# ON THE INTEGRABILITY OF POLYNOMIAL VECTOR FIELDS IN THE PLANE BY MEANS OF PICARD-VESSIOT THEORY 

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#### Abstract

We study the integrability of polynomial vector fields using Galois theory of linear differential equations when the associated foliations is reduced to a Riccati type foliation. In particular we obtain integrability results for some families of quadratic vector fields, Liénard equations and equations related with special functions such as Hypergeometric and Heun ones. The Poincaré problem for some families is also approached.


## Introduction

Given a polynomial differential system in $\mathbb{C}^{2}$,

$$
\begin{equation*}
\frac{d x}{d t}=\dot{x}=P(x, y), \quad \frac{d y}{d t}=\dot{y}=Q(x, y) \tag{1}
\end{equation*}
$$

with $P, Q \in \mathbb{C}[x, y]$, we consider its associated differential vector field

$$
\begin{equation*}
X=P(x, y) \frac{\partial}{\partial x}+Q(x, y) \frac{\partial}{\partial y}, \tag{2}
\end{equation*}
$$

whose integral curves are intimately related to the solutions of system (1). These solutions, taken as curves on the plane and leaving for a while its time-dependence, constitute its socalled foliation and satisfy the first order differential equation

$$
\begin{equation*}
y^{\prime}=\frac{d y}{d x}=\frac{Q(x, y)}{P(x, y)} . \tag{3}
\end{equation*}
$$

This expression (3) is often written as a Pfaff equation

$$
\begin{equation*}
\Omega=0, \tag{4}
\end{equation*}
$$

where $\Omega=Q(x, y) d x-P(x, y) d y$ is the corresponding differential 1-form. The connection between integral curves of the vector field $X$ and solutions of $\Omega=0$ is clear:

- geometrically, it is given by $\Omega \cdot X=0$, which means that the vector field $X$ is tangent to the leaves of the foliation (the orbits) defined by (4);

[^0]- dynamically, the general solution of equation (4), $H(x, y)=c t t$, is given by a first integral $H$ of the original vector field $X$, that is, a non-constant scalar function which remains constant along any of its solutions $(x(t), y(t))$. Since $\Omega \cdot X=0$, this is equivalent to say that $X(H)=0$ and to the existence of a suitable scalar function $f$ such that $\Omega=f d H$. In fact, $1 / f$ is an integrating factor of the field $X$, or of the 1 -form $\Omega$. Later we will say something more about the geometrical meaning of the integrating factors (Remark 1.3).
From a geometrical point of view, we focus our attention on invariant algebraic curves, that is, polynomial integral curves of the vector field $X$. This is the natural framework where Darboux Theory can be applied. Dynamically speaking, we restrict ourselves to first integrals $H(x, y)$ which are Liouvillian, i.e., written as a combination of algebraic functions, quadratures and exponential of quadratures in $\mathbb{C}(x, y)$, the field of rational functions on $x, y$. As we will see later, Galois Theory provides very useful and powerful tools to approach it.

Two classical problems remain still open for complex polynomial fields:
(i) Concerning the existence of invariant algebraic curves of system (1) or, equivalently, of algebraic solutions of the foliation equation (3).
(ii) About the existence of Liouvillian first integrals for systems (1) (or, in other words, to determine when the general solution of equation (3) is Liouvillian).
For general polynomial vector fields, problems (i) and (ii) are very difficult and we are still far from obtaining an effective method to decide whether a given arbitrary polynomial field has or not an invariant curve or admits a Liouvillian first integral. In fact, problem $(i)$ is connected with the (also unsolved) classical Poincaré problem, which seeks for a bound of the degree of the invariant algebraic curves as a function of the degree of the vector field (or of the associated foliation defined by (3)). It is known that Darboux Theory and adjacent results, as the ones due to Prelle-Singer and Singer [48, 50], provide connections between this two problems.

In this work we are concerned with the study of the (Darboux, Galois) integrability of some families of equations of type (4) inside the complex analytic category, that is, when the original vector field $X$ defining $\Omega$ is complex polynomial or can be reduced to it. Precisely, we restrict ourselves to those systems which can be reduced to a Riccati type equation

$$
\begin{equation*}
v^{\prime}=a_{0}(x)+a_{1}(x) v+a_{3}(x) v^{2}, \tag{5}
\end{equation*}
$$

$a_{1}, a_{2}$ and $a_{3}$ being rational functions with complex coefficients. For Riccati equations there is a very nice theory of integrability in the context of the Galois theory of its associated second order linear differential equation. This is, in our opinion, a natural framework where several known results concerning integrability of Riccati equations (5) should be considered. Kovacic provided in 1986 (see [30]) an effective algorithm which allows to decide whether an equation(5) has got an algebraic solution or not. And, additionally, a theorem of Liouville (see [31]) proved that the existence of an algebraic solution is exactly the definition of the integrability for (5) in the context of the Galois theory for linear differential equations. Thus, for foliations of type (5) problems (i) and (ii) are equivalent and Kovacic algorithm becomes an extremely powerful tool to approach them.

In some sense, this work can be considered as a very particular case of the Malgrange approach to the Galois theory of codimension-1 foliations [39, 40, 9], that is, for Riccati codimension- 1 foliations on the complex plane. Our target is not to obtain general theoretical
classification results, but several effective criteria of integrability for such foliations. We will provide integrability criteria for some families of polynomial quadratic vector fields and some Liénard equations involving special functions, allowing this to recover previous results established by several authors. For instance, we will solve completely the integrability problem for the family of Liénard type equations of the form (see [46])

$$
y y^{\prime}=\left(a(2 m+k) x^{2 k}+b(2 m-k) x^{m-k-1}\right) y-\left(a^{2} m x^{4 k}+c x^{2 k}+b^{2} m\right) x^{2 m-2 k-1},
$$

with $a, b, c, m, k$ complex parameters. It is known that when this equation comes from a polynomial vector field the constants $m$ and $k$ must be rational numbers.

We also approach the Poincaré problem for some particular families of systems (see Proposition 3.1 and Theorem 3.4).

The paper is structured as follows. To make it as self-contained as possible, we introduce in Section 1 a basic background about Galois theory of linear differential equations and Darboux theory of integrability of polynomial vector fields. In Section 2 we remind some useful properties concerning Riccati equations and Section 3 is devoted to applications. For completeness we include two Appendixes about Kovacic algorithm and some special functions.

## 1. Two notions of integrability for planar polynomial vector fields

1.1. Darboux theory of Integrability. We give a very brief overview of Darboux's integrability ideas [20], his terminology and some essential results.

Let us consider a vector field (2) and an irreducible polynomial $f \in \mathbb{C}[x, y]$. The curve $f=0$ is called an invariant algebraic curve of vector field (2) if it satisfies

$$
\left.\dot{f}\right|_{f=0}=0
$$

This condition is equivalent to the existence of a polynomial $K \in \mathbb{C}[x, y]$, called cofactor, such that

$$
X(f(x, y))=P(x, y) \frac{\partial f}{\partial x}+Q(x, y) \frac{\partial f}{\partial y}=K(x, y) f(x, y)
$$

or, equivalently, that

$$
\begin{equation*}
\frac{X(f)}{f}=X(\log (f))=K \tag{6}
\end{equation*}
$$

From this expression it follows that the curve $f=0$ is formed by leaves and critical points of the vector field $X=(P, Q)$ defined by (2). If the polynomial system (1) has degree $d$, that is $d=\max \{\operatorname{deg} P, \operatorname{deg} Q\}$, then we have that $\operatorname{deg} K \leq d-1$, independently of the degree of the curve $f(x, y)=0$. From definition (6) it follows that if the cofactor $K$ vanishes identically then the polynomial $f$ is a first integral of the vector field $X$. In terms of the associated foliation, this invariant curve $f=0$ is a particular solution of $y^{\prime}=Q / P$ and $Q d x-P d y=0$.

An analytic $\mathbb{C}$-valued non-constant function $\mu$ is called an integrating factor of system (1) if the expression

$$
\frac{X(\mu)}{\mu}=X(\log (\mu))=-\nabla \cdot X
$$

holds, where $\nabla \cdot X=(\partial P / \partial x)+(\partial Q / \partial y)$ is the divergence of the vector field $X=(P, Q)$. In case that the domain of definition of $X$ is simply connected, from the integrating factor $\mu$
it follows that

$$
H(x, y)=\int \mu(x, y) P(x, y) d y+\varphi(x)
$$

is a first integral of $X$, provided that $\partial H / \partial y=-\mu Q$.
To ensure the existence of a first integral for a system (1) is, in general, a very difficult problem. In [20], Darboux introduced a method to detect and construct first integrals using invariant algebraic curves. Namely, he proved that any planar polynomial differential system of degree $d$ having, at least, $d(d+1) / 2$ invariant algebraic curves, admits a first integral or an integrating factor which can be obtained from them (see also Jouanolou [27] for a study in a general context for codimension-1 foliations). Furthermore, Darboux's original ideas have been improved by taking into account the multiplicity of the invariant algebraic curves (see [19] for more details). Related to them some other invariant objects have been introduced (see [14]). They are the so-called exponential factors: given $h, g \in \mathbb{C}[x, y]$ relatively prime, the function $F=\exp (g / h)$ is called an exponential factor of the polynomial system (1) if there exists a polynomial $\widetilde{K} \in \mathbb{C}[x, y]$ (also called cofactor) that satisfies the equation

$$
\begin{equation*}
\frac{X(F)}{F}=X\left(\frac{g}{h}\right)=\widetilde{K} \tag{7}
\end{equation*}
$$

It is known that if $h$ is not a constant polynomial then $h=0$ is an invariant algebraic curve of (1) of cofactor $K_{h}$ satisfying that $X(g)=g K_{h}+h \widetilde{K}$.

The following theorem (coming originally from Darboux) shows how the construction of first integrals and integrating factors of (2) can be carried out from its invariant algebraic curves.

Theorem 1.1. Let consider a planar polynomial system (1) of degree $m$, having

- $p$ invariant algebraic curves $f_{i}=0$ with cofactors $K_{i}$, for $i=1, \ldots, p$ and
- $q$ exponential factors $F_{j}=\exp \left(g_{j} / h_{j}\right)$ with cofactors $\widetilde{K}_{j}, j=1, \ldots, q$.

Then the following assertions hold:
(a) There exist constants $\lambda_{i}, \tilde{\lambda}_{j} \in \mathbb{C}$ not all vanishing such that

$$
\sum_{i=1}^{p} \lambda_{i} K_{i}+\sum_{j=1}^{q} \widetilde{\lambda}_{j} \widetilde{K}_{j}=0
$$

if and only if the multivalued function

$$
\begin{equation*}
f_{1}^{\lambda_{1}} \ldots f_{p}^{\lambda_{p}} F_{1}^{\widetilde{\lambda}_{1}} \ldots F_{q}^{\tilde{\lambda}_{q}} \tag{8}
\end{equation*}
$$

is a (Darboux) first integral of system (1).
(b) There exist constants $\lambda_{i}, \widetilde{\lambda}_{j} \in \mathbb{C}$ not all vanishing such that

$$
\sum_{i=1}^{p} \lambda_{i} K_{i}+\sum_{j=1}^{q} \widetilde{\lambda}_{j} \widetilde{K}_{j}+\operatorname{div} X=0
$$

if and only if the function defined by (8) is a (Darboux) integrating factor of $X$.
For more recent versions of Theorem 1.1 see [37, 38] and for some generalizations see $[34,7]$.

Functions of the form (8) are called Darboux functions. We say that the polynomial system (1) is Darboux integrable if it admits a first integral or an integrating factor which is given by a Darboux function.

Remark 1.2. Prelle and Singer [48] showed that if system (1) admits an elementary first integral then it admits an integrating factor which is the $n$-th root of a rational function (a slightly improved version of this result can be found in [33, Corollary 4]). Later, Singer in [50] showed that if system (1) admits a Liouvillian first integral then it has an integrating factor which is given by a Darboux function. This is an important argument to motivate sentences like "Darboux functions capture Liouvillian integrability" or "Liouvillian first integrals are either Darboux first integrals or integrals coming from a Darboux integrating factor".

Given a polynomial system (1) of degree $m$, the computation of all its invariant algebraic curves becomes a complicated problem since nothing is known a priori about the maximum degree of these curves. This makes necessary to impose additional conditions either on the structure of the system (1) or on the nature of such curves (see for instance, [45, 8, 10, 17] or references therein). This difficulty has motivated the study of different types of inverse problems of the Darboux theory of integrability [43, 17, 15, 18, 16, 35].

We finish this subsection with a remark about a geometrical meaning of the integrating factor, pretty known to people coming from the Sophus Lie mathematical community. It is an established fact in Fluid Dynamics that integrating factors arise as a density in stationary planar regimes: the equation $\nabla \cdot(\mu X)=0$ is the continuity equation for the field of velocities $X$, with density function $\mu=\mu(x, y)$ (considered in the context of the Symplectic geometry and Hamiltonian dynamics).

