu)=\frac{2 \sqrt{2}}{\sqrt{a+b+c}} \operatorname{arctanh}\left(\frac{2 a+b-\sqrt{b^{2}-4 a c}}{2 \sqrt{a(a+b+c)}}\right) .
$$

Then the following holds:
(a) If $\mu \in U_{1}$ then $P(s ; \mu)=\Delta_{0}(\mu)+\Delta_{1}(\mu) s+s f_{1}(s ; \mu)$, where $f_{1} \in \mathcal{I}\left(U_{1}\right)$ and

$$
\begin{aligned}
\Delta_{1}(\mu)= & \frac{-1 / \sqrt{2 a}}{\left(p_{2}-p_{1}\right)\left(1-p_{1}\right)} \\
& \times\left\{2-\int_{0}^{1}\left(u^{-\frac{1}{\lambda}}\left(\frac{1-p_{2}}{1-p_{1}}(u-1)+1\right)^{1+\frac{1}{\lambda}}-1\right) \frac{d u}{(1-u)^{3 / 2}}\right\}
\end{aligned}
$$

(b) If $\mu \in U_{2}$ then $P(s ; \mu)=\Delta_{0}(\mu)+\Delta_{2}(\mu) s^{\lambda}+s^{\lambda} f_{2}(s ; \mu)$, where $f_{2} \in \mathcal{I}\left(U_{2}\right)$ and

$$
\Delta_{2}(\mu)=\sqrt{\frac{2 \pi}{a}} \frac{\lambda\left(p_{2}-p_{1}\right)^{\lambda}}{\left(1-p_{1}\right)^{2 \lambda+1}} \frac{\Gamma\left(\frac{1}{2(1-F)}\right)}{\Gamma\left(\frac{F-2}{2(F-1)}\right)}
$$

(c) If $\mu \in U_{3}$ then $P(s ; \mu)=\Delta_{0}(\mu)+\Delta_{3}(\mu) s \omega(s ; \lambda)+\Delta_{4}(\mu) s+s f_{3}(s ; \mu)$, where $f_{3} \in \mathcal{I}\left(U_{3}\right)$ and the functions $\Delta_{3}(\mu)$ and $\Delta_{4}(\mu)$ are analytic on $U_{3}$. Furthermore, if $\lambda\left(\mu_{0}\right)=1$ then

$$
\Delta_{3}\left(\mu_{0}\right)=-\frac{p_{2}-p_{1}}{\sqrt{2 a}\left(1-p_{1}\right)^{3}} .
$$

Notice that, in the Poincaré disc, the outer boundary of the $\mathcal{P}_{\mu}$ is a polycycle with two hyperbolic saddles located at infinity (see Fig. 4). In addition, taking advantage of the symmetry of the Loud's family with respect to the $x$-axes, in order to prove Theorem 3.6 it is enough to study half of the period. Consequently we must only study the time function associated to the passage through one of these saddles. To this end we shall use a result which appears in [13]. In that paper, given an analytic family of vector fields in $\mathbb{R}^{2}$ having a saddle point, we studied the asymptotic development of the time function along the union of two separatrices. Next, for the sake of completeness, we state this result (see Proposition 3.9) and we explain the related definitions.

Let $W$ be an open set of $\mathbb{R}^{m}$ and let $\left\{\widetilde{X}_{\mu}, \mu \in W\right\}$ be an analytic family of vector fields defined on some open set $V$ of $\mathbb{R}^{2}$. Assume that each vector field $\widetilde{X}_{\mu}$ has a hyperbolic saddle $p_{\mu}$ as the unique critical point inside $V$. In this situation it is well known that there exist exactly two analytic transverse invariant curves $\mathcal{S}_{\mu}$ and $\mathcal{T}_{\mu}$, the stable and unstable manifolds, passing through $p_{\mu}$ (depending also analytically on $\mu$ ). We consider an analytic family of meromorphic vector fields $X_{\mu}$ proportional to $\widetilde{X}_{\mu}$ with a pole of order $n>0$ along $\mathcal{T}_{\mu}$. We can take a coordinate system $(u, v, \mu)$ on $V \times W \subset \mathbb{R}^{2+m}$ such that $p_{\mu}=(0,0, \mu), \mathcal{S}_{\mu}=\{(u, v, \mu): u=0\}$ and $\mathcal{T}_{\mu}=\{(u, v, \mu):$ $v=0\}$. In these coordinates the family $\left\{X_{\mu}, \mu \in W\right\}$ can be written as

$$
\begin{equation*}
X_{\mu}(u, v)=\frac{1}{v^{n}}\left(u P(u, v ; \mu) \partial_{u}+v Q(u, v ; \mu) \partial_{v}\right), \tag{7}
\end{equation*}
$$

where $P$ and $Q$ are analytic functions such that $P(u, 0 ; \mu)>0$ and $Q(0, v ; \mu)<0$ for any $(0, v, \mu) \in \mathcal{S}_{\mu}$ and $(u, 0, \mu) \in \mathcal{T}_{\mu}$. Moreover, by hypothesis, we have that

$$
\lambda(\mu):=-\frac{Q(0,0 ; \mu)}{P(0,0 ; \mu)}>0 .
$$

The family $\left\{X_{\mu}, \mu \in W\right\}$ can be thought of as a single vector field $X$ defined on $V \times W \subset \mathbb{R}^{2+m}$ whose trajectories are contained inside the submanifolds $\{\mu=$ const $\}$. Let $\sigma: I \times W \longrightarrow \Sigma_{\sigma}$ and $\tau: I \times W \longrightarrow \Sigma_{\tau}$ be two analytic transverse sections to $X$


Fig. 6. Definition of $T$ and $R$ in Proposition 3.9.
defined by

$$
\sigma(s ; \mu)=\left(\sigma_{1}(s ; \mu), \sigma_{2}(s ; \mu) ; \mu\right) \quad \text { and } \quad \tau(s ; \mu)=\left(\tau_{1}(s ; \mu), \tau_{2}(s ; \mu) ; \mu\right)
$$

such that $\sigma(0 ; \mu) \in \mathcal{S}_{\mu}$ and $\tau(0 ; \mu) \in \mathcal{T}_{\mu}$. Here $I$ denotes a small interval of $\mathbb{R}$ containing 0 .

We denote the Dulac and time mappings between the transverse sections $\Sigma_{\sigma}$ and $\Sigma_{\tau}$ by $R$ and $T$ respectively. More precisely (see Fig. 6), if $\varphi\left(t,\left(u_{0}, v_{0}\right) ; \mu\right)$ is the solution of $X_{\mu}$ passing through $\left(u_{0}, v_{0}\right)$ at $t=0$, for each $s>0$ we define $R(s ; \mu)$ and $T(s ; \mu)$ by means of the relation

$$
\begin{equation*}
\varphi(T(s ; \mu), \sigma(s) ; \mu)=\tau(R(s ; \mu)) \tag{8}
\end{equation*}
$$

Definition 3.7. We will say that $\left\{X_{\mu}, \mu \in W\right\}$ verifies the family linearization property (FLP) if there exist an open set $U \subset \mathbb{R}^{2}$ containing the origin and an analytic local diffeomorphism $\Phi: U \times W \rightarrow V \times W$ of the form $\Phi(x, y ; \mu)=(x+$ h.o.t., $y+$ h.o.t., $\mu)$ such that

$$
X_{\mu}=\Phi_{*}\left(\frac{1}{f(x, y ; \mu)}\left(x \partial_{x}-\lambda(\mu) y \partial_{y}\right)\right)
$$

where $f$ is an analytic function on $U \times W$.
Remark 3.8. It is easy to show that the family of meromorphic vector fields $\left\{X_{\mu}, \mu \in\right.$ $W$ \} defined in (7) verifies FLP if it has a Darboux first integral

$$
H_{\mu}(x, y)=f_{1}(x, y ; \mu)^{\beta_{1}(\mu)} \cdots f_{k}(x, y ; \mu)^{\beta_{k}(\mu)}
$$

where $f_{j}$ and $\beta_{j}$ are analytic functions on $V \times W$ and $W$, respectively.
Recall that $H_{\mu}(x, y)=(1-x)^{-2 F}\left(\frac{1}{2} y^{2}-q_{\mu}(x)\right)$ is a Darboux first integral for $X_{\mu}$ if $F(F-1)(2 F-1) \neq 0$, so the FLP is verified in these cases.

In order to simplify the expressions that appear in the statement of the next result we introduce the functions

$$
\begin{aligned}
& L(u ; \mu):=\exp \left(\int_{\sigma_{2}(0)}^{u}\left(\frac{P(0, y)}{Q(0, y)}+\frac{1}{\lambda}\right) \frac{d y}{y}\right), \\
& M(u ; \mu):=\exp \left(\int_{0}^{u}\left(\frac{Q(x, 0)}{P(x, 0)}+\lambda\right) \frac{d x}{x}\right)
\end{aligned}
$$

and the covering of the parameter space $W$ given by the open subsets

$$
\begin{aligned}
& W_{1}:=\left\{\mu \in W: \lambda>\frac{1}{n}\right\}, \quad W_{2}:=\left\{\mu \in W: \lambda<\frac{1}{n}\right\} \quad \text { and } \\
& W_{3}:=\left\{\mu \in W: \frac{1}{n+1}<\lambda<\frac{2}{n}\right\} .
\end{aligned}
$$

Proposition 3.9. Let $\left\{X_{\mu}, \mu \in W\right\}$ be the family of vector fields defined in (7) and assume that it verifies FLP. Let $R$ and $T$ be, respectively, the Dulac map and the time function associated to the transverse sections $\Sigma_{\sigma}$ and $\Sigma_{\tau}$ as introduced in (8). Denote

$$
\rho(\mu)=\frac{\sigma_{1}^{\prime}(0)^{\lambda} \sigma_{2}(0)}{\tau_{2}^{\prime}(0) \tau_{1}(0)^{\lambda}} L(0)^{\lambda} M\left(\tau_{1}(0)\right) \quad \text { and } \quad \Delta_{0}(\mu)=\int_{\sigma_{2}(0)}^{0} \frac{v^{n-1}}{Q(0, v)} d v
$$

Then $R(s ; \mu)=\rho(\mu) s^{\lambda}+s^{\lambda} f_{0}(s ; \mu)$ with $f_{0} \in \mathcal{I}(W)$. In addition, the time function $T(s ; \mu)$ verifies the following:
(a) If $\mu \in W_{1}$ then $T(s ; \mu)=\Delta_{0}(\mu)+\Delta_{1}(\mu) s+s f_{1}(s ; \mu)$, where $f_{1} \in \mathcal{I}\left(W_{1}\right)$ and

$$
\Delta_{1}(\mu)=-\frac{\sigma_{2}^{\prime}(0) \sigma_{2}(0)^{n-1}}{Q\left(0, \sigma_{2}(0)\right)}+\sigma_{1}^{\prime}(0) \sigma_{2}(0)^{1 / \lambda} \int_{0}^{\sigma_{2}(0)} \frac{Q_{u}(0, v) L(v) v^{n-1 / \lambda}}{Q(0, v)^{2}} \frac{d v}{v}
$$

(b) If $\mu \in W_{2}$ then $T(s ; \mu)=\Delta_{0}(\mu)+\Delta_{2}(\mu) s^{\lambda n}+s^{\lambda n} f_{2}(s ; \mu)$, where $f_{2} \in \mathcal{I}\left(W_{2}\right)$ and

$$
\Delta_{2}(\mu)=\sigma_{1}^{\prime}(0)^{\lambda n} \sigma_{2}(0)^{n} L(0)^{\lambda n}\left\{\frac{\tau_{1}(0)^{-\lambda n}}{n Q(0,0)}+\int_{0}^{\tau_{1}(0)}\left(\frac{M(u)^{n}}{P(u, 0)}-\frac{M(0)^{n}}{P(0,0)}\right) \frac{d u}{u^{\lambda n+1}}\right\}
$$

(c) If $\mu \in W_{3}$ then $T(s ; \mu)=\Delta_{0}(\mu)+\Delta_{3}(\mu) s \omega(s ; \lambda n)+\Delta_{4}(\mu) s+s f_{3}(s ; \mu)$, where $f_{3} \in \mathcal{I}\left(W_{3}\right)$ and the functions $\Delta_{3}(\mu)$ and $\Delta_{4}(\mu)$ are analytic on $W_{3}$. Furthermore,
if $\lambda\left(\mu_{0}\right)=1 / n$ then

$$
\Delta_{3}\left(\mu_{0}\right)=-n \sigma_{1}^{\prime}(0) \sigma_{2}(0)^{n} L(0) \frac{Q_{u}(0,0)}{P(0,0)^{2}}
$$

Proposition 3.9 constitutes the main ingredient in the proof of Theorem 3.6. However, in order to apply it we must first perform a change of coordinates that sends each separatrix of the saddle at infinity to a straight line. This will raise some technical complications because the coordinate transformation that we use is singular and it creates a line of critical points. To bypass this problem we will have to split up the time function and to introduce an additional parameter associated to the new transverse sections. This makes the proof more complicated than one could expect. In particular we shall need the following result to study the remainder terms. Its proof can be found in [13] and, for the sake of brevity, in the statement we denote $\mathcal{I}(W)$ and $\mathcal{I}_{0}(W)$ by $\mathcal{I}$ and $\mathcal{I}_{0}$, respectively (see Definition 3.4).

Lemma 3.10. Assume that $a(\mu), k(\mu)$ and $r(\mu)$ are positive analytic functions.
(a) If $g(s ; \mu)$ and $f(s ; \mu)$ belong to $\mathcal{I}_{0}$ and $\mathcal{I}$, respectively, then $g \circ f \in \mathcal{I}$.
(b) If $f(s ; \mu)$ belongs to $\mathcal{I}$ (respectively $\left.\mathcal{I}_{0}\right)$ and $\varphi:=s^{r}(a+f)$ then $s^{k} \circ \varphi-a^{k} s^{k r}$ belongs to $s^{k r} \mathcal{I}$ (respectively $s^{k r} \mathcal{I}_{0}$ ).
(c) If $f(s ; \mu)$ and $g(s ; \mu)$ belong to $\mathcal{I}$ and $\varphi:=s^{r}(a+f)$ then $\left(s^{k} g\right) \circ \varphi$ belongs to $s^{k r} \mathcal{I}$.
(d) If $g(s ; \mu)$ belongs to $\mathcal{I}_{0}$ then $g \omega(s ; r) \in \mathcal{I}$.
(e) If $g(s ; \mu)$ belongs to $\mathcal{I}_{0}$ then $(s \omega(s ; r)) \circ(s(a+g))=a^{r} s \omega(s ; r)+a \omega(a ; r) s+\Psi$ with $\Psi \in s \mathcal{I}$.

Proof of Theorem 3.6. Note first that, since the transformation $(x, y, t) \longmapsto(x,-y,-t)$ preserves the Loud normal form, it is enough to study half of the period. More concretely, denoting the solution of $X_{\mu}$ passing through $\left(x_{0}, y_{0}\right)$ at $t=0$ by $\varphi\left(t,\left(x_{0}, y_{0}\right) ; \mu\right)$, for each $s \in\left(0, p_{1}\right)$ we define $T(s ; \mu)$ as the minimum positive number so that $\varphi_{2}\left(T(s ; \mu),\left(p_{1}-s, 0\right) ; \mu\right)=0$.

