# Exactly solvable Schrödinger operators 

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Abstract

We systematically describe and classify 1-dimensional Schrödinger equations that can be solved in terms of hypergeometric type functions. Beside the well-known families, we explicitly describe 2 new classes of exactly solvable Schrödinger equations that can be reduced to the Hermite equation.

## 1 Introduction

Exactly solvable 1-dimensional Schrödinger equations play an important role in quantum physics. The best known are the harmonic oscillator and the radial equation for the hydrogen atom, which are covered in every course of quantum mechanics. A number of other examples were discovered in the 30's of the last century. They include (known also under other names) the Eckart potential [Eck], the trigonometric and hyperbolic Pöschl-Teller potencjal $\overline{\mathrm{PT}}$ and the Rosen potential MR . All these examples can be reduced to the hypergeometric equation (see eg. WW). One should also mention the Morse potential that leads to the confluent equation. Problems involving these potential are often used in classes on quantum mechanics, see eg. the well-known problem book of Flügge [Flü. A number of techniques have been developed to study their properties, such as the factorization method [HI] and the method of the superpotential (see for instance CKS).

In later years J. N. Ginocchio Gin discovered that these examples can be generalized to a larger class of potentials equivalent to the hypergeometric or confluent equation. Later this class was extended by G. A. Natanzon [Nat and further generalized by R. Milson Mil]. These classes, besides the hitherto known potentials, are not very practical in applications, since they are not given by explicit expressions.

There exist even more general classes of potentials that can be called exactly solvable, found by A. Khare and U. Sukhatme [KS]. They are, however, expressed in exotic special functions.

In our paper we will systematically describe basic classes of exactly solvable Schrödinger equations. Most of the material of our paper is scattered in the literature, notably in [Nat and Mil] (see also [Bos]). Our treatment is somewhat more systematic than what one can find in the literature. For instance, we distinguish between the complex classification and the real classification. It also seems that for the first time we explicitly describe 2 new classes of exactly solvable Schrödinger equations, which can be reduced to the Hermite equation (see Subsects 5.2 and 5.3).

By a (stationary 1-dimensional) Schrödinger equation we will mean an equation of the form

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V(r)-E\right) \phi(r)=0 \tag{1.1}
\end{equation*}
$$

This equation can be interpreted as the eigenvalue problem for the operator

$$
\begin{equation*}
H:=-\partial_{r}^{2}+V(r) \tag{1.2}
\end{equation*}
$$

An operator of the form (1.2) will be called a (1-dimensional) Schrödinger operator. $V(r)$ will be called a potential and the parameter $E$ an energy.
(1.2) can be interpreted as an operator in a number of ways. If $V$ is a holomorphic function of the complex variable on some open $\Omega \subset \mathbb{C}$, then $r$ can be interpreted as a complex variable and (1.2) can be viewed as an operator on holomorphic functions. In this case it is natural to allow $E$ to be a complex parameter. The corresponding eigenvalue equation (1.1) will be then called a Schrödinger equation of complex variable. Let us note that, if $a \neq 0, b$ are complex constants, a complex affine transformation $r \mapsto a r+b$ preserves the class of Schrödinger equations of the complex variable.

One can also interpret $r$ as a real variable in some open $I \subset \mathbb{R}$. The operator (1.2) is then viewed as acting on functions on $I$. Understood in this way, (1.1) will be called a Schrödinger equation of the real variable. Of special interest is then the case of real $V$.

Our paper is organized as follows. First, we will briefly dicuss some general facts related to 2nd order linear differential equations. In particular, we will describe basic ingredients of the so-called BoseNatanzon method, which permits to obtain a class of Schrödinger equations equivalent to a chosen linear equation. It is worthful to point out, that it applies to more general situations than those described in the literature ( Nat, Mil]). We will formulate a criterion which determines when the Bose-Natanzon method can be used.

Next, we will focus on hypergeometric type equations. Recall that equations of the form

$$
\begin{equation*}
\left(\sigma(z) \partial_{z}^{2}+\tau(z) \partial_{z}+\eta\right) f(z)=0 \tag{1.3}
\end{equation*}
$$

where $\sigma(z), \tau(z)$ are polynomials with

$$
\operatorname{deg} \sigma \leq 2, \operatorname{deg} \tau \leq 1
$$

and $\eta$ is a number, are occasionally called hypergeometric type equations (see A.F. Nikiforov and V.B. Uvarov (NU). Solutions of (1.3) are called hypergeometric type functions. Hypergeometric type functions are very well understood and include the hypergeomeric function, the confluent function, Hermite, Laguerre and Jacobi polynomials, etc. Traditionally, one classifies hypergeometric type equations into several distinct classes invariant under complex affine transformations of the form $z \mapsto a z+b(a \neq 0)$. In each class one chooses a simple representative, to which the whole class can be reduced. Such representatives are the hypergeometric equation, the confluent equation, etc.

In our analysis, we will concentrate on 3 specific equations:
(1) the hypergeometric equation

$$
\begin{equation*}
\left(z(1-z) \partial_{z}^{2}+(c-(a+b+1) z) \partial_{z}-a b\right) f(z)=0 \tag{1.4}
\end{equation*}
$$

(2) the generalized confluent equation

$$
\begin{equation*}
\left(z \partial_{z}^{2}+(c-\gamma z) \partial_{z}-a\right) f(z)=0 \tag{1.5}
\end{equation*}
$$

(essentially just the rescaled confluent equation);
(3) the generalized harmonic oscillator

$$
\begin{equation*}
\left(-\partial_{z}^{2}+\theta^{2} z^{2}+\rho z+\lambda\right) f(z)=0 \tag{1.6}
\end{equation*}
$$

(essentially just the translated harmonic oscillator equation).
Note that all hypergeometric type equations with a nonzero $\sigma$ are related with one of the above equations by elementary operations (multiplication of $f(z)$ and of both sides of the equation by a function) and by a change of variables not depending on the parameters $a, b, c$ (respectively $a, c, \gamma$ or $\theta, \rho, \lambda$ ).
(1) If $\sigma$ is 2 nd order and has 2 distinct roots, an affine transformation reduces (1.3) to (1.4).
(2) If $\sigma$ is 2 nd order and has only one root, then a transformation involving $z \mapsto z^{-1}$ reduces (1.3) to (1.5) with $\gamma \neq 0$.
(3) If $\sigma$ is 1 st order, then an affine transformation reduces (1.3) to (1.5).
(4) If $\sigma$ is 0 th order and nonzero, then (1.3) is equivalent to (1.6) with $\theta \neq 0$.

Obviously, (1.4) and (1.5) are themselves equations of the hypergeometric type. (1.6) is not, but can be reduced to a hypergeometric type equation. This reduction is depends on whether $\theta \neq 0$ or $\theta=0$. The former case leads to a hypergeometric type equation with a constant $\sigma$, as mentioned above. In the case $\theta=0,(1.6)$ is the Airy equation, which can be reduced to a special case of the hypergeometric type equation by a transformation involving $z \mapsto z^{3}$.

Equations (1.4), (1.5) and (1.6) are very well understood and have well known solutions. We will describe classes of solvable potentials reducible to one of the equations (1.4), (1.5) or (1.6). We will describe both the complex and the real classification of such potentials.

All solvable potentials that we consider are holomorphic on the complex plane, apart from some isolated singularities. Therefore, it is natural to consider first a classification of Schrödinger equations of the complex variable. Obviously, we can always use a complex affine transformation to put the equation in a convenient form. One can also move a complex constant from the potential to the energy.

Suppose we are given a Schrödinger equations of the complex variable. Embedding a real line (or its subset) in the complex plane in an arbitrary way yields a Schrödinger equation of the real variable. In order to obtain an interesting and useful real classification we will demand in addition that the potential is real. Thus, we will look for an open, possibly infinite interval $I \subset \mathbb{R}$, on which the potential is real for a selected class of parameters. It turns out, that in all cases that we consider $I$, can be chosen so that it ends either at $-\infty,+\infty$ or at a singularity of the potential, and no singularities lie inside $I$. Such an interval $I$ will be called a natural real domain. Again, a real affine transformation and moving a real constant from the potential to the energy can be used to put the equation in a convenient form.

Clearly, if $I$ is a natural real domain, then the operator $H=-\partial_{r}^{2}+V(r)$ is hermitian on $C_{\mathrm{c}}^{\infty}(I)$. It is then natural to ask about self-adjoint extensions of $H$. We will not discuss this question here. We plan to consider it in later papers.

In Equations (1.4), (1.5) and (1.6) we have 3 arbitrary parameters. Therefore, in all the solvable cases that we will describe, the potential depends on 2 parameters, since the third parameter is responsible for the energy $E$.

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## 2 Second order homogeneous differential equations

The main object of our paper are ordinary homogeneous 2 nd order linear differential equations, that is equations of the form

$$
\begin{equation*}
\left(a(r) \partial_{r}^{2}+b(r) \partial_{r}+c(r)\right) \phi(r)=0 \tag{2.7}
\end{equation*}
$$

It will be convenient to treat (2.7) as the problem of finding the kernel of the operator

$$
\begin{equation*}
\mathcal{A}\left(r, \partial_{r}\right):=a(r) \partial_{r}^{2}+b(r) \partial_{r}+c(r) \tag{2.8}
\end{equation*}
$$

We will then say that the equation (2.7) is given by the operator (2.8).
We will treat $r$ either as a complex or a real variable. In the complex case we will usually assume that the coefficients are analytic.

In this section we describe some general facts related to the theory of equations of the form (2.7) and their reduction to the Schrödinger equation.

### 2.1 The Bose invariant of a second order differential equation

Suppose that we are given a second order differential equation (2.7). Then we can always eliminate the first order term as follows. We divide from the left by $a(r)$, set

$$
h(r):=\exp \left(\int^{r} \frac{b(t)}{2 a(t)} \mathrm{d} t\right)
$$

We check that

$$
h(r) \mathcal{A}\left(r, \partial_{r}\right) h(r)^{-1}=\partial_{r}^{2}+I(r)
$$

where, using the notation $a^{\prime}:=\frac{\mathrm{d}}{\mathrm{d} r} a(r)$,

$$
\begin{equation*}
I=\frac{4 a c-2 a b^{\prime}+2 b a^{\prime}-b^{2}}{4 a^{2}} \tag{2.9}
\end{equation*}
$$

Clearly, $\check{\phi}(r):=h(r) \phi(r)$ solves the equation

$$
\begin{equation*}
\left(\partial_{r}^{2}+I(r)\right) \check{\phi}(r)=0 \tag{2.10}
\end{equation*}
$$

(2.10) will be called the canonical form of the equation (2.7). (2.9) will be called the Bose invariant of (2.7) Bos. (The name "Bose invariant" is used eg. by Milson Mil, however the object itself was clearly known before Bose).

It is easy to check that for any functions $f, g$ the equation

$$
\begin{equation*}
f(r) \mathcal{A}\left(r, \partial_{r}\right) g(r) \phi(r)=0 \tag{2.11}
\end{equation*}
$$

has the same Bose invariant as (2.7). Conversely, if an equation has the same Bose invariant as (2.7), it equals (2.11) for some $f, g$.

### 2.2 Schwarz derivative

The Schwarz derivative of $r \mapsto r(y)$ appears naturally in the context of 2 nd order ordinary differential equations and was known to XIX century mathematicians. It was mentioned already by Kummer in 1836 Kum. The univaersally adopted name Schwarz derivative was introduced by A. Cayley, who derived its basic properties in 1880 Cay. For more historical comments together with numerous applications we refer the reader to a book by Osgood Osg.