Remark 1.3. Let $\mu=\mu(x, y)$ be an integrating factor of the vector field (2) defined in some domain of the plane. Then the vector field $X$ is a hamiltonian vector field with respect to the symplectic form $\omega=\mu d x \wedge d y$ (this form degenerates only at the zeros of $\mu$ ). In fact from $i_{X} \omega=d H$, we obtain

$$
-\mu Q=\omega\left(X, \frac{\partial}{\partial x}\right)=\frac{\partial H}{\partial x}, \quad \mu P=\omega\left(X, \frac{\partial}{\partial y}\right)=\frac{\partial H}{\partial y} .
$$

Hence, the vector field (2) can be rewritten

$$
X=P \frac{\partial}{\partial x}+Q \frac{\partial}{\partial y}=\frac{1}{\mu}\left(\frac{\partial H}{\partial y} \frac{\partial}{\partial x}-\frac{\partial H}{\partial x} \frac{\partial}{\partial y}\right)
$$

which is hamiltonian with Hamilton function the first integral H. It is straightforward to verify that the symplectic form $\omega$ is invariant under the action of the flow of $X$, i.e.,

$$
L_{X} \omega=d i_{X} \omega+i_{X} d \omega=d i_{X} \omega=d d H=0 .
$$

From this point of view, the dynamics of the vector fields in the plane can be formally considered as an Ergodic theory problem: the existence of an invariant measure, the one defined by the associated integrating factor of the flow.
1.2. Picard-Vessiot theory. Picard-Vessiot theory is the Galois theory of linear differential equations. We will just remind here some of its main definitions and results but we refer the reader to [49] for a wide theoretical background.

We start recalling some basic notions on algebraic groups and, afterwards, Picard-Vessiot theory will be introduced.

An algebraic group of matrices $2 \times 2$ is a subgroup $G \subset G L(2, \mathbb{C})$ defined by means of algebraic equations in its matrix elements and in the inverse of its determinant. That is, there exists a set of polynomials $P_{i} \in \mathbb{C}\left[x_{1}, \ldots, x_{5}\right]$, for $i \in I$, such that $A \in \operatorname{GL}(2, \mathbb{C})$ given by

$$
A=\left(\begin{array}{ll}
x_{11} & x_{12} \\
x_{21} & x_{22}
\end{array}\right)
$$

belongs to $G$ if and only if $P_{i}\left(x_{11}, x_{12}, x_{21}, x_{22},(\operatorname{det} A)^{-1}\right)=0$ for all $i \in I$ and where $\operatorname{det} A=$ $x_{11} x_{22}-x_{21} x_{12}$. It is said that $G$ is an algebraic manifold endowed with a group structure.

Recall that a group $G$ is called solvable if and only if there exists a chain of normal subgroups

$$
e=G_{0} \triangleleft G_{1} \triangleleft \ldots \triangleleft G_{n}=G
$$

satisfying that the quotient $G_{i} / G_{j}$ is abelian for all $n \geq i \geq j \geq 0$.
It is well known that any algebraic group $G$ has a unique connected normal algebraic subgroup $G^{0}$ of finite index. In particular, the identity connected component $G^{0}$ of $G$ is defined as the largest connected algebraic subgroup of $G$ containing the identity. In case that $G=G^{0}$ we say that $G$ is a connected group. Moreover, if $G^{0}$ is solvable we say that $G$ is virtually solvable.

The following result provides the relation between virtual solvability of an algebraic group and its structure.

Theorem 1.4 (Lie-Kolchin). Let $G \subseteq \mathrm{GL}(2, \mathbb{C})$ be a virtually solvable group. Then, $G^{0}$ is triangularizable, that is, it is conjugate to a subgroup of upper triangular matrices.

Now, we briefly introduce Picard-Vessiot Theory.
First, we say that ( $\left.\mathcal{K},{ }^{\prime}\right)$ - or, simply, $\mathcal{K}$ - is a differential field if $\mathcal{K}$ is a commutative field of characteristic zero, depending on $x$ and ' is a derivation on $\mathcal{K}$ (that is, satisfying that $(a+b)^{\prime}=a^{\prime}+b^{\prime}$ and $(a \cdot b)^{\prime}=a^{\prime} \cdot b+a \cdot b^{\prime}$ for all $\left.a, b \in \mathcal{K}\right)$. We denote by $\mathcal{C}$ the field of constants of $\mathcal{K}$, defined as $\mathcal{C}=\left\{c \in \mathcal{K} \mid c^{\prime}=0\right\}$.

We will deal with second order linear homogeneous differential equations, that is, equations of the form

$$
\begin{equation*}
y^{\prime \prime}+b_{1} y^{\prime}+b_{0} y=0, \quad b_{1}, b_{0} \in \mathcal{K}, \tag{9}
\end{equation*}
$$

and we will be concerned with the algebraic structure of their solutions. Moreover, along this work, we will refer to the current differential field as the smallest one containing the field of coefficients of this differential equation.

Let us suppose that $y_{1}, y_{2}$ is a basis of solutions of equation (9), i.e., $y_{1}, y_{2}$ are linearly independent over $\mathcal{K}$ and every solution is a linear combination over $\mathcal{C}$ of these two. Let $\mathcal{L}=\mathcal{K}\left\langle y_{1}, y_{2}\right\rangle=\mathcal{K}\left(y_{1}, y_{2}, y_{1}^{\prime}, y_{2}^{\prime}\right)$ be the differential extension of $\mathcal{K}$ such that $\mathcal{C}$ is the field of constants for $\mathcal{K}$ and $\mathcal{L}$. In this terms, we say that $\mathcal{L}$, the smallest differential field containing $\mathcal{K}$ and $\left\{y_{1}, y_{2}\right\}$, is the Picard-Vessiot extension of $\mathcal{K}$ for the differential equation (9).

The group of all the differential automorphisms of $\mathcal{L}$ over $\mathcal{K}$ that commute with the derivation ' is called the Galois group of $\mathcal{L}$ over $\mathcal{K}$ and is denoted by $\operatorname{Gal}(\mathcal{L} / \mathcal{K})$. This means, in particular, that for any $\sigma \in \operatorname{Gal}(\mathcal{L} / \mathcal{K}), \sigma\left(a^{\prime}\right)=(\sigma(a))^{\prime}$ for all $a \in \mathcal{L}$ and that $\sigma(a)=a$ for all $a \in \mathcal{K}$. Thus, if $\left\{y_{1}, y_{2}\right\}$ is a fundamental system of solutions of (9) and $\sigma \in \operatorname{Gal}(\mathcal{L} / \mathcal{K})$
then $\left\{\sigma y_{1}, \sigma y_{2}\right\}$ is also a fundamental system. This implies the existence of a non-singular constant matrix

$$
A_{\sigma}=\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right) \in \mathrm{GL}(2, \mathbb{C})
$$

such that

$$
\sigma\left(\begin{array}{ll}
y_{1} & y_{2}
\end{array}\right)=\left(\begin{array}{ll}
\sigma\left(y_{1}\right) & \sigma\left(y_{2}\right)
\end{array}\right)=\left(\begin{array}{ll}
y_{1} & y_{2}
\end{array}\right) A_{\sigma} .
$$

This fact can be extended in a natural way to a system

$$
\sigma\left(\begin{array}{ll}
y_{1} & y_{2} \\
y_{1}^{\prime} & y_{2}^{\prime}
\end{array}\right)=\left(\begin{array}{ll}
\sigma\left(y_{1}\right) & \sigma\left(y_{2}\right) \\
\sigma\left(y_{1}^{\prime}\right) & \sigma\left(y_{2}^{\prime}\right)
\end{array}\right)=\left(\begin{array}{ll}
y_{1} & y_{2} \\
y_{1}^{\prime} & y_{2}^{\prime}
\end{array}\right) A_{\sigma},
$$

which leads to a faithful representation $\operatorname{Gal}(\mathcal{L} / \mathcal{K}) \rightarrow \mathrm{GL}(2, \mathbb{C})$ and makes possible to consider $\operatorname{Gal}(\mathcal{L} / \mathcal{K})$ as a subgroup of $\mathrm{GL}(2, \mathbb{C})$ depending (up to conjugacy) on the choice of the fundamental system $\left\{y_{1}, y_{2}\right\}$.

One of the fundamental results of the Picard-Vessiot Theory is the following theorem (see [28, 31]).
Theorem 1.5. The Galois group $\operatorname{Gal}(\mathcal{L} / \mathcal{K})$ is an algebraic subgroup of $\operatorname{GL}(2, \mathbb{C})$.
We say that equation (9) is integrable if the Picard-Vessiot extension $\mathcal{L} \supset \mathcal{K}$ is obtained as a tower of differential fields $\mathcal{K}=\mathcal{L}_{0} \subset \mathcal{L}_{1} \subset \cdots \subset \mathcal{L}_{m}=\mathcal{L}$ such that $\mathcal{L}_{i}=\mathcal{L}_{i-1}(\eta)$ for $i=1, \ldots, m$, where either
(i) $\eta$ is algebraic over $\mathcal{L}_{i-1}$, that is $\eta$ satisfies a polynomial equation with coefficients in $\mathcal{L}_{i-1}$.
(ii) $\eta$ is primitive over $\mathcal{L}_{i-1}$, that is $\eta^{\prime} \in \mathcal{L}_{i-1}$.
(iii) $\eta$ is exponential over $\mathcal{L}_{i-1}$, that is $\eta^{\prime} / \eta \in \mathcal{L}_{i-1}$.

Usually in terms of Differential Algebra's terminology we say that equation (9) is integrable if the corresponding Picard-Vessiot extension is Liouvillian. Moreover, the following theorem holds.

Theorem 1.6 (Kolchin). Equation (9) is integrable if and only if $\operatorname{Gal}(\mathcal{L} / \mathcal{K})$ is virtually solvable, that is, its identity component $(\operatorname{Gal}(\mathcal{L} / \mathcal{K}))^{0}$ is solvable.

For instance, for the case $a=0$ in equation (9), i.e. $y^{\prime \prime}+b y=0$, it is very well known [28, 31, 49] that $\operatorname{Gal}(\mathcal{L} / \mathcal{K})$ is an algebraic subgroup of $\operatorname{SL}(2, \mathbb{C})$ (remind that $A \in \operatorname{SL}(2, \mathbb{C}) \Leftrightarrow$ $A \in \mathrm{GL}(2, \mathbb{C})$ and $\operatorname{det} A=1)$. For a more detailed study see the Appendix A.

## 2. Some remarks about Riccati equation

Ricatti equation is probably one of the most studied equations in Dynamical Systems. Its rôle in the study of the Darboux and Picard-Vessiot integrability leads us to devote this section to some of its properties. Even though these results are known, their proofs have been included for completeness. We divide these properties in two types: the first one (see Subsection 2.1) concerning transformations leading a general second order differential equation into a Riccati equation (written in the so-called reduced form). Remind that this becomes the starting point of the celebrated Kovacic algorithm (see Appendix A). A second type, more Darboux-like, that studies first integrals and integrating factors for a Riccati equation (Subsection 2.2).
2.1. Transformations related to Riccati equations. It is known that any second order differential equation can be led into a general Riccati equation through a classical logarithmic change of variable (see, for instance, $[46,26]$ ). The following proposition recall it and summarises some other related transformations.
Proposition 2.1. Let $\mathcal{K}$ be a differential field and let consider functions $a_{0}(x), a_{1}(x), a_{2}(x)$, $r(x), \rho(x), b_{0}(x), b_{1}(x)$ belonging to $\mathcal{K}$ that, for simplicity, will be denoted without their explicit dependence on $x$. Consider now the following forms associated to any second order differential equation (ode) and Riccati equation:
(i) Second order ode (in general form):

$$
\begin{equation*}
y^{\prime \prime}+b_{1} y^{\prime}+b_{0} y=0 \tag{10}
\end{equation*}
$$

(ii) Second order ode (in reduced form):

$$
\begin{equation*}
\xi^{\prime \prime}=\rho \xi \tag{11}
\end{equation*}
$$

(iii) Riccati equation (in general form):

$$
\begin{equation*}
v^{\prime}=a_{0}+a_{1} v+a_{2} v^{2}, \quad a_{2} \neq 0 \tag{12}
\end{equation*}
$$

(iv) Riccati equation (in reduced form):

$$
\begin{equation*}
w^{\prime}=r-w^{2}, \tag{13}
\end{equation*}
$$

Then, there exist transformations $\mathcal{T}, \mathcal{B}, \mathcal{S}$ and $\mathcal{R}$ leading some of these equations into the other ones, as showed in the following diagram:


The new independent variables are defined by means of

$$
\begin{array}{ll}
\mathcal{T}: v=-\left(\frac{a_{2}^{\prime}}{2 a_{2}^{2}}+\frac{a_{1}}{2 a_{2}}\right)-\frac{1}{a_{2}} w, & \mathcal{B}: v=-\frac{1}{a_{2}} \frac{y^{\prime}}{y}, \\
\mathcal{S}: y=\xi \mathrm{e}^{-\frac{1}{2} \int b_{1} d x}, & \mathcal{R}: w=\frac{\xi^{\prime}}{\xi}
\end{array}
$$

and the functions $r, \rho, b_{0}$ and $b_{1}$ are given by

$$
\begin{align*}
r & =\frac{1}{\beta}\left(a_{0}+a_{1} \alpha+a_{2} \alpha^{2}-\alpha^{\prime}\right)  \tag{14}\\
\alpha & =-\left(\frac{a_{2}^{\prime}}{2 a_{2}^{2}}+\frac{a_{1}}{2 a_{2}}\right), \quad \beta=-\frac{1}{a_{2}},  \tag{15}\\
b_{1} & =-\left(a_{1}+\frac{a_{2}^{\prime}}{a_{2}}\right), \quad b_{0}=a_{0} a_{2},  \tag{16}\\
\rho & =r=\frac{b_{1}^{2}}{4}+\frac{b_{1}^{\prime}}{2}-b_{0} . \tag{17}
\end{align*}
$$