Thus we only need to obtain the coefficients of the asymptotic development of $T(s ; \mu)$ at $s=0$, which involves only one passage through a saddle at infinity. Clearly the coefficients of $P(s ; \mu)$ at $s=0$ will follow then using that $P(s ; \mu)=2 T(s ; \mu)$.

Notice now (see Fig. 7) that $T(s ; \mu)$ is the time function associated to the transverse sections $\Sigma_{1}$ and $\Sigma_{2}$, which are given, respectively, by $\alpha^{1}(s)=\left(p_{1}-s, 0\right)$ and $\alpha^{2}(s)=$ $(-1 / s, 0)$. In order to study $T(s ; \mu)$ we introduce two auxiliary transverse sections, say $\Sigma_{1}^{\eta}$ and $\Sigma_{2}^{\eta}$, on the straight line $y=\eta\left(p_{2}-x\right)$, where $\eta \in(0, \varepsilon)$. To this end, let $\left(x_{\eta}, \eta\left(p_{2}-x_{\eta}\right)\right)$ be the intersection point between this straight line and the hyperbola $\left\{\frac{1}{2} y^{2}-q_{\mu}(x)=0\right\}$. We parametrize $\Sigma_{1}^{\eta}$ and $\Sigma_{2}^{\eta}$ by

$$
\alpha_{\eta}^{1}(s)=\left(x_{\eta}-s, \eta\left(p_{2}-x_{\eta}+s\right)\right) \quad \text { and } \quad \alpha_{\eta}^{2}(s)=\left(-1 / s, \eta\left(p_{2}+1 / s\right)\right)
$$



Fig. 7. Auxiliary transverse sections.
respectively. Let us denote the time function between $\Sigma_{1}$ and $\Sigma_{1}^{\eta}$ by $T_{1}(s ; \mu, \eta)$, the one between $\Sigma_{1}^{\eta}$ and $\Sigma_{2}^{\eta}$ by $T_{2}(s ; \mu, \eta)$, and the one between $\Sigma_{2}^{\eta}$ and $\Sigma_{2}$ by $T_{3}(s ; \mu, \eta)$. Then

$$
T(s ; \mu)=T_{1}(s ; \mu, \eta)+T_{2}\left(R_{1}(s ; \mu, \eta) ; \mu, \eta\right)+T_{3}\left(R_{2}(s ; \mu, \eta) ; \mu, \eta\right)
$$

where $R_{1}(s ; \mu, \eta)$ is the Poincare mapping between $\Sigma_{1}$ and $\Sigma_{1}^{\eta}$ and $R_{2}(s ; \mu, \eta)$ is the one between $\Sigma_{1}$ and $\Sigma_{2}^{\eta}$.

It is well known that $T_{1}, T_{3}$ and $R_{1}$ can be extended analytically to $s=0$, and it is also clear that they are analytical for $(\mu, \eta) \in U \times(-\varepsilon, \varepsilon)$. Hence

$$
\begin{align*}
& T_{1}(s ; \mu, \eta)=\Delta_{0}^{1}(\mu, \eta)+\Delta_{1}^{1}(\mu, \eta) s+s f_{1}(s ; \mu, \eta) \quad \text { with } f_{1} \in \mathcal{I}_{0}(U \times(-\varepsilon, \varepsilon))  \tag{9}\\
& T_{3}(s ; \mu, \eta)=\Delta_{1}^{3}(\mu, \eta) s+s f_{3}(s ; \mu, \eta) \quad \text { with } f_{3} \in \mathcal{I}_{0}(U \times(-\varepsilon, \varepsilon))  \tag{10}\\
& R_{1}(s ; \mu, \eta)=\rho_{1}(\mu, \eta) s+s g_{1}(s ; \mu, \eta) \quad \text { with } g_{1} \in \mathcal{I}_{0}(U \times(-\varepsilon, \varepsilon)) \tag{11}
\end{align*}
$$

Notice moreover that $T_{1}(s ; \mu, \eta) \longrightarrow 0, T_{3}(s ; \mu, \eta) \longrightarrow 0$ and $R_{1}(s ; \mu, \eta) \longrightarrow s$ as $\eta \longrightarrow 0$. Therefore

$$
\begin{equation*}
\lim _{\eta \rightarrow 0} \Delta_{0}^{1}(\mu, \eta)=\lim _{\eta \rightarrow 0} \Delta_{1}^{1}(\mu, \eta)=\lim _{\eta \rightarrow 0} \Delta_{1}^{3}(\mu, \eta)=0 \quad \text { and } \quad \lim _{\eta \rightarrow 0} \rho_{1}(\mu, \eta)=1 \tag{12}
\end{equation*}
$$

The asymptotic developments of $T_{2}$ and $R_{2}$, which correspond to the passage through a saddle at infinity, are more delicate. To obtain them we shall apply Proposition 3.9, and to this end we must first perform a coordinate transformation that sends the separatrices of the saddle to straight lines. We thus consider the singular change of variables given
by

$$
(z, w)=\phi(x, y):=\left(\frac{2 q_{\mu}(x)-y^{2}}{2 a\left(p_{2}-x\right)^{2}}, \frac{p_{2}-p_{1}}{p_{2}-x}\right) .
$$

Setting $k_{1}:=p_{2}-p_{1}$ and $k_{2}:=1 / \sqrt{2 a}$ for the sake of shortness, some computations show that it brings (2) to the system given by the vector field

$$
X_{\mu}=\frac{1}{w}\left(z P(z, w ; \mu) \partial_{z}+w Q(z, w ; \mu) \partial_{w}\right)
$$

where

$$
P(z, w ; \mu)=\frac{2}{k_{2}} \sqrt{1-z-w}\left(k_{1}(F-1)+\left(p_{2}-1\right) w\right)
$$

and

$$
Q(z, w ; \mu)=\frac{1}{k_{2}} \sqrt{1-z-w}\left(-k_{1}+\left(p_{2}-1\right) w\right)
$$

One can verify (see Fig. 8) that $\Sigma_{\sigma}:=\phi\left(\Sigma_{1}^{\eta}\right)$ and $\Sigma_{\tau}:=\phi\left(\Sigma_{2}^{\eta}\right)$ are in the straight line $z+w=1-k_{2}^{2} \eta^{2}$. We parameterize them with the transferred parameterizations of $\Sigma_{1}^{\eta}$ and $\Sigma_{2}^{\eta}$. More concretely, we consider

$$
\sigma(s ; \mu, \eta):=\phi\left(\alpha_{\eta}^{1}(s)\right)=\left(\frac{s\left(1-k_{2}^{2} \eta^{2}\right)^{2}}{k_{1}+s\left(1-k_{2}^{2} \eta^{2}\right)}, \frac{k_{1}\left(1-k_{2}^{2} \eta^{2}\right)}{k_{1}+s\left(1-k_{2}^{2} \eta^{2}\right)}\right)
$$

and

$$
\tau(s ; \mu, \eta):=\phi\left(\alpha_{\eta}^{2}(s)\right)=\left(\frac{s p_{1}+1}{s p_{2}+1}-k_{2}^{2} \eta^{2}, \frac{s k_{1}}{s p_{2}+1}\right)
$$

Note therefore that $T_{2}(s ; \mu, \eta)$ is precisely the time function between $\Sigma_{\sigma}$ and $\Sigma_{\tau}$, and that, on the other hand, the vector field $X_{\mu}$ is meromorphic on the region under consideration. (We point out that Proposition 3.9 cannot be applied to compute $T(s ; \mu)$ directly because $\phi\left(\Sigma_{1}\right)$ and $\phi\left(\Sigma_{2}\right)$ are in the straight line $z+w=1$ and the vector field $X_{\mu}$ is not meromorphic there.) Moreover, since $H_{\mu}\left(\phi^{-1}(z, w)\right)$ is a Darboux first integral of $X_{\mu}$, from Remark 3.8 it follows that $\left\{X_{\mu}, \mu \in U\right\}$ is a family of vector fields verifying FLP. Consequently we can apply Proposition 3.9 to compute the asymptotic development of $T_{2}$ at $s=0$. As a matter of fact, to be precise, since the transverse


Fig. 8. Passage through the saddle at infinity in $(z, w)$-coordinates.
sections depend also on $\eta$, we shall apply it to the family $\left\{X_{(\mu, \eta)},(\mu, \eta) \in U \times(0, \varepsilon)\right\}$. Thus, following the notation of Proposition 3.9, since

$$
\lambda(\mu):=-\frac{Q(0,0 ; \mu)}{P(0,0 ; \mu)}=\frac{1}{2(F-1)}
$$

does not depend on $\eta$, it turns out that $W_{i}=U_{i} \times(0, \varepsilon)$ where

$$
\begin{aligned}
& U_{1}=\{\mu \in U: F<3 / 2\}, \quad U_{2}=\{\mu \in U: F>3 / 2\} \quad \text { and } \\
& U_{3}=\{\mu \in U: 5 / 4<F<2\} .
\end{aligned}
$$

In addition we can assert that

$$
\begin{align*}
& T_{2}(s ; \mu, \eta)=\Delta_{0}^{2}(\mu, \eta)+\Delta_{1}^{2}(\mu, \eta) s+s f_{2}^{1}(s ; \mu, \eta) \quad \text { if } \mu \in U_{1},  \tag{13}\\
& T_{2}(s ; \mu, \eta)=\Delta_{0}^{2}(\mu, \eta)+\Delta_{2}^{2}(\mu, \eta) s^{\lambda}+s^{\lambda} f_{2}^{2}(s ; \mu, \eta) \quad \text { if } \mu \in U_{2},  \tag{14}\\
& T_{2}(s ; \mu, \eta)=\Delta_{0}^{2}(\mu, \eta)+\Delta_{3}^{2}(\mu, \eta) s \omega(s ; \lambda)+\Delta_{4}^{2}(\mu, \eta) s+s f_{2}^{3}(s ; \mu, \eta) \\
& \quad \text { if } \mu \in U_{3} \tag{15}
\end{align*}
$$

where $f_{2}^{i} \in{ }_{1}\left(U_{i} \times(0, \varepsilon)\right)$. Some computations show that

$$
\begin{equation*}
\Delta_{0}^{2}(\mu, \eta)=\int_{\sigma_{2}(0)}^{0} \frac{d w}{Q(0, w)}=\int_{1-k_{2}^{2} \eta^{2}}^{0} \frac{k_{2}}{\left(p_{2}-1\right) w-k_{1}} \frac{d w}{\sqrt{1-w}} \tag{16}
\end{equation*}
$$

Let us compute next the coefficient $\Delta_{1}^{2}(\mu, \eta)$. From Proposition 3.9 we know that it is given by

$$
\Delta_{1}^{2}(\mu, \eta)=-\frac{\sigma_{2}^{\prime}(0)}{Q\left(0, \sigma_{2}(0)\right)}+\sigma_{1}^{\prime}(0) \sigma_{2}(0)^{1 / \lambda} \int_{0}^{\sigma_{2}(0)} \frac{Q_{z}(0, w) L(w)}{Q(0, w)^{2}} \frac{d w}{w^{1 / \lambda}}
$$

One can verify that

$$
\begin{aligned}
& L(w)=\left(\frac{\left(1-p_{2}\right) w+k_{1}}{1-p_{1}-k_{2}^{2} \eta^{2}\left(1-p_{2}\right)}\right)^{2 F} \text { and } \\
& \frac{Q_{z}(0, w)}{Q(0, w)^{2}}=\frac{k_{2}}{2(1-w)^{3 / 2}\left(\left(1-p_{2}\right) w+k_{1}\right)}
\end{aligned}
$$

Consequently, using also that

$$
\frac{\sigma_{2}^{\prime}(0)}{Q\left(0, \sigma_{2}(0)\right)}=\frac{\left(1-k_{2}^{2} \eta^{2}\right)^{2}}{k_{1} \eta\left(1-p_{1}-k_{2}^{2} \eta^{2}\left(1-p_{2}\right)\right)} \quad \text { and } \quad \sigma_{1}^{\prime}(0) \sigma_{2}(0)^{1 / \lambda}=\frac{\left(1-k_{2}^{2} \eta^{2}\right)^{2 F}}{k_{1}}
$$

it turns out that

$$
\begin{align*}
\Delta_{1}^{2}(\mu, \eta)= & \frac{-\left(1-k_{2}^{2} \eta^{2}\right)^{2}}{k_{1} \eta\left(1-p_{1}-k_{2}^{2} \eta^{2}\left(1-p_{2}\right)\right)} \\
& +\frac{k_{2}\left(1-k_{2}^{2} \eta^{2}\right)^{2 F}}{2 k_{1}\left(1-p_{1}-k_{2}^{2} \eta^{2}\left(1-p_{2}\right)\right)^{2 F}} \int_{0}^{1-k_{2}^{2} \eta^{2}} \frac{G(w)}{(1-w)^{3 / 2}} d w, \tag{17}
\end{align*}
$$

where $G(w):=w^{-1 / \lambda}\left(\left(1-p_{2}\right) w+k_{1}\right)^{2 F-1}$. Let us turn now to the computation of $\Delta_{2}^{2}(\mu, \eta)$, which is given by

$$
\Delta_{2}^{2}(\mu, \eta)=\sigma_{1}^{\prime}(0)^{\lambda} \sigma_{2}(0) L(0)^{\lambda}\left\{\frac{\tau_{1}(0)^{-\lambda}}{Q(0,0)}+\int_{0}^{\tau_{1}(0)}\left(\frac{M(z)}{P(z, 0)}-\frac{M(0)}{P(0,0)}\right) \frac{d z}{z^{\lambda+1}}\right\}
$$

In this case, since one can show that $M(z) \equiv 1$,

$$
\begin{aligned}
& \sigma_{1}^{\prime}(0)^{\lambda} \sigma_{2}(0) L(0)^{\lambda}=\frac{1}{k_{1}^{\lambda}}\left(\frac{k_{1}\left(1-k_{2}^{2} \eta^{2}\right)}{1-p_{1}-k_{2}^{2} \eta^{2}\left(1-p_{2}\right)}\right)^{2 \lambda F} \text { and } \\
& \frac{\tau_{1}(0)^{-\lambda}}{Q(0,0)}=\frac{-k_{2}}{k_{1}\left(1-k_{2}^{2} \eta^{2}\right)^{\lambda}}
\end{aligned}
$$

we conclude that

$$
\begin{align*}
\Delta_{2}^{2}(\mu, \eta)= & \frac{k_{2} k_{1}^{\lambda}\left(1-k_{2}^{2} \eta^{2}\right)^{2 \lambda F}}{\left(1-p_{1}-k_{2}^{2} \eta^{2}\left(1-p_{2}\right)\right)^{2 \lambda F}}\left\{\frac{-1}{\left(1-k_{2}^{2} \eta^{2}\right)^{\lambda}}\right. \\
& \left.+\lambda \int_{0}^{1-k_{2}^{2} \eta^{2}}\left(\frac{1}{\sqrt{1-z}}-1\right) \frac{d z}{z^{\lambda+1}}\right\} \tag{18}
\end{align*}
$$