The Schwarz derivative is defined as

$$
\begin{equation*}
\{r, y\}=\left(\frac{r^{\prime \prime}(y)}{r^{\prime}(y)}\right)^{\prime}-\frac{1}{2}\left(\frac{r^{\prime \prime}(y)}{r^{\prime}(y)}\right)^{2} \tag{2.12}
\end{equation*}
$$

Here are some basic properties of the Schwarz derivative:
Proposition 2.1 1. Let $y \mapsto s(y), s \mapsto r(s)$ be two functions. Then

$$
\begin{equation*}
\{r, y\}=\{r, s\} s^{\prime}(y)^{2}+\{s, y\} \tag{2.13}
\end{equation*}
$$

2. Let $r \mapsto y(r)$ be the inverse of $y \mapsto r(y)$. Then

$$
\begin{equation*}
\{y, r\}=-\left[r^{\prime}(y)\right]^{-2}\{r, y\} \tag{2.14}
\end{equation*}
$$

3. $\{r, y\} \equiv 0 \Leftrightarrow r(y)=\frac{a y+b}{c y+d}, a d-b c \neq 0$

We will need the following fact:
Proposition 2.2 Suppose that $z^{\prime}=z^{p}(1-z)^{q}$. Then

$$
\begin{equation*}
\{z, r\}=\frac{1}{2} z^{2 p-2}(1-z)^{2 q-2}\left(\left(p^{2}-2 p\right)(1-z)^{2}+\left(q^{2}-2 q\right) z^{2}-2 p q z(1-z)\right) \tag{2.15}
\end{equation*}
$$

### 2.3 The Liouville transformation

Consider a 2 nd order equation in the canonical form, that is

$$
\begin{equation*}
\left(\partial_{r}^{2}+I(r)\right) \phi(r)=0 . \tag{2.16}
\end{equation*}
$$

Let us make the transformation $r=r(y)$ in this equation. We obtain

$$
\begin{equation*}
\frac{1}{\left[r^{\prime}(y)\right]^{2}} \partial_{r}^{2} \phi(r(y))-\frac{r^{\prime \prime}(y)}{\left[r^{\prime}(y)\right]^{3}} \partial_{r} \phi(r(y))+I(r(y)) \phi(r(y))=0 \tag{2.17}
\end{equation*}
$$

The resulting equation we transform again to its canonical form:

$$
\begin{equation*}
\left(\partial_{y}^{2}+J(y)\right) \psi(y)=0, \tag{2.18}
\end{equation*}
$$

where

$$
\begin{gather*}
\psi(y)=\frac{1}{\left[r^{\prime}(y)\right]^{2}} \phi(r(y)),  \tag{2.19}\\
J(y)=\left[r^{\prime}(y)\right]^{2} I(r(y))+\frac{1}{2}\{r, y\}, \tag{2.20}
\end{gather*}
$$

The above procedure is called the Liouville transformation of (2.16) by the change of variables $r=r(y)$ Lio. One can check using (2.17), that the composition of two Liouville transformations, first by the change of variables $r=r(y)$, then $y=y(z)$, is the Liouville transformation by the change of variables $r=r(y(z))$.

### 2.4 Natanzon's problem

Let $\mathcal{A}\left(a_{1}, \ldots, a_{k} ; z, \partial_{z}\right)$ be a 2 nd order differential operator depending on $k$ parameters, such that its Bose invariant can be written as

$$
\begin{equation*}
I\left(b_{1}, \ldots, b_{n} ; z\right)=b_{1} I_{1}(z)+\ldots+b_{n} I_{n}(z) \tag{2.21}
\end{equation*}
$$

for some (at least one) $n \leq k$ linearly independent functions $I_{1}, \ldots, I_{n}$ (independent of $a_{1}, \ldots, a_{k}$ ), and some $n$ numbers $b_{i}\left(a_{1}, \ldots, a_{k}\right)$. Then one can solve the following problem (which we will call Natanzon's problem for $\left.\mathcal{A}\left(z, \partial_{z}\right)\right)$ :

Find all potentials $V(r)$ such that the 1-dimensional stationary Schrödinger equation

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V(r)-E\right) \phi(E, r)=0 \tag{2.22}
\end{equation*}
$$

can be transformed to the equation given by $\mathcal{A}\left(a_{1}, \ldots, a_{k} ; z, \partial_{z}\right)$ for some $a_{1}(E), \ldots, a_{k}(E)$. We allow the following operations:
(1) multiplication of both sides of the equation by some $f(E, r)$;
(2) substitution of $\check{\phi}(E, r):=g(E, r) \phi(E, r)$ for some $g(E, r)$;
(3) change of coordinates $r \mapsto z$ independent of $E$.

This is the problem solved by Natanzon in the case of the hypergeometric equation and by Milson in the general case of hypergeometric type equations. In the following, we recall the construction they used, in a slightly generalized form.

Let us consider an arbitrary equation of the form (2.22). Obviously, its Bose invariant equals $E-V(r)$. Clearly, the transformations (1) and (2) allow us to transform $\mathcal{A}\left(a_{1}, \ldots, a_{k} ; z, \partial_{z}\right)$ to its canonical form

$$
\begin{equation*}
-\partial_{z}^{2}-I\left(b_{1}, \ldots, b_{n} ; z\right) . \tag{2.23}
\end{equation*}
$$

Thus (2.22) can be transformed to $\mathcal{A}\left(a_{1}, \ldots, a_{k} ; z, \partial_{z}\right)$ if

$$
\begin{equation*}
I\left(b_{1}, \ldots, b_{n} ; z\right)=\left(r^{\prime}(z)\right)^{2}\left(E-V(r(z))+\frac{1}{2}\{r, z\}\right. \tag{2.24}
\end{equation*}
$$

that is, if the two Bose invariants are related by a Liouville transformation. Using (2.21) we rewrite (2.24) as

$$
\begin{equation*}
b_{1}(E) I_{1}(z)+\ldots+b_{n}(E) I_{n}(z)=\left(r^{\prime}(z)\right)^{2}\left(E-V(r(z))+\frac{1}{2}\{r, z\}\right. \tag{2.25}
\end{equation*}
$$

By assumption, $b_{i}$ depend on $E$ and $I_{i}(z)$ are independent of $E$. Thus, the dependence of $b_{i}$ on $E$ is linear. Therefore, by transforming linearly $\left(b_{1}, \ldots, b_{n}\right)$ into $\left(\tilde{b}_{1}, \ldots, \tilde{b}_{n-1}, E\right)$, we can assume that

$$
I\left(b_{1}, \ldots, b_{n} ; z\right)=\tilde{b}_{1} \tilde{I}_{1}(z)+\ldots+\tilde{b}_{n-1} \tilde{I}_{n-1}(z)+E \tilde{I}_{n}(z)
$$

for some functions $\tilde{I}_{i}(z)$, which are linear combinations of $I_{i}(z)$. Then we can write

$$
\begin{equation*}
E \tilde{I}_{n}(z)+\tilde{b}_{1} \tilde{I}_{1}(z)+\ldots+\tilde{b}_{n-1} \tilde{I}_{n-1}(z)=\left(r^{\prime}(z)\right)^{2}\left(E-V(r(z))+\frac{1}{2}\{r, z\}\right. \tag{2.26}
\end{equation*}
$$

Therefore, Natanzon's problem is solved by the following pair of equations

$$
\begin{align*}
V\left(\tilde{b}_{1}, \ldots, \tilde{b}_{n-1}, r(z)\right) & :=\left(r^{\prime}(z)\right)^{-2}\left(\tilde{b}_{1} \tilde{I}_{1}(z)+\ldots+\tilde{b}_{n-1} \tilde{I}_{n-1}(z)-\frac{1}{2}\{r, z\}\right) \\
\left(r^{\prime}(z)\right)^{2} & =\tilde{I}_{n}(z) \tag{2.27}
\end{align*}
$$

We can rewrite (2.27) in an equivalent way as

$$
\begin{align*}
V\left(\tilde{b}_{1}, \ldots, \tilde{b}_{n-1}, r\right) & :=\left(z^{\prime}(r)\right)^{2}\left(\tilde{b}_{1} \tilde{I}_{1}(z(r))+\ldots+\tilde{b}_{n-1} \tilde{I}_{n-1}(z(r))\right)+\frac{1}{2}\{z, r\} \\
\left(z^{\prime}(r)\right)^{-2} & =\tilde{I}_{n}(z(r)) \tag{2.28}
\end{align*}
$$

If desired, we can renormalize $E$ by substracting a constant.
Note that in practical applications the value of the above method depends on whether it is possible to invert the relation $r \mapsto z$ expressing it in terms of standard functions. In the next sections, we will describe such solutions of the Natanzon problem for the three classes of equations described in the introduction.

## 3 Schrödinger operators reducible to the hypergeometric equation

### 3.1 Hypergeometric equation

The hypergeometric equation is given by the operator

$$
\begin{equation*}
\mathcal{F}\left(a, b ; c ; z, \partial_{z}\right):=z(1-z) \partial_{z}^{2}+(c-(a+b+1) z) \partial_{z}-a b \tag{3.29}
\end{equation*}
$$

where $a, b, c$ are arbitrary complex parameters. Its Bose invariant equals

$$
{ }_{2} I_{1}(a, b ; c ; z)=\frac{\left(1-a^{2}-b^{2}-2 a b\right) z^{2}+2(-2 a b+a c+b c-c) z+2 c-c^{2}}{4 z^{2}(1-z)^{2}}
$$

It satisfies condition (2.21) for $n=3$. The functions $I_{i}(z), i=1,2,3$ can be taken to be $\frac{1}{4 z^{2}(1-z)}$ $\frac{1}{4 z(1-z)^{2}}$ and $\frac{1}{4 z(1-z)}$. It is natural to demand that $\left(r^{\prime}(z)\right)^{2}=\tilde{I}_{3}(z)$ is a function of the form $z^{2 p}(1-z)^{2 q}$. It is easy to see that this gives the following possibilities:

$$
\begin{align*}
& \text { (1) } \frac{1}{4 z(1-z)}, \quad \text { (2) } \frac{1}{4 z^{2}(1-z)}, \quad \text { (3) } \frac{1}{4 z(1-z)^{2}} \\
& \text { (4) } \frac{1}{4 z^{2}(1-z)^{2}}, \quad \text { (5) } \frac{1}{4 z^{2}}, \quad(6) \frac{1}{4(1-z)^{2}}
\end{align*}
$$

We will see that each of these ansatzes leads to exactly solvable potentials considered in the literature. Note also that the formula for the Schwarz derivative (2.15) will be handy.