Proof. The proof is quite standard.
$[\mathcal{T}]$ : Applying the change $v=\alpha+\beta w$ we get the equation

$$
\alpha^{\prime}+\beta^{\prime} w+\beta w^{\prime}=a_{0}+a_{1} \alpha+a_{1} \beta w+a_{2} \alpha^{2}+2 a_{2} \alpha \beta w+a_{2} \beta^{2} w^{2}
$$

that, regrouping terms, leads to

$$
w^{\prime}=\frac{1}{\beta}\left(a_{0}+a_{1} \alpha+a_{2} \alpha^{2}-\alpha^{\prime}\right)+\left(a_{1}+2 a_{2} \alpha-\frac{\beta^{\prime}}{\beta}\right) w+a_{2} \beta w^{2} .
$$

Since $a_{2} \neq 0$ we can take $\beta=-1 / a_{2}$ and, therefore, $a_{2} \beta=-1$. Having this into account, the value of $\alpha$ satisfying that the coefficient in $w$ vanishes is given by

$$
\alpha=\frac{1}{2 a_{2}}\left(\frac{\beta^{\prime}}{\beta}-a_{1}\right) .
$$

The expressions for $\alpha, \beta$ and $r$ follow straightforwardly,

$$
r=\frac{1}{\beta}\left(a_{0}+a_{1} \alpha+a_{2} \alpha^{2}-\alpha^{\prime}\right), \quad \alpha=-\left(\frac{a_{2}^{\prime}}{2 a_{2}^{2}}+\frac{a_{1}}{2 a_{2}}\right), \quad \beta=-\frac{1}{a_{2}} .
$$

Moreover, it is clear that $\alpha, \beta$ and $r$ belong to $\mathcal{K}$.
$[\mathcal{B}]:$ Imposing $\alpha=0$ and taking $\beta=-1 / a_{2}$ in transformation $\mathcal{T}$ we have $v=-w / a_{2}$ and we obtain the Riccati equation

$$
w^{\prime}=-a_{0} a_{2}+\left(a_{1}+\frac{a_{2}^{\prime}}{a_{2}}\right) w-w^{2}
$$

Performing now the change of variables $w=(\log y)^{\prime}$ (or, equivalently, $v=-a_{2} y^{\prime} / y$ ) we obtain the differential equation $y^{\prime \prime}+b_{1} y^{\prime}+b_{0} y=0$ with

$$
b_{1}=-\left(a_{1}+\frac{a_{2}^{\prime}}{a_{2}}\right), \quad b_{0}=a_{0} a_{2}
$$

Obviously, $b_{0}$ and $b_{1}$ belong to $\mathcal{K}$.
[ $\mathcal{S}$ ]: The change of variable $y=\mu \xi$, with $\mu=\mu(x)$ and $\xi=\xi(x)$, lead us to

$$
\xi^{\prime \prime}+\left(2 \frac{\mu^{\prime}}{\mu}+b_{0}\right) \xi^{\prime}+\left(\frac{\mu^{\prime \prime}}{\mu}+b_{0} \frac{\mu^{\prime}}{\mu}+b_{1}\right) \xi=0
$$

In order to obtain the equation $\xi^{\prime \prime}=\rho \xi$ we need to impose

$$
2 \frac{\mu^{\prime}}{\mu}+b_{0}=0, \quad \frac{\mu^{\prime \prime}}{\mu}+b_{0} \frac{\mu^{\prime}}{\mu}+b_{1}=-\rho,
$$

which gives rise to

$$
\mu=\mathrm{e}^{-\frac{1}{2} \int \mathrm{~b}_{0}}, \quad \rho=\frac{\mathrm{b}_{0}^{2}}{4}+\frac{\mathrm{b}_{0}^{\prime}}{2}-\mathrm{b}_{1} .
$$

Moreover, it is straightforward to check that $\rho \in \mathcal{K}$.
$[\mathcal{R}]:$ This is a particular case of transformation $[\mathcal{B}]$ with the particular choice $a_{0}=r$, $a_{1}=0$ and $a_{2}=-1$.

Finally, composing the transformations provided by $[\mathcal{B}],[\mathcal{R}]$ and $[\mathcal{S}]$ :

$$
-a_{2} v=\frac{y^{\prime}}{y}, \quad y=\xi e^{-\frac{1}{2} \int b_{0}}, \quad b_{0}=-\left(a_{1}+\frac{a_{2}^{\prime}}{a_{2}}\right) \quad \frac{\xi^{\prime}}{\xi}=w
$$

we recover the result given by $[\mathcal{T}]$,

$$
v=-\left(\frac{a_{1}}{2 a_{2}}+\frac{a_{2}^{\prime}}{2 a_{2}^{2}}\right)-\frac{1}{a_{2}} w=\alpha+\beta w,
$$

which implies that, in some sense and taking $\rho=r$, the diagram commutes: $\mathcal{S} \circ \mathcal{B}=\mathcal{R} \circ \mathcal{T}$.

From this Lemma, it follows that the function $v$ is algebraic over $\mathcal{K}$ if and only if the function $w$ is also algebraic over $\mathcal{K}$. Furthermore, in such case, the degree over $\mathcal{K}$ of both functions $v$ and $w$ is the same.

It is known that a Riccati equation (12) has an algebraic solution over $\mathcal{K}$ if and only if the differential equation (10) is integrable in a Picard-Vessiot sense. In this situation we say that the Riccati equation is integrable over $\mathcal{K}$. We notice that Kovacic algorithm (see Appendix A) starts from an equation in form (11).
2.2. Integrating factor and first integrals for Riccati vector fields. We briefly show some relations between the existence of invariant curves of a certain type of vector fields and the integrability, via Kovacic algorithm (see Appendix A), of its associated Riccati foliation.

From Singer [50] (see Remark 1.2) we know that if a planar polynomial vector field (2) admits a Liouvillian first integral then it has also an integrating factor given by a Darboux function. However, very few results are known about the relation between the existence of an algebraic invariant curve of a general planar vector field and the Liouvillian integrability of its foliation.

Let us consider a family of planar vector fields of the form

$$
\begin{equation*}
X=\left(p(x)-q(x) w^{2}\right) \frac{\partial}{\partial w}+q(x) \frac{\partial}{\partial x}, \tag{18}
\end{equation*}
$$

with $p(x), q(x) \in \mathbb{C}[x]$ complex polynomials. Introducing an independent variable $t$, usually called time, we can associate to them the following system of differential equations

$$
\begin{aligned}
\dot{w} & =p(x)-q(x) w^{2}, \\
\dot{x} & =q(x),
\end{aligned}
$$

where we denote by $\cdot=d / d t$. Its foliation, governed by the equation

$$
w^{\prime}=\frac{d w}{d x}=\frac{p(x)-q(x) w^{2}}{q(x)}=\frac{p(x)}{q(x)}-w^{2},
$$

is a Riccati equation given in reduced form $w^{\prime}=r(x)-w^{2}$ with $r=p / q \in \mathbb{C}(x)$. Next lemma shows that the integrability of this "Riccati foliation" is closely related to the existence of an algebraic invariant curve of its vector field (18). A similar approach for this problem can be found in $[24,25]$.

Lemma 2.2. Let $w_{1}=w_{1}(x)$ be a solution of a Riccati equation in reduced form

$$
w^{\prime}=r(x)-w^{2},
$$

with $r(x)=p(x) / q(x) \in \mathbb{C}[x]$. Then the associated vector field (18) has an integrating factor given by

$$
\begin{equation*}
\mu_{1}(w, x)=\frac{e^{-2 \int w_{1}(x) d x}}{\left(-w+w_{1}(x)\right)^{2}} \tag{19}
\end{equation*}
$$

Proof. It is straightforward to check that if $w_{1}(x)$ is a solution of $w^{\prime}=p / q-w^{2}$ then it holds $X\left(f_{1}\right)=K_{1} f_{1}$ with $f_{1}(w, x)=-w+w_{1}(x)=0$ and $K_{1}=-q\left(w+w_{1}(x)\right)$. In addition, $F_{1}(x)=\mathrm{e}^{-\int \omega_{1}(\mathrm{x}) \mathrm{dx}}$ satisfies $X\left(F_{1}\right)=L_{1} F_{1}$ with $L_{1}=-q \omega_{1}$. Notice that $X$ has divergence $\operatorname{div} X=-2 q w$, and additionally it holds $-2 K_{1}+2 L_{1}+\operatorname{div} X=0$. Proceeding as in Theorem 1.1, vector field (18) admits the integrating factor

$$
\mu_{1}(w, x)=\frac{F_{1}^{2}}{f_{1}^{2}}=\frac{e^{-2 \int w_{1}(x) d x}}{\left(-w+w_{1}(x)\right)^{2}},
$$

as it was claimed.

Remark 2.3. We stress that the result in Lemma 2.2 is independent of the nature of the solution $w=w_{1}(x)$. In general, the integral $\int \omega_{1}(x) d x$ is an abelian integral.

The important fact is that, conversely, Picard-Vessiot theory and in particular, Kovacic algorithm, supply information about first integrals and integrating factors of the equation $w^{\prime}=r(x)-w^{2}$ from the knowledge of some of its solutions, $w_{1}, w_{2}, w_{3}$. Indeed, from the first three cases in Kovacic algorithm [30] (the integrable ones) one obtains the following types of first integrals (see Weil [54] and Zoła̧dek [57]).
Proposition 2.4. The following statements hold.
Case 1: One has two possibilities:

- If only $w_{1} \in \mathbb{C}(x)$ then $X$ admits a first integral of Darboux-Schwarz-Christoffel type.
- If both $w_{1}, w_{2} \in \mathbb{C}(x)$ then $X$ admits a first integral of Darboux type. In particular, from Lemma 2.2 we can construct two integrating factors $\mu_{1}$ and $\mu_{2}$ and so $\mu_{1} / \mu_{2}$ is a first integral of $X$. Thus, we have

$$
H(w, x)=\frac{\left(-w+w_{2}(x)\right)}{\left(-w+w_{1}(x)\right)} e^{\int\left[\left(w_{2}(x)-w_{1}(x)\right) d x\right]}
$$

Case 2: If $w_{1}$ is a solution of a quadratic polynomial then $X$ admits a first integral of hyperelliptic type.
Case 3: If all $w_{1}, w_{2}, w_{3}$ are algebraic over $\mathbb{C}(x)$ then $X$ admits a rational first integral.
The following result characterises the rational integrability of the polynomial vector fields that we are considering in this work.

Corollary 2.5. The Galois group of (11) is finite if and only if its corresponding planar polynomial vector field has a rational first integral.

Proof. Notice that $\operatorname{Gal}(\mathcal{L} / \mathcal{K})$ is finite if and only if we fall in case 3 of Kovacic Algorithm or in case 1 of the type

$$
\operatorname{Gal}(\mathcal{L} / \mathcal{K})=\left\{\left(\begin{array}{cc}
c & 0 \\
0 & c^{-1}
\end{array}\right), \quad c^{n}=1\right\} .
$$

Hence, only remains to study this last case (the cyclic one). Let $\xi_{1}, \xi_{2}$ be solutions of $\xi^{\prime \prime}=\rho \xi$. Then there exists $g \in \mathbb{C}(x)$ such that $\xi_{1}=g^{\frac{1}{n}}$ and $\xi_{2}=g^{-\frac{1}{n}}$. We define $\omega_{1}=\xi_{1}^{\prime} / \xi_{1}, \omega_{2}=\xi_{2}^{\prime} / \xi_{2}$, and we obtain

$$
\omega_{1}=\frac{1}{n} \frac{g^{\prime}}{g}, \quad \omega_{2}=-\frac{1}{n} \frac{g^{\prime}}{g} .
$$

Thus, the corresponding vector field (11) admits the first integral

$$
H_{1}(\omega, x)=\frac{-\omega+\omega_{2}}{-\omega+\omega_{1}} e^{\int\left(\omega_{2}-\omega_{1}\right) d x}=\frac{-\omega-\frac{1}{n} \frac{g^{\prime}}{g}}{-\omega+\frac{1}{n} \frac{g^{\prime}}{g}} g^{\frac{-2}{n}},
$$

and so also admits the first integral

$$
H(\omega, x)=H_{1}(\omega, x)^{n}=\frac{1}{g^{2}}\left(\frac{-n g \omega-g^{\prime}}{-n g \omega+g^{\prime}}\right)^{n} \in \mathbb{C}(\omega, x),
$$

which completes the proof.

Remark 2.6. Let $P \in \mathbb{C}[x]$ be a polynomial of odd degree. It is known that the planar polynomial vector field $\dot{x}=1, \dot{y}=P(x)+y^{2}$ (with associated foliation $y^{\prime}=P(x)+y^{2}$ ) is not integrable, that is, it has no invariant curves, since it falls in case 4 of Kovacic Algorithm [30].
2.3. Riccati foliations. Let us recall some well-known geometrical properties of Riccati foliations defined by planar polynomials vector fields (see, for instance, [32]). Although we are not using these properties along the paper, we include them for completeness.

Let

$$
\begin{equation*}
\Omega=q(x) d y-\left(p_{1}(x)+p_{2}(x) y+p_{3}(x) y^{2}\right) d x, \quad p_{i}, q \in \mathbb{C}[x] \tag{20}
\end{equation*}
$$

be the 1 -form defining a Riccati foliation on the complex plane. Since the equation $\Omega=0$ is the projective version of the corresponding second order linear differential equation defined over the vector bundle $\mathbf{P}^{1} \times \mathbb{C}^{2}$ (i.e., the fibre $\mathbb{C}^{2}$ is projectivized to $\mathbf{P}^{1}$, see subsection 2.1), the Riccati equation $\Omega=0$ is defined in a natural way over $(x, y) \in \mathbf{P}^{1} \times \mathbf{P}^{1}$.