Concerning the coefficient $\Delta_{3}^{2}(\mu, \eta)$, we know that if $\lambda\left(\mu_{0}\right)=1$ then

$$
\Delta_{3}^{2}\left(\mu_{0}, \eta\right)=-\sigma_{1}^{\prime}(0) \sigma_{2}(0) L(0) \frac{Q_{z}(0,0)}{P(0,0)^{2}}
$$

In our situation, using that $\lambda\left(\mu_{0}\right)=1$ corresponds to $F=3 / 2$, some computations show that

$$
\begin{equation*}
\Delta_{3}^{2}\left(\mu_{0}, \eta\right)=\frac{-k_{1} k_{2}\left(1-k_{2}^{2} \eta^{2}\right)^{3}}{2\left(1-p_{1}-k_{2}^{2} \eta^{2}\left(1-p_{2}\right)\right)^{3}} \tag{19}
\end{equation*}
$$

In order to study $R_{2}(s ; \mu, \eta)$ we will also use $(z, w)$-coordinates. Notice (see Fig. 8) that, taking the transferred parametrizations, it is precisely the Dulac map between $\phi\left(\Sigma_{1}\right)$ and $\phi\left(\Sigma_{2}^{\eta}\right)$. We point out that $\phi\left(\Sigma_{1}\right)$ is in the straight line $z+w=1$. However, in this case, this is not a problem for our purpose. Indeed, in order to study the Poincaré mapping we can apply Proposition 3.9 with the polynomial vector field

$$
\tilde{X}_{\mu}:=\frac{w}{\sqrt{1-z-w}} X_{\mu}
$$

which provides the same foliation as $X_{\mu}$ and it is obviously analytic. So we can assert that

$$
\begin{equation*}
R_{2}(s ; \mu, \eta)=\rho_{2}(\mu, \eta) s^{\lambda}+s^{\lambda} g_{2}(s ; \mu, \eta) \tag{20}
\end{equation*}
$$

where $g_{2} \in 1(U \times(-\varepsilon, \varepsilon))$ and $\rho_{2}(\mu, \eta)$ is an analytic function on $U \times(-\varepsilon, \varepsilon)$. For values of $\mu$ such that $\lambda(\mu) \approx 1$ we need more information about the remainder term in $R_{2}$. In fact, if $\mu \in U_{3}$ then

$$
\begin{equation*}
s^{\lambda} g_{2}(s ; \mu, \eta)=s \tilde{g}_{2}(s ; \mu, \eta) \quad \text { with } \tilde{g}_{2} \in 1\left(U_{3} \times(-\varepsilon, \varepsilon)\right) \tag{21}
\end{equation*}
$$

This fact does not follow from Proposition 3.9 but it is easy to show and so, for the sake of brevity, we do not prove it here. As we shall see later on, we do not need the
concrete expression of $\rho_{2}(\mu, \eta)$. We shall only use that it is convergent as $\eta \longrightarrow 0$ and this follows from its analyticity at $\eta=0$.

We can now study the composition $T_{3}\left(R_{2}(s ; \mu, \eta) ; \mu, \eta\right)$. Thus, on account of (10) and (20), by applying Lemma 3.10 we can assert that

$$
\begin{equation*}
T_{3}\left(R_{2}(s)\right)=\Delta_{1}^{3} \rho_{2} s^{\lambda}+s^{\lambda} h_{1} \quad \text { with } h_{1} \in \mathcal{I}(U \times(0, \varepsilon)) . \tag{22}
\end{equation*}
$$

In case that $\mu \in U_{3}$ we must be sharper. Since $f_{3} \in \mathcal{I}_{0}(U \times(-\varepsilon, \varepsilon))$, we have that $f_{3}=s \widehat{f_{3}}$ where $\widehat{f_{3}}$ is an analytic function on $s=0$. Hence, from (10) and (21), it follows that

$$
\begin{aligned}
T_{3}\left(R_{2}(s)\right) & =\Delta_{1}^{3}\left(\rho_{2} s^{\lambda}+s \widetilde{g}_{2}\right)+\left(\rho_{2} s^{\lambda}+s \widetilde{g}_{2}\right) f_{3}\left(\rho_{2} s^{\lambda}+s \widetilde{g}_{2}\right) \\
& =\Delta_{1}^{3}\left(\rho_{2} s^{\lambda}+s \widetilde{g}_{2}\right)+s\left(\rho_{2} s^{\lambda-1 / 2}+s^{1 / 2} \widetilde{g}_{2}\right)^{2} \widehat{f}_{3}\left(\rho_{2} s^{\lambda}+s \widetilde{g}_{2}\right)
\end{aligned}
$$

Since $\lambda(\mu)>1 / 2$ for $\mu \in U_{3}$, note that $\rho_{2} s^{\lambda-1 / 2}$ belongs to $1\left(U_{3} \times(0, \varepsilon)\right)$. Consequently the expression above shows that $T_{3}\left(R_{2}(s)\right)=\Delta_{1}^{3} \rho_{2} s^{\lambda}+s \widetilde{h}_{1}$ with $\widetilde{h}_{1} \in \mathcal{I}\left(U_{3} \times(0, \varepsilon)\right)$. Hence, using that $s^{\lambda}=(\lambda-1) s \omega(s ; \lambda)+s$, we obtain

$$
\begin{equation*}
T_{3}\left(R_{2}(s)\right)=\Delta_{1}^{3} \rho_{2}(\lambda-1) s \omega(s ; \lambda)+\Delta_{1}^{3} \rho_{2} s+s \widetilde{h}_{1} \quad \text { with } \widetilde{h}_{1} \in \mathcal{I}\left(U_{3} \times(0, \varepsilon)\right) \tag{23}
\end{equation*}
$$

We have now all the necessary ingredients to study $T(s ; \mu)$. Let us consider first the case $\mu \in U_{1}$. In this case, from (11) and (13), by applying Lemma 3.10 we obtain

$$
T_{2}\left(R_{1}(s)\right)=\Delta_{0}^{2}+\Delta_{1}^{2} \rho_{1} s+s h_{2} \quad \text { with } h_{2} \in \mathcal{I}\left(U_{1} \times(0, \varepsilon)\right)
$$

Therefore, taking (11) and (22) also into account, we get

$$
T(s ; \mu)=\Delta_{0}^{1}+\Delta_{0}^{2}+\left(\Delta_{1}^{1}+\Delta_{1}^{2} \rho_{1}\right) s+s\left(h_{2}+f_{1}+s^{\lambda-1}\left(\Delta_{1}^{3} \rho_{2}+h_{1}\right)\right)
$$

Then, using that $\lambda(\mu)>1$ for $\mu \in U_{1}$, we conclude that

$$
T(s ; \mu)=\Delta_{0}^{1}(\mu, \eta)+\Delta_{0}^{2}(\mu, \eta)+\left(\Delta_{1}^{1}(\mu, \eta)+\Delta_{1}^{2}(\mu, \eta) \rho_{1}(\mu, \eta)\right) s+s h_{3}(s ; \mu, \eta)
$$

with $h_{3} \in \mathcal{I}\left(U_{1} \times(0, \varepsilon)\right)$. At this point we stress that the coefficients

$$
\Delta_{0}(\mu):=\Delta_{0}^{1}(\mu, \eta)+\Delta_{0}^{2}(\mu, \eta) \quad \text { and } \quad \Delta_{1}(\mu):=\Delta_{1}^{1}(\mu, \eta)+\Delta_{1}^{2}(\mu, \eta) \rho_{1}(\mu, \eta)
$$

depend only on $\mu$ because $T(s ; \mu)$ does not depend on $\eta$. This proves in particular that $h_{3} \in \mathcal{I}\left(U_{1}\right)$. In order to compute explicitly these coefficients we take advantage of (12). For the first one we get

$$
\Delta_{0}(\mu)=\lim _{\eta \rightarrow 0}\left(\Delta_{0}^{1}(\mu, \eta)+\Delta_{0}^{2}(\mu, \eta)\right)=\lim _{\eta \rightarrow 0} \Delta_{0}^{2}(\mu, \eta)
$$

Thus, by applying the Dominate Convergence Theorem to the expression of $\Delta_{0}^{2}(\mu, \eta)$ given in (16) we obtain

$$
\Delta_{0}(\mu)=\int_{0}^{1} \frac{k_{2}}{\left(1-p_{2}\right) w+k_{1}} \frac{d w}{\sqrt{1-w}}=\frac{\sqrt{2}}{\sqrt{a+b+c}} \operatorname{arctanh}\left(\frac{2 a+b-\sqrt{b^{2}-4 a c}}{2 \sqrt{a(a+b+c)}}\right)
$$

The last equality above follows from direct integration and using the relation of $p_{2}, k_{1}$ and $k_{2}$ with the coefficients of $q_{\mu}(x)=a x^{2}+b x+c$. On the other hand, taking (12) into account again,

$$
\Delta_{1}(\mu)=\lim _{\eta \rightarrow 0}\left(\Delta_{1}^{1}(\mu, \eta)+\Delta_{1}^{2}(\mu, \eta) \rho_{1}(\mu, \eta)\right)=\lim _{\eta \rightarrow 0} \Delta_{1}^{2}(\mu, \eta)
$$

The computation of this limit is more delicate because $\Delta_{1}^{2}(\mu, \eta)$, which is given in (17), contains two terms that considered separately diverge as $\eta \longrightarrow 0$. To show that these cancel each other we proceed as follows.

$$
\begin{aligned}
\Delta_{1}^{2}(\mu, \eta)= & \frac{-\left(1-k_{2}^{2} \eta^{2}\right)^{2}}{k_{1} \eta\left(1-p_{1}-k_{2}^{2} \eta^{2}\left(1-p_{2}\right)\right)}+\frac{k_{2}\left(1-k_{2}^{2} \eta^{2}\right)^{2 F}}{2 k_{1}\left(1-p_{1}-k_{2}^{2} \eta^{2}\left(1-p_{2}\right)\right)^{2 F}} \\
& \times\left\{2 G(1) \frac{1-k_{2} \eta}{k_{2} \eta}+\int_{0}^{1-k_{2}^{2} \eta^{2}} \frac{G(w)-G(1)}{(1-w)^{3 / 2}} d w\right\} \\
= & \frac{\left(1-k_{2}^{2} \eta^{2}\right)^{2}}{k_{1} \eta\left(1-p_{1}-k_{2}^{2} \eta^{2}\left(1-p_{2}\right)\right)}\left\{-1+\frac{\left(1-k_{2}^{2} \eta^{2}\right)^{1 / \lambda}\left(1-k_{2} \eta\right)\left(1-p_{1}\right)^{2 F-1}}{\left(1-p_{1}-k_{2}^{2} \eta^{2}\left(1-p_{2}\right)\right)^{2 F-1}}\right\} \\
& +\frac{k_{2}\left(1-k_{2}^{2} \eta^{2}\right)^{2 F}}{2 k_{1}\left(1-p_{1}-k_{2}^{2} \eta^{2}\left(1-p_{2}\right)\right)^{2 F}} \int_{0}^{1-k_{2}^{2} \eta^{2}} \frac{G(w)-G(1)}{(1-w)^{3 / 2}} d w .
\end{aligned}
$$

Now, to compute the limit we apply L'Hôpital's rule to the first term and the Dominate Convergence Theorem to the second one. It can be shown in this way that

$$
\begin{aligned}
\Delta_{1}(\mu) & =\lim _{\eta \rightarrow 0} \Delta_{1}^{2}(\mu, \eta)=\frac{k_{2}}{k_{1}\left(p_{1}-1\right)}+\frac{k_{2}}{2 k_{1}\left(1-p_{1}\right)^{2 F}} \int_{0}^{1} \frac{G(w)-G(1)}{(1-w)^{3 / 2}} d w \\
& =\frac{k_{2}}{2 k_{1}\left(p_{1}-1\right)}\left\{2-\int_{0}^{1}\left(w^{-\frac{1}{\lambda}}\left(\frac{1-p_{2}}{1-p_{1}}(w-1)+1\right)^{2 F-1}-1\right) \frac{d w}{(1-w)^{3 / 2}}\right\},
\end{aligned}
$$

and this concludes the proof of the assertion in (a).
Let us turn now to the case $\mu \in U_{2}$. In this case, from (11) and (14), by applying Lemma 3.10 we obtain that

$$
T_{2}\left(R_{1}(s)\right)=\Delta_{0}^{2}+\Delta_{2}^{2} \rho_{1}^{\lambda} s^{\lambda}+s^{\lambda} h_{4} \quad \text { with } h_{4} \in \mathcal{I}\left(U_{2} \times(0, \varepsilon)\right)
$$

The combination of (11) and (22) with the above expression shows that

$$
T(s ; \mu)=\Delta_{0}^{1}+\Delta_{0}^{2}+\left(\Delta_{1}^{3} \rho_{2}+\Delta_{2}^{2} \rho_{1}^{\lambda}\right) s^{\lambda}+s^{\lambda} h_{5}
$$

where $h_{5}:=h_{1}+h_{4}+s^{1-\lambda}\left(\Delta_{1}^{1}+f_{1}\right)$ is a function that belongs to $\mathcal{I}\left(U_{2} \times(0, \varepsilon)\right)$. This assertion follows from using that $\lambda(\mu)<1$ for $\mu \in U_{2}$. Note in addition that, since $T(s ; \mu)$ and $s^{\lambda(\mu)}$ do not depend on $\eta$, the coefficients

$$
\Delta_{0}(\mu):=\Delta_{0}^{1}(\mu, \eta)+\Delta_{0}^{2}(\mu, \eta) \quad \text { and } \quad \Delta_{2}(\mu):=\Delta_{1}^{3}(\mu, \eta) \rho_{2}(\mu, \eta)+\Delta_{2}^{2}(\mu, \eta) \rho_{1}(\mu, \eta)^{\lambda(\mu)}
$$

depend only on $\mu$. On the other hand, since $\rho_{2}(\mu, \eta)$ is analytic at $\eta=0$, from (12) it turns out that

$$
\Delta_{2}(\mu)=\lim _{\eta \rightarrow 0}\left(\Delta_{1}^{3}(\mu, \eta) \rho_{2}(\mu, \eta)+\Delta_{2}^{2}(\mu, \eta) \rho_{1}(\mu, \eta)^{\lambda(\mu)}\right)=\lim _{\eta \rightarrow 0} \Delta_{2}^{2}(\mu, \eta)
$$

Thus, by applying the Dominate Convergence Theorem to the expression of $\Delta_{2}^{2}(\mu, \eta)$ given in (18), one can easily show that

$$
\Delta_{2}(\mu)=\frac{k_{2} k_{1}^{\lambda}}{\left(1-p_{1}\right)^{2 \lambda F}}\left\{\lambda \int_{0}^{1}\left(\frac{1}{\sqrt{1-z}}-1\right) \frac{d z}{z^{\lambda+1}}-1\right\}=\frac{k_{2} k_{1}^{\lambda}}{\left(1-p_{1}\right)^{2 \lambda F}} \frac{\lambda \sqrt{\pi} \Gamma(-\lambda)}{\Gamma\left(\frac{1}{2}-\lambda\right)}
$$