### 3.2 The Riemann equation

It is well known that it is useful to consider the hypergeometric equation as a special case of the so-called Riemann equation. The Riemann equation is defined as the class of equations on the Riemann sphere $\overline{\mathbb{C}}$ having 3 regular singular points. The following theorem summarizes the basic theory of these equations:
Theorem 3.1 (1) Suppose that we are given a $2 n d$ order differential equation on the Riemann sphere having 3 singular points $z_{1}, z_{2}, z_{3}$, all of them regular singular points with the following indices

$$
\begin{array}{ll}
z_{1}: & \rho_{1}, \tilde{\rho}_{1} \\
z_{2}: & \rho_{2}, \tilde{\rho}_{2} \\
z_{3}: & \rho_{3}, \tilde{\rho}_{3}
\end{array}
$$

Then the following condition is satisfied:

$$
\begin{equation*}
\rho_{1}+\tilde{\rho}_{1}+\rho_{2}+\tilde{\rho}_{2}+\rho_{3}+\tilde{\rho}_{3}=1 \tag{3.31}
\end{equation*}
$$

Such an equation, normalized to have coefficient 1 at the 2nd derivative, is always equal to

$$
\mathcal{P}\left(\begin{array}{llll}
z_{1} & z_{2} & z_{3} &  \tag{3.32}\\
\rho_{1} & \rho_{2} & \rho_{3} & z, \partial_{z} \\
\tilde{\rho}_{1} & \tilde{\rho}_{2} & \tilde{\rho}_{3} &
\end{array}\right) \phi(z)=0
$$

where

$$
\begin{aligned}
& \mathcal{P}\left(\begin{array}{llll}
z_{1} & z_{2} & z_{3} \\
\rho_{1} & \rho_{2} & \rho_{3} & z, \partial_{z} \\
\tilde{\rho}_{1} & \tilde{\rho}_{2} & \tilde{\rho}_{3}
\end{array}\right):=\partial_{z}^{2}-\left(\frac{\rho_{1}+\tilde{\rho}_{1}-1}{z-z_{1}}+\frac{\rho_{2}+\tilde{\rho}_{2}-1}{z-z_{2}}+\frac{\rho_{3}+\tilde{\rho}_{3}-1}{z-z_{3}}\right) \partial_{z} \\
&+\frac{\rho_{1} \tilde{\rho}_{1}\left(z_{1}-z_{2}\right)\left(z_{1}-z_{3}\right)}{\left(z-z_{1}\right)^{2}\left(z-z_{2}\right)\left(z-z_{3}\right)}+\frac{\rho_{2} \tilde{\rho}_{2}\left(z_{2}-z_{3}\right)\left(z_{2}-z_{1}\right)}{\left(z-z_{2}\right)^{2}\left(z-z_{3}\right)\left(z-z_{1}\right)}+\frac{\rho_{3} \tilde{\rho}_{3}\left(z_{3}-z_{1}\right)\left(z_{3}-z_{2}\right)}{\left(z-z_{3}\right)^{2}\left(z-z_{1}\right)\left(z-z_{2}\right)}
\end{aligned}
$$

(2) Let $z \mapsto w=h(z)=\frac{a z+b}{c z+d}$. (Transformations of this form are called homographies or Möbius transformations). We can always assume that $a d-b c=1$. Then

$$
\left.\mathcal{P}\left(\begin{array}{cccc}
h\left(z_{1}\right) & h\left(z_{2}\right) & h\left(z_{3}\right) \\
\rho_{1} & \rho_{2} & \rho_{3} & \\
\tilde{\rho}_{1} & \tilde{\rho}_{2} & \tilde{\rho}_{3} &
\end{array}\right), \partial_{w}\right)=(c z+d)^{4} \mathcal{P}\left(\begin{array}{cccc}
z_{1} & z_{2} & z_{3} & \\
\rho_{1} & \rho_{2} & \rho_{3} & z, \partial_{z} \\
\tilde{\rho}_{1} & \tilde{\rho}_{2} & \tilde{\rho}_{3} &
\end{array}\right)
$$

(3)

$$
\begin{aligned}
& \left(z-z_{1}\right)^{-\lambda}\left(z-z_{2}\right)^{\lambda} \mathcal{P}\left(\begin{array}{cccc}
z_{1} & z_{2} & z_{3} & \\
\rho_{1} & \rho_{2} & \rho_{3} & z, \partial_{z} \\
\tilde{\rho}_{1} & \tilde{\rho}_{2} & \tilde{\rho}_{3} &
\end{array}\right)\left(z-z_{1}\right)^{\lambda}\left(z-z_{2}\right)^{-\lambda} \\
& =\mathcal{P}\left(\begin{array}{cccc}
z_{1} & z_{2} & z_{3} & \\
\rho_{1}-\lambda & \rho_{2}+\lambda & \rho_{3} & z, \partial_{z} \\
\tilde{\rho}_{1}-\lambda & \tilde{\rho}_{2}+\lambda & \tilde{\rho}_{3} &
\end{array}\right)
\end{aligned}
$$

Clearly, in all above formulas one of $z_{i}$ can equal $\infty$, with an obvious meaning of various expressions. For convenience we give the expression for the Riemann operator with $z_{3}=\infty$ :

$$
\begin{align*}
& \mathcal{P}\left(\begin{array}{ccc}
z_{1} & z_{2} & \infty \\
\rho_{1} & \rho_{2} & \rho_{3} \\
\tilde{\rho}_{1} & \tilde{\rho}_{2} & \tilde{\rho}_{3}
\end{array}\right) \\
= & \partial_{z}^{2}-\left(\frac{\rho_{1}+\tilde{\rho}_{1}-1}{z-z_{1}}+\frac{\rho_{2}+\tilde{\rho}_{2}-1}{z-z_{2}}\right) \partial_{z} \\
& +\frac{\rho_{1} \tilde{\rho}_{1}\left(z_{1}-z_{2}\right)}{\left(z-z_{1}\right)^{2}\left(z-z_{2}\right)}+\frac{\rho_{2} \tilde{\rho}_{2}\left(z_{2}-z_{1}\right)}{\left(z-z_{2}\right)^{2}\left(z-z_{1}\right)}+\frac{\rho_{3} \tilde{\rho}_{3}}{\left(z-z_{1}\right)\left(z-z_{2}\right)} \tag{3.33}
\end{align*}
$$

The hypergeometric equation is a special case of the Riemann equation, since

$$
\begin{aligned}
\mathcal{F}\left(a, b ; c ; z, \partial_{z}\right) & =z(1-z) \mathcal{P}\left(\begin{array}{cccc}
0 & 1 & \infty & \\
0 & 0 & a & z, \partial_{z} \\
1-c & c-a-b & b &
\end{array}\right) \\
& =z(1-z) \partial_{z}^{2}+(c-(a+b+1) z) \partial_{z}-a b
\end{aligned}
$$

Every Riemann equation can be brought to the form of the hypergeometric equation by applying (1) and (2) of Theorem 3.1

From Theorem 3.1 we also see that symmetries of the hypergeometric equation are better visible if we replace $a, b, c$ with $\alpha, \beta, \mu$ :

$$
\begin{array}{lll}
\alpha=c-1 & \beta=a+b-c, & \mu=a-b \\
a=\frac{1+\alpha+\beta-\mu}{2}, & b=\frac{1+\alpha+\beta+\mu}{2}, & c=1+\alpha
\end{array}
$$

In fact, the new parameters coincide with the differences of the indices of the points $0,1, \infty$ :

$$
\alpha=\rho_{1}-\tilde{\rho}_{1}, \beta=\rho_{2}-\tilde{\rho}_{2}, \quad \mu=\rho_{3}-\tilde{\rho}_{3} .
$$

In the parameters $\alpha, \beta, \mu$, the Bose invariant of the hypergeometric equation has a more symmetric form:

$$
{ }_{2} I_{1}(\alpha, \beta, \mu ; z)=\frac{\left(1-\alpha^{2}\right)(1-z)+\left(1-\beta^{2}\right) z+\left(\mu^{2}-1\right)(1-z) z}{4 z^{2}(1-z)^{2}}
$$

We can summarize the relation between the hypergeometric equation and its canonical form by

$$
\begin{align*}
& -z^{\frac{\alpha}{2}+\frac{1}{2}-1}(1-z)^{\frac{\beta}{2}+\frac{1}{2}-1} \mathcal{F}\left(\frac{\alpha+\beta+\mu+1}{2}, \frac{\alpha+\beta-\mu+1}{2} ; 1+\alpha ; z, \partial_{z}\right) z^{-\frac{\alpha}{2}-\frac{1}{2}}(1-z)^{-\frac{\beta}{2}-\frac{1}{2}} \\
= & -\mathcal{P}\left(\begin{array}{ccc}
0 & 1 & \infty \\
\frac{\alpha}{2}+\frac{1}{2} & \frac{\beta}{2}+\frac{1}{2} & \frac{\mu}{2}-\frac{1}{2} \\
-\frac{\alpha}{2}+\frac{1}{2} & -\frac{\beta}{2}+\frac{1}{2} & -\frac{\mu}{2}-\frac{1}{2}
\end{array}\right) \\
= & -\partial_{z}^{2}+\left(\alpha^{2}-1\right) \frac{1}{4 z^{2}(1-z)}+\left(\beta^{2}-1\right) \frac{1}{4 z(1-z)^{2}}-\left(\mu^{2}-1\right) \frac{1}{4 z(1-z)} . \tag{3.34}
\end{align*}
$$

### 3.3 Trigonometric Pöschl-Teller potentials

In this subsection we consider Ansatz (1) of (3.30). We set

$$
\begin{equation*}
z=\sin ^{2} r, \quad \text { which solves } z^{\prime}=2 z^{\frac{1}{2}}(1-z)^{\frac{1}{2}} \tag{3.35}
\end{equation*}
$$

This leads to the Schrödinger equation

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V_{\alpha, \beta}(r)-\mu^{2}\right) \phi(r)=0 \tag{3.36}
\end{equation*}
$$

where

$$
\begin{align*}
V_{\alpha, \beta}(r) & :=\left(\alpha^{2}-\frac{1}{4}\right) \frac{1}{\sin ^{2} r}+\left(\beta^{2}-\frac{1}{4}\right) \frac{1}{\cos ^{2} r}  \tag{3.37}\\
& =\left(\frac{\alpha^{2}+\beta^{2}}{2}-\frac{1}{4}\right) \frac{4}{\sin ^{2} 2 r}+\frac{\left(\alpha^{2}-\beta^{2}\right)}{2} \frac{4 \cos 2 r}{\sin ^{2} 2 r}
\end{align*}
$$

This potential was proposed and solved by G. Pöschl and E. Teller [PT]. It is usually called the PöschlTeller potential, sometimes also the Pöschl-Teller potential of the first kind.

A natural real domain for this potential is $] 0,2 \pi[$. The potential is real if $r \in] 0,2 \pi\left[\right.$ and $\alpha^{2}, \beta^{2}$ are real.