The singular points of $\Omega=0$ are the zeros of the polynomial $q(x)$ and, possibly, the point at infinity $x=\infty \in \mathbf{P}^{1}$. Moreover, these singular points are exactly the poles of the coefficients of the associated second order linear differential equation. We define $d:=$ $\max \left(\operatorname{deg}\left(p_{1}\right), \operatorname{deg}\left(p_{2}\right), \operatorname{deg}\left(p_{3}\right), \operatorname{deg}(q)-2\right)$. Thus the point $x=\infty$ is a singular point if and only if $\operatorname{deg}(q)$ is less than $d+2$. In Kovacic algorithm, which applies to the reduced form of the second order linear differential equation (see Appendix A), this set of singular points is denoted by $\Gamma$. Therefore it seems natural to call it in the same way also here, that is, $\Gamma=\left\{x_{1}, \ldots, x_{r}\right\}$. Thus the Riccati foliation is holomorphic on $\left(\mathbf{P}^{1}-\Gamma\right) \times \mathbf{P}^{1}$, because the singular points of the associated linear differential equation are exactly the set $\Gamma$, given by the singularities of their coefficients. We notice that the ("singular") sets $\left\{x_{i}\right\} \times P^{1}$ are invariant
by the foliation and are usually called invariant fibres because they are already fibres of the fibration $\pi: \mathbf{P}^{1} \times \mathbf{P}^{1} \rightarrow \mathbf{P}^{1},(x, y) \mapsto x$. This fibration is transversal to the Riccati foliation since fibres $T(x):=\pi^{-1}(x)=\{x\} \times \mathbf{P}^{1}$, with $x$ non singular, are global and transversal to the foliation, i.e., transversal to all the leaves. Over any of these transversals the holonomy group of the foliation is defined as a representation

$$
\pi_{1}\left(\mathbf{P}^{1}-\Gamma, x_{0}\right) \rightarrow \operatorname{Diff}\left(T\left(x_{0}\right)\right)
$$

(where $\operatorname{Diff}\left(T\left(x_{0}\right)\right)$ is the group of diffeomorphisms on the transversal), given by lifting the loops in the fundamental group to the leaves of the foliation, this is, by solving the Riccati equation with initial conditions and final points on the transversal $T\left(x_{0}\right)$. As the Riccati equation is the projectivization of a second order linear differential equation, the holonomy group must be the projectivization of the monodromy group of the linear second order equation acting on the vector space fibre of the meromorphic vector bundle $\mathbf{P}^{1} \times \mathbb{C}^{2},\{x\} \times \mathbb{C}^{2} \approx \mathbf{C}^{2}$. By fixing a base of fundamental solutions, this can be considered as the space of solutions of the linear differential equation. Hence, as the monodromy group is represented by the linear group $G L(2, \mathbb{C})$, the holonomy group is represented by the projective linear group $\operatorname{PGL}(2, \mathbb{C})$, the Möbius transformations

$$
\pi_{1}\left(\mathbf{P}^{1}-\Gamma, x_{0}\right) \rightarrow P G L(2, \mathbb{C})
$$

The Riccati foliations are the most well-known class of a family of foliations, the projective foliations, with holonomy group represented in the projective group. For Riccati foliations, the holonomy group is contained in the Galois group of the foliation either in the Malgrange approach [39, 40, 9] or in the Lie-Vessiot-Kolchin approach [6, 5]. In fact, if the singular points of the associated linear differential equations are singular regular ones, then the Zariski adherence of the holonomy group is the Galois group of the foliation, because in this case the Zariski adherence of the monodromy group of the linear differential equation is the Galois group of the associated linear differential equation.

The critical points of the associated vector field, i.e., zeroes of $q(x)$ and of $p_{1}(x)+p_{2}(x) y+$ $p_{3}(x) y^{2}$ are obviously contained in the invariant fibres. For general Riccati foliations there are two critical points on the invariant fibre for each point in $\Gamma$. For some special Riccati foliations there are no critical points. This is the case, for example, of Riccati foliations given in reduced form by $\Omega=q(x) d y-\left(p(x)+q(x) y^{2}\right) d x$ (with $p$ and $q$ relatively primes), corresponding to the field (18).

## 3. Applications

In this section we analyse some examples involving integrability and non-integrability of some families of Riccati planar vector fields or planar vector fields whose foliation can be reduced into a Riccati form.
3.1. Quadratic polynomials fields. The study of the integrability of the quadratic polynomial vector field

$$
\begin{aligned}
& \dot{x}=a_{20} x^{2}+a_{11} x y+a_{02} y^{2}+a_{10} x+a_{01} y+a_{00}, \\
& \dot{y}=b_{20} x^{2}+b_{11} x y+b_{02} y^{2}+b_{10} x+b_{01} y+b_{00},
\end{aligned}
$$

with $a_{i j}, b_{i, j} \in \mathbb{C}$ is, in its general form, a hard problem. One of its possible approaches is the so-called linear-quadratic case, when one of the two components is a polynomial of degree
one. In [36, Prop.3] it is proved that its study around a finite equilibrium point (the origin) can be reduced to consider two families of systems. Using the notation introduced therein, we refer to these families as (S1)-type,

$$
\begin{align*}
& \dot{x}=x, \\
& \dot{y}=\varepsilon x+\lambda y+b_{20} x^{2}+b_{11} x y+b_{02} y^{2} \tag{S1}
\end{align*}
$$

and (S2)-type,

$$
\begin{align*}
& \dot{x}=y, \\
& \dot{y}=\varepsilon x+\lambda y+b_{20} x^{2}+b_{11} x y+b_{02} y^{2} . \tag{S2}
\end{align*}
$$

In [36], the authors prove that the linear-quadratic systems having a global analytic first integral are those satisfying:

$$
\begin{aligned}
& \left(a_{1}\right) b_{02}=\lambda=0 . \\
& \left(b_{1}\right) b_{02}=0 \text { and } \lambda=-p / q \in \mathbb{Q}^{-},
\end{aligned}
$$

in the case of (S1)-type systems and
$\left(a_{2}\right) b_{20}=b_{02}=\lambda=0$ and $\varepsilon b_{11} \neq 0$.
$\left(b_{2}\right) b_{20}=b_{11}=\lambda=0$ and $\varepsilon b_{02} \neq 0$.
( $\left.c_{2}\right) b_{11}=\lambda=0$ and $b_{20} \neq 0$,
for (S2)-type systems. Furthermore, they also provide the explicit form of the corresponding first integrals. It is important to notice that all of them are of Darboux type and, therefore, Liouvillian.

Our aim in this example is to show that these results can be recovered using arguments coming from the Galois theory of linear differential equations. We start first with the (S1)case, whose associated foliation is given by the Riccati equation:

$$
\begin{equation*}
\frac{d y}{d x}=\left(\varepsilon+b_{20} x\right)+\left(\frac{\lambda+b_{11} x}{x}\right) y+\frac{b_{02}}{x} y^{2} . \tag{21}
\end{equation*}
$$

By Lemma 2.1 this equation can be transformed into the reduced form $w^{\prime}=r(x)-w^{2}$, with

$$
\begin{equation*}
r(x)=\frac{1}{4}-\frac{\kappa}{x}+\frac{4 \mu^{2}-1}{4 x^{2}}, \quad \kappa=\frac{1}{\sqrt{b_{11}^{2}-4 b_{20} b_{02}}}\left(b_{02} \varepsilon+\frac{b_{11}}{2}(1-\lambda)\right), \quad \mu=\frac{\lambda}{2}, \tag{22}
\end{equation*}
$$

provided $b_{11}^{2}-4 b_{20} b_{02} \neq 0$, and into the form $\xi^{\prime \prime}=r(x) \xi$. This equation is a Whittaker equation (see Appendix B) to which one can apply the Martinet-Ramis Theorem (see theorem B.2). This Theorem asserts that such Whittaker equation is integrable if and only if at least one of the following conditions is verified:

$$
\pm \kappa \pm \mu \in \frac{1}{2}+\mathbb{N}
$$

or, equivalently (and more suitable for the expressions derived of $\kappa$ and $\mu$ )

$$
2(\kappa \pm \mu) \in 2 \mathbb{Z}+1
$$

In our case one has that

$$
2(\kappa \pm \mu)=\frac{2 b_{20} \varepsilon+b_{11}(1-\lambda)}{\sqrt{b_{11}^{2}-4 b_{20} b_{02}}} \pm \lambda
$$

so for (S1)-type systems conditions $\left(a_{1}\right)$ and ( $b_{1}$ ) read, respectively, $2(\kappa \pm \mu)=1 \in 2 \mathbb{Z}+1$ and $2(\kappa+\mu)=(1+(p / q))+(-p / q)=1 \in 2 \mathbb{Z}+1$. Therefore Galois Theory recovers the integrability result asserted in [36, Thm.1].

Let us consider now a (S2)-type system, namely,

$$
\begin{aligned}
& \dot{x}=y \\
& \dot{y}=\varepsilon x+\lambda y+b_{20} x^{2}+b_{11} x y+b_{02} y^{2}
\end{aligned}
$$

with foliation given by the differential equation

$$
\begin{equation*}
\frac{d y}{d x}=\left(\lambda+b_{11} x\right)+\left(\varepsilon x+b_{20} x^{2}\right) \frac{1}{y}+b_{02} y . \tag{23}
\end{equation*}
$$

This equation falls in one of the following situations:
(i) $\lambda=b_{11}=0$ yields to a Bernoulli equation

$$
\frac{d y}{d x}=\left(\varepsilon x+b_{20} x^{2}\right) \frac{1}{y}+b_{02} y
$$

which corresponds to cases $\left(b_{2}\right)$ and $\left(c_{2}\right)$.
(ii) If $\varepsilon=b_{20}=0$ we obtain the linear equation (and, of course, integrable in a Liouville sense)

$$
\frac{d y}{d x}=\left(\lambda+b_{11} x\right)+b_{02} y
$$

This possibility is not taken into account by Llibre and Valls [36] since this equation is not, strictly speaking, in Riccati form.
(iii) $b_{20}=b_{02}=\lambda=0$ and $\varepsilon b_{11} \neq 0$ (case $\left.\left(a_{2}\right)\right)$ gives rise to $d y / d x=b_{11} x+\varepsilon x y^{-1}$, which is a separable equation (and a Bernoulli as well) and whose solutions are all Liouvillian.
(iv) If $b_{02}=0$ we obtain a Liénard equation,

$$
y \frac{d y}{d x}=\left(\lambda+b_{11} x\right) y+\left(\varepsilon x+b_{20} x^{2}\right)
$$

that will be considered more deeply in a forthcoming section.
3.2. Families of orthogonal polynomials. Recall that the Hypergeometric equation, including confluences, is a particular case of the differential equation

$$
\begin{equation*}
y^{\prime \prime}+\frac{L}{Q} y^{\prime}+\frac{\lambda}{Q} y=0, \quad \lambda \in \mathbb{C}, \quad L=a_{0}+a_{1} x, \quad Q=b_{0}+b_{1} x+b_{2} x^{2} \tag{24}
\end{equation*}
$$

It is well known (see, for example, [13]) that classical orthogonal and Bessel polynomials are solutions of equation (24) for suitable values of $a_{j}, b_{j}$ and $\lambda$. Namely,

- Hermite $H_{n}$,
- Chebyshev of first kind $T_{n}$,
- Chebyshev of second kind $U_{n}$,
- Legendre $P_{n}$,
- Laguerre $L_{n}$,
- associated Laguerre $L_{n}^{(m)}$,
- Gegenbauer $C_{n}^{(m)}$,
- Jacobi $\mathcal{P}_{n}^{(m, \nu)}$
- Bessel $B_{n}$,
where

| Family | $Q$ | $L$ | $\lambda$ |
| :---: | :---: | :---: | :---: |
| $H_{n}$ | 1 | $-2 x$ | $2 n$ |
| $T_{n}$ | $1-x^{2}$ | $-x$ | $n^{2}$ |
| $U_{n}$ | $1-x^{2}$ | $-3 x$ | $n(n+2)$ |
| $P_{n}$ | $1-x^{2}$ | $-2 x$ | $n(n+1)$ |
| $L_{n}$ | $x$ | $1-x$ | $n$ |
| $L_{n}^{(m)}$ | $x$ | $m+1-x$ | $n$ |
| $C_{n}^{(m)}$ | $1-x^{2}$ | $-(2 m+1) x$ | $n(n+2 m)$ |
| $\mathcal{P}_{n}^{(m, \nu)}$ | $1-x^{2}$ | $\nu-m-(m+\nu+2) x$ | $n(n+1+m+\nu)$ |
| $B_{n}$ | $x^{2}$ | $2(x+1)$ | $-n(n+1)$ |

Integrability conditions and solutions of equation (24) can be obtained applying Kovacic algorithm (Case 1 of the algorithm). Besides, they can also be achieved via Kimura and Martinet-Ramis Theorems and the parabolic cylinder equation (see [22, 29, 41]).