The last equality above follows from direct integration. This proves (b).
Let us study finally the case $\mu \in U_{3}$. In this case, from (11) and (15), by applying Lemma 3.10 we obtain

$$
T_{2}\left(R_{1}(s)\right)=\Delta_{0}^{2}+\left(\Delta_{3}^{2} \rho_{1}^{\lambda}\right) s \omega(s ; \lambda)+\left(\Delta_{3}^{2} \rho_{1} \omega\left(\rho_{1} ; \lambda\right)+\Delta_{4}^{2} \rho_{1}\right) s+s h_{6}
$$

with $h_{6} \in \mathcal{I}\left(U_{3} \times(0, \varepsilon)\right)$. Therefore, taking (11) and (23) also into account, we get

$$
\begin{aligned}
T(s ; \mu)= & \Delta_{0}^{1}+\Delta_{0}^{2}+\left(\Delta_{3}^{2} \rho_{1}^{\lambda}+\Delta_{1}^{3} \rho_{2}(\lambda-1)\right) s \omega(s ; \lambda) \\
& +\left(\Delta_{1}^{1}+\Delta_{3}^{2} \rho_{1} \omega\left(\rho_{1} ; \lambda\right)+\Delta_{4}^{2} \rho_{1}+\Delta_{1}^{3} \rho_{2}\right) s+s h_{7}
\end{aligned}
$$

where $h_{7}:=f_{1}+\tilde{h}_{1}+h_{6}$ is a function that belongs to $\mathcal{I}\left(U_{3} \times(0, \varepsilon)\right)$. On the other hand, since $T(s ; \mu)$ and $\omega(s ; \lambda(\mu))$ depend only on $\mu$, the coefficients

$$
\begin{aligned}
& \Delta_{0}(\mu):=\Delta_{0}^{1}(\mu, \eta)+\Delta_{0}^{2}(\mu, \eta) \\
& \Delta_{3}(\mu):=\Delta_{3}^{2}(\mu, \eta) \rho_{1}(\mu, \eta)^{\lambda(\mu)}+\Delta_{1}^{3}(\mu, \eta) \rho_{2}(\mu, \eta)(\lambda(\mu)-1)
\end{aligned}
$$

and

$$
\begin{aligned}
\Delta_{4}(\mu):= & \Delta_{1}^{1}(\mu, \eta)+\Delta_{3}^{2}(\mu, \eta) \rho_{1}(\mu, \eta) \omega\left(\rho_{1}(\mu, \eta) ; \lambda(\mu)\right) \\
& +\Delta_{4}^{2}(\mu, \eta) \rho_{1}(\mu, \eta)+\Delta_{1}^{3}(\mu, \eta) \rho_{2}(\mu, \eta)
\end{aligned}
$$

do not depend on $\eta$. This implies in particular that $h_{7}(s ; \mu, \eta)$ does not depend on $\eta$, and so we can assert that $h_{7} \in \mathcal{I}\left(U_{3}\right)$. Finally, if we consider some $\mu_{0} \in U_{3}$ such that $\lambda\left(\mu_{0}\right)=1$, then $\Delta_{3}\left(\mu_{0}\right)=\Delta_{3}^{2}\left(\mu_{0}, \eta\right) \rho_{1}\left(\mu_{0}, \eta\right)$ and consequently, from (12),

$$
\Delta_{3}\left(\mu_{0}\right)=\lim _{\eta \rightarrow 0} \Delta_{3}^{2}\left(\mu_{0}, \eta\right) \rho_{1}\left(\mu_{0}, \eta\right)=\lim _{\eta \rightarrow 0} \Delta_{3}^{2}\left(\mu_{0}, \eta\right)
$$

Thus, on account of (19), it follows that

$$
\Delta_{3}\left(\mu_{0}\right)=\frac{k_{1} k_{2}}{2\left(p_{1}-1\right)^{3}}
$$

This proves (c) and concludes the proof of the result.
It is clear that the sign of $P_{S}(s ; \mu)$ for small positive $s$ determines the monotonicity of the period function near the outer boundary of the period annulus. So we need to study the coefficients of the second monomial in the asymptotic development given in Theorem 3.6. To this end we introduce the sets

$$
\begin{align*}
& \Gamma_{1}:=\left\{\mu \in U_{1}: \Delta_{1}(\mu)=0\right\} \\
& \Gamma_{2}:=\left\{\mu \in U_{2}: \Delta_{2}(\mu)=0\right\} \\
& \Gamma_{3}:=\left\{\mu \in U_{3}: \Delta_{3}(\mu)=0 \text { with } \lambda(\mu)=1\right\} . \tag{24}
\end{align*}
$$

One can easily verify that $\Gamma_{2}=\left\{\mu \in U_{2}: F=2\right\}$ and that $\Gamma_{3}$ is empty. Fig. 5 shows the set $\Gamma_{1}$ computed numerically. We are now in position to prove the main result of this subsection:

Proof of Theorem 3.3. Fix some $\mu^{\star} \in U \backslash\left\{\Gamma_{1} \cup \Gamma_{2}\right\}$ and note that, taking (6) and (24) into account, there are three different situations to consider:
(a) $\mu^{\star} \in U_{1} \backslash \Gamma_{1}$,
(b) $\mu^{\star} \in U_{2} \backslash \Gamma_{2}$,
(c) $\mu^{\star} \in U_{3}$ such that $\lambda\left(\mu^{\star}\right)=1$.

The fact that $\mu^{\star}$ is a local regular value in the cases (a) and (b) follows exactly the same way as in the proof of Theorem 5.1 in [13]. So let us consider only the case (c), which corresponds to the values of $\mu \in U$ such that $F=3 / 2$. Note first of all that, from (c) in Theorem 3.6, we can assert that if $\mu \in U_{3}$ then

$$
P_{s}(s ; \mu)=\Delta_{3}(\mu)(\lambda \omega(s ; \lambda)+1)+\Delta_{4}(\mu)+s f_{3}^{\prime}(s ; \mu)+f_{3}(s ; \mu),
$$

where $f_{3} \in \mathcal{I}\left(U_{3}\right)$. Here we used that, on account of Definition 3.5, $s \omega_{s}=(\lambda-1) \omega+1$. On the other hand, since $\lambda(\mu) \longrightarrow 1$ as $\mu \longrightarrow \mu^{\star}$, it is clear that $\omega(s ; \lambda) \longrightarrow-\infty$ as $(s, \mu) \longrightarrow\left(0, \mu^{\star}\right)$. Consequently, using also that $f_{3} \in \mathcal{I}\left(U_{3}\right)$, from the above equality we obtain that

$$
\frac{P_{s}(s ; \mu)}{\lambda \omega(s ; \lambda)+1} \longrightarrow \Delta_{3}\left(\mu^{\star}\right) \quad \text { as }(s, \mu) \longrightarrow\left(0, \mu^{\star}\right)
$$

Therefore, since one can easily verify that $\Delta_{3}\left(\mu^{\star}\right)<0$, we can assert that there exists a neighborhood $U^{\star}$ of $\mu^{\star}$ and $\varepsilon>0$ such that $P_{s}(s ; \mu)>0$ for all $s \in(0, \varepsilon)$ and $\mu \in U^{\star}$. According to (b) in Remark 2.6, this proves that $\mu^{\star}$ is a local regular value. It also shows that the period function is monotonous decreasing on the outer boundary of $\mathcal{P}_{\mu}$. Indeed, $P(s ; \mu)$ is by definition the period of the periodic orbit of $X_{\mu}$ passing through the point $\left(p_{1}-s, 0\right)$, which approaches to the outer boundary as $s$ decreases.

The assertions concerning the monotonicity in the cases (a) and (b) follow exactly the same way taking into account the sign of $\Delta_{1}\left(\mu^{\star}\right)$ and $\Delta_{2}\left(\mu^{\star}\right)$, respectively.

The rest of the subsection is devoted to show some properties of $\Gamma_{1}$. We prove the following:

Proposition 3.11. The set $\Gamma_{1}$ is the graphic of an analytic function $D=\mathcal{G}(F)$ defined for $F \in(1,3 / 2)$ that has the following properties:
(a) $-F<\mathcal{G}(F)<-1 / 2$ for all $F \in(1,3 / 2)$,
(b) $\mathcal{G}(F) \longrightarrow-3 / 2$ as $F \nearrow 3 / 2$,
(c) $\mathcal{G}(F) \longrightarrow-1 / 2$ as $F \searrow 1$,
(d) $\mathcal{G}(5 / 4)=-1$.

This result follows almost directly from Lemma 3.13. However, in order to prove Lemma 3.13 we shall need a previous result concerning a general property of the coefficients in Proposition 3.9.

Lemma 3.12. Under the hypothesis of Proposition 3.9, let $\left\{\mu_{k}\right\}$ be a sequence of parameters in $W_{1}$ (respectively $W_{2}$ ) so that $\mu_{k} \longrightarrow \widehat{\mu}$ with $\lambda(\widehat{\mu})=1 / n$ and $\Delta_{3}(\widehat{\mu}) \neq 0$.
(a) If $\Delta_{3}(\widehat{\mu})>0$ then $\Delta_{1}\left(\mu_{k}\right)$ (respectively $\Delta_{2}\left(\mu_{k}\right)$ ) tends to $-\infty$ as $\mu_{k} \longrightarrow \widehat{\mu}$.
(b) If $\Delta_{3}(\widehat{\mu})<0$ then $\Delta_{1}\left(\mu_{k}\right)$ (respectively $\Delta_{2}\left(\mu_{k}\right)$ ) tends to $+\infty$ as $\mu_{k} \longrightarrow \widehat{\mu}$.

Proof. We shall prove (a) and (b) for a sequence $\left\{\mu_{k}\right\}$ in $W_{1}$ (the other case follows exactly the same way). Notice first that, on account of $\mu_{k} \in W_{1}$, we have $\lambda\left(\mu_{k}\right)>1 / n$ and

$$
\begin{equation*}
T\left(s ; \mu_{k}\right)=\Delta_{0}\left(\mu_{k}\right)+\Delta_{1}\left(\mu_{k}\right) s+s f_{1}\left(s ; \mu_{k}\right) \quad \text { with } f_{1} \in \mathcal{I}\left(W_{1}\right) \tag{25}
\end{equation*}
$$

On the other hand, note that $\mu_{k} \in W_{3}$ for $k$ large enough because $\mu_{k} \longrightarrow \widehat{\mu} \in W_{3}$. Therefore

$$
\begin{equation*}
T\left(s ; \mu_{k}\right)=\Delta_{0}\left(\mu_{k}\right)+\Delta_{3}\left(\mu_{k}\right) s \omega\left(s ; \lambda\left(\mu_{k}\right) n\right)+\Delta_{4}\left(\mu_{k}\right) s+s f_{3}\left(s ; \mu_{k}\right) \quad \text { with } f_{3} \in \mathcal{I}\left(W_{3}\right) . \tag{26}
\end{equation*}
$$

By definition

$$
s \omega\left(s ; \lambda\left(\mu_{k}\right) n\right)=\frac{s^{\lambda\left(\mu_{k}\right) n}-s}{\lambda\left(\mu_{k}\right) n-1}
$$

and consequently, from (26),

$$
T\left(s ; \mu_{k}\right)=\Delta_{0}\left(\mu_{k}\right)+\left(\Delta_{4}\left(\mu_{k}\right)-\frac{\Delta_{3}\left(\mu_{k}\right)}{\lambda\left(\mu_{k}\right) n-1}\right) s+s g\left(s ; \mu_{k}\right),
$$

where

$$
g(s ; \mu):=f_{3}(s ; \mu)+\frac{\Delta_{3}(\mu)}{\lambda(\mu) n-1} s^{\lambda(\mu) n-1}
$$

is a function that belongs to $\mathcal{I}\left(W_{1}\right)$. Consequently the combination of this expression for $T\left(s ; \mu_{k}\right)$ and the one in (25) shows that

$$
\begin{equation*}
\Delta_{1}\left(\mu_{k}\right)=\Delta_{4}\left(\mu_{k}\right)-\frac{\Delta_{3}\left(\mu_{k}\right)}{\lambda\left(\mu_{k}\right) n-1} \tag{27}
\end{equation*}
$$

Note that, since $\Delta_{3}$ and $\Delta_{4}$ are analytic on $W_{3}, \Delta_{3}\left(\mu_{k}\right) \longrightarrow \Delta_{3}(\widehat{\mu})$ and $\Delta_{4}\left(\mu_{k}\right) \longrightarrow \Delta_{4}(\widehat{\mu})$ as $\mu_{k} \longrightarrow \widehat{\mu}$. In addition, due to $\mu_{k} \in W_{1}$, it turns out that $\lambda\left(\mu_{k}\right) n-1 \searrow 0$ as $\mu_{k} \longrightarrow \widehat{\mu}$. Hence, from (27), we conclude that

$$
\lim _{\mu_{k} \rightarrow \widehat{\mu}} \Delta_{1}\left(\mu_{k}\right)=-\infty \text { if } \Delta_{3}(\widehat{\mu})>0 \quad \text { and } \quad \lim _{\mu_{k} \rightarrow \widehat{\mu}} \Delta_{1}\left(\mu_{k}\right)=+\infty \text { if } \Delta_{3}(\widehat{\mu})<0
$$

as claimed.
In what follows we shall use the notation $k_{1}=p_{2}-p_{1}$ and $k_{2}=1 / \sqrt{2 a}$ introduced in the proof of Theorem 3.6. Let us also define $\Psi(\mu)$ by means of the relation $\Delta_{1}(\mu)=$ $\frac{k_{2}}{2 k_{1}\left(p_{1}-1\right)} \Psi(\mu)$, that is,

$$
\begin{gather*}
\Psi(\mu):=2-\int_{0}^{1}\left(u^{2(1-F)}((u-1) \kappa+1)^{2 F-1}-1\right) \frac{d u}{(1-u)^{3 / 2}} \\
\text { where } \kappa(\mu):=\frac{1-p_{2}}{1-p_{1}} . \tag{28}
\end{gather*}
$$

Concerning this function we prove the following:
Lemma 3.13. If $(D, F) \in U_{1}$ then the following holds:
(a) $\Psi_{D}(D, F)<0$,
(b) $\Psi(D, F)<0$ for $D \geqslant-1 / 2$,
(c) $\Psi(-1,5 / 4)=0$,
(d) $\Psi(D, F) \longrightarrow 4$ as $(D, F) \longrightarrow(-q, q)$ with $1<q<3 / 2$,
(e) $\Psi(D, F) \longrightarrow-\infty$ as $(D, F) \longrightarrow(q, 3 / 2)$ with $-3 / 2<q<0$,
(f) $\frac{\Psi\left(D, 1+(D+1 / 2)^{2}\right)}{(D+1 / 2)^{2}} \longrightarrow 4(4-\pi)$ as $D \quad \nearrow-1 / 2$.