Explicitly, the reduction of (3.36) to the hypergeometric equation is derived as follows:

$$
\begin{align*}
& -4 z^{\frac{\alpha}{2}+\frac{1}{4}}(1-z)^{\frac{\beta}{2}+\frac{1}{4}} \mathcal{F}\left(\frac{\alpha+\beta+\mu+1}{2}, \frac{\alpha+\beta-\mu+1}{2} ; 1+\alpha ; z, \partial_{z}\right) z^{-\frac{\alpha}{2}-\frac{1}{4}}(1-z)^{-\frac{\beta}{2}-\frac{1}{4}} \\
= & \left.-4 z(1-z) \mathcal{P}\left(\begin{array}{ccc}
0 & 1 & \infty \\
\frac{\alpha}{2}+\frac{1}{4} & \frac{\beta}{2}+\frac{1}{4} & \frac{\mu}{2} \\
-\frac{\alpha}{2}+\frac{1}{4} & -\frac{\beta}{2}+\frac{1}{4} & -\frac{\mu}{2}
\end{array}\right), \partial_{z}\right) \\
= & -4 z(1-z)\left(\partial_{z}^{2}+\left(\frac{1}{2 z}-\frac{1}{2(1-z)}\right) \partial_{z}\right) \\
& +\left(\alpha^{2}-\frac{1}{4}\right) \frac{1}{z}+\left(\beta^{2}-\frac{1}{4}\right) \frac{1}{1-z}-\mu^{2}  \tag{3.38}\\
= & -\partial_{r}^{2}+\left(\alpha^{2}-\frac{1}{4}\right) \frac{1}{\sin ^{2} r}+\left(\beta^{2}-\frac{1}{4}\right) \frac{1}{\cos ^{2} r}-\mu^{2} .
\end{align*}
$$

### 3.4 Hyperbolic Pöschl-Teller potentials

We continue with (1) of (3.30). We set

$$
\begin{equation*}
z=-\sinh ^{2} r, \quad \text { which solves } z^{\prime}=-2(-z)^{\frac{1}{2}}(1-z)^{\frac{1}{2}} \tag{3.39}
\end{equation*}
$$

This leads to the Schrödinger equation

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V_{\alpha, \beta}(r)-\mu^{2}\right) \phi(r)=0 \tag{3.40}
\end{equation*}
$$

where

$$
\begin{align*}
V_{\alpha, \beta}(r) & :=\left(\alpha^{2}-\frac{1}{4}\right) \frac{1}{\sinh ^{2} r}-\left(\eta^{2}-\frac{1}{4}\right) \frac{1}{\cosh ^{2} r}  \tag{3.41}\\
& =\left(\frac{\alpha^{2}+\beta^{2}}{2}-\frac{1}{4}\right) \frac{4}{\sinh ^{2} 2 r}+\frac{\left(\alpha^{2}-\beta^{2}\right)}{2} \frac{4 \cosh 2 r}{\sinh ^{2} 2 r}
\end{align*}
$$

This potential was also proposed and solved by G. Pöschl and E. Teller $\overline{\mathrm{PT}}$. In the literature it is known as hyperbolic, generalized Pöschl-Teller potential or the Pöschl-Teller potential of the second kind.

A natural real domain for this potential is $] 0, \infty[$. The potential is real if $r \in] 0, \infty\left[\right.$ and $\alpha^{2}, \beta^{2}$ are real.
To see that (3.42) can be solved in terms of the hypergeometric equation, we first repeat the computations leading to (3.38), and then set $z=-\sinh ^{2} r$.

Note that the trigonometric and hyperbolic Pöschl-Teller potentials are related with a complex linear transformation $r \mapsto \mathrm{i} r$. Therefore, these two potentials are two real cases of the same complex case.

### 3.5 Alternative approach to Pöschl-Teller potentials

Ansatzes (2) and (3) lead to the same classes of exactly solvable potentials as Ansatz (1). To see this it is enough to consider Ansatz (3). In fact, we first repeat computations analogous to (3.38):

$$
\begin{aligned}
& -4 z(1-z)^{2} \mathcal{P}\left(\begin{array}{cccc}
0 & 1 & \infty & \\
\frac{\alpha}{2}+\frac{1}{4} & \frac{\beta}{2} & \frac{\mu}{2}+\frac{1}{4} & z, \partial_{z} \\
-\frac{\alpha}{2}+\frac{1}{4} & -\frac{\beta}{2} & -\frac{\mu}{2}+\frac{1}{4}
\end{array}\right) \\
= & -4 z(1-z)^{2}\left(\partial_{z}^{2}+\left(\frac{1}{2 z}-\frac{1}{1-z}\right) \partial_{z}\right) \\
& +\left(\alpha^{2}-\frac{1}{4}\right) \frac{1-z}{z}+\beta^{2} \frac{1}{1-z}-\left(\mu^{2}-\frac{1}{4}\right)(1-z)
\end{aligned}
$$

We set

$$
\begin{equation*}
z=\operatorname{tgh}^{2} r, \quad \text { which solves } z^{\prime}=2 z^{\frac{1}{2}}(1-z) \tag{3.42}
\end{equation*}
$$

obtaining the trigonometric Pöschl-Teller potentials, or

$$
\begin{equation*}
z=-\operatorname{tg}^{2} r, \quad \text { which solves } \quad z^{\prime}=-2(-z)^{\frac{1}{2}}(1-z) \tag{3.43}
\end{equation*}
$$

obtaining the hyperbolic Pöschl-Teller potentials.

### 3.6 Rosen/Scarf/Woods-Saxon potentials

Let us consider Ansatz (3) of (3.30). We set

$$
\begin{equation*}
z=\frac{1}{1+\mathrm{e}^{2 r}}, \quad \text { which solves } \quad z^{\prime}=2 z(z-1) \tag{3.44}
\end{equation*}
$$

We also set $\kappa=\beta^{2}-\alpha^{2}$. This leads to the Schrödinger equation

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V_{\kappa, \mu}(r)+\alpha^{2}\right) \phi(r)=0 \tag{3.45}
\end{equation*}
$$

where

$$
V_{\kappa, \mu}(y):=\kappa \frac{\mathrm{e}^{r}}{2 \cosh r}+\left(\frac{\mu^{2}}{4}-\frac{1}{4}\right) \frac{1}{\cosh ^{2} r}
$$

This potential was proposed and solved by M. F. Manning and N. Rosen MR. In the literature it is also called the Scarf [Sca] or Woods-Saxon [SW] potential (for instance in [Flü]).

A natural real domain for this potential is $]-\infty, \infty[$. The potential is real if $r \in]-\infty, \infty\left[\right.$ and $\kappa, \beta^{2}$ are real.

Here is a derivation of (3.45) from the hypergeometric equation:

$$
\begin{aligned}
& -4 z^{1+\frac{\alpha}{2}}(1-z)^{1+\frac{\beta}{2}} \mathcal{F}\left(\frac{\alpha+\beta+\mu+1}{2}, \frac{\alpha+\beta-\mu+1}{2} ; 1+\alpha ; z, \partial_{z}\right) z^{-\frac{\alpha}{2}}(1-z)^{-\frac{\beta}{2}} \\
= & -4 z^{2}(z-1)^{2} \mathcal{P}\left(\begin{array}{cccc}
0 & 1 & \infty \\
\frac{\alpha}{2} & \frac{\beta}{2} & \mu+\frac{1}{2} & z, \partial_{z} \\
-\frac{\alpha}{2} & -\frac{\beta}{2} & -\mu+\frac{1}{2} &
\end{array}\right)
\end{aligned}
$$

$$
\begin{align*}
& =-4 z^{2}(1-z)^{2}\left(\partial_{z}^{2}+\left(\frac{1}{z}-\frac{1}{1-z}\right) \partial_{z}\right)+\alpha^{2}(1-z)+\beta^{2} z-\left(\mu^{2}-1\right) z(z-1)  \tag{3.46}\\
& =-\partial_{r}^{2}+\alpha^{2}+\left(-\alpha^{2}+\beta^{2}\right) \frac{1}{1+\mathrm{e}^{2 r}}-\left(\mu^{2}-1\right) \frac{\mathrm{e}^{2 r}}{\left(1+\mathrm{e}^{2 r}\right)^{2}} \tag{3.47}
\end{align*}
$$

### 3.7 Eckart/Hulthen potentials

We still consider Ansatz (3) of (3.30).
We set

$$
\begin{equation*}
z=\frac{1}{1-\mathrm{e}^{-2 r}}, \quad \text { which solves } \quad z^{\prime}=2 z(1-z) \tag{3.48}
\end{equation*}
$$

We also set $\kappa=\beta^{2}-\alpha^{2}$. This leads to the Schrödinger equation

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V_{\kappa, \mu}(r)+\alpha^{2}\right) \phi(r)=0 \tag{3.49}
\end{equation*}
$$

where

$$
V_{\kappa, \mu}(y):=\kappa \frac{\mathrm{e}^{r}}{2 \sinh r}+\left(\frac{\mu^{2}}{4}-\frac{1}{4}\right) \frac{1}{\sinh ^{2} r}
$$

This potential was proposed and solved by C. Eckart Eck. In the literature (for instance [Flü]) it is also called the Hulthen potential Hul, sometimes also the generalized Morse potential, because of its similarity to the Morse potential, see Subsect. 4.4

A natural real domain for this potential is $] 0, \infty[$. If $r \in] 0, \infty\left[\right.$ and $\kappa, \beta^{2}$ are real, then the potential is real.

To check these computations, we first repeat the computations leading to (3.46), and then set $z=$ $\frac{1}{1-\mathrm{e}^{-2 r}}$.

Note that from the Rosen and Eckart potentials are related with a complex affine transformation $r \mapsto \mathrm{i} \pi-r$. Therefore, these two potentials can be treated as two real cases of the same complex case.

### 3.8 Alternative approach to Rosen/Scarf/Wood-Saxon and Eckart/Hulthen potentials

Ansatzes (5) and (6) lead to the same classes of exactly solvable potentials as Ansatz (4). To see this it is enough to consider Ansatz (5). In fact, we first repeat computations analogous to (3.46):

$$
\begin{aligned}
& -4 z^{-\frac{\alpha+1}{2}}(1-z)^{1+\frac{\beta}{2}} \mathcal{F}\left(\frac{\alpha+\beta+\mu+1}{2}, \frac{\alpha+\beta+-\mu+1}{2} ; 1+\alpha ; z, \partial_{z}\right) z^{-\frac{\alpha+1}{2}}(1-z)^{-\frac{\beta}{2}} \\
= & -4(z-1)^{2} \mathcal{P}\left(\begin{array}{cccc}
0 & 1 & \infty \\
\alpha+\frac{1}{2} & \frac{\beta}{2} & \frac{\mu}{2} & z, \partial_{z} \\
-\alpha+\frac{1}{2} & -\frac{\beta}{2} & -\frac{\mu}{2}
\end{array}\right) \\
= & -4(1-z)^{2}\left(\partial_{z}^{2}-\frac{1}{1-z} \partial_{z}\right) \\
& +\left(\alpha^{2}-1\right) \frac{1-z}{z^{2}}+\beta^{2} \frac{1}{z}-\mu^{2} \frac{1-z}{z} .
\end{aligned}
$$

We set

$$
\begin{equation*}
z=1-\mathrm{e}^{2 r}, \quad \text { which solves } z^{\prime}=2(z-1) \tag{3.50}
\end{equation*}
$$

obtaining the Rosen potential, or

$$
\begin{equation*}
z=1+\mathrm{e}^{2 r}, \quad \text { which solves } \quad z^{\prime}=-2(z-1) \tag{3.51}
\end{equation*}
$$

obtaining the Eckart potential.