As a consequence, from this table, we obtain the following result.
Proposition 3.1. We consider $Q, L$ and $\lambda$ as in the previous table. Then, for any $\mu \neq 0$, the planar quadratic polynomial vector field

$$
\begin{align*}
& \frac{d v}{d t}=\frac{\lambda}{\mu} Q+\left(Q^{\prime}-L\right) v+\mu v^{2}  \tag{25}\\
& \frac{d x}{d t}=Q
\end{align*}
$$

has invariant algebraic curves of the form $\mu v+\frac{Q(x) \mathcal{P}_{n}^{\prime}(x)}{\mathcal{P}_{n}(x)}=0$ where $\mathcal{P}_{n}(x)$ is any orthogonal polynomial associated to $\lambda, Q, L$ and $n \in \mathbb{N}$.
Proof. The associated (Riccati) foliation of system (25) is written into the form

$$
\begin{equation*}
\frac{d v}{d x}=\frac{\lambda}{\mu}+\frac{Q^{\prime}-L}{Q} v+\frac{\mu}{Q} v^{2} . \tag{26}
\end{equation*}
$$

Performing the change $\tilde{v}=\mu v$ equation (26) becomes

$$
\frac{d \tilde{v}}{d x}=\lambda+\frac{Q^{\prime}-L}{Q} \tilde{v}+\frac{1}{Q} \tilde{v}^{2},
$$

and can be led into the form (24) through the transformation $\tilde{v}=-Q y^{\prime} / y$ (see transformation $\mathcal{B}$ of Proposition 2.1). For fixed $\lambda, Q$ and $L$, let $\mathcal{P}_{n}$ be any orthogonal polynomial solution of equation (24). Then $v_{n}=-\left(Q \mathcal{P}_{n}^{\prime}\right) /\left(\mu \mathcal{P}_{n}\right)$ is a rational solution of equation (26). Therefore the curve $v-v_{n}=0$ is an invariant curve of the vector field (25) for any $n \in \mathbb{N}$.

According to Proposition 2.1 equation (26) can be reduced to the form $\xi^{\prime \prime}=\rho \xi$ with

$$
\rho=\frac{1}{2}\left(\frac{L}{Q}\right)^{\prime}-\frac{\lambda}{Q}+\left(\frac{L}{2 Q}\right)^{2}
$$

and $\xi=\mathcal{P}_{n} e^{\int \frac{L}{2 Q}}$ is a solution for any $n \in \mathbb{N}$. Additionally, we notice that we fall in Case 1 of Kovacic algorithm.
3.3. Liénard equation. Let us consider first order differential equations whose associated foliation can be expressed into the Liénard form

$$
\begin{equation*}
y y^{\prime}=f(x) y+g(x) \tag{27}
\end{equation*}
$$

with $y=y(x)$ and rational functions $f(x)$ and $g(x)$. We are concerned with the problem of obtaining criteria on $f(x)$ and $g(x)$ such that equation (27) can be led into a Riccati equation.

This is a difficult problem and, as far as the authors know, only partial answers have been given to it. In what follows we give some examples of such results coming from the handbook [46] and papers [11, 12].

A first example is given by the 5 -parametric family (1.3.3.11 in [46])

$$
\begin{equation*}
y y^{\prime}=\left(a(2 m+k) x^{2 k}+b(2 m-k) x^{m-k-1}\right) y-\left(a^{2} m x^{4 k}+c x^{2 k}+b^{2} m\right) x^{2 m-2 k-1} \tag{28}
\end{equation*}
$$

$a, b, c, m, k$ being complex parameters. In order that (28) come from a polynomial vector field we have that $m$ and $k$ must be rational numbers (see [46]). The change $w=x^{k}$, $y=x^{m}\left(z+a x^{k}+b x^{-k}\right)$ leads (28) into the Riccati form

$$
\begin{equation*}
\left(-m z^{2}+2 a b m-c\right) w^{\prime}(z)=b k+k z w+a k w^{2} \tag{29}
\end{equation*}
$$

whose associated second order linear equation is a Riemann equation. More precisely, by Lemma 2.1, it can be written as a Legendre equation

$$
\begin{equation*}
\left(1-t^{2}\right) u^{\prime \prime}(t)-2 t u^{\prime}(t)+\left(\nu(\nu+1)-\frac{\mu^{2}}{1-t^{2}}\right) u(t)=0 \tag{30}
\end{equation*}
$$

with

$$
\mu=-\frac{m+k}{2 m}
$$

and $\nu$ being a solution of

$$
\nu^{2}+\nu+\frac{m^{2}-k^{2}}{4 m^{2}}-\frac{a b k^{2}}{m c-2 a b m^{2}}=0 .
$$

The difference of exponents in (30) is $\mu, \mu$ and $2 \nu-1$ and, therefore, we are under the hypotheses of Kimura's Theorem (see Appendix B.1.1).

Proposition 3.2. Legendre equation (30) is integrable if and only if, either
(1) $\mu \pm \nu \in \mathbb{Z}$ or $\nu \in \mathbb{Z}$, or
(2) $\pm \mu, \pm \mu, \pm(2 \nu+1)$ belong to one of the following seven families

| Case | $\mu \in$ | $\nu \in$ | $\mu+\nu \in$ |
| :---: | :---: | :---: | :---: |
| $(a)$ | $\mathbb{Z}+\frac{1}{2}$ | $\mathbb{C}$ |  |
| $(b)$ | $\mathbb{Z} \pm \frac{1}{3}$ | $\frac{1}{2} \mathbb{Z} \pm \frac{1}{3}$ | $\mathbb{Z}+\frac{1}{6}$ |
| $(c)$ | $\mathbb{Z} \pm \frac{2}{5}$ | $\frac{1}{2} \mathbb{Z} \pm \frac{1}{5}$ | $\mathbb{Z}+\frac{1}{0}$ |
| $(d)$ | $\mathbb{Z} \pm \frac{1}{3}$ | $\frac{1}{2} \mathbb{Z} \pm \frac{2}{5}$ | $\mathbb{Z}+\frac{1}{10}$ |
| $(e)$ | $\mathbb{Z} \pm \frac{1}{5}$ | $\frac{1}{2} \mathbb{Z} \pm \frac{2}{5}$ | $\mathbb{Z}+\frac{1}{10}$ |
| $(f)$ | $\mathbb{Z} \pm \frac{2}{5}$ | $\frac{1}{2} \mathbb{Z} \pm \frac{1}{3}$ | $\mathbb{Z}+\frac{1}{6}$ |

Proof. In Kimura's Theorem, the difference of exponents $\mu, \mu$ and $2 \nu+1$ correspond to the possibilities listed above. Indeed, they are cases $(i),(i i .1),(i i .3),(i i .11),(i i .12),(i i .13)$ and (ii.15) of Kimura's table (see Appendix B.1.1). For the case ( $i$ ) we have

- $\mu+\mu+2 \nu+1 \in 2 \mathbb{Z}+1 \Rightarrow \mu+\nu \in \mathbb{Z}$,
- $-\mu+\mu+2 \nu+1 \in 2 \mathbb{Z}+1 \Rightarrow \nu \in \mathbb{Z}$,
- $\mu-\mu+2 \nu+1 \in 2 \mathbb{Z}+1 \Rightarrow \nu \in \mathbb{Z}$,
- $\mu+\mu-2 \nu-1 \in 2 \mathbb{Z}+1 \Rightarrow \mu-\nu \in \mathbb{Z}$.

The rest of the cases can be proven in a similar way.
(ii.1) We see that $\pm \mu \in \frac{1}{2}+\mathbb{Z}$ and $\pm(2 \nu+1) \in \mathbb{C}$ and therefore $\mu \in \mathbb{Z}+\frac{1}{2}$ and $\nu \in \mathbb{C}$.
(ii.3) We consider that $\pm \mu=l+\frac{2}{3}$, $\pm \mu=m+\frac{1}{3}, \pm(2 \nu+1)=q+\frac{1}{3}$, being $l, m, q \in \mathbb{Z}$. Take for instance $\mu \in \mathbb{Z} \pm \frac{1}{3}$ and $\nu \in \frac{1}{2} \mathbb{Z} \pm \frac{1}{3}$. Furthermore, $l+m+q$ must be even and in consequence we obtain $\mu+\nu \in \mathbb{Z}+\frac{1}{6}$.
(ii.11) We have that $\pm \mu=l+\frac{2}{5}, \pm \mu=m+\frac{2}{5}, \pm(2 \nu+1)=q+\frac{2}{5}$ with $l, m, q \in \mathbb{Z}$. For example we consider $\mu \in \mathbb{Z} \pm \frac{2}{5}$ and $\nu \in \frac{1}{2} \mathbb{Z} \pm \frac{1}{5}$. Moreover $l+m+q$ must be even and therefore we have that $\mu+\nu \in \mathbb{Z}+\frac{1}{10}$.
(ii.12) Now we consider $\pm \mu=l+\frac{2}{3}, \pm \mu=m+\frac{1}{3}, \pm(2 \nu+1)=q+\frac{1}{5}$, being $l, m, q \in \mathbb{Z}$. We take for instance $\mu \in \mathbb{Z} \pm \frac{1}{3}$ and $\nu \in \frac{1}{2} \mathbb{Z} \pm \frac{2}{5}$. Additionally $l+m+q$ must be even and so $\mu+\nu \in \mathbb{Z}+\frac{1}{10}$.
(ii.13) Let be $\pm \mu=l+\frac{2}{3}, \pm \mu=m+\frac{1}{3}, \pm(2 \nu+1)=q+\frac{1}{5}$ with $l, m, q \in \mathbb{Z}$. Consider for example $\mu \in \mathbb{Z} \pm \frac{1}{5}$ and $\nu \in \frac{1}{2} \mathbb{Z} \pm \frac{2}{5}$. Furthermore, $l+m+q$ must be even and therefore $\mu+\nu \in \mathbb{Z}+\frac{1}{10}$.
(ii.15) Consider $\pm \mu=l+\frac{3}{5}, \pm \mu=m+\frac{2}{5}, \pm(2 \nu+1)=q+\frac{1}{3}$, being $l, m, q \in \mathbb{Z}$, take for instance $\mu \in \mathbb{Z} \pm \frac{2}{5}$ and $\nu \in \frac{1}{2} \mathbb{Z} \pm \frac{1}{3}$. Moreover, $l+m+q$ must be even and in consequence $\mu+\nu \in \mathbb{Z}+\frac{1}{6}$.
Finally, observe that the difference exponents $\mu, \mu$ and $2 \nu+1$ do not satisfy the conditions (ii.2), (ii.4), (ii.5), (ii.6), (ii.7), (ii.8), (ii.9), (ii.10) and (ii.14)

Now we deal with an example from $[11,12]$. We consider the equation

$$
\begin{equation*}
\frac{d x}{d w}=A(x)+B(x) w \tag{31}
\end{equation*}
$$

where $x$ is the dependent variable and $w$ is the independent one. Now, for $B \not \equiv 0$, by means of the change of variable

$$
w=y-\frac{A}{B}
$$

and changing $(w, x) \longrightarrow(x, y)$ (that is, we consider now $x$ as the independent variable and $y$ as the dependent one) equation (31) is transformed into the Liénard equation

$$
\begin{equation*}
y \frac{d y}{d x}=\frac{1}{B}+\frac{d}{d x}\left(\frac{A}{B}\right) y, \tag{32}
\end{equation*}
$$

for any functions $A$ and $B \not \equiv 0$. In particular, for

$$
A=A(x)=a+b x+c x^{2}, \quad B=B(x)=\alpha+\beta x+\gamma x^{2}
$$

equation (31) falls into the Riccatti form

$$
\frac{d x}{d w}=(a+\alpha w)+(b+\beta w) x+(c+\gamma w) x^{2}
$$

By Proposition 2.1, applying the transformation $\mathcal{T}$ it follows the reduced Riccatti equation (13) and through the transformation $\mathcal{R}$ the normalized second order differential equation $\xi^{\prime \prime}=\rho(x) \xi$ with

$$
\begin{equation*}
\rho(x)=\frac{\beta^{2}-4 \alpha \gamma}{4} x^{2}-\frac{2 a \gamma+2 \alpha c-b \beta}{2} x-\frac{4 a c-b^{2}}{4}+\frac{b \gamma-\beta c}{2(\gamma x+c)}+\frac{3 \gamma^{2}}{4(\gamma x+c)^{2}} . \tag{33}
\end{equation*}
$$

Introducing the change $\tau=\gamma x+c$ we get

$$
\begin{equation*}
\xi^{\prime \prime}=\rho(\tau) \xi \tag{34}
\end{equation*}
$$

where

$$
\begin{aligned}
\rho(\tau)= & \frac{\beta^{2}-4 \alpha \gamma}{4 \gamma^{2}} \tau^{2}-\frac{2 a \gamma^{2}-2 \alpha c \gamma-\beta b \gamma+\beta^{2} c}{2 \gamma^{2}} \tau+\frac{b^{2} \gamma^{2}-2 b \beta c \gamma+\beta^{2} c^{2}}{4 \gamma^{2}}+ \\
& \frac{b \gamma-\beta c}{2 \tau}+\frac{3 \gamma^{2}}{4 \tau^{2}}
\end{aligned}
$$

and performing $z=\sqrt[4]{\frac{\beta^{2}-4 \alpha \gamma}{4 \gamma^{2}}} \tau$ we arrive to

$$
\begin{equation*}
\psi^{\prime \prime}=\phi(z) \psi, \quad \phi(z)=z^{2}+\delta_{1} z+\frac{\delta_{1}^{2}}{4}-\delta_{2}+\frac{\delta_{3}}{2 z}+\frac{\delta_{0}^{2}-1}{4 z^{2}} \tag{35}
\end{equation*}
$$

with $\delta_{i}$ being algebraic functions in $a, b, c, \alpha, \beta$ and $\gamma$. This equation (35) is exactly the biconfluent Heun equation whose integrability is analysed in Appendix B.2.