Proof. Some computations show that

$$
\begin{equation*}
\kappa=\frac{(2 D+1) \sqrt{F(F-1)}+\sqrt{(F+D)(F-D-1)}}{(2 D+1) \sqrt{F(F-1)}-\sqrt{(F+D)(F-D-1)}} \tag{29}
\end{equation*}
$$

and

$$
\frac{d \kappa}{d D}=\frac{-(2 F-1)^{2}}{((2 D+1) \sqrt{F(F-1)}-\sqrt{(F+D)(F-D-1)})^{2}} \sqrt{\frac{F(F-1)}{(F+D)(F-D-1)}} .
$$

Thus, from the last expression above it follows that

$$
\Psi_{D}(\mu)=(2 F-1) \frac{d \kappa}{d D} \int_{0}^{1}\left(\frac{(u-1) \kappa+1}{u}\right)^{2(F-1)} \frac{d u}{(1-u)^{1 / 2}}
$$

is negative for $\mu \in U_{1}$. This proves (a). Let us turn next to the assertion in (b). Notice first that

$$
\begin{equation*}
\Psi(\mu)=\int_{0}^{1} \frac{2-u-u^{2(1-F)}((u-1) \kappa+1)^{2 F-1}}{(1-u)^{3 / 2}} d u=\int_{0}^{1} \frac{2-u(h(u ; \mu)+1)}{(1-u)^{3 / 2}} d u \tag{30}
\end{equation*}
$$

where

$$
h(u ; \mu):=\left(\frac{(u-1) \kappa+1}{u}\right)^{2 F-1}
$$

On the other hand one can verify that $\kappa \leqslant-1$ for $D \geqslant-1 / 2$. Taking this into account it is easy to show that

$$
\frac{(u-1) \kappa+1}{u}>\frac{2-u}{u}>1 \quad \text { for } u \in(0,1) .
$$

Therefore, since $F>1$, we have that

$$
h(u ; \mu)>\left(\frac{2-u}{u}\right)^{2 F-1}>\frac{2-u}{u}
$$

and this, on account of (30), proves (b). The assertion in (c) is straightforward because $\kappa(-1, F)=0$ and direct integration yields

$$
\Psi(-1,5 / 4)=2-\int_{0}^{1} \frac{u^{-1 / 2}-1}{(1-u)^{3 / 2}}=0 .
$$

To show (d) note first that $\kappa(-q, q)=1$. Thus, if $(D, F) \longrightarrow(-q, q)$ with $1<q<3 / 2$, then

$$
\Psi(D, F) \longrightarrow 2+\int_{0}^{1} \frac{d u}{\sqrt{1-u}}=4
$$

Here we apply the Dominate Convergence Theorem and to do so it is necessary that $1<q<3 / 2$. In order to prove (e) we shall apply Lemma 3.12 to the results obtained in Theorem 3.6. To this end notice first that the parameter $\widehat{\mu}:=(q, 3 / 2)$ satisfies $\lambda(\widehat{\mu})=1$. Hence (c) in Theorem 3.6 shows that

$$
\Delta_{3}(\widehat{\mu})=\frac{k_{1} k_{2}}{2\left(p_{1}-1\right)^{3}},
$$

which is negative because $0<p_{1}<1$ and $k_{i}>0$. Thus, if we consider any sequence $\left\{\mu_{k}\right\}$ in $U_{1}$ with $\mu_{k} \longrightarrow \widehat{\mu}$ then, by applying Lemma 3.12, it follows that

$$
\begin{equation*}
\Delta_{1}\left(\mu_{k}\right) \longrightarrow+\infty \quad \text { as } \mu_{k} \longrightarrow \widehat{\mu} . \tag{31}
\end{equation*}
$$

Recall at this point that, by definition,

$$
\begin{equation*}
\Delta_{1}(\mu)=\frac{k_{2}}{2 k_{1}\left(p_{1}-1\right)} \Psi(\mu) \tag{32}
\end{equation*}
$$

One can verify moreover that if $(D, F) \longrightarrow(q, 3 / 2)$ with $-3 / 2<q<0$, then

$$
\frac{k_{2}}{2 k_{1}\left(p_{1}-1\right)} \longrightarrow \frac{18}{6 q+3-\sqrt{9-12 q-12 q^{2}}}\left(\frac{2 q^{3}}{12 q^{2}+12 q-9}\right)^{1 / 2}<0
$$

Hence, on account of (31) and (32), we get that $\Psi(D, F) \longrightarrow-\infty$ as $(D, F) \longrightarrow(q, 3 / 2)$. This proves (e). Finally, to show $(f)$ we shall take advantage of the expression of $\Psi(\mu)$ given in (30). Let us define

$$
g(u ; D):=h\left(u ;\left(D, 1+(D+1 / 2)^{2}\right)\right)
$$

The idea will be to use the Taylor's development of $g(u ; D)$ at $D=-1 / 2$. We point out, however, that the term $\sqrt{F-1}$ in $\kappa$, see (29), makes that $\kappa\left(D, 1+(D+1 / 2)^{2}\right)$ is not smooth enough at $D=-1 / 2$. Nevertheless, this will not be a problem for our purpose because we only need to study its behavior as $D \nearrow-1 / 2$, and it is clear that this function coincides on $(-1 / 2-\varepsilon,-1 / 2]$ with an analytic function at $D=-1 / 2$. In the sequel, when we study $g(u ; D)$, we will use this analytic function instead of the original $\kappa\left(D, 1+(D+1 / 2)^{2}\right)$. Keeping this in mind, to avoid cumbersome notation we shall maintain the name of the functions. Now, since one can check that

$$
g(u ; D)=\frac{2-u}{u}+\frac{1}{2} \frac{d^{2}}{d D^{2}} g\left(u ; \xi_{D}\right)(D+1 / 2)^{2} \quad \text { with } \xi_{D} \in(D,-1 / 2)
$$

from (30) we obtain that

$$
\begin{equation*}
\frac{\Psi\left(D, 1+(D+1 / 2)^{2}\right)}{(D+1 / 2)^{2}}=-\frac{1}{2} \int_{0}^{1} u \frac{d^{2}}{d D^{2}} g\left(u ; \xi_{D}\right) \frac{d u}{(1-u)^{3 / 2}} \tag{33}
\end{equation*}
$$

Lengthy computations, which are not included here for the sake of brevity, allow to verify that, for all $D \in(-1 / 2-\varepsilon,-1 / 2]$,

$$
\left|\frac{u}{(1-u)^{3 / 2}} \frac{d^{2}}{d D^{2}} g(u ; D)\right|<f(u) \quad \text { with } f \in L^{1}((0,1))
$$

Therefore, by applying the Dominate Convergence Theorem, from (33) it turns out that

$$
\frac{\Psi\left(D, 1+(D+1 / 2)^{2}\right)}{(D+1 / 2)^{2}} \longrightarrow-\frac{1}{2} \int_{0}^{1} u \frac{d^{2}}{d D^{2}} g(u ;-1 / 2) \frac{d u}{(1-u)^{3 / 2}} \quad \text { as } D \nearrow-1 / 2
$$

Finally, since one can check that

$$
\frac{d^{2}}{d D^{2}} g(u ;-1 / 2)=\frac{4}{u}\left((2-u) \ln \left(\frac{2-u}{u}\right)+4(u-1)\right)
$$

and

$$
-2 \int_{0}^{1}\left((2-u) \ln \left(\frac{2-u}{u}\right)+4(u-1)\right) \frac{d u}{(1-u)^{3 / 2}}=4(4-\pi)
$$

the result follows.
Proof of Proposition 3.11. Recall that, by definition,

$$
\Delta_{1}(\mu)=\frac{k_{2}}{2 k_{1}\left(p_{1}-1\right)} \Psi(\mu)
$$

Then, since one can verify that $\frac{k_{2}}{2 k_{1}\left(p_{1}-1\right)}$ does not vanish on $U_{1}$, it suffices to study $\left\{\mu \in U_{1}: \Psi(\mu)=0\right\}$. Notice first that, by applying the Implicit Function Theorem, (a) and (c) in Lemma 3.13 show that this set is the graphic of an analytic function $D=\mathcal{G}(F)$ with $\mathcal{G}(5 / 4)=-1$. The fact that $\mathcal{G}$ is defined for all $F \in(1,3 / 2)$ and that $-F<\mathcal{G}(F)<-1 / 2$ follow from (b) and (d) in Lemma 3.13. On the other hand, by applying (e) in Lemma 3.13 we can assert that $\mathcal{G}(F) \longrightarrow-3 / 2$ as $F \nearrow 3 / 2$. So it only remains to prove (c). To this end note that $(f)$ in Lemma 3.13 implies that

$$
\Psi\left(D, 1+(D+1 / 2)^{2}\right)>0 \quad \text { for all } D \in(-1 / 2-\varepsilon,-1 / 2)
$$

In addition, from (d) in Lemma 3.13, we have that $\Psi(-1 / 2, F)<0$ for all $F \in$ $(1,1+\varepsilon)$. Consequently, by Bolzano's Theorem, in any neighborhood $V$ of $(D, F)=$ $(-1 / 2,1)$ there exists some $\mu \in U_{1} \cap V$ such that $\Psi(\mu)=0$. This shows (c) and concludes the proof of the result.
3.2.2. The case $0<F<1$ and $-1<D<0$.

The aim of this subsection is only to recall the results that we obtain in [13] concerning the period function of the center at the origin of $X_{\mu}$ in case (see Fig. 9) that $\mu$ belongs to

$$
W:=\left\{(D, F) \in \mathbb{R}^{2}:-1<D<0 \text { and } 0<F<1\right\} .
$$



Fig. 9. Monotonicity of the period function at the outer boundary of $\mathcal{P}_{\mu}$.

Setting

$$
\Gamma_{3}:=\left\{\mu \in W: D=-\frac{1}{2}, F \in\left(\frac{1}{2}, 1\right)\right\} \quad \text { and } \quad \Gamma_{4}:=\left\{\mu \in W: F=\frac{1}{2}\right\}
$$

from Theorem 5.1 and Proposition 5.2 in [13] it follows the next result:
Theorem 3.14. Let $\left\{X_{\mu}, \mu \in \mathbb{R}^{2}\right\}$ be the family of vector fields in (2) and consider the period function of the center at the origin. Then the open set $W \backslash\left\{\Gamma_{3} \cup \Gamma_{4}\right\}$ corresponds to local regular values of the period function at the outer boundary of the period annulus. Moreover, for these parameters, the period function is monotonous near the outer boundary and the corresponding character is shown in Fig. 9.

### 3.3. Bounded period annulus

In this section we study the period function of the center at the origin for the parameter values $\mu$ such that $\mathcal{P}_{\mu}$ is bounded. Notice that among these parameters there are two main situations to consider (see Fig. 4). The first one are those parameters such that the outer boundary of $P_{\mu}$ is a saddle loop, which corresponds to

$$
M:=\left\{\mu \in \mathbb{R}^{2}: D<-1 \text { and } F+D<0\right\} \cup\left\{\mu \in \mathbb{R}^{2}: D>0 \text { and } F+D>0\right\} .
$$

The second one are those parameters such that it is a bicycle, which corresponds to

$$
N:=\left\{\mu \in \mathbb{R}^{2}:-1<D<0 \text { and } F<0\right\} \cup\left\{\mu \in \mathbb{R}^{2}: D>0 \text { and } F+D<0\right\} .
$$



Fig. 10. Parameters corresponding to bounded period annulus.

Now, with these definitions (see Fig. 10), our goal is to prove the following result:
Theorem 3.15. Let $\left\{X_{\mu}, \mu \in \mathbb{R}^{2}\right\}$ be the family of vector fields in (2) and consider the period function of the center at the origin. Then the open set $M \cup N$ corresponds to local regular values of the period function at the outer boundary of the period annulus. Moreover, for these parameters, the period function is monotonous increasing near the outer boundary.

Proof. Let us begin by recalling some facts about the normal form of a family of vector fields with a hyperbolic saddle. So let us consider a $\mathcal{C}^{\infty}$-family of vector fields $\left\{X_{\mu}: \mu \in W\right\}$ with a hyperbolic saddle $p_{\mu}$. Let the eigenvalues of $X_{\mu}$ at $p_{\mu}$ be $\lambda_{1}(\mu)$ and $\lambda_{2}(\mu)$, with $\lambda_{2}<0<\lambda_{1}$, and let $r(\mu):=-\frac{\lambda_{2}(\mu)}{\lambda_{1}(\mu)}$ be its ratio of hyperbolicity. Fix some $\mu_{0}$ and assume first that $r\left(\mu_{0}\right)$ is rational, i.e., $r\left(\mu_{0}\right)=\frac{p}{q}$ with $(p, q)=1$. Then, for each $k \in \mathbb{N}$, there exists a $\mathcal{C}^{k}$-diffeomorphism $\Phi$ such that, in some neighborhood of $p_{\mu}$ and for $\mu \approx \mu_{0}$,

$$
\begin{equation*}
X_{\mu}=\Phi_{*}\left(\frac{1}{f(u ; \mu)}\left(x \partial_{x}-y g(u ; \mu) \partial_{y}\right)\right) \tag{34}
\end{equation*}
$$

where $u=x^{p} y^{q}$ and

$$
\begin{aligned}
& f(u ; \mu)=\frac{1}{\lambda_{1}(\mu)}+\beta_{1}(\mu) u+\cdots+\beta_{n_{k}} u^{n_{k}}, \\
& g(u ; \mu)=r(\mu)+\alpha_{1}(\mu) u+\cdots+\alpha_{n_{k}} u^{n_{k}} .
\end{aligned}
$$

If $r\left(\mu_{0}\right) \notin \mathbb{Q}$ then the above results holds with $\alpha_{i}(\mu) \equiv 0$ and $\beta_{i}(\mu) \equiv 0$ for all $i$.