## 4 Schrödinger operators reducible to the confluent equation

One of basic exactly solvable equations is the confluent equation, given by the operator

$$
\begin{equation*}
\mathcal{F}\left(a ; c ; z, \partial_{z}\right):=z \partial_{z}^{2}+(c-z) \partial_{z}-a . \tag{4.52}
\end{equation*}
$$

It is convenient to consider (4.52) in parallel with the equation given by the operator

$$
\begin{equation*}
\mathcal{F}\left(a, b ;-; z, \partial_{z}\right):=z^{2} \partial_{z}^{2}+(-1+(1+a+b) z) \partial_{z}+a b . \tag{4.53}
\end{equation*}
$$

The equation given by (4.53) is sometimes called the ${ }_{2} F_{0}$ equation. Note that

$$
z^{a} \mathcal{F}\left(a, b ;-; z, \partial_{z}\right) z^{-a}=w \mathcal{F}\left(a ; 1+a-b ; w, \partial_{w}\right), \quad w=-z^{-1}, \quad z=-w^{-1}
$$

Hence the ${ }_{2} F_{0}$ equation is equivalent to the confluent equation. The relationship between the parameters is

$$
c=1+a-b, \quad b=1+a-c
$$

Another exactly solvable equation that we will consider in this section is sometimes called the ${ }_{0} F_{1}$ equation. It is given by

$$
\begin{equation*}
\mathcal{F}\left(c ; z, \partial_{z}\right):=z \partial_{z}^{2}+(c-z) \partial_{z}-a, \tag{4.54}
\end{equation*}
$$

and is equivalent to the Bessel equation.
Clearly, the confluent, ${ }_{2} F_{0}$ and ${ }_{0} F_{1}$ equations belong to the class of hypergeometric type equations.
The basic equation of this section will be (1.5), given by

$$
\begin{equation*}
\mathcal{F}\left(a ; c ; \gamma ; z ; \partial_{z}\right):=z \partial_{z}^{2}+(c-\gamma z) \partial_{z}-a \tag{4.55}
\end{equation*}
$$

It will be called the generalized confluent equation. Note that in the case $\gamma=1$ 4.55) coincides with the confluent equation. If $\gamma \neq 0$, 4.55) can be reduced to the confluent equation by rescaling (and hence also to the ${ }_{2} F_{0}$ equation). If $\gamma=0$, 4.55) coincides with the ${ }_{0} F_{1}$ equation.

The Bose invariant of the generalized confluent equation equals

$$
{ }_{1} I_{1}(a, b ; \gamma ; z)=\frac{-\gamma^{2} z^{2}+2(c \gamma-2 a) z+2 c-c^{2}}{4 z^{2}}
$$

### 4.1 Symmetries of the generalized confluent equation

Let us first describe symmetries of the generalized confluent equation.

$$
\begin{align*}
z^{c-1} \mathcal{F}\left(a ; c ; \gamma ; z, \partial_{z}\right) z^{1-c} & =\mathcal{F}\left(a-c \gamma+\gamma ; 2-c ; \gamma ; z, \partial_{z}\right)  \tag{4.56}\\
\mathrm{e}^{-\gamma z} \mathcal{F}\left(a ; c ; \gamma ; z, \partial_{z}\right) \mathrm{e}^{\gamma z} & =-\mathcal{F}\left(c-a ; c ; \gamma ; w, \partial_{w}\right), \quad z=-w ; \tag{4.57}
\end{align*}
$$

Besides, the scaling acts as follows:

$$
\mathcal{F}\left(a ; c ; \gamma ; z, \partial_{z}\right)=\gamma \mathcal{F}\left(a / \gamma ; c ; 1 ; w, \partial_{w}\right), \quad w=\gamma z
$$

It is convenient to introduce new parameters $\alpha, \nu$ :

$$
\begin{array}{ll}
a=\frac{\gamma+\alpha \gamma-\nu}{2}, & c=1+\alpha \\
\alpha=c-1=a-b, & \nu=c \gamma-2 a=1-a-b .
\end{array}
$$

In the new parameters the Bose invariant of the generalized confluent equation has a more symmetric form:

$$
{ }_{1} I_{1}(\alpha, \nu, \gamma ; z)=-\frac{\gamma^{2}}{4}+\frac{\nu}{2 z}+\frac{1-\alpha^{2}}{4 z^{2}}
$$

Thus the starting point for the further analysis will be the equation

$$
\begin{equation*}
\left(-\partial_{z}^{2}+\frac{\gamma^{2}}{4}-\frac{\nu}{2 z}+\left(\frac{\alpha^{2}}{4}-\frac{1}{4}\right) \frac{1}{z^{2}}\right) \phi(z)=0 \tag{4.58}
\end{equation*}
$$

We will describe 3 classes of Schrödinger operators solved using the confluent equation corresponding to three obvious choices for $\left(r^{\prime}(z)\right)^{2}$ :

$$
\begin{equation*}
\text { (1) } \frac{1}{4}, \quad \text { (2) } \frac{1}{4 z}, \quad \text { (3) } \frac{1}{4 z^{2}} \text {. } \tag{4.59}
\end{equation*}
$$

### 4.2 Hydrogen atom

We consider (1) of (4.59). We set

$$
\begin{equation*}
z=2 r \tag{4.60}
\end{equation*}
$$

This leads to the Schrödinger equation

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V_{\alpha, \nu}(r)+\gamma^{2}\right) \phi(r)=0 \tag{4.61}
\end{equation*}
$$

with the potential

$$
V_{\alpha, \nu}(r):=-\frac{\nu}{r}+\left(\frac{\alpha^{2}}{4}-\frac{1}{4}\right) \frac{1}{r^{2}}
$$

(4.61) is the radial part of the Schrödinger equation for the hydrogen atom.
$] 0, \infty\left[\right.$ is a natural real domain for this equation. If $\nu, \alpha^{2}$ are real, then so is $V_{\alpha, \nu}$ on $] 0, \infty[$. The derivation of (4.61) from (4.58) is immediate:

$$
\begin{aligned}
& 4\left(-\partial_{z}^{2}+\frac{\gamma^{2}}{4}-\frac{\nu}{2 z}+\left(\frac{\alpha^{2}}{4}-\frac{1}{4}\right) \frac{1}{z^{2}}\right) \\
= & -\partial_{r}^{2}+\gamma^{2}-\frac{\nu}{r}+\left(\frac{\alpha^{2}}{4}-\frac{1}{4}\right) \frac{1}{r^{2}}
\end{aligned}
$$

### 4.3 Rotationally symmetric harmonic oscillator

We consider (2) of (4.59). We set

$$
\begin{equation*}
z=r^{2}, \quad \text { which solves } \quad z^{\prime}=2 \sqrt{z} \tag{4.62}
\end{equation*}
$$

This leads to the Schrödinger equation

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V_{\alpha, \gamma}(r)-2 \nu\right) \phi(r)=0 \tag{4.63}
\end{equation*}
$$

with the potential

$$
V_{\alpha, \gamma}(r):=\gamma^{2} r^{2}+\left(\alpha^{2}-\frac{1}{4}\right) \frac{1}{r^{2}}
$$

$] 0, \infty\left[\right.$ is a natural real domain. For real $\gamma^{2}, \alpha^{2}$, the potential is real on $] 0, \infty[$. By scaling, the only differet real cases are $\gamma^{2}=1,0,-1$.

For $\gamma^{2}=1$ the equation is the radial part of the rotationally symmetric harmonic oscillator. For $\gamma^{2}=0$ the equation is the radial part of the Helmholtz equation.

Here is an explicit derivation of (4.63) from (4.58):

$$
\begin{aligned}
& 4 z^{1-\frac{1}{4}}\left(-\partial_{z}^{2}+\frac{\gamma^{2}}{4}-\frac{\nu}{2 z}+\left(\frac{\alpha^{2}}{4}-\frac{1}{4}\right) \frac{1}{z^{2}}\right) z^{\frac{1}{4}} \\
= & -4 z \partial_{z}^{2}-2 \partial_{z}+\gamma^{2} z-2 \nu+\left(\alpha^{2}-\frac{1}{4}\right) \frac{1}{z} \\
= & -\partial_{r}^{2}+\gamma^{2} r^{2}-2 \nu+\left(\alpha^{2}-\frac{1}{4}\right) \frac{1}{r^{2}} .
\end{aligned}
$$

### 4.4 Morse potential

We consider (3) of (4.59). We set

$$
\begin{equation*}
z=\mathrm{e}^{-r}, \quad \text { which solves } \quad z^{\prime}=-z \tag{4.64}
\end{equation*}
$$

This leads to the Schrödinger equation

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V_{\nu, \gamma}(r)+\frac{\alpha^{2}}{4}\right) \phi(r)=0 \tag{4.65}
\end{equation*}
$$

where

$$
V_{\nu, \gamma}(r)=\frac{\gamma^{2}}{4} \mathrm{e}^{-2 r}-\frac{\nu}{2} \mathrm{e}^{-r}
$$

$]-\infty, \infty\left[\right.$ is a natural real domain. For real $\nu, \gamma^{2}$, the potential is real on $]-\infty, \infty[$. By translation, the only different real cases are $\gamma^{2}=1,0,-1$.

Here is an explicit derivation of (4.65) from (4.58):

$$
\begin{aligned}
& z^{2-\frac{1}{2}}\left(-\partial_{z}^{2}+\frac{\gamma^{2}}{4}-\frac{\nu}{2 z}+\left(\frac{\alpha^{2}}{4}-\frac{1}{4}\right) \frac{1}{z^{2}}\right) z^{\frac{1}{2}} \\
= & -z^{2} \partial_{z}^{2}-z \partial_{z}+\frac{\gamma^{2}}{4} z^{2}-\frac{\nu}{2}+\frac{\alpha^{2}}{4} \\
= & -\partial_{r}^{2}+\frac{\gamma^{2}}{4} \mathrm{e}^{-2 r}-\frac{\nu}{2} \mathrm{e}^{-r}+\frac{\alpha^{2}}{4} .
\end{aligned}
$$

## 5 Schrödinger equations reducible to the generalized harmonic oscillator

The last family of exactly solvable Schrödinger equations that we describe is (1.6), given by

$$
\begin{equation*}
-\partial_{z}^{2}+\theta^{2} z^{2}+\rho z+\lambda \tag{5.66}
\end{equation*}
$$

We will call (5.66) the generalized harmonic oscillator equation. It is already in the canonical form.

If $\theta^{2} \neq 0$, then (5.66) is the Schrödinger equation for translated harmonic oscillator. It is then equivalent to the Hermite equation, given by

$$
\mathcal{G}\left(a, y, \partial_{y}\right):=\left(\partial_{y}^{2}-2 y \partial_{y}-2 a\right)
$$

In fact, we have

$$
\begin{aligned}
& \theta \mathrm{e}^{-\frac{y^{2}}{2}} \mathcal{G}\left(a, y, \partial_{y}\right) \mathrm{e}^{\frac{y^{2}}{2}} \\
= & \theta\left(-\partial_{y}^{2}+y^{2}+2 a-1\right) \\
= & -\partial_{z}^{2}+\theta^{2} z^{2}+\rho z+\frac{\rho^{2}}{2 \theta^{2}}+2 a-1
\end{aligned}
$$

where $y=\sqrt{\theta}\left(z+\frac{\rho}{2 \theta^{2}}\right)$.
For $\theta=0, \rho \neq 0$, (5.66) is equivalent to the Airy equation

$$
\left(\partial_{y}^{2}+y\right) \psi(y)=0
$$

which in turn is equivalent to a special case of the ${ }_{0} F_{1}$ equation:

$$
\begin{aligned}
& -\partial_{z}^{2}+\rho z+\lambda \\
= & \rho^{\frac{2}{3}}\left(-\partial_{y}^{2}+y\right) \\
= & (\rho 3)^{\frac{2}{3}} w^{\frac{1}{3}} \mathcal{F}\left(\frac{2}{3} ; w, \partial_{w}\right),
\end{aligned}
$$

where we set $y=\rho^{\frac{1}{3}} z+\rho^{-\frac{2}{3}} \lambda$, and $w=3^{-2} y^{3}$.
We will describe 3 classes of Schrödinger operators that can be reduced to (5.66) corresponding to three obvious choices for $\left(r^{\prime}(z)\right)^{2}$ :

$$
\begin{array}{ll}
(1) & (2)  \tag{5.67}\\
, & (3) z^{2}
\end{array}
$$

### 5.1 Generalized harmonic oscillator

We consider Ansatz (1) of (5.67), which corresponds to the most obvious choice of the energy that is $-\lambda$. We rename the variable $z=r$. This leads to the Schrödinger equation