Assuming $\beta=\gamma=0$ we obtain a Liénard equation which is transformable into a reduced second order differential equation with $r \in \mathbb{C}[x]$ and $\operatorname{deg}(r)=1$. This means that the equation is not integrable (see [30] and section 3.4). As a particular case, we have a Liénard equation that can be reduced to the Riccati equation given in [?, equation 1.3.2.1],

$$
\begin{equation*}
2 y y^{\prime}=(a x+b) y+1 \tag{36}
\end{equation*}
$$

Moreover, under some restrictions over the parameters of the biconfluent Heun equation one can obtain the Whittaker equation. For instance, the Liénard equation

$$
y \frac{d y}{d x}=\left(\lambda+b_{11} x\right) y+\left(\varepsilon x+b_{20} x^{2}\right),
$$

falls into a Whittaker equation for some special values of the parameters.
Remark 3.3. We would like to stress the following facts.
(a) It is well-known that via the change $z(x)=\int f(x) d x$ (with inverse $x=x(z)$ ), the Liénard equation (27) can be led into the equation

$$
\begin{equation*}
y y^{\prime}(z)=y+h(z), \tag{37}
\end{equation*}
$$

with

$$
h(z):=\frac{g(x(z))}{f(x(z))} .
$$

In a similar way, the change $z(x)=\int g(x) d x$ reduces (27) to

$$
\begin{equation*}
y y^{\prime}(z)=h(z) y+1, \tag{38}
\end{equation*}
$$

with

$$
h(z):=\frac{f(x(z))}{g(x(z))} .
$$

However, in general, these transformations do not preserve the differential field of the coefficients. This is why we do not use them to reduce equation (27).
(b) Sometimes equation (27) is called Abel equation of second kind since through the change $y=1 / w$ it is reduced to an Abel equation

$$
w^{\prime}=-f(x) w^{2}-g(x) w^{3}
$$

In Cheb-Terrab [11, 12] the authors consider some families of Abel equations which are reducible to Riccati equations.
3.4. Other families. Here we consider some special cases of Riccati equations.
a) Polynomial Riccati equations

The Riccati equation

$$
w^{\prime}=r(x)-w^{2}, \quad r(x) \in \mathbb{C}[x]
$$

has been studied by several authors (see, for instance [2,57,53]). The Galois group of its associated second order linear differential equation is connected and can be either $\operatorname{SL}(2, \mathbb{C})$ or the Borel group [3, 2]. In the first case the tangent field associated to its Riccati equation has no invariant curves. In the second one, there is no rational first integral for its vector field. As an example, the reduced form for the triconfluent Heun equation is of this type and is given by

$$
\begin{equation*}
\xi^{\prime \prime}=\rho(x) \xi, \quad \rho(x)=\frac{9 x^{4}}{4}+\frac{3}{2} \delta_{2} x^{2}-\delta_{1} x+\frac{\delta_{2}^{2}}{4}-\delta_{0} . \tag{39}
\end{equation*}
$$

b) Equations with finite Galois group

Consider the polynomial Riccati vector field

$$
\begin{align*}
& \dot{x}=-x^{2}(x-1)^{2}  \tag{40}\\
& \dot{y}=a(x-1)^{2}+b x^{2}+c x(x-1)+x^{2}(x-1)^{2} v^{2} .
\end{align*}
$$

Its foliation was studied in [55] and is related to the following differential equation

$$
\begin{equation*}
\xi^{\prime \prime}=-\left(\frac{1-\lambda^{2}}{4 x^{2}}+\frac{1-\mu^{2}}{4(x-1)^{2}}+\frac{1-\nu^{2}+\lambda^{2}+\mu^{2}}{4 x(x-1)}\right) \xi \tag{41}
\end{equation*}
$$

where $\lambda, \mu$ and $\nu$ are the differences of the exponents at 0,1 and $\infty$ of the hypergeometric equation, see also Appendix B. When the equation (41) is integrable their solutions are Legendre functions. For $(\lambda, \mu, \nu)=\left(\frac{1}{2}, \frac{1}{2}, \frac{1}{n}\right)$, the solutions of equation (41) are given by

$$
\begin{aligned}
& \xi_{1}=\sqrt[4]{x^{2}-x}\left(2 x-1+2 \sqrt{x^{2}-x}\right)^{\frac{1}{2} \sqrt{\frac{1-2 n^{2}}{n^{2}}}} \\
& \xi_{2}=\sqrt[4]{x^{2}-x}\left(2 x-1+2 \sqrt{x^{2}-x}\right)^{-\frac{1}{2} \sqrt{\frac{1-2 n^{2}}{n^{2}}}}
\end{aligned}
$$

and its differential Galois group is the Dihedral Group $D_{n}$ (case 2 of Kovacic algorithm). One can also obtain the tetrahedral, octahedral and icosahedral groups for $(\lambda, \mu, \nu)=\left(\frac{1}{3}, \frac{1}{2}, \frac{1}{3}\right)$, $(\lambda, \mu, \nu)=\left(\frac{1}{3}, \frac{1}{2}, \frac{1}{4}\right)$ and $(\lambda, \mu, \nu)=\left(\frac{1}{3}, \frac{1}{2}, \frac{1}{5}\right)$, respectively (case 3 of Kovacic algorithm).

Related to (41) we have the equation (see [55])

$$
\begin{equation*}
y^{\prime \prime}+\frac{7 x-4}{6 x(x-1)} y^{\prime}-\frac{36 \nu^{2}-1}{144 x(x-1)} y=0, \tag{42}
\end{equation*}
$$

where the differential Galois group can be tetrahedral, octahedral and icosahedral depending on the values of $\nu$ (indeed, for $1 / 3,1 / 4$ and $1 / 5$ respectively). Even though explicit solutions of equations (41) and (42) are difficult to get in general, it can be proved the existence of rational first integrals for the associated vector field for values of the parameters $\nu \in\left\{\frac{1}{3}, \frac{1}{4}, \frac{1}{5}\right\}$. For instance, solutions of equation (42) for $\nu=\frac{1}{3}$ (tetrahedral group) are

$$
\begin{aligned}
& y_{1}=\sqrt[4]{-\sqrt[3]{x}+1} \sqrt[8]{\frac{2 \sqrt{\sqrt[3]{x^{2}}+\sqrt[3]{x}+1}-\sqrt{3} \sqrt[3]{x}-\sqrt{3}}{-\sqrt{3} \sqrt[3]{x}-\sqrt{3}-2 \sqrt{\sqrt[3]{x^{2}}+\sqrt[3]{x}+1}}} \\
& y_{2}=\sqrt[4]{-\sqrt[3]{x}+1} \sqrt[8]{\frac{2 \sqrt{\sqrt[3]{x^{2}}+\sqrt[3]{x}+1}-\sqrt{3} \sqrt[3]{x}-\sqrt{3}}{-\sqrt{3} \sqrt[3]{x}-\sqrt{3}-2 \sqrt{\sqrt[3]{x^{2}}+\sqrt[3]{x}+1}}}
\end{aligned}
$$

In the case $\nu=\frac{1}{4}$ (octahedral group) we obtain the solutions $y_{1}=\frac{f_{1} f_{2}}{f_{3}}$ and $y_{2}=\frac{f_{1}}{f_{2} f_{3}}$ where

$$
\begin{aligned}
& f_{1}=\sqrt[6]{\sqrt[3]{\frac{\sqrt{1-x}+1}{\sqrt{1-x}-1}}-1} \\
& f_{2}=\sqrt[8]{-\frac{-2 \sqrt{\left(\frac{\sqrt{1-x}+1}{\sqrt{1-x}-1}\right)^{2 / 3}+\sqrt[3]{\frac{\sqrt{1-x}+1}{\sqrt{1-x-1}}}+\sqrt{3}} \sqrt[3]{\sqrt{\sqrt{1-x}+1}}+\sqrt{3}}{\left.\sqrt{3} \sqrt[3]{\frac{\sqrt{1-x}+1}{\sqrt{1-x-1}}+\sqrt{3}+2} \sqrt{\left(\frac{\sqrt{1-x}+1}{\sqrt{1-x}-1}\right.}\right)^{2 / 3}+\sqrt[3]{\frac{\sqrt{1-x}+1}{\sqrt{1-x-1}}}}} \\
& f_{3}=\left(\left(\frac{\sqrt{1-x}+1}{\sqrt{1-x}-1}\right)^{2 / 3}+\sqrt[3]{\frac{\sqrt{1-x}+1}{\sqrt{1-x}-1}}+1\right)^{\frac{1}{12}}
\end{aligned}
$$

Finally, setting $\nu=\frac{1}{5}$ (icosahedral group), we get the solutions $y_{1}=e^{\int \omega_{1}}$ and $y_{2}=e^{\int \omega_{2}}$, where $\omega_{1}=\omega_{1}(x)$ and $\omega_{2}=\omega_{2}(x)$ are roots of the polynomial $\sum_{k=0}^{12} a_{k}(x) \omega^{k}$ and

$$
\begin{aligned}
a_{0}(x)= & 102400-11264 x-11 x^{2} \\
a_{1}(x)= & -3686400 x+3679200 x^{2}+7200 x^{3} \\
a_{2}(x)= & 479001600 x^{2}-476863200 x^{3}-2138400 x^{4} \\
a_{3}(x)= & 30412800000 x^{2}-60445440000 x^{3}+29652480000 x^{4}+380160000 x^{5} \\
a_{4}(x)= & -821145600000 x^{3}+1597384800000 x^{4}-731332800000 x^{5}-44906400000 x^{6} \\
a_{5}(x)= & -3695155200000 x^{4}+11085465600000 x^{5}+3695155200000 x^{7}- \\
& 11085465600000 x^{6}
\end{aligned}
$$

$$
\begin{aligned}
a_{6}(x)= & -492687360000000 x^{4}++1693612800000000 x^{5}-2124714240000000 x^{6}+ \\
& 1139339520000000 x^{7}-215550720000000 x^{8} \\
a_{7}(x)= & 8868372480000000 x^{5}-35473489920000000 x^{6}+53210234880000000 x^{7}- \\
& 35473489920000000 x^{8}+8868372480000000 x^{9}, \\
a_{8}(x)= & -249422976000000000 x^{6}+997691904000000000 x^{7}-1496537856000000000 x^{8}+ \\
& 997691904000000000 x^{9}-249422976000000000 x^{10} \\
a_{9}(x)= & -4434186240000000000 x^{6}+22170931200000000000 x^{7}- \\
& 44341862400000000000 x^{8}+44341862400000000000 x^{9}- \\
& 22170931200000000000 x^{10}+4434186240000000000 x^{11} \\
a_{10}(x)= & 39907676160000000000 x^{7}-199538380800000000000 x^{8}+ \\
& 399076761600000000000 x^{9}-399076761600000000000 x^{10}+ \\
& 199538380800000000000 x^{11}-39907676160000000000 x^{12} \\
a_{11}(x)= & 0 \\
a_{12}(x)= & 2176782336000000000000 x^{8}-13060694016000000000000 x^{9}+ \\
& 32651735040000000000000 x^{10}-43535646720000000000000 x^{11}+ \\
& 32651735040000000000000 x^{12}-13060694016000000000000 x^{13}+ \\
& 2176782336000000000000 x^{14} .
\end{aligned}
$$

c) Lamé families.

Let us consider in the Lamé equation

$$
\begin{equation*}
\frac{d^{2} y}{d x^{2}}+\frac{f^{\prime}(x)}{2 f(x)} \frac{d y}{d x}-\frac{n(n+1) x+B}{f(x)} y=0 \tag{43}
\end{equation*}
$$

the case where $f(x)=4 x^{3}-g_{2} x-g_{3}$ and parameters $n, B, g_{2}, g_{3}$ such that the discriminant of $f$, namely $27 g_{3}^{2}-g_{2}^{3}$, is non-zero (see Appendix B). Performing the transformation

$$
\begin{equation*}
v=-\frac{y^{\prime}}{c y}, \tag{44}
\end{equation*}
$$

with $c=c(x)$ being any non-zero arbitrary rational function, it follows the family of Riccati equations associated to the Lamé equation,

$$
\begin{equation*}
v^{\prime}=-\frac{n(n+1) x+B}{c f}-\left(\frac{f^{\prime}(x)}{2 f(x)}+\frac{c^{\prime}}{c}\right) v+c v^{2} . \tag{45}
\end{equation*}
$$

For the Lamé case, (i.1) of B. 3 with $B=B_{i}$, and from Remark B. 5 it turns out the existence of a polynomial solution of equation (45) despite of the general solution of this equation is not algebraic. Furthermore, it is clear that for a fixed $n$ we have associated Riccati equations (45) with arbitrary degree, because the non-zero rational function $c=c(x)$ at the transformation above is arbitrary. The Lamé functions correspond here to algebraic solutions of (45) and moving $n$ along the integers, we obtain algebraic invariants curves

$$
v+\frac{E^{\prime}(x)}{c(x) E(x)}=0
$$

of unbounded degree.
Hence, we have prove the following result.
Theorem 3.4. For any fixed degree in the Riccati family (45) associated to the Lamé equation (43) there exist invariant algebraic curves of any tangent field $X$ of the corresponding foliation with unbounded degree. Furthermore, the first integral of $X$ is not rational.

## Appendix

## Appendix A. Kovacic Algorithm

This algorithm is devoted to solve the reduced linear differential equation (RLDE) $\xi^{\prime \prime}=\rho \xi$ and is based on the algebraic subgroups of $\operatorname{SL}(2, \mathbb{C})$. For more details see [30]. Although improvements for this algorithm are given in $[22,52]$, we follow the original version given by Kovacic in [30].

Theorem A.1. Let $G$ be an algebraic subgroup of $\operatorname{SL}(2, \mathbb{C})$. Then one of the following four cases can occur.
(1) $G$ is triangularizable.
(2) $G$ is conjugate to a subgroup of infinite dihedral group (also called meta-abelian group) and case 1 does not hold.
(3) Up to conjugation $G$ is one of the following finite groups: Tetrahedral group, Octahedral group or Icosahedral group, and cases 1 and 2 do not hold.
(4) $G=\operatorname{SL}(2, \mathbb{C})$.