Let us first take a parameter $\mu^{\star}$ such that the period annulus is bounded by a saddle loop (i.e., $\mu \in M$ ). Notice that, according to Remark 3.1, there is a neigbourhood $U^{\star}$ of $\mu^{\star}$ such that the saddle loop persists and it is the outer boundary of $\mathcal{P}_{\mu}$ for all $\mu \in U^{\star}$. Moreover in Section 3.1 we showed that the saddle is located at $\left(-\frac{1}{D}, 0\right)$ and has eigenvalues

$$
\lambda_{1}(\mu)=\sqrt{1+\frac{1}{D}} \quad \text { and } \quad \lambda_{2}(\mu)=-\sqrt{1+\frac{1}{D}}
$$

and so $r(\mu) \equiv 1$. Since in particular $r\left(\mu^{\star}\right)=1$, there exists a local diffeomorphism $\Phi$ such that it holds (34) with $p=q=1$. We introduce two transversal sections $\Sigma_{\sigma}$ and $\Sigma_{\tau}$ given by $s \longmapsto \Phi(s, 1)$ and $s \longmapsto \Phi(1, s)$, respectively. Let us also define $P(s ; \mu)$ as the period of the periodic orbit of $X_{\mu}$ passing through $\Phi(s, 1)$. We split it as

$$
P(s ; \mu)=T_{1}(s ; \mu)+T_{2}(s ; \mu)
$$

where $T_{1}$ is the time function for $-X_{\mu}$ from $\Sigma_{\sigma}$ to $\Sigma_{\tau}$ and $T_{2}$ is the time function for $X_{\mu}$ from $\Sigma_{\sigma}$ to $\Sigma_{\tau}$ (i.e., the passage through the saddle). To be more precise, $T_{1}(s ; \mu)$ is the minimum positive time necessary so that the solution of $-X_{\mu}$ passing through $\Phi(s, 1) \in \Sigma_{\sigma}$ reaches $\Sigma_{\tau}$. Note that $T_{1}(s ; \mu)$ is a smooth function on $s=0$. The time function associated to the passage through a saddle with $r\left(\mu^{\star}\right)=1$ has already been studied. Indeed, Lemma 2 in [1] shows that

$$
\begin{equation*}
T_{2}(s ; \mu)=-\frac{1}{\lambda_{1}(\mu)} \log s-\beta_{1}(\mu) s \omega(s ; r(\mu))+\psi(s ; \mu) \tag{35}
\end{equation*}
$$

where $\psi$ is a $\mathcal{C}^{1}$ function at $(s, \mu)=\left(0, \mu^{\star}\right)$ and 1-flat at $s=0$ for all $\mu$. In our situation $\omega(s ; r(\mu))=\log s$ because in fact $r(\mu) \equiv 1$. From the above expression we obtain that

$$
s \frac{d}{d s} T_{2}(s ; \mu) \longrightarrow-1 / \lambda_{1}\left(\mu^{\star}\right) \quad \text { as }(s, \mu) \longrightarrow\left(0, \mu^{\star}\right)
$$

This shows, since $T_{1}$ is smooth at $s=0$, that $s P_{s}(s ; \mu) \longrightarrow-1 / \lambda_{1}\left(\mu^{\star}\right)$ as $(s, \mu) \longrightarrow$ $\left(0, \mu^{\star}\right)$. Therefore, due to $\lambda_{1}\left(\mu^{\star}\right)>0$, we can assert that there exists $\varepsilon>0$ such that $P_{s}(s ; \mu)<0$ for all $s \in(0, \varepsilon)$ and $\mu \approx \mu^{\star}$. Consequently, by (b) in Remark 2.6, $\mu^{\star}$ is a local regular value of the period function at the outer boundary. We can conclude in addition, noting that $\Phi(s, 1)$ approaches to the saddle loop as $s$ decreases, that the period function is increasing near the outer boundary.

Let us consider next a parameter $\mu^{\star}$ such that the outer boundary of $\mathcal{P}_{\mu^{\star}}$ is a bicycle (i.e., $\mu \in N$ ). As before, according to Remark 3.1, the bicycle persists and it is the outer boundary of $\mathcal{P}_{\mu}$ for $\mu \approx \mu^{\star}$. We shall take advantage of the symmetry of the Loud's systems with respect to $\{y=0\}$ to study only the passage through one of
the saddles. For instance, let us consider the saddle in $\{y>0\}$. Recall (see Section 3.1) that this saddle is located at

$$
p_{\mu}:=\left(1, \sqrt{-\frac{(D+1)}{F}}\right) \quad \text { with } \lambda_{1}(\mu)=\sqrt{-\frac{(D+1)}{F}} \quad \text { and } \quad \lambda_{2}(\mu)=2 F \sqrt{-\frac{(D+1)}{F}} \text {, }
$$

and so its ratio of hyperbolicity is $r(\mu)=-2 F$. We take now a local diffeomorphism $\Phi$ that conjugates $X_{\mu}$ for $\mu \approx \mu^{\star}$ with its normal form, which depends on $r\left(\mu^{\star}\right) \in \mathbb{Q}$ or $r\left(\mu^{\star}\right) \notin \mathbb{Q}$. As before, we shall take two transversal sections $\Sigma_{\sigma}$ and $\Sigma_{\tau}$ given by $s \longmapsto \Phi(s, 1)$ and $s \longmapsto \Phi(1, s)$, respectively. Denote by $P(s ; \mu)$ the period of the periodic orbit of $X_{\mu}$ passing through $\Phi(s, 1)$. We decompose it as

$$
P(s ; \mu)=2 T_{1}(s ; \mu)+2 T_{2}(s ; \mu)+2 T_{3}(R(s ; \mu) ; \mu)
$$

where $T_{1}$ is the time function for $-X_{\mu}$ from $\Sigma_{\sigma}$ to $\{y=0\}, R$ and $T_{2}$ are, respectively, the Dulac map and the time function for $X_{\mu}$ from $\Sigma_{\sigma}$ to $\Sigma_{\tau}$ and, finally, $T_{3}$ is the time function for $X_{\mu}$ from $\Sigma_{\tau}$ to $\{y=0\}$. It is clear that $T_{1}$ and $T_{3}$ are smooth functions on $s=0$. On the other hand, it is well known (see $[14,16]$ for instance) that

$$
R(s ; \mu)=s^{r(\mu)}\left(\rho(\mu)+\psi_{1}(s ; \mu)\right) \quad \text { where } \psi_{1} \in \mathcal{I}\left(U^{\star}\right)
$$

for some neighborhood $U^{\star}$ of $\mu^{\star}$. Concerning $T_{2}$, if $r\left(\mu^{\star}\right) \notin \mathbb{Q}$ then the normal form (34) is linear and one can easily verify that $T_{2}(s ; \mu)=-\frac{1}{\lambda_{1}(\mu)} \log s$. The expression is not so easy when $r\left(\mu^{\star}\right)=\frac{p}{q}$. The case $p=q=1$ is treated in [1] and, as we already mentioned, one obtains the expression given in (35). In the general case, following the same approach it can be shown (see [8, Proposition 23]) that

$$
T_{2}(s ; \mu)=-\frac{1}{\lambda_{1}(\mu)} \log s-\frac{1}{p} \beta_{1}(\mu) s^{p} \omega\left(s^{p} ; \frac{q}{p} r(\mu)\right)+\psi_{2}(s ; \mu)
$$

where $\psi_{2}$ is $\mathcal{C}^{1}$ at $(s, \mu)=\left(0, \mu^{\star}\right)$ and 1 -flat at $s=0$ for all $\mu$. Some computations show that

$$
\frac{d}{d s} T_{2}(s ; \mu)=-\frac{1}{\lambda_{1}(\mu)} \frac{1}{s}-\beta_{1}(\mu) s^{p-1}\left(\frac{q}{p} r(\mu) \omega\left(s^{p} ; \frac{q}{p} r(\mu)\right)+1\right)+\psi_{2}^{\prime}(s ; \mu)
$$

Therefore we can assert that in both cases, $r\left(\mu^{\star}\right)$ rational or irrational, it holds

$$
s \frac{d}{d s} T_{2}(s ; \mu) \longrightarrow-1 / \lambda_{1}\left(\mu^{\star}\right) \quad \text { as }(s, \mu) \longrightarrow\left(0, \mu^{\star}\right)
$$

Finally this implies that $s P_{s}(s ; \mu) \longrightarrow-2 / \lambda_{1}\left(\mu^{\star}\right)$ as $(s, \mu) \longrightarrow\left(0, \mu^{\star}\right)$ because $T_{1}$ and $T_{3}$ are smooth at $s=0$ and $\psi_{1} \in \mathcal{I}\left(U^{\star}\right)$. Exactly as in the saddle loop case, this proves
that $\mu^{\star}$ is a local regular value at the outer boundary and that the period function is monotonous increasing there.

Remark 3.16. From the proof of Theorem 3.15 it follows that if $\mu \in M \cup N$ then the period function of $X_{\mu}$ tends to $+\infty$ as we approach to the outer boundary of $\mathcal{P}_{\mu}$.

### 3.4. Proof of the main result

Proof of Theorem A. The fact that the parameters in $\mathbb{R}^{2} \backslash\left\{\Gamma_{\mathrm{B}} \cup \Gamma_{\mathrm{U}}\right\}$ are local regular values follows from the application of Theorems 3.3, 3.14 and 3.15 in the corresponding regions that cover. These theorems also show the assertions in (a) and (b) concerning the monotonicity near the outer boundary. Consider now a parameter $\mu_{0} \in \Gamma_{\mathrm{B}}$ and note that any neighborhood of $\mu_{0}$ intersects $\mathcal{D}_{\mathrm{B}} \backslash \Gamma_{\mathrm{U}}$ and $\mathcal{I}_{\mathrm{B}} \backslash \Gamma_{\mathrm{U}}$. Consequently, any neighborhood of $\mu_{0}$ contains two parameters $\mu^{+}$and $\mu^{-}$such that the respective period functions have different monotonicity in the outer boundary. (Here we use, recall Remark 2.5, that the character increasing or decreasing does not depend on the particular parametrization of the period function used.) This clearly implies that $\mu_{0}$ is a local bifurcation value at the outer boundary and so the result is proved.

## 4. Bifurcation in the interior

In this section we determine some local bifurcation values of the period function in the interior of the period annulus of the dehomogenized Loud's systems (2). To be more precise, we prove that there are three parameter values, namely

$$
\begin{align*}
& L_{1}=\left(-\frac{3}{2}, \frac{5}{2}\right), \quad L_{2}=\left(\frac{-11+\sqrt{105}}{20}, \frac{15-\sqrt{105}}{20}\right) \text { and } \\
& L_{3}=\left(\frac{-11-\sqrt{105}}{20}, \frac{15+\sqrt{105}}{20}\right) \tag{36}
\end{align*}
$$

such that at each $L_{i}$ there exists a germ of analytic curve corresponding to this type of bifurcation. We describe moreover the relative position of this curve with respect to other bifurcation curves. The result is based on the work of Chicone and Jacobs in [4].

Setting $\mu=(D, F)$ as usual, let us denote by $P(s ; \mu)$ the period of the periodic orbit of system (2) passing through the point $(s, 0)$. Note that $P(s ; \mu)$ is a well defined analytic function for $(s, \mu) \in(0, \varepsilon) \times \mathbb{R}^{2}$ because the center is nondegenerate. Moreover, since the eigenvalues of the linear part of $X_{\mu}$ at the origin are $\pm i$, it can be extended analytically to $s=0$ by setting $P(0 ; \mu):=2 \pi$. We can thus consider the Taylor expansion of $P(s ; \mu)$ at $s=0$,

$$
\begin{equation*}
P(s ; \mu)=2 \pi+P_{2}(\mu) s^{2}+P_{3}(\mu) s^{3}+P_{4}(\mu) s^{4}+P_{5}(\mu) s^{5}+P_{6}(\mu) s^{6}+\cdots \tag{37}
\end{equation*}
$$

The coefficients $P_{k}(\mu)$, which are real polynomials in the parameters of the system, are called the period constants of the center. For instance (see [4]),

$$
\begin{aligned}
P_{2}(D, F)= & \frac{\pi}{12}\left(10 D^{2}+10 D F-D+4 F^{2}-5 F+1\right) \\
P_{4}(D, F)= & \frac{\pi}{1152}\left(1540 D^{4}+4040 D^{3} F+1180 D^{3}+4692 D^{2} F^{2}+1992 D^{2} F\right. \\
& +453 D^{2}+2768 D F^{3}+228 D F^{2}+318 D F-2 D+784 F^{4} \\
& \left.-616 F^{3}-63 F^{2}-154 F+49\right)
\end{aligned}
$$

Chicone and Jacobs prove in [4] that the ideals generated by the period constants verify that

$$
\begin{equation*}
\left(P_{2}\right)=\left(P_{2}, P_{3}\right) \subsetneq\left(P_{2}, P_{4}\right)=\left(P_{2}, P_{4}, P_{5}\right) \subsetneq\left(P_{2}, P_{4}, P_{6}\right)=\left(P_{i}, i \in \mathbb{N}\right) \tag{38}
\end{equation*}
$$

They also show that the ideal $\left(P_{2}, P_{4}, P_{6}\right)$ determines the points

$$
S_{1}=(-1 / 2,1 / 2), \quad S_{2}=(0,1), \quad S_{3}=(0,1 / 4), \quad S_{4}=(-1 / 2,2)
$$

which correspond to the four nonlinear quadratic isochronous centers. The ideal $\left(P_{2}, P_{4}\right)$ determines, apart from the four isochronous centers, the three weak centers $L_{i}$ given in (36). Note in particular that these seven parameter values are over the conic $\Gamma_{\mathrm{C}}:=$ $\left\{\mu \in \mathbb{R}^{2}: P_{2}(\mu)=0\right\}$ (see Fig. 11).

In [4] the authors use a notion of bifurcation which differs from the one introduced in Section 2. Indeed, they say that $k$ critical periods bifurcate from the center corresponding to the parameter $\mu_{0}$ if for every $\varepsilon>0$ and every neighborhood $U$ of $\mu_{0}$ there is a point $\mu_{1} \in U$ such that the equation $P^{\prime}\left(s ; \mu_{1}\right)=0$ has $k$ solutions in the interval $(0, \varepsilon)$. With this definition and the notation introduced above we can now summarize their result concerning the dehomogenized Loud's systems:

Theorem 4.1 (Chicone-Jacobs). The maximal number of critical periods bifurcating from the center at the origin of the dehomogenized Loud's family is two. In addition,
(a) If $\mu \notin \Gamma_{\mathrm{C}}$ then no critical period bifurcates from the center.
(b) If $\mu \in \Gamma_{\mathrm{C}} \backslash\left\{L_{1}, L_{2}, L_{3}\right\}$ then at most one critical period bifurcates from the center and there are perturbations with exactly one critical period.
(c) If $\mu \in\left\{L_{1}, L_{2}, L_{3}\right\}$ then at most two critical periods bifurcate from the center and there are perturbations with exactly one and exactly two critical periods.

Remark 4.2. It is clear, on account of (37), that if $\mu \notin \Gamma_{\mathrm{C}}$ then the monotonicity of $P(s ; \mu)$ for $s>0$ small enough is given by the sign of $P_{2}(\mu)$. Note that $\Gamma_{\mathrm{C}}$ is a Jordan curve. We denote the bounded and unbounded component of $\mathbb{R}^{2} \backslash \Gamma_{\mathrm{C}}$ by $\mathcal{D}_{\mathrm{C}}$ and $\mathcal{I}_{\mathrm{C}}$, respectively. One can check then that $P_{2}(\mu)$ is positive for $\mu \in \mathcal{I}_{\mathrm{C}}$ and negative


Fig. 11. The ellipse $\Gamma_{\mathrm{C}}=\left\{\mu \in \mathbb{R}^{2}: P_{2}(\mu)=0\right\}$.
for $\mu \in \mathcal{D}_{\mathrm{C}}$. Thus, if $\mu$ belongs to $\mathcal{I}_{\mathrm{C}}$ (respectively, $\mathcal{D}_{\mathrm{C}}$ ) then the period function of $X_{\mu}$ is monotonous increasing (respectively, decreasing) at the inner boundary (i.e., the center). Therefore, according to (a) in Remark 2.6, the parameters in $\Gamma_{\mathrm{C}}$ are local bifurcation values of the period function at the inner boundary. On the other hand, since the expansion in (37) is uniform with respect to $\mu$, from (b) in Remark 2.6 we can assert that any parameter $\mu \notin \Gamma_{\mathrm{C}}$ is a local regular value of the period function at the inner boundary. In short, the regular and bifurcation values at the center are the same with the definition in [4] and Definition 2.4.