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V_{\theta, \rho}(r)+\lambda\right) \phi(r)=0 \tag{5.68}
\end{equation*}
$$

where

$$
V_{\theta, \rho}(r)=\theta^{2} r^{2}+\rho r
$$

### 5.2 Special potential I

We can choose the energy to be $-\rho$, that is Ansatz (2) of (5.67). This corresponds to the substitution

$$
z=\left(\frac{3 r}{2}\right)^{\frac{2}{3}}, \quad \text { which solves } \quad z^{\prime}=z^{-\frac{1}{2}}
$$

and leads to the Schrödinger equation

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V_{\theta, \lambda}(r)+\rho\right) \phi(r)=0 \tag{5.69}
\end{equation*}
$$

where

$$
V_{\theta, \lambda}(r)=\theta^{2}\left(\frac{3 r}{2}\right)^{\frac{2}{3}}+\lambda\left(\frac{2}{3 r}\right)^{\frac{2}{3}}-\frac{5}{36} \frac{1}{r^{2}}
$$

In fact,

$$
\begin{aligned}
& z^{\frac{1}{4}-1}\left(-\partial_{z}^{2}+\theta^{2} z^{2}+\rho z+\lambda\right) z^{-\frac{1}{4}} \\
= & -\frac{1}{z} \partial_{z}^{2}+\frac{1}{2 z^{2}} \partial_{z}+\theta^{2} z+\rho+\frac{\lambda}{z}-\left(\frac{1}{4}+\frac{1}{4^{2}}\right) \frac{1}{z^{3}} \\
= & -\partial_{r}^{2}+\theta^{2}\left(\frac{3 r}{2}\right)^{\frac{2}{3}}+\rho+\lambda\left(\frac{2}{3 r}\right)^{\frac{2}{3}}-\left(\frac{1}{4}+\frac{1}{4^{2}}\right)\left(\frac{2}{3 r}\right)^{2} .
\end{aligned}
$$

Note the following intriguing feature of the above potential: the coefficient at $r^{-2}$ is fixed and one cannot change it by rescaling the variable $r$.

### 5.3 Special potential II

We can choose the energy to be $-\theta^{2}$, that is Ansatz (3) of (5.67). This corresponds to the substitution

$$
z=(2 r)^{\frac{1}{2}}, \quad \text { which solves } \quad z^{\prime}=z^{-1}
$$

and leads to the Schrödinger equation

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V_{\rho, \lambda}(r)+\theta^{2}\right) \phi(r)=0 \tag{5.70}
\end{equation*}
$$

where

$$
V_{\rho, \lambda}(r)=\frac{\rho}{(2 r)^{\frac{1}{2}}}+\frac{\lambda}{2 r}-\frac{3}{16} \frac{1}{r^{2}}
$$

In fact,

$$
\begin{aligned}
& z^{\frac{1}{2}-2}\left(-\partial_{z}^{2}+\theta^{2} z^{2}+\rho z+\lambda\right) z^{-\frac{1}{2}} \\
= & -\frac{1}{z^{2}} \partial_{z}^{2}+\frac{1}{z^{3}} \partial_{z}+\theta^{2}+\frac{\rho}{z}+\frac{\lambda}{z^{2}}-\left(\frac{1}{2}+\frac{1}{2^{2}}\right) \frac{1}{z^{4}} \\
= & -\partial_{r}^{2}+\theta^{2}+\frac{\rho}{(2 r)^{\frac{1}{2}}}+\frac{\lambda}{2 r}-\left(\frac{1}{2}+\frac{1}{2^{2}}\right) \frac{1}{2^{2} r^{2}} .
\end{aligned}
$$

Again, the coefficient at $r^{-2}$ is fixed and one cannot change it by rescaling the variable $r$.

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# Exactly solvable Schrödinger operators 

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Abstract

We systematically describe and classify 1-dimensional Schrödinger equations that can be solved in terms of hypergeometric type functions. Beside the well-known families, we explicitly describe 2 new classes of exactly solvable Schrödinger equations that can be reduced to the Hermite equation.

## 1 Introduction

Exactly solvable 1-dimensional Schrödinger equations play an important role in quantum physics. The best known are the harmonic oscillator and the radial equation for the hydrogen atom, which are covered in every course of quantum mechanics. A number of other examples were discovered in the 30's of the last century. They include (known also under other names) trigonometric and hyperbolic Pöschl-Teller potentials [PT], the Scarf potential [Sca, the Eckart potential [Eck], the Manning-Rosen potential MR] and the Rosen-Morse potential. All these examples can be reduced to the hypergeometric equation (see eg. WW). One should also mention the Morse potential that leads to the confluent equation. Problems involving these potentials are often used in classes on quantum mechanics, see eg. the well-known problem book of Flügge [Flü]. A number of techniques have been developed to study their properties, such as the factorization method [HI] and the closely related method of the superpotential (see for instance CKS and Cot ).

In later years J. N. Ginocchio Gin discovered that these examples can be generalized to a larger class of potentials equivalent to the hypergeometric or confluent equation. Later, this class was extended by G. A. Natanzon [Nat] and further generalized by R. Milson Mil]. These classes, besides the hitherto known potentials, are not very practical in applications, since they are not given by explicit expressions.

There exist even more general classes of potentials that can be called exactly solvable, found by A. Khare and U. Sukhatme [KS. They are, however, expressed in exotic special functions.

The literature on exactly solvable Schrödinger equations is very large. The subject is in fact very useful for applications, especially in quantum physics. We also believe that it is quite beautiful, which in the existing literature is perhaps not so easy to see.

In our paper we would like to systematically describe basic classes of exactly solvable Schrödinger equations. Most of the material of our paper is scattered in the literature, notably in (Nat, Mil] and CKS (see also Bos). Our treatment is, however, somewhat more systematic than what we could find in the literature. For instance, we distinguish between the complex classification and the real classification. It also seems that for the first time we explicitly describe 2 new classes of exactly solvable Schrödinger equations, which can be reduced to the Hermite equation (see Subsects 5.2 and 5.3).

By a (stationary 1-dimensional) Schrödinger equation we will mean an equation of the form

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V(r)-E\right) \phi(r)=0 \tag{1.1}
\end{equation*}
$$

This equation can be interpreted as the eigenvalue problem for the operator

$$
\begin{equation*}
H:=-\partial_{r}^{2}+V(r) \tag{1.2}
\end{equation*}
$$

An operator of the form (1.2) will be called a (1-dimensional) Schrödinger operator. $V(r)$ will be called a potential and the parameter $E$ an energy.
(1.2) can be interpreted as an operator in a number of ways. If $V$ is a holomorphic function of the complex variable on some open $\Omega \subset \mathbb{C}$, then $r$ can be interpreted as a complex variable and (1.2) can be
viewed as an operator on holomorphic functions. In this case it is natural to allow $E$ to be a complex parameter. The corresponding eigenvalue equation (1.1) will be then called the Schrödinger equation of the complex variable. Let us note that, complex affine transformations preserve the class of Schrödinger equations of the complex variable. (By a complex affine transformations we mean $r \mapsto a r+b$, where $a \neq 0, b$ are complex constants).

One can also interpret $r$ as a real variable in some open $I \subset \mathbb{R}$. The operator (1.2) is then viewed as acting on functions on $I$. Understood in this way, (1.1) will be called the Schrödinger equation of the real variable. Of special interest is then the case of real potential. Clearly, real affine transformations preserve the class of real Schrödinger equations of the real variable.

Our paper is organized as follows. First, we will briefly dicuss some general facts related to 2nd order linear differential equations. In particular, we will describe basic ingredients of the so-called BoseNatanzon method, which permits to obtain a class of Schrödinger equations equivalent to a chosen linear equation. Note that it applies to more general situations than those described in the literature ( Nat, Mil]. We will formulate a criterion that determines when the Bose-Natanzon method can be used.

Next, we will focus on hypergeometric type equations. Recall that equations of the form

$$
\begin{equation*}
\left(\sigma(z) \partial_{z}^{2}+\tau(z) \partial_{z}+\eta\right) f(z)=0 \tag{1.3}
\end{equation*}
$$

where $\sigma(z), \tau(z)$ are polynomials with

$$
\operatorname{deg} \sigma \leq 2, \operatorname{deg} \tau \leq 1
$$

and $\eta$ is a number, are sometimes called hypergeometric type equations (see A.F. Nikiforov and V.B. Uvarov [NU]). Solutions of (1.3) are called hypergeometric type functions. Hypergeometric type functions are very well understood and include the hypergeomeric function, the confluent function, Hermite, Laguerre and Jacobi polynomials, etc. Traditionally, one classifies hypergeometric type equations into several distinct classes invariant under complex affine transformations. In each class one chooses a simple representative, to which the whole class can be reduced. Such representatives are the hypergeometric equation, the confluent equation, etc.

In our analysis, we will concentrate on 3 specific equations:
(1) the hypergeometric equation

$$
\begin{equation*}
\left(z(1-z) \partial_{z}^{2}+(c-(a+b+1) z) \partial_{z}-a b\right) f(z)=0 \tag{1.4}
\end{equation*}
$$

(2) the rescaled confluent equation

$$
\begin{equation*}
\left(z \partial_{z}^{2}+(c-\gamma z) \partial_{z}-a\right) f(z)=0 \tag{1.5}
\end{equation*}
$$

(3) the translated harmonic oscillator

$$
\begin{equation*}
\left(-\partial_{z}^{2}+\theta^{2} z^{2}+\rho z+\lambda\right) f(z)=0 \tag{1.6}
\end{equation*}
$$

Note that all hypergeometric type equations with a nonzero $\sigma$ are related with one of the above equations by elementary operations

- multiplication of $f(z)$ by a function, multiplication of the equation by a function,
- a change of variables not depending on the parameters $a, b, c$ (respectively $a, c, \gamma$ or $\theta, \rho, \lambda$ ).
(1) If $\sigma$ is 2 nd order and has 2 distinct roots, an affine transformation reduces (1.3) to (1.4).
(2) If $\sigma$ is 2 nd order and has only one root, then a transformation involving $z \mapsto z^{-1}$ reduces (1.3) to (1.5) with $\gamma \neq 0$.
(3) If $\sigma$ is 1st order, then an affine transformation reduces (1.3) to (1.5).
(4) If $\sigma$ is 0 th order and nonzero, then (1.3) is equivalent to (1.6) with $\theta \neq 0$.

Obviously, (1.4) and (1.5) are themselves equations of the hypergeometric type. (1.6) is not, but can be reduced to a hypergeometric type equation. This reduction is depends on whether $\theta \neq 0$ or $\theta=0$. The former case leads to a hypergeometric type equation with a constant $\sigma$, as mentioned above. In the case $\theta=0$, (1.6) is the Airy equation, which can be reduced to a special case of the hypergeometric type equation by a transformation involving $z \mapsto z^{3}$.

Equations (1.4), (1.5) and (1.6) are very well understood and have well known solutions. We will describe classes of solvable potentials reducible to one of the equations (1.4), (1.5) or (1.6). We will describe both the complex and the real classification of such potentials.

It is natural to consider first a classification of exactly solvable Schrödinger equations of the complex variable. Note that all solvable potentials that we consider are holomorphic on the complex plane, apart from some isolated singularities. Obviously, we can always use a complex affine transformation to put the equation in a convenient form. One can also move a complex constant from the potential to the energy.