Each case in Kovacic algorithm is related to each one of the algebraic subgroups of $\mathrm{SL}(2, \mathbb{C})$ and its associated Riccatti equation

$$
\theta^{\prime}=r-\theta^{2}=(\sqrt{r}-\theta)(\sqrt{r}+\theta), \quad \theta=\frac{\xi^{\prime}}{\xi}, r=\rho .
$$

According to Theorem A. 1 we obtain four cases. Only for cases 1,2 and 3 one can solve the differential equation RLDE and in case 4 one has not Liouvillian solutions for it. Kovacic algorithm can possibly provide one solution $\left(\xi_{1}\right)$, so the second one $\left(\xi_{2}\right)$ can be got through

$$
\begin{equation*}
\xi_{2}=\xi_{1} \int \frac{d x}{\xi_{1}^{2}} \tag{46}
\end{equation*}
$$

Notations. For the RLDE given by

$$
\xi^{\prime \prime}=\rho \xi=r \xi, \quad r=\frac{s}{t}, \quad s, t \in \mathbb{C}[x]
$$

we use:
(1) Denote by $\Gamma^{\prime}$ be the set of (finite) poles of $r, \Gamma^{\prime}=\{c \in \mathbb{C}: t(c)=0\}$.
(2) Denote by $\Gamma=\Gamma^{\prime} \cup\{\infty\}$.
(3) By the order of $r$ at $c \in \Gamma^{\prime}, \circ\left(r_{c}\right)$, we mean the multiplicity of $c$ as a pole of $r$.
(4) By the order of $r$ at $\infty, \circ\left(r_{\infty}\right)$, we mean the order of $\infty$ as a zero of $r$. That is $\circ\left(r_{\infty}\right)=\operatorname{deg}(t)-\operatorname{deg}(s)$.
A.1. The four cases. Case 1. In this case $[\sqrt{r}]_{c}$ and $[\sqrt{r}]_{\infty}$ means the Laurent series of $\sqrt{r}$ at $c$ and the Laurent series of $\sqrt{r}$ at $\infty$ respectively. Furthermore, we define $\varepsilon(p)$ as follows: if $p \in \Gamma$, then $\varepsilon(p) \in\{+,-\}$. Finally, the complex numbers $\alpha_{c}^{+}, \alpha_{c}^{-}, \alpha_{\infty}^{+}, \alpha_{\infty}^{-}$will be defined in the first step. If the differential equation has not poles it only can fall in this case.

Step 1. Search for each $c \in \Gamma^{\prime}$ and for $\infty$ the corresponding situation as follows:
$\left(c_{0}\right)$ If $\circ\left(r_{c}\right)=0$, then

$$
[\sqrt{r}]_{c}=0, \quad \alpha_{c}^{ \pm}=0
$$

$\left(c_{1}\right)$ If $\circ\left(r_{c}\right)=1$, then

$$
[\sqrt{r}]_{c}=0, \quad \alpha_{c}^{ \pm}=1
$$

$\left(c_{2}\right)$ If $\circ\left(r_{c}\right)=2$, and

$$
\begin{aligned}
& r=\cdots+b(x-c)^{-2}+\cdots, \quad \text { then } \\
& {[\sqrt{r}]_{c}=0, \quad \alpha_{c}^{ \pm}=\frac{1 \pm \sqrt{1+4 b}}{2}}
\end{aligned}
$$

( $c_{3}$ ) If $\circ\left(r_{c}\right)=2 v \geq 4$, and

$$
\begin{gathered}
r=\left(a(x-c)^{-v}+\ldots+d(x-c)^{-2}\right)^{2}+b(x-c)^{-(v+1)}+\cdots, \quad \text { then } \\
\quad[\sqrt{r}]_{c}=a(x-c)^{-v}+\ldots+d(x-c)^{-2}, \quad \alpha_{c}^{ \pm}=\frac{1}{2}\left( \pm \frac{b}{a}+v\right) .
\end{gathered}
$$

$\left(\infty_{1}\right)$ If $\circ\left(r_{\infty}\right)>2$, then

$$
[\sqrt{r}]_{\infty}=0, \quad \alpha_{\infty}^{+}=0, \quad \alpha_{\infty}^{-}=1
$$

$\left(\infty_{2}\right)$ If $\circ\left(r_{\infty}\right)=2$, and $r=\cdots+b x^{2}+\cdots$, then

$$
[\sqrt{r}]_{\infty}=0, \quad \alpha_{\infty}^{ \pm}=\frac{1 \pm \sqrt{1+4 b}}{2}
$$

$\left(\infty_{3}\right)$ If $\circ\left(r_{\infty}\right)=-2 v \leq 0$, and

$$
\begin{aligned}
r & =\left(a x^{v}+\ldots+d\right)^{2}+b x^{v-1}+\cdots, \quad \text { then } \\
{[\sqrt{r}]_{\infty} } & =a x^{v}+\ldots+d, \quad \text { and } \quad \alpha_{\infty}^{ \pm}=\frac{1}{2}\left( \pm \frac{b}{a}-v\right) .
\end{aligned}
$$

Step 2. Find $D \neq \emptyset$ defined by

$$
D=\left\{m \in \mathbb{Z}_{+}: m=\alpha_{\infty}^{\varepsilon(\infty)}-\sum_{c \in \Gamma^{\prime}} \alpha_{c}^{\varepsilon(c)}, \forall(\varepsilon(p))_{p \in \Gamma}\right\}
$$

If $D=\emptyset$, then we should start with the case 2 . Now, if $\# D>0$, then for each $m \in D$ we search $\omega \in \mathbb{C}(x)$ such that

$$
\omega=\varepsilon(\infty)[\sqrt{r}]_{\infty}+\sum_{c \in \Gamma^{\prime}}\left(\varepsilon(c)[\sqrt{r}]_{c}+\alpha_{c}^{\varepsilon(c)}(x-c)^{-1}\right) .
$$

Step 3. For each $m \in D$, search for a monic polynomial $P_{m}$ of degree $m$ with

$$
P_{m}^{\prime \prime}+2 \omega P_{m}^{\prime}+\left(\omega^{\prime}+\omega^{2}-r\right) P_{m}=0
$$

If one successes then $\xi_{1}=P_{m} e^{\int \omega}$ is a solution of the differential equation RLDE. Else, Case 1 cannot hold.

Case 2. Search for each $c \in \Gamma^{\prime}$ and for $\infty$ the corresponding situation as follows.
Step 1. Search for each $c \in \Gamma^{\prime}$ and $\infty$ the sets $E_{c} \neq \emptyset$ and $E_{\infty} \neq \emptyset$. For each $c \in \Gamma^{\prime}$ and for $\infty$ we define $E_{c} \subset \mathbb{Z}$ and $E_{\infty} \subset \mathbb{Z}$ as follows:
$\left(c_{1}\right)$ If $\circ\left(r_{c}\right)=1$, then $E_{c}=\{4\}$.
$\left(c_{2}\right)$ If $\circ\left(r_{c}\right)=2$, and $r=\cdots+b(x-c)^{-2}+\cdots$, then

$$
E_{c}=\{2+k \sqrt{1+4 b}: k=0, \pm 2\}
$$

( $c_{3}$ ) If $\circ\left(r_{c}\right)=v>2$, then $E_{c}=\{v\}$.
$\left(\infty_{1}\right)$ If $\circ\left(r_{\infty}\right)>2$, then $E_{\infty}=\{0,2,4\}$.
$\left(\infty_{2}\right)$ If $\circ\left(r_{\infty}\right)=2$, and $r=\cdots+b x^{2}+\cdots$, then

$$
E_{\infty}=\{2+k \sqrt{1+4 b}: k=0, \pm 2\} .
$$

$\left(\infty_{3}\right)$ If $\circ\left(r_{\infty}\right)=v<2$, then $E_{\infty}=\{v\}$.
Step 2. Find $D \neq \emptyset$ defined by

$$
D=\left\{m \in \mathbb{Z}_{+}: \quad m=\frac{1}{2}\left(e_{\infty}-\sum_{c \in \Gamma^{\prime}} e_{c}\right), \forall e_{p} \in E_{p}, p \in \Gamma\right\} .
$$

If $D=\emptyset$, then we should start the case 3 . Now, if $\# D>0$, then for each $m \in D$ we search a rational function $\theta$ defined by

$$
\theta=\frac{1}{2} \sum_{c \in \Gamma^{\prime}} \frac{e_{c}}{x-c} .
$$

Step 3. For each $m \in D$, search a monic polynomial $P_{m}$ of degree $m$, such that

$$
P_{m}^{\prime \prime \prime}+3 \theta P_{m}^{\prime \prime}+\left(3 \theta^{\prime}+3 \theta^{2}-4 r\right) P_{m}^{\prime}+\left(\theta^{\prime \prime}+3 \theta \theta^{\prime}+\theta^{3}-4 r \theta-2 r^{\prime}\right) P_{m}=0 .
$$

If $P_{m}$ does not exist, then Case 2 cannot hold. If such a polynomial is found, set $\phi=\theta+P^{\prime} / P$ and let $\omega$ be a solution of

$$
\omega^{2}+\phi \omega+\frac{1}{2}\left(\phi^{\prime}+\phi^{2}-2 r\right)=0 .
$$

Then $\xi_{1}=e^{\int \omega}$ is a solution of the differential equation RLDE.
Case 3. Search for each $c \in \Gamma^{\prime}$ and for $\infty$ the corresponding situation as follows:
Step 1. Search for each $c \in \Gamma^{\prime}$ and $\infty$ the sets $E_{c} \neq \emptyset$ and $E_{\infty} \neq \emptyset$. For each $c \in \Gamma^{\prime}$ and for $\infty$ we define $E_{c} \subset \mathbb{Z}$ and $E_{\infty} \subset \mathbb{Z}$ as follows:
$\left(c_{1}\right)$ If $\circ\left(r_{c}\right)=1$, then $E_{c}=\{12\}$.
$\left(c_{2}\right)$ If $\circ\left(r_{c}\right)=2$, and $r=\cdots+b(x-c)^{-2}+\cdots$, then

$$
E_{c}=\{6+k \sqrt{1+4 b}: \quad k=0, \pm 1, \pm 2, \pm 3, \pm 4, \pm 5, \pm 6\} .
$$

$(\infty)$ If $\circ\left(r_{\infty}\right)=v \geq 2$, and $r=\cdots+b x^{2}+\cdots$, then

$$
E_{\infty}=\left\{6+\frac{12 k}{n} \sqrt{1+4 b}: \quad k=0, \pm 1, \pm 2, \pm 3, \pm 4, \pm 5, \pm 6\right\}, \quad n \in\{4,6,12\} .
$$

Step 2. Find $D \neq \emptyset$ defined by

$$
D=\left\{m \in \mathbb{Z}_{+}: \quad m=\frac{n}{12}\left(e_{\infty}-\sum_{c \in \Gamma^{\prime}} e_{c}\right), \forall e_{p} \in E_{p}, p \in \Gamma\right\}
$$

In this case we start with $n=4$ to obtain the solution, afterwards $n=6$ and finally $n=12$. If $D=\emptyset$, then the differential equation has not Liouvillian solution because it falls in the case 4 . Now, if $\# D>0$, then for each $m \in D$ with its respective $n$, search a rational function

$$
\theta=\frac{n}{12} \sum_{c \in \Gamma^{\prime}} \frac{e_{c}}{x-c},
$$

and a polynomial $S$ defined as

$$
S=\prod_{c \in \Gamma^{\prime}}(x-c)
$$

Step 3. Search for each $m \in D$, with its respective $n$, a monic polynomial $P_{m}=P$ of degree $m$, such that its coefficients can be determined recursively by

$$
\begin{gathered}
P_{-1}=0, \quad P_{n}=-P \\
P_{i-1}=-S P_{i}^{\prime}-\left((n-i) S^{\prime}-S \theta\right) P_{i}-(n-i)(i+1) S^{2} r P_{i+1},
\end{gathered}
$$

where $i \in\{0,1 \ldots, n-1, n\}$. If $P$ does not exist, then the differential equation has not Liouvillian solution because it falls in Case 4. Now, if $P$ exists search $\omega$ such that

$$
\sum_{i=0}^{n} \frac{S^{i} P}{(n-i)!} \omega^{i}=0
$$

then a solution of the differential equation the RLDE is given by

$$
\xi=e^{\int \omega},
$$

where $\omega$ is solution of the previous polynomial of degree $n$.

## Appendix B. Some Special Functions

## B.1. Hypergeometric families.

B.1.1. Kimura's Theorem. The hypergeometric (or Riemann) equation is the more general second order linear differential equation over the Riemann sphere with three regular singular singularities. If we place the singularities at $x=0,1, \infty$ it is given by

$$
\begin{align*}
\frac{d^{2} y}{d x^{2}}+\left(\frac{1-\alpha-\alpha^{\prime}}{x}\right. & \left.+\frac{1-\gamma-\gamma^{\prime}}{x-1}\right) \frac{d y}{d x}  \tag{47}\\
& +\left(\frac{\alpha \alpha^{\prime}}{x^{2}}+\frac{\gamma \gamma^{\prime}}{(x-1)^{2}}+\frac{\beta \beta^{\prime}-\alpha \alpha^{\prime} \gamma \gamma^{\prime}}{x(x-1)}\right) y=0,
\end{align*}
$$

where $\left(\alpha, \alpha^{\prime}\right),\left(\gamma, \gamma^{\prime}\right),\left(\beta, \beta^{\prime}\right)$ are the exponents at the singular points and must satisfy the Fuchs relation $\alpha+\alpha^{\prime}+\gamma+\gamma^{\prime}+\beta+\beta^{\prime}=1$.