One can easily verify that $\Gamma_{\mathrm{C}}=\left\{\mu \in \mathbb{R}^{2}: P_{2}(\mu)=0\right\}$ and $\left\{\mu \in \mathbb{R}^{2}: P_{4}(\mu)=0\right\}$ are analytic curves that intersect transversally at each $L_{i}$. We shall prove next the following result.

Theorem 4.3. For each $i=1,2,3$ there exist a neighborhood $U_{i}$ of $L_{i}$ and an analytic curve $\delta_{i}$ which is tangent to $\Gamma_{\mathrm{C}}$ at $L_{i}$ such that the arc $\delta_{i} \cap\left\{\mu \in U_{i}: P_{4}(\mu)<0\right\}$ corresponds to local bifurcation values of the period function in the interior. Moreover this arc is inside $\left\{\mu \in \mathbb{R}^{2}: P_{2}(\mu)>0\right\}$.

Proof. To study the period function near the center it is more convenient to parametrize the periodic orbits by means of the first integral $H_{\mu}(x, y)-H_{\mu}(0,0)$. (Recall that $H_{\mu}$ is given in Section 3.1.) This will eliminate the rather artificial property that only the even coefficients in (37) are significant, which is due to the fact that each periodic
orbit intersects twice the $x$-axis. So for each $h>0$ denote by $\widehat{P}(h ; \mu)$ the period of the periodic orbit of $X_{\mu}$ inside the energy level $\left\{H_{\mu}(x, y)-H_{\mu}(0,0)=h\right\}$. Thus, since one can verify that

$$
h=H_{\mu}(s, 0)-H(0,0)=\frac{s^{2}}{2}+(1+D+2 F) \frac{s^{3}}{3}+(1+D+F)(1+2 F) \frac{s^{4}}{4}+\cdots,
$$

from (37) it follows that $\widehat{P}(h ; \mu)=2 \pi+Q_{0}(\mu) h+Q_{1}(\mu) \frac{h^{2}}{2}+Q_{2}(\mu) \frac{h^{3}}{3}+\cdots$ with

$$
\begin{align*}
& Q_{0}(\mu)=2 P_{2}(\mu) \\
& Q_{1}(\mu)=8 P_{4}(\mu)-4(1+D+F)(1+2 F) P_{2}(\mu) \\
& Q_{2}(\mu)=24 P_{6}(\mu) \quad \bmod \left(P_{2}(\mu), P_{4}(\mu)\right) \tag{39}
\end{align*}
$$

It is clear moreover that the critical periods coincide with the positive zeros of the function

$$
\begin{equation*}
Z(h ; \mu):=\widehat{P}_{h}(h ; \mu)=Q_{0}(\mu)+Q_{1}(\mu) h+Q_{2}(\mu) h^{2}+\cdots \tag{40}
\end{equation*}
$$

Since $P_{2}\left(L_{i}\right)=P_{4}\left(L_{i}\right)=0$, from (39) and using also the expression for $P_{6}$ given in [4], it follows that $Q_{2}\left(L_{i}\right)=24 P_{6}\left(L_{i}\right)>0$ for each $i=1,2,3$. One can also verify, taking (39) into account again, that the gradients $\nabla Q_{0}\left(L_{i}\right)$ and $\nabla Q_{1}\left(L_{i}\right)$ are linearly independent for each $i=1,2,3$. From now on let us only consider $L_{1}$ for the sake of simplicity in the exposition.

By applying the Weierstrass Preparation Theorem, there exist a neighborhood $U$ of $L_{1}$, two analytic functions $a_{0}$ and $a_{1}$ with $a_{i}: U \longrightarrow \mathbb{R}$ and a positive analytic function $K:(-\varepsilon, \varepsilon) \times U \longrightarrow \mathbb{R}$ such that

$$
\begin{equation*}
Z(h ; \mu)=K(h ; \mu)\left(h^{2}+a_{1}(\mu) h+a_{0}(\mu)\right) . \tag{41}
\end{equation*}
$$

Accordingly, if $K(h ; \mu)=k_{0}(\mu)+k_{1}(\mu) h+\mathrm{o}(h)$ then $k_{0}\left(L_{1}\right)>0$ and, from (40),

$$
\begin{equation*}
a_{0}(\mu) k_{0}(\mu)=Q_{0}(\mu) \quad \text { and } \quad a_{0}(\mu) k_{1}(\mu)+a_{1}(\mu) k_{0}(\mu)=Q_{1}(\mu) \tag{42}
\end{equation*}
$$

So it turns out that $a_{0}\left(L_{1}\right)=a_{1}\left(L_{1}\right)=0$. Therefore

$$
\nabla Q_{0}\left(L_{1}\right)=k_{0}\left(L_{1}\right) \nabla a_{0}\left(L_{1}\right) \quad \text { and } \quad \nabla Q_{1}\left(L_{1}\right)=k_{1}\left(L_{1}\right) \nabla a_{0}\left(L_{1}\right)+k_{0}\left(L_{1}\right) \nabla a_{1}\left(L_{1}\right)
$$

and we can thus assert that the gradients $\nabla a_{0}\left(L_{1}\right)$ and $\nabla a_{1}\left(L_{1}\right)$ are linearly independent. Consequently $\psi(\mu):=\left(a_{0}(\mu), a_{1}(\mu)\right)$ is a local diffeomorphism between $U$ and some neighborhood $V$ of $(0,0)$. We define $\delta_{1}$ as the preimage by $\psi$ of the analytic curve $\left\{\left(a_{0}, a_{1}\right) \in V: a_{1}^{2}-4 a_{0}=0\right\}$. Note in particular that $L_{1}=\psi^{-1}(0,0)$ belongs to $\delta_{1}$.

On the other hand, since $\Gamma_{\mathrm{C}} \cap U=\left\{\mu \in U: Q_{0}(\mu)=0\right\}$ is the preimage by $\psi$ of $\left\{\left(a_{0}, a_{1}\right) \in V: a_{0}=0\right\}$, we conclude that $\delta_{1}$ is tangent to $\Gamma_{\mathrm{C}}$ at $L_{1}$.

Now we shall prove, recall Definition 2.4, that any parameter in $\delta_{1} \cap\{\mu \in U$ : $\left.P_{4}(\mu)<0\right\}$ is a local bifurcation value of the period function in the interior. To this end fix some $\mu^{\star} \in \delta_{1} \cap U$ with $P_{4}\left(\mu^{\star}\right)<0$ and define $h^{\star}:=-\frac{a_{1}\left(\mu^{\star}\right)}{2}$. Observe that if $\mu \in \delta_{1}$ then

$$
\begin{aligned}
8 P_{4}(\mu) & =Q_{1}(\mu)+2(1+D+F)(1+2 F) Q_{0}(\mu) \\
& =a_{1}(\mu) k_{0}(\mu)+a_{0}(\mu)\left(k_{1}(\mu)+2(1+D+F)(1+2 F) k_{0}(\mu)\right) \\
& =a_{1}(\mu)\left\{k_{0}(\mu)+\frac{1}{4} a_{1}(\mu)\left(k_{1}(\mu)+2(1+D+F)(1+2 F) k_{0}(\mu)\right)\right\} .
\end{aligned}
$$

Here we use (39) in the first equality, (42) in the second one and that $\mu \in \delta_{1}$ in the third one. Therefore, since $k_{0}\left(L_{1}\right)>0$ and $a_{1}\left(L_{1}\right)=0$, the above equality shows that $P_{4}(\mu)$ and $a_{1}(\mu)$ have the same sign for $\mu \approx L_{1}$. Thus, shrinking the neighborhood $U$ of $L_{1}$ if necessary, we have that $h^{\star}>0$. Consider now any neighborhood $U^{\star}$ of $\mu^{\star}$. Due to $\mu^{\star} \in \delta_{1}$, there exist $\bar{\mu} \in U^{\star}$ such that $a_{1}^{2}(\bar{\mu})-4 a_{0}(\bar{\mu})>0$. Hence (41) implies that $Z\left(-a_{1}(\bar{\mu}) / 2 ; \bar{\mu}\right)<0$. Note also that $Z\left(h ; \mu^{\star}\right)>0$ for $h \neq h^{\star}$. Therefore, since $-\frac{a_{1}(\bar{\mu})}{2} \longrightarrow h^{\star}$ as $\bar{\mu} \longrightarrow \mu^{\star}$, it turns out that relation (3) cannot be verified in any neighborhood of $\mu^{\star}$. So $\mu^{\star}$ is a local bifurcation value in the interior because Definition 2.4 is not fulfilled for $c=h^{\star}$. Finally, from (39) and (42) and taking $\mu^{\star} \in \delta_{1}$ into account,

$$
2 P_{2}\left(\mu^{\star}\right)=a_{0}\left(\mu^{\star}\right) k_{0}\left(\mu^{\star}\right)=a_{1}\left(\mu^{\star}\right)^{2} k_{0}\left(\mu^{\star}\right) / 4
$$

and this implies that $P_{2}\left(\mu^{\star}\right)>0$. The proof is completed.
Remark 4.4. In fact the preceding proof provides a stronger result than Theorem 4.3. Namely, that for each $i=1,2,3$ there exists a neighborhood $U_{i}$ of $L_{i}$ and a local analytic equivalence $\psi_{i}: U_{i} \longrightarrow V_{i}$ between the local bifurcation diagrams of the families $\widehat{P}_{h}(h ; \mu)$ with $\mu \in U_{i}$ and $\mathcal{C}(h ; a):=h^{2}+a_{1} h+a_{0}$ with $a \in V_{i}$ for $h \in(0, \varepsilon)$.

On the other hand it is also possible to obtain the asymptotic expansion of the curves $\delta_{i}$. To do so it is enough to find the expansion of the functions $a_{0}(\mu)$ and $a_{1}(\mu)$ in terms of the coefficients $Q_{k}(\mu)$ and substitute them into the equation $a_{1}^{2}-4 a_{0}=0$. In this way one can obtain for instance the second-order expansion of the curve $\delta_{1}$ at the point $L_{1}=(-3 / 2,5 / 2)$,

$$
\begin{aligned}
& \frac{189}{10}\left(D+\frac{3}{2}\right)-\frac{63}{2}\left(D+\frac{3}{2}\right)^{2}-\frac{63}{2}\left(F-\frac{5}{2}\right)\left(D+\frac{3}{2}\right)+\frac{17}{5}\left(F-\frac{5}{2}\right)^{2} \\
& \quad+\mathrm{O}\left(\left\|\left(D+\frac{3}{2}, F-\frac{5}{2}\right)\right\|^{3}\right)=0
\end{aligned}
$$

## 5. Existence of critical periods

In this section we determine two subsets of the parameter space such that the corresponding period function has at least one critical period and at least two critical periods, respectively. We shall use the notation $\mathcal{D}_{\mathrm{B}}$ and $\mathcal{I}_{\mathrm{B}}$, introduced in Section 1, for the bounded and unbounded components of $\mathbb{R}^{2} \backslash \Gamma_{\mathrm{B}}$, and the notation $\mathcal{D}_{\mathrm{C}}$ and $\mathcal{I}_{\mathrm{C}}$, introduced in Remark 4.2, for the bounded and unbounded components of $\mathbb{R}^{2} \backslash \Gamma_{\mathrm{C}}$. Let us note that $\mathcal{D}$ and $\mathcal{I}$ stand for decreasing and increasing, respectively.

Theorem 5.1. Consider a parameter $\mu_{0}$ inside $\mathcal{I}_{\mathrm{C}} \cap \mathcal{D}_{\mathrm{B}}$ or $\mathcal{D}_{\mathrm{C}} \cap \mathcal{I}_{\mathrm{B}}$. If $\mu_{0} \notin \Gamma_{\mathrm{U}}$ then the period function of $X_{\mu_{0}}$ has at least one critical period.

Proof. Let us prove for instance the assertion concerning $\mathcal{I}_{\mathrm{C}} \cap \mathcal{D}_{\mathrm{B}}$ (the other one follows exactly the same way). So consider some $\mu_{0} \in \mathcal{I}_{\mathrm{C}} \cap \mathcal{D}_{\mathrm{B}} \backslash \Gamma_{\mathrm{U}}$ and let $P_{\mu_{0}}:(0,1) \longrightarrow \mathbb{R}$ be a parametrization of the period function of $X_{\mu_{0}}$. Then, on account of Remark 4.2, we have that $P_{\mu_{0}}^{\prime}$ is positive near $s=0$ because $\mu_{0} \in \mathcal{I}_{\mathrm{C}}$. On the other hand, using that $\mu_{0} \in \mathcal{D}_{\mathrm{B}}$, by applying Theorem A it follows that $P_{\mu_{0}}^{\prime}$ is negative near $s=1$. Therefore, by Bolzano's Theorem, we can assert that there exists $s_{0} \in(0,1)$ such that $P_{\mu_{0}}^{\prime}\left(s_{0}\right)=0$.

Consider now the subsets $U, W, M$ and $N$ introduced in Section 3 and let $\left\{P_{\mu}\right.$ : $\left.(0,1) \longrightarrow \mathbb{R}, \mu \in \mathbb{R}^{2}\right\}$ be any parametrization of the period function. It follows then that $L(\mu):=\lim _{s \rightarrow 1} P_{\mu}(s)$ is a well defined function on $U \cup W \cup M \cup N \backslash \Gamma_{4}$, where $\Gamma_{4}$ is the segment represented in Fig. 9. Indeed, Theorem 3.6 shows that

$$
L(\mu)=\frac{2 \sqrt{2}}{\sqrt{a+b+c}} \operatorname{arctanh}\left(\frac{2 a+b-\sqrt{b^{2}-4 a c}}{2 \sqrt{a(a+b+c)}}\right) \quad \text { if } \mu \in U
$$

Next, Proposition 5.2 in [13] shows that $L(\mu)=\frac{\pi}{\sqrt{F(D+1)}}$ if $\mu \in W \backslash \Gamma_{4}$ and, finally, Remark 3.16 shows that $L(\mu)=+\infty$ if $\mu \in M \cup N$. Some easy computations, that are not included here for the sake of brevity, show that the set

$$
\left\{\mu \in U \cup W \backslash \Gamma_{4}: L(\mu)-2 \pi=0\right\}
$$

together with the points $(-3 / 4,1)$ and $(-1 / 2,1 / 2)$ and the segment $\{0\} \times[1 / 4,1]$ form a Jordan curve, say $\Gamma_{0}$ (see Fig. 12). We can thus consider the bounded and unbounded components of $\mathbb{R}^{2} \backslash \Gamma_{0}$, which we denote by $\mathcal{J}_{-}$and $\mathcal{J}_{+}$, respectively. The subscripts are chosen in this way because one can verify that $L(\mu)-2 \pi$ is negative for $\mu \in \mathcal{J}_{-}$ and positive for $\mu \in \mathcal{J}_{+}$. With this notation we obtain the following result:

Theorem 5.2. Consider a parameter $\mu_{0}$ inside $\mathcal{I}_{\mathrm{C}} \cap \mathcal{I}_{\mathrm{B}} \cap \mathcal{J}_{-}$or $\mathcal{D}_{\mathrm{C}} \cap \mathcal{D}_{\mathrm{B}} \cap \mathcal{J}_{+}$. If $\mu_{0} \notin \Gamma_{\mathrm{U}}$ then the period function of $X_{\mu_{0}}$ has at least two critical periods.