Consider a family of exactly solvable holomorphic potentials. Suppose that for a selected subfamily of parameters the potential is real if restricted to an open, possibly infinite interval $I \subset \mathbb{C}$. By elementary properties of holomorphic functions, if we extend the interval $I$, the potential is still real until we hit a singularity. A real affine transformation can be used to put the equation in a convenient form. In particular, we can always assume that $I$ is a subset of the real line. We can also move a real constant from the potential to the energy.

The above discussion motivates the following definition: We will say that an open interval $I \subset \mathbb{R}$ is a natural real domain for a certain family of potentials if it ends either at $-\infty,+\infty$ or at a singularity of the potential, and no singularities lie inside $I$. When describing the real classification, we will always specify a natural real domain of the potential.

Clearly, if $I$ is a natural real domain, then the operator $H=-\partial_{r}^{2}+V(r)$ is hermitian on $C_{\mathrm{c}}^{\infty}(I)$. It is then natural to ask about self-adjoint extensions of $H$. We will not discuss this question here. We plan to consider it in later papers.

In Equations (1.4), (1.5) and (1.6) we have 3 arbitrary parameters. Therefore, in all the solvable cases that we will describe, the potential depends on 2 parameters, since the third parameter is responsible for the energy $E$.

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## 2 Second order homogeneous differential equations

The main object of our paper are ordinary homogeneous 2 nd order linear differential equations, that is equations of the form

$$
\begin{equation*}
\left(a(r) \partial_{r}^{2}+b(r) \partial_{r}+c(r)\right) \phi(r)=0 \tag{2.7}
\end{equation*}
$$

It will be convenient to treat (2.7) as the problem of finding the kernel of the operator

$$
\begin{equation*}
\mathcal{A}\left(r, \partial_{r}\right):=a(r) \partial_{r}^{2}+b(r) \partial_{r}+c(r) \tag{2.8}
\end{equation*}
$$

We will then say that the equation (2.7) is given by the operator (2.8).
We will treat $r$ either as a complex or a real variable. In the complex case we will usually assume that the coefficients are analytic.

In this section we describe some general facts related to the theory of equations of the form (2.7) and their reduction to the Schrödinger equation.

### 2.1 The Bose invariant of a second order differential equation

Suppose that we are given a second order differential equation (2.7). Then we can always eliminate the first order term as follows. We divide from the left by $a(r)$, set

$$
h(r):=\exp \left(\int^{r} \frac{b(t)}{2 a(t)} \mathrm{d} t\right)
$$

We check that

$$
h(r) \mathcal{A}\left(r, \partial_{r}\right) h(r)^{-1}=\partial_{r}^{2}+I(r)
$$

where, using the notation $a^{\prime}:=\frac{\mathrm{d}}{\mathrm{d} r} a(r)$,

$$
\begin{equation*}
I=\frac{4 a c-2 a b^{\prime}+2 b a^{\prime}-b^{2}}{4 a^{2}} \tag{2.9}
\end{equation*}
$$

Clearly, $\check{\phi}(r):=h(r) \phi(r)$ solves the equation

$$
\begin{equation*}
\left(\partial_{r}^{2}+I(r)\right) \check{\phi}(r)=0 \tag{2.10}
\end{equation*}
$$

(2.10) will be called the canonical form of the equation (2.7). (2.9) will be called the Bose invariant of (2.7) Bos. (The name "Bose invariant" is used eg. by Milson Mil, however the object itself was clearly known before Bose).

It is easy to check that for any functions $f, g$ the equation

$$
\begin{equation*}
f(r) \mathcal{A}\left(r, \partial_{r}\right) g(r) \phi(r)=0 \tag{2.11}
\end{equation*}
$$

has the same Bose invariant as (2.7). Conversely, if an equation has the same Bose invariant as (2.7), it equals (2.11) for some $f, g$.

### 2.2 Schwarz derivative

The Schwarz derivative of $r \mapsto r(y)$ appears naturally in the context of 2 nd order ordinary differential equations and was known to XIX century mathematicians. It was mentioned already by Kummer in 1836 Kum. The universally adopted name Schwarz derivative was introduced by A. Cayley, who derived its basic properties in 1880 Cay. For more historical comments together with numerous applications we refer the reader to a book by Osgood Osg.

The Schwarz derivative is defined as

$$
\begin{equation*}
\{r, y\}=\left(\frac{r^{\prime \prime}(y)}{r^{\prime}(y)}\right)^{\prime}-\frac{1}{2}\left(\frac{r^{\prime \prime}(y)}{r^{\prime}(y)}\right)^{2} \tag{2.12}
\end{equation*}
$$

Here are some basic properties of the Schwarz derivative:
Proposition 2.1 1. Let $y \mapsto s(y), s \mapsto r(s)$ be two functions. Then

$$
\begin{equation*}
\{r, y\}=\{r, s\} s^{\prime}(y)^{2}+\{s, y\} \tag{2.13}
\end{equation*}
$$

2. Let $r \mapsto y(r)$ be the inverse of $y \mapsto r(y)$. Then

$$
\begin{equation*}
\{y, r\}=-\left[r^{\prime}(y)\right]^{-2}\{r, y\} \tag{2.14}
\end{equation*}
$$

3. $\{r, y\} \equiv 0 \Leftrightarrow r(y)=\frac{a y+b}{c y+d}, a d-b c \neq 0$

We will need the following fact:
Proposition 2.2 Suppose that $z^{\prime}=z^{p}(1-z)^{q}$. Then

$$
\begin{equation*}
\{z, r\}=\frac{1}{2} z^{2 p-2}(1-z)^{2 q-2}\left(\left(p^{2}-2 p\right)(1-z)^{2}+\left(q^{2}-2 q\right) z^{2}-2 p q z(1-z)\right) \tag{2.15}
\end{equation*}
$$

### 2.3 The Liouville transformation

Consider a 2 nd order equation in the canonical form, that is

$$
\begin{equation*}
\left(\partial_{r}^{2}+I(r)\right) \phi(r)=0 . \tag{2.16}
\end{equation*}
$$

Let us make the transformation $r=r(y)$ in this equation. We obtain

$$
\begin{equation*}
\frac{1}{\left[r^{\prime}(y)\right]^{2}} \partial_{r}^{2} \phi(r(y))-\frac{r^{\prime \prime}(y)}{\left[r^{\prime}(y)\right]^{3}} \partial_{r} \phi(r(y))+I(r(y)) \phi(r(y))=0 \tag{2.17}
\end{equation*}
$$

The resulting equation we transform again to its canonical form:

$$
\begin{equation*}
\left(\partial_{y}^{2}+J(y)\right) \psi(y)=0, \tag{2.18}
\end{equation*}
$$

where

$$
\begin{gather*}
\psi(y)=\frac{1}{\left[r^{\prime}(y)\right]^{2}} \phi(r(y)),  \tag{2.19}\\
J(y)=\left[r^{\prime}(y)\right]^{2} I(r(y))+\frac{1}{2}\{r, y\}, \tag{2.20}
\end{gather*}
$$

The above procedure is called the Liouville transformation of (2.16) by the change of variables $r=r(y)$ Lio. One can check using (2.17), that the composition of two Liouville transformations, first by the change of variables $r=r(y)$, then $y=y(z)$, is the Liouville transformation by the change of variables $r=r(y(z))$.

### 2.4 Natanzon's problem

Let $\mathcal{A}\left(a_{1}, \ldots, a_{k} ; z, \partial_{z}\right)$ be a 2 nd order differential operator depending on $k$ parameters, such that its Bose invariant can be written as

$$
\begin{equation*}
I\left(b_{1}, \ldots, b_{n} ; z\right)=b_{1} I_{1}(z)+\ldots+b_{n} I_{n}(z) \tag{2.21}
\end{equation*}
$$

for some (at least one) $n \leq k$ linearly independent functions $I_{1}, \ldots, I_{n}$ (independent of $a_{1}, \ldots, a_{k}$ ), and some $n$ numbers $b_{i}\left(a_{1}, \ldots, a_{k}\right)$. Then one can solve the following problem (which we will call Natanzon's problem for $\left.\mathcal{A}\left(z, \partial_{z}\right)\right)$ :

Find all potentials $V(r)$ such that the 1-dimensional stationary Schrödinger equation

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V(r)-E\right) \phi(E, r)=0 \tag{2.22}
\end{equation*}
$$

can be transformed to the equation given by $\mathcal{A}\left(a_{1}, \ldots, a_{k} ; z, \partial_{z}\right)$ for some $a_{1}(E), \ldots, a_{k}(E)$. We allow the following operations:
(1) multiplication of both sides of the equation by some $f(E, r)$;
(2) substitution of $\check{\phi}(E, r):=g(E, r) \phi(E, r)$ for some $g(E, r)$;
(3) change of coordinates $r \mapsto z$ independent of $E$.

This is the problem solved by Natanzon in the case of the hypergeometric equation and by Milson in the general case of hypergeometric type equations. In the following, we recall the construction they used, in a slightly generalized form.

Let us consider an arbitrary equation of the form (2.22). Obviously, its Bose invariant equals $E-V(r)$. Clearly, the transformations (1) and (2) allow us to transform $\mathcal{A}\left(a_{1}, \ldots, a_{k} ; z, \partial_{z}\right)$ to its canonical form

$$
\begin{equation*}
-\partial_{z}^{2}-I\left(b_{1}, \ldots, b_{n} ; z\right) . \tag{2.23}
\end{equation*}
$$

Thus (2.22) can be transformed to $\mathcal{A}\left(a_{1}, \ldots, a_{k} ; z, \partial_{z}\right)$ if

$$
\begin{equation*}
I\left(b_{1}, \ldots, b_{n} ; z\right)=\left(r^{\prime}(z)\right)^{2}\left(E-V(r(z))+\frac{1}{2}\{r, z\}\right. \tag{2.24}
\end{equation*}
$$

that is, if the two Bose invariants are related by a Liouville transformation. Using (2.21) we rewrite (2.24) as

$$
\begin{equation*}
b_{1}(E) I_{1}(z)+\ldots+b_{n}(E) I_{n}(z)=\left(r^{\prime}(z)\right)^{2}\left(E-V(r(z))+\frac{1}{2}\{r, z\}\right. \tag{2.25}
\end{equation*}
$$

By assumption, $b_{i}$ depend on $E$ and $I_{i}(z)$ are independent of $E$. Thus, the dependence of $b_{i}$ on $E$ is linear. Therefore, by transforming linearly $\left(b_{1}, \ldots, b_{n}\right)$ into $\left(\tilde{b}_{1}, \ldots, \tilde{b}_{n-1}, E\right)$, we can assume that

$$
I\left(b_{1}, \ldots, b_{n} ; z\right)=\tilde{b}_{1} \tilde{I}_{1}(z)+\ldots+\tilde{b}_{n-1} \tilde{I}_{n-1}(z)+E \tilde{I}_{n}(z)
$$

for some functions $\tilde{I}_{i}(z)$, which are linear combinations of $I_{i}(z)$. Then we can write

$$
\begin{equation*}
E \tilde{I}_{n}(z)+\tilde{b}_{1} \tilde{I}_{1}(z)+\ldots+\tilde{b}_{n-1} \tilde{I}_{n-1}(z)=\left(r^{\prime}(z)\right)^{2}\left(E-V(r(z))+\frac{1}{2}\{r, z\}\right. \tag{2.26}
\end{equation*}
$$

Therefore, Natanzon's problem is solved by the following pair of equations

$$
\begin{align*}
V\left(\tilde{b}_{1}, \ldots, \tilde{b}_{n-1}, r(z)\right) & :=\left(r^{\prime}(z)\right)^{-2}\left(\tilde{b}_{1} \tilde{I}_{1}(z)+\ldots+\tilde{b}_{n-1} \tilde{I}_{n-1}(z)-\frac{1}{2}\{r, z\}\right) \\
\left(r^{\prime}(z)\right)^{2} & =\tilde{I}_{n}(z) \tag{2.27}
\end{align*}
$$

We can rewrite (2.27) in an equivalent way as

$$
\begin{align*}
V\left(\tilde{b}_{1}, \ldots, \tilde{b}_{n-1}, r\right) & :=\left(z^{\prime}(r)\right)^{2}\left(\tilde{b}_{1} \tilde{I}_{1}(z(r))+\ldots+\tilde{b}_{n-1} \tilde{I}_{n-1}(z(r))\right)+\frac{1}{2}\{z, r\} \\
\left(z^{\prime}(r)\right)^{-2} & =\tilde{I}_{n}(z(r)) \tag{2.28}
\end{align*}
$$

If desired, we can renormalize $E$ by substracting a constant.
Note that in practical applications the value of the above method depends on whether it is possible to invert the relation $r \mapsto z$ expressing it in terms of standard functions. In the next sections, we will describe such solutions of the Natanzon problem for the three classes of equations described in the introduction.