Now, we will briefly describe Kimura's Theorem that provides necessary and sufficient conditions for the integrability of the hypergeometric equation. Let be $\lambda=\alpha-\alpha^{\prime}, \mu=\beta-\beta^{\prime}$ and $\nu=\gamma-\gamma^{\prime}$.

Theorem B. 1 (Kimura, [29]). The hypergeometric equation (47) is integrable if and only if either
(i) At least one of the four numbers $\lambda+\mu+\nu,-\lambda+\mu+\nu, \lambda-\mu+\nu, \lambda+\mu-\nu$ is an odd integer, or
(ii) The numbers $\lambda$ or $-\lambda, \mu$ or $-\mu$ and $\nu$ or $-\nu$ belong (in an arbitrary order) to some of the following fifteen families

| 1 | $1 / 2+l$ | $1 / 2+m$ | arbitrary complex number |  |
| :---: | :---: | :---: | :---: | :---: |
| 2 | $1 / 2+l$ | $1 / 3+m$ | $1 / 3+q$ |  |
| 3 | $2 / 3+l$ | $1 / 3+m$ | $1 / 3+q$ | $l+m+q$ even |
| 4 | $1 / 2+l$ | $1 / 3+m$ | $1 / 4+q$ |  |
| 5 | $2 / 3+l$ | $1 / 4+m$ | $1 / 4+q$ | $l+m+q$ even |
| 6 | $1 / 2+l$ | $1 / 3+m$ | $1 / 5+q$ |  |
| 7 | $2 / 5+l$ | $1 / 3+m$ | $1 / 3+q$ | $l+m+q$ even |
| 8 | $2 / 3+l$ | $1 / 5+m$ | $1 / 5+q$ | $l+m+q$ even |
| 9 | $1 / 2+l$ | $2 / 5+m$ | $1 / 5+q$ | $l+m+q$ even |
| 10 | $3 / 5+l$ | $1 / 3+m$ | $1 / 5+q$ | $l+m+q$ even |
| 11 | $2 / 5+l$ | $2 / 5+m$ | $2 / 5+q$ | $l+m+q$ even |
| 12 | $2 / 3+l$ | $1 / 3+m$ | $1 / 5+q$ | $l+m+q$ even |
| 13 | $4 / 5+l$ | $1 / 5+m$ | $1 / 5+q$ | $l+m+q$ even |
| 14 | $1 / 2+l$ | $2 / 5+m$ | $1 / 3+q$ | $l+m+q$ even |
| 15 | $3 / 5+l$ | $2 / 5+m$ | $1 / 3+q$ | $l+m+q$ even |

Here $l, m$ and $q$ are integers.
B.1.2. Confluent hypergeometric. The confluent Hypergeometric equation is a degenerate form of the Hypergeometric differential equation where two of the three regular singularities merge into an irregular singularity. The following are two classical forms.

- Kummer's form

$$
\begin{equation*}
y^{\prime \prime}+\frac{c-x}{x} y^{\prime}-\frac{a}{x} y=0, \quad a, c \in \mathbb{C} \tag{48}
\end{equation*}
$$

- Whittaker's form

$$
\begin{equation*}
y^{\prime \prime}=\left(\frac{1}{4}-\frac{\kappa}{x}+\frac{4 \mu^{2}-1}{4 x^{2}}\right) y \tag{49}
\end{equation*}
$$

where the parameters of the two equations are linked by $\kappa=\frac{c}{2}-a$ and $\mu=\frac{c}{2}-\frac{1}{2}$. Furthermore, using the expression (10), we can see that the Whittaker's equation is the reduced form of the Kummer's equation (48). The Galoisian structure of these equations has been deeply studied in [41, 22].
Theorem B. 2 (Martinet \& Ramis, [41]). The Whittaker's differential equation (49) is integrable if and only if either, $\kappa+\mu \in \frac{1}{2}+\mathbb{N}$, or $\kappa-\mu \in \frac{1}{2}+\mathbb{N}$, or $-\kappa+\mu \in \frac{1}{2}+\mathbb{N}$, or $-\kappa-\mu \in \frac{1}{2}+\mathbb{N}$.

The Bessel's equation is a particular case of the confluent Hypergeometric equation and is given by

$$
\begin{equation*}
y^{\prime \prime}+\frac{1}{x} y^{\prime}+\frac{x^{2}-n^{2}}{x^{2}} y=0 . \tag{50}
\end{equation*}
$$

Under a suitable transformation, the reduced form of the Bessel's equation is a particular case of the Whittaker's equation (49).
Corollary B.3. The Bessel's differential equation (50) is integrable if and only if $n \in \frac{1}{2}+\mathbb{Z}$. B.2. Heun's families. The Heun's equation is the generic differential equation with four regular singular points at $0,1, c$ and $\infty$. In its reduced form, the Heun's equation is $y^{\prime \prime}=$ $r(x) y$, where

$$
\begin{gather*}
r(x)=\frac{A}{x}+\frac{B}{x-1}+\frac{C}{x-c}+\frac{D}{x^{2}}+\frac{E}{(x-1)^{2}}+\frac{F}{(x-c)^{2}},  \tag{51}\\
A=-\frac{\alpha \beta}{2}-\frac{\alpha \gamma}{2 c}+\frac{\delta \eta h}{c}, \quad B=\frac{\alpha \beta}{2}-\frac{\beta \gamma}{2(c-1)}-\frac{\delta \eta(h-1)}{c-1}, \\
C=\frac{\alpha \gamma}{2 c}+\frac{\beta \gamma}{2(c-1)}-\frac{\delta \eta(c-h)}{c(c-1)}, \quad D=\frac{\alpha}{2}\left(\frac{\alpha}{2}-1\right), \quad E=\frac{\beta}{2}\left(\frac{\beta}{2}-1\right), \\
F=\frac{\gamma}{2}\left(\frac{\gamma}{2}-1\right), \quad \text { with } \quad \alpha+\beta+\gamma-\delta-\eta=1 .
\end{gather*}
$$

To our purposes we write the determinant $\Pi_{d+1}(a, b, u, v, \xi, w)$ as in [22]:

$$
\left|\begin{array}{ccccccc}
w & u & 0 & 0 & 0 & \ldots & 0 \\
d \xi w+1 & v & 2(u+b) & 0 & 0 & \cdots & 0 \\
0 & (d-1) \xi & w+2(v+a) & 3(u+2 b) & 0 & \cdots & 0 \\
0 & 0 & (d-2) \xi & w+3(v+2 a) & 4(u+3 b) & \cdots & 0 \\
\vdots & & & & & \ldots & \\
0 & \ldots & & \ldots & 2 \xi & w+(d-1)(v+(d-2) a) & d(u+(d-1) b) \\
0 & \cdots & & \cdots & 0 & \xi & w+d(v+(d-1) a)
\end{array}\right|
$$

B.2.1. Biconfluent Heun. The equation

$$
\begin{equation*}
\xi^{\prime \prime}=\left(x^{2}+\delta_{1} x+\frac{\delta_{1}^{2}}{4}-\delta_{2}+\frac{\delta_{3}}{2 x}+\frac{\delta_{0}^{2}-1}{4 x^{2}}\right) \xi \tag{52}
\end{equation*}
$$

is the well known biconfluent Heun equation which has been deeply analyzed by Duval and Loday-Richaud in [22, p. 236].
Theorem B.4. [22]. The biconfluent Heun equation (52) has Liouvillian solutions if and only if it falls in Case 1 of Kovacic algorithm and one of the following conditions is fulfilled.
(1) $\delta_{0}^{2}=1, \delta_{3}=0$ and $\delta_{2} \in 2 \mathbb{Z}+1$.
(2) $\delta_{0}^{2}=1, \delta_{3} \neq 0$ and $\delta_{2} \in 2 \mathbb{Z}^{*}+1$ with $\left|\delta_{2}\right| \geq 3$, and if $\varepsilon=\operatorname{sign} \delta_{2}$, then

$$
\Pi_{\left(\left|\delta_{2}\right|-1\right) / 2}\left(0,1,2, \varepsilon \delta_{1},-2 \varepsilon, \varepsilon \delta_{1}-\frac{\delta_{3}}{2}\right)=0 .
$$

(3) $\delta_{0} \neq \pm 1, \pm \delta_{0} \pm \delta_{2} \in 2 \mathbb{Z}^{*}$ and if $\varepsilon_{0}, \varepsilon_{\infty} \in\{ \pm 1\}$ are such that $\varepsilon_{\infty} \delta_{2}-\varepsilon_{0} \delta_{0}=2 d^{*} \in 2 \mathbb{N}^{*}$ then

$$
\Pi_{d^{*}}\left(0,1,1+\varepsilon_{0} \delta_{0}, \varepsilon_{\infty} \delta_{1},-2 \varepsilon_{\infty}, \frac{1}{2}\left(\varepsilon_{\infty} \delta_{1}\left(1+\varepsilon_{0} \delta_{0}\right)-\delta_{3}\right)\right)=0 .
$$

B.3. Lamé equation. The algebraic form of the Lamé Equation is [47, 56]

$$
\begin{equation*}
\frac{d^{2} y}{d x^{2}}+\frac{f^{\prime}(x)}{2 f(x)} \frac{d y}{d x}-\frac{n(n+1) x+B}{f(x)} y=0, \tag{53}
\end{equation*}
$$

where $f(x)=4 x^{3}-g_{2} x-g_{3}$, with $n, B, g_{2}$ and $g_{3}$ parameters such that the discriminant of $f, \Delta=27 g_{3}^{2}-g_{2}^{3}$, is non-zero and, therefore, it has no multiple roots. This equation is a Fuchsian differential equation with four singular points over the Riemann sphere: the roots $e_{1}, e_{2}, e_{3}$ of $f$ and the point of the infinity.

The mutually-exclusive known cases of solutions in closed form of the Lamé equation (53) are the following.
(i) The Lamé-Hermite case (see $[47,56]$ ). We have $n \in \mathbf{N}$ and arbitrary parameters $B, g_{2}$ and $g_{3}$.
(ii) The Brioschi-Halphen-Crawford case (see [23, 47]). We have $n \in \mathbb{N}$ such that $m:=n+$ $\frac{1}{2} \in \mathbb{N}$ and parameters $B, g_{2}$ and $g_{3}$ satisfying an algebraic condition $Q_{m}\left(g_{2} / 4, g_{3} / 4, B\right)=$ 0 , where $Q_{m} \in \mathbb{Z}\left[g_{2} / 4, g_{3} / 4, B\right]$ is a polynomial of degree $m$ in $B$, known as the Brioschi determinant.
(iii) The Baldassarri case (see [4]). One asks $n$ to satisfy that $n+\frac{1}{2} \in \frac{1}{3} \mathbb{Z} \cup \frac{1}{4} \mathbb{Z} \cup \frac{1}{5} \mathbb{Z}-\mathbb{Z}$ besides some additional (involved) algebraic restrictions on the other parameters.
It is possible to prove that the only integrable cases of the Lamé equation are cases (i)-(iii) above, see [42]. In (ii) and (iii) the general solution of (53) is algebraic and the Galois group is finite. Case $(i)$ splits in the following two subcases [47, 56].
(i.1) The Lamé case. For a fixed integer $n$, this equation admits a solution (called Lamé function) of the form

$$
\begin{equation*}
E(x)=\prod_{i=1}^{3}\left(x-e_{i}\right)^{k_{i}} P_{m}(x), \tag{54}
\end{equation*}
$$

being $P_{m}$ a monic polynomial of degree $m=n / 2-\left(k_{1}+k_{2}+k_{3}\right)$ and $k_{i} \in\left\{0, \frac{1}{2}\right\}$, $i=1,2,3$. Since $m \in \mathbb{N}$, eight different possibilities regarding $n$ appear: If $n$ is even, we have $k_{1}=k_{2}=k_{3}=0$ or just one zero $k_{i}$; if $n$ is odd, we could have all non-zero $k_{i}$ 's or combinations with exactly one non-zero $k_{i}$. All these possibilities give rise to classes of Lamé functions. Concerning parameter $B$, it must be one of the $m+1$ different roots $B_{1}, \ldots, B_{m+1}$ of certain irreducible polynomial of degree $m+1$, with all its roots real and simple [23]. Furthermore, the numbers $B_{i}$ are reals.
(i.2) The Hermite case. Here we are not in case (i.1) and $n$ is an arbitrary natural number. We are also fall in case 1 of Kovacic algorithm, but with a diagonal Galois group.

Remark B.5. We notice that the polynomial $P_{m}$ in (i.1) satisfies a second order linear differential equation similar to the one that appears in the first case of Kovacic algorithm. In fact it is possible to obtain the above passing into normal form and applying Kovacic algorithm. Therefore the second linear independent solution is not algebraic and the associated Riccati equation has no rational first integral.

## Acknowledgements

Remark 1.3 comes from a question posed by J. Llibre to the third author some years ago. The authors are indebted to J.A. Weil for several discussion and suggestions. The last author also thanks J. Llibre for many interesting discussions on several subjects of this work.

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[^0]:    2010 Mathematics Subject Classification. Primary: 12H05. Secondary: 32S65.
    Key words and phrases. Differential Galois Theory, Darboux theory of Integrability, Poincaré problem, Rational first integral, Integrating factor, Riccati equation, Liénard Equation, Liouvillian solution.

    All the authors are partially supported by the MICIIN/FEDER grant number MTM2009-06973 and by the Generalitat de Catalunya grant number 2009SGR859. C.P. is additionally partial supported by the MICIIN/FEDER grant MTM2008- 03437.