Fig. 12. Numerical drawing of the regions in Theorem 5.1 (left) and Theorem 5.2 (right).

Proof. Let us prove for instance the assertion concerning $\mathcal{I}_{\mathrm{C}} \cap \mathcal{I}_{\mathrm{B}} \cap \mathcal{J}_{-}$(the other one follows exactly the same way). Fix $\mu_{0} \notin \Gamma_{\mathrm{U}}$ and consider a parametrization $P_{\mu_{0}}:(0,1) \longrightarrow \mathbb{R}$ of the period function of $X_{\mu_{0}}$. Then, using that the eigenvalues of the linear part of $X_{\mu_{0}}$ are $\pm i$, it follows that $\lim _{s \rightarrow 0} P_{\mu_{0}}(s)=2 \pi$. We have in addition, recall Remark 4.2, that $P_{\mu_{0}}$ is increasing near $s=0$ because $\mu_{0} \in \mathcal{I}_{\mathrm{C}}$. Then, since $\lim _{s \rightarrow 1} P_{\mu_{0}}(s)<2 \pi$ due to $\mu_{0} \in \mathcal{J}_{-}$, there exists $\bar{s} \in(0,1)$ such that $P_{\mu_{0}}(\bar{s})=2 \pi$ and $P_{\mu_{0}}^{\prime}(\bar{s}) \leqslant 0$. Thus, since $P_{\mu_{0}}^{\prime}$ is positive near $s=0$ and, on account of $\mu_{0} \in \mathcal{I}_{\mathrm{B}}$, also near $s=1$ by Theorem A, we conclude that there exist $s_{1} \in(0, \bar{s})$ and $s_{2} \in[\bar{s}, 1)$ such that $P_{\mu_{0}}^{\prime}\left(s_{i}\right)=0$.

Unfortunately it is very difficult to provide an explicit and simple analytical description of the regions in Theorems 5.1 and 5.2. We prefer instead to make a numerical drawing of them using the exact analytic expressions that define the curves $\Gamma_{\mathrm{B}}, \Gamma_{\mathrm{C}}$, and $\Gamma_{0}$. The picture on the left in Fig. 12 shows several regions, which can be clearly observed, with at least one critical period. The picture on the right shows two regions, one of them very tiny, with at least two critical periods. Both regions corresponds to the set $\mathcal{I}_{\mathrm{C}} \cap \mathcal{I}_{\mathrm{B}} \cap \mathcal{J}_{-}$, the other one seems to be empty. We refer here to the result of Chicone and Dumortier in [3]. They proved that there exists some $D^{\star} \approx-1.47$ such if $\mu \in\left(D^{\star},-1.4\right) \times\{2\}$ then the period function of $X_{\mu}$ has at least one critical point. Observe in Fig. 12 that this is the horizontal segment in the boundary of the biggest component of $\mathcal{I}_{\mathrm{C}} \cap \mathcal{I}_{\mathrm{B}} \cap \mathcal{J}_{-}$. Their results follows from the fact that this segment is inside $\mathcal{I}_{\mathrm{C}} \cap \mathcal{J}_{-}$.

## 6. Conjectures and open problems

In this section we explain how we come to the conjecture posed in the introduction and we comment on various steps that should be done in order to prove it. Some of them seem feasible, while others seem out of reach for the moment.

### 6.1. Geometrical picture

Consider the parametrization of the set of periodic orbits in the period annulus that provides the first integral $H_{\mu}$ given in (4). Let us assume that we normalize it in order that $h=0$ corresponds to the center (i.e., the inner boundary) and $h=1$ to the polycycle (i.e., the outer boundary). Then the parametrization of the period function $P(h ; \mu)$ that we obtain is defined on $(0,1)$ for all $\mu \in \mathbb{R}^{2}$. Recall (see Section 4) that it can be extended analytically to $h=0$ by setting $P(0 ; \mu):=2 \pi$ and that

$$
P^{\prime}(h ; \mu)=Q_{0}(\mu)+Q_{1}(\mu) h+Q_{2}(\mu) h^{2}+\cdots \quad \text { for } h \approx 0
$$

Let $M \subset[0,1] \times \mathbb{R}^{2}$ be the set of points $(h, \mu)$ verifying $P^{\prime}(h ; \mu)=0$ for $h \in$ $[0,1)$ and extended to $h=1$ by continuity using an asymptotic development at the polycycle. We conjecture that $M$ is a smooth surface for $h \neq 1$ and that it is fibered by simple closed curves $M_{h}$ given by $\left\{\mu \in \mathbb{R}^{2}: P^{\prime}(h ; \mu)=0\right\}$ for each fixed $h \in[0,1)$. This last condition is equivalent to require that $\frac{\partial P^{\prime}(h ; \mu)}{\partial D}$ and $\frac{\partial P^{\prime}(h ; \mu)}{\partial F}$ do not vanish simultaneously. Thus $M_{0}$ is the ellipse $\Gamma_{\mathrm{C}}$ in Fig. 11 and $M_{1}$ should correspond to the local bifurcation values at the outer boundary, that we conjecture to be the curve $\Gamma_{\mathrm{B}}$ in Theorem A together with the segment $\{0\} \times[0,1 / 2]$. In fact the strange kidney-shaped curve that appears in the numerical bifurcation diagram of Chicone and Jacobs (see Fig. 1) would correspond approximately to some curve $M_{h}$ with $h \approx 1$. The curves that correspond to local bifurcation values of the period function in the interior are obtained by the projection of $M$ on the $\mu$-plane. More precisely, these curves would be given as envelopes of the family $\left\{M_{h}, h \in[0,1]\right\}$. We obtained our conjectural bifurcation diagram by trying to interpolate a continuous family of curves $M_{h}$ starting at $\Gamma_{\mathrm{C}}$ for $h=0$ and ending for $h=1$ at our conjectural bifurcation diagram at the polycycle. Fig. 13 shows two intermediate curves $M_{h}$ and $M_{h^{\prime}}$ with $0<h<h^{\prime}<1$. Note that every curve $M_{h}$ must pass through $S_{1}, S_{2}, S_{3}$ and $S_{4}$, the parameters corresponding to the four isochronous centers of the family. On the other hand, according to [4], $Q_{0}(\mu)$, $Q_{1}(\mu)$ and $Q_{2}(\mu)$ generate the ideal of the coefficients of $P^{\prime}(h ; \mu)$ at $h=0$. This gives that for $\mu \approx S_{i}$ and $h \approx 0$ the family of curves $M_{h}$ is approximately of the form of the pencil $Q_{0}(\mu)+Q_{1}(\mu) h=0$. Since the curves $Q_{0}(\mu)=0$ and $Q_{1}(\mu)=0$ are transverse at $S_{2}, S_{3}$ and $S_{4}$, it follows that, in a neighborhood of these three parameters and at least for $h$ small, the curves $M_{h}$ look like a pencil of straight lines passing through $S_{i}$. At the other isochronous center $S_{1}$, since $Q_{0}(\mu)=0$ and $Q_{1}(\mu)=0$ have quadratic contact, the curves $M_{h}$ look like a pencil of parabolas tangent to $\Gamma_{\mathrm{C}}$. Consider finally the curves $\delta_{i}$ passing through the three weak centers $L_{i}$ that we obtain in Theorem 4.3. From Remark 4.4 it follows that, in a neighborhood of each $L_{i}$, the curves $M_{h}$ are


Fig. 13. The intermediate curves $M_{h}$.
tangent to $\delta_{i}$ for $h \approx 0$ and that these curves correspond to parameters in which two critical periods collapse disappearing in the interior. In other words, near each $L_{i}$ and for $h \approx 0$ the curves $\delta_{i}$ are double bifurcation curves in the interior. We conjecture that this behavior holds for the entire curve $\delta_{i}$ and that they separate the regions with two critical periods from the region in which the period function is globally monotonous increasing. Since these bifurcation curves $\delta_{i}$ begin at the weak centers $L_{i} \in \Gamma_{\mathrm{C}}$, which correspond to double bifurcations at the inner boundary, we presume that they end at three special points in the curve $\Gamma_{\mathrm{B}} \cup(\{0\} \times[0,1 / 2])$, which would play the role of double bifurcation parameters at the outer boundary (see Fig. 3).

### 6.2. Bifurcation at the outer boundary

Our main result determines two sets $\Gamma_{\mathrm{B}}$ and $\Gamma_{\mathrm{U}}$ such that the parameters in $\Gamma_{\mathrm{B}}$ are local bifurcation values of the period function at the outer boundary and the ones in $\mathbb{R}^{2} \backslash\left(\Gamma_{\mathrm{B}} \cup \Gamma_{\mathrm{U}}\right)$ are local regular values. The character of the parameters in $\Gamma_{\mathrm{U}}$ remains unspecified in our work. The first natural problem that raises is to determine the character of the parameters in $\Gamma_{\mathrm{U}}$. As we already mention, we conjecture that they are all regular values except for the segment $\{0\} \times[0,1 / 2]$, whose conjectural bifurcation is described below.

The set $\Gamma_{\mathrm{U}}$ is stratified as a union of open segments and a few points at the intersection of them. Probably one has to treat first the open segments and next the points (higher codimension strata). Our main tool, Proposition 3.9, does not apply along some curves of $\Gamma_{\mathrm{U}}$ because the singular points at infinity of the polycycle are saddle-nodes or resonant nonlinearizable saddles. It seems reasonable to think that an analogue of Proposition 3.9 can be developed in all these cases and that one could determine the behavior of the period function at the polycycle in a neighborhood of these curves. Next, a specific study should be done for each of the codimension two points.

Our study of the bifurcation diagram of the period function does not deal with higherorder bifurcation parameters. The so called weak centers of order one and two that appear in the study of the period function near the center (see [4]) have a counterpart near the polycycle. The determination of these parameters and their study requires the knowledge of at least one more coefficient in the asymptotic expansion of the period function near the polycycle. There is no theoretical obstacle in doing so, but the technicalities seem prohibitive. Once these parameters are determined, using the derivation-division process as in the study of Chebyshev systems (see [12] for instance), one should be able to prove the equivalence of the local bifurcation at the outer boundary with some polynomial model.

Let us say a few words about the study at $D=0$. In this case the polycycle in the boundary of the period annulus has a degenerate singularity at infinity. Blowingup this singularity and applying a still unpublished generalization of Proposition 3.9, we hope to obtain the beginning of the asymptotic expansion of the period function $P(h ; \mu)$ at the outer boundary. It seems feasible to prove in this way that the segment $\{0\} \times(0,1 / 2)$ consists of local bifurcation values at the outer boundary. On the other hand, we have a polynomial family of functions that we think is a good model for the family $P^{\prime}(h ; D, F)$ near $S_{3}=(0,1 / 4)$. To be more precise, we conjecture that the families

$$
\begin{aligned}
Z(h ; D, F)= & P^{\prime}(h ; D, F), \quad h \in[0,1),(D, F) \in \mathbb{R}^{2} \\
\widehat{Z}(h ; \widehat{D}, \widehat{F})= & (\widehat{F}+\widehat{D}-\widehat{F} \widehat{D}) h^{2}+\left(\widehat{D} \widehat{F}-2 \widehat{F}-\widehat{D}-\widehat{D}^{2}\right) h+\widehat{F} \\
& h \in[0,1],(\widehat{D}, \widehat{F}) \in \mathbb{R}^{2}
\end{aligned}
$$

have locally equivalent bifurcation diagrams near the points $(0 ; 0,1 / 4)$ and $(0 ; 0,0)$ respectively. By this equivalence the curve $\Gamma_{\mathrm{C}}$ would correspond to $\{\widehat{F}=0\},\{D=0\}$ to $\{\widehat{D}=0\}$ and the curve $\delta_{2}$ to the curve defined by $\widehat{F}^{2}+\widehat{D}^{2}-2 \widehat{F} \widehat{D}+2 \widehat{D}+2 \widehat{F}+1=0$. (Note that $\widehat{D}^{2}\left(\widehat{F}^{2}+\widehat{D}^{2}-2 \widehat{F} \widehat{D}+2 \widehat{D}+2 \widehat{F}+1\right)$ is the discriminant of $\widehat{Z}(h ; \widehat{D}, \widehat{F})$ with respect to $h$.) Furthermore the intermediate curves $M_{h}=\{Z(h ; D, F)=0\}$ would correspond to the hyperbolic branches $\widehat{M}_{h}=\{\widehat{Z}(h ; \widehat{D}, \widehat{F})=0\}$. In Fig. 14 we show the bifurcation diagram of this polynomial model, which consists of the straight lines $\widehat{D}=0$ and $\widehat{F}=0$ together with the curve that joins the points $(-1,0)$ and $(0,-1)$, and we draw the curves $\widehat{M}_{h}$ for $h=i / 10, i=1, \ldots, 9$.


Fig. 14. Conjectured polynomial model near $D=0$.

### 6.3. Bifurcation in the interior

The study of the local bifurcation values of the period function in the interior is equivalent to the study of the behavior of the zeros of $P^{\prime}(h ; \mu)$ for $h \in(0,1)$. Here we assume again that $h$ is, up to a normalization, the energy of the first integral $H_{\mu}$ given in (4). In order to study these zeros one hopes to apply methods which proved successful for the abelian integrals. However, in this case the situation is more complicated for two reasons. The first one is that the first integral is not rational but only of Darboux-type. The methods used for abelian integrals have not yet been successfully adapted to this situation despite several efforts [9,21]. The second reason is that in the usual setting of the abelian integrals, the parameters enter linearly in the study as linear coefficients of the form that one integrates. In our situation the dependence on the parameters is highly nonlinear.

Due to the form of the first integral it can be verified that the complex fibers $\left\{H_{\mu}(x, y)=h\right\}$ given by a fixed $(h, \mu)$ are generically of infinite genus. This suggests complicated study unless, for some reason (the symmetry of the system for instance), one can project to some smaller space. The study of asymptotic cycles probably plays some role too.

Let us finally refer a method developed in [6]. The authors obtain a general formula for the derivative of the period function that can be applied to determine the critical periods that persists after the perturbation of an isochronous center. This formula can be viewed as an analogous of the first Melnikov function used to study limit cycles. As
an example of application they study the isochronous center $S_{2}=(0,1)$. Proposition 5 in [6] shows that for each closed interval $I$ inside $(0,1)$ there exists a neighborhood $U$ of $S_{2}$ such that if $\mu_{0} \in U$ then $P\left(h ; \mu_{0}\right)$ has at most one critical period in $I$ and that this critical period exists only in the case $\frac{1-F_{0}}{D_{0}}>3$. (Observe in Fig. 12 that this region is precisely the "linear approximation" at $S_{2}$ of the region in Theorem 5.1.) This implies, since one can also verify that the critical period is simple, that in a punctured neighborhood of $S_{2}$ there are no local bifurcation values of the period function in the interior.

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