## 3 Schrödinger operators reducible to the hypergeometric equation

### 3.1 Hypergeometric equation

The hypergeometric equation is given by the operator

$$
\begin{equation*}
\mathcal{F}\left(a, b ; c ; z, \partial_{z}\right):=z(1-z) \partial_{z}^{2}+(c-(a+b+1) z) \partial_{z}-a b \tag{3.29}
\end{equation*}
$$

where $a, b, c$ are arbitrary complex parameters. Its Bose invariant equals

$$
{ }_{2} I_{1}(a, b ; c ; z)=\frac{\left(1-a^{2}-b^{2}-2 a b\right) z^{2}+2(-2 a b+a c+b c-c) z+2 c-c^{2}}{4 z^{2}(1-z)^{2}}
$$

It satisfies condition (2.21) for $n=3$. The functions $I_{i}(z), i=1,2,3$ can be taken to be $\frac{1}{4 z^{2}(1-z)}$ $\frac{1}{4 z(1-z)^{2}}$ and $\frac{1}{4 z(1-z)}$. It is natural to demand that $\left(r^{\prime}(z)\right)^{2}=\tilde{I}_{3}(z)$ is a function of the form $z^{2 p}(1-z)^{2 q}$. It is easy to see that this gives the following possibilities for $\tilde{I}_{3}(z)$ :
(1) $\frac{1}{4 z(1-z)}$,
(2) $\frac{1}{4 z^{2}(1-z)}$,
(3) $\frac{1}{4 z(1-z)^{2}}$,
(4) $\frac{1}{4 z^{2}(1-z)^{2}}$,
(5) $\frac{1}{4 z^{2}}, \quad(6) \frac{1}{4(1-z)^{2}}$.

We will see that each of these ansatzes leads to exactly solvable potentials considered in the literature. Note also that the formula for the Schwarz derivative (2.15) will be handy.

### 3.2 The Riemann equation

It is well known that it is useful to consider the hypergeometric equation as a special case of the so-called Riemann equation. The Riemann equation is defined as the class of equations on the Riemann sphere $\overline{\mathbb{C}}$ having 3 regular singular points. The following theorem summarizes the basic theory of these equations:

Theorem 3.1 (1) Suppose that we are given a 2nd order differential equation on the Riemann sphere having 3 singular points $z_{1}, z_{2}, z_{3}$, all of them regular singular points with the following indices

$$
\begin{array}{ll}
z_{1}: & \rho_{1}, \tilde{\rho}_{1} \\
z_{2}: & \rho_{2}, \tilde{\rho}_{2} \\
z_{3}: & \rho_{3}, \tilde{\rho}_{3}
\end{array}
$$

Then the following condition is satisfied:

$$
\begin{equation*}
\rho_{1}+\tilde{\rho}_{1}+\rho_{2}+\tilde{\rho}_{2}+\rho_{3}+\tilde{\rho}_{3}=1 \tag{3.31}
\end{equation*}
$$

Such an equation, normalized to have coefficient 1 at the 2nd derivative, is always equal to

$$
\mathcal{P}\left(\begin{array}{llll}
z_{1} & z_{2} & z_{3} &  \tag{3.32}\\
\rho_{1} & \rho_{2} & \rho_{3} & z, \partial_{z} \\
\tilde{\rho}_{1} & \tilde{\rho}_{2} & \tilde{\rho}_{3} &
\end{array}\right) \phi(z)=0
$$

where

$$
\begin{aligned}
& \mathcal{P}\left(\begin{array}{lll}
z_{1} & z_{2} & z_{3} \\
\rho_{1} & \rho_{2} & \rho_{3} \\
\tilde{\rho}_{1} & \tilde{\rho}_{2} & \tilde{\rho}_{3}
\end{array}\right), \partial_{z} \\
&+\frac{\rho_{1} \tilde{\rho}_{1}\left(z_{1}-z_{2}\right)\left(z_{1}-z_{3}\right)}{\left(z-z_{1}\right)^{2}\left(z-z_{2}\right)\left(z-z_{3}\right)}+\frac{\rho_{2} \tilde{\rho}_{2}\left(z_{2}-z_{3}\right)\left(z_{2}-z_{1}\right)}{\left(z-z_{2}\right)^{2}\left(z-z_{3}\right)\left(z-z_{1}\right)}+\frac{\rho_{3} \tilde{\rho}_{3}\left(z_{3}-z_{1}\right)\left(z_{3}-z_{2}\right)}{\left(z-z_{3}\right)^{2}\left(z-z_{1}\right)\left(z-z_{2}\right)}
\end{aligned}
$$

(2) Let $z \mapsto w=h(z)=\frac{a z+b}{c z+d}$. (Transformations of this form are called homographies or Möbius transformations). We can always assume that $a d-b c=1$. Then

$$
\mathcal{P}\left(\begin{array}{cccc}
h\left(z_{1}\right) & h\left(z_{2}\right) & h\left(z_{3}\right) & \\
\rho_{1} & \rho_{2} & \rho_{3} & w, \partial_{w} \\
\tilde{\rho}_{1} & \tilde{\rho}_{2} & \tilde{\rho}_{3} &
\end{array}\right)=(c z+d)^{4} \mathcal{P}\left(\begin{array}{cccc}
z_{1} & z_{2} & z_{3} & \\
\rho_{1} & \rho_{2} & \rho_{3} & z, \partial_{z} \\
\tilde{\rho}_{1} & \tilde{\rho}_{2} & \tilde{\rho}_{3} &
\end{array}\right)
$$

(3)

$$
\begin{aligned}
& \left(z-z_{1}\right)^{-\lambda}\left(z-z_{2}\right)^{\lambda} \mathcal{P}\left(\begin{array}{cccc}
z_{1} & z_{2} & z_{3} & \\
\rho_{1} & \rho_{2} & \rho_{3} & z, \partial_{z} \\
\tilde{\rho}_{1} & \tilde{\rho}_{2} & \tilde{\rho}_{3} &
\end{array}\right)\left(z-z_{1}\right)^{\lambda}\left(z-z_{2}\right)^{-\lambda} \\
& =\mathcal{P}\left(\begin{array}{cccc}
z_{1} & z_{2} & z_{3} & \\
\rho_{1}-\lambda & \rho_{2}+\lambda & \rho_{3} & z, \partial_{z} \\
\tilde{\rho}_{1}-\lambda & \tilde{\rho}_{2}+\lambda & \tilde{\rho}_{3} &
\end{array}\right)
\end{aligned}
$$

Clearly, in all above formulas one of $z_{i}$ can equal $\infty$, with an obvious meaning of various expressions. For convenience we give the expression for the Riemann operator with $z_{3}=\infty$ :

$$
\begin{align*}
& \mathcal{P}\left(\begin{array}{ccc}
z_{1} & z_{2} & \infty \\
\rho_{1} & \rho_{2} & \rho_{3} \\
\tilde{\rho}_{1} & \tilde{\rho}_{2} & \tilde{\rho}_{3}
\end{array}\right) \\
= & \partial_{z}^{2}-\left(\frac{\rho_{1}+\tilde{\rho}_{1}-1}{z-z_{1}}+\frac{\rho_{2}+\tilde{\rho}_{2}-1}{z-z_{2}}\right) \partial_{z} \\
& +\frac{\rho_{1} \tilde{\rho}_{1}\left(z_{1}-z_{2}\right)}{\left(z-z_{1}\right)^{2}\left(z-z_{2}\right)}+\frac{\rho_{2} \tilde{\rho}_{2}\left(z_{2}-z_{1}\right)}{\left(z-z_{2}\right)^{2}\left(z-z_{1}\right)}+\frac{\rho_{3} \tilde{\rho}_{3}}{\left(z-z_{1}\right)\left(z-z_{2}\right)} \tag{3.33}
\end{align*}
$$

The hypergeometric equation is a special case of the Riemann equation, since

$$
\begin{aligned}
\mathcal{F}\left(a, b ; c ; z, \partial_{z}\right) & =z(1-z) \mathcal{P}\left(\begin{array}{cccc}
0 & 1 & \infty & \\
0 & 0 & a & z, \partial_{z} \\
1-c & c-a-b & b
\end{array}\right) \\
& =z(1-z) \partial_{z}^{2}+(c-(a+b+1) z) \partial_{z}-a b
\end{aligned}
$$

Every Riemann equation can be brought to the form of the hypergeometric equation by applying (1) and (2) of Theorem 3.1

From Theorem 3.1 we also see that symmetries of the hypergeometric equation are better visible if we replace $a, b, c$ with $\alpha, \beta, \mu$ :

$$
\begin{array}{lll}
\alpha=c-1 & \beta=a+b-c, & \mu=a-b \\
a=\frac{1+\alpha+\beta-\mu}{2}, & b=\frac{1+\alpha+\beta+\mu}{2}, & c=1+\alpha
\end{array}
$$

In fact, the new parameters coincide with the differences of the indices of the points $0,1, \infty$ :

$$
\alpha=\rho_{1}-\tilde{\rho}_{1}, \beta=\rho_{2}-\tilde{\rho}_{2}, \quad \mu=\rho_{3}-\tilde{\rho}_{3} .
$$

In the parameters $\alpha, \beta, \mu$, the Bose invariant of the hypergeometric equation has a more symmetric form:

$$
{ }_{2} I_{1}(\alpha, \beta, \mu ; z)=\frac{\left(1-\alpha^{2}\right)(1-z)+\left(1-\beta^{2}\right) z+\left(\mu^{2}-1\right)(1-z) z}{4 z^{2}(1-z)^{2}}
$$

We can summarize the relation between the hypergeometric equation and its canonical form by

$$
\begin{align*}
& -z^{\frac{\alpha}{2}+\frac{1}{2}-1}(1-z)^{\frac{\beta}{2}+\frac{1}{2}-1} \mathcal{F}\left(\frac{\alpha+\beta+\mu+1}{2}, \frac{\alpha+\beta-\mu+1}{2} ; 1+\alpha ; z, \partial_{z}\right) z^{-\frac{\alpha}{2}-\frac{1}{2}}(1-z)^{-\frac{\beta}{2}-\frac{1}{2}} \\
= & -\mathcal{P}\left(\begin{array}{ccc}
0 & 1 & \infty \\
\frac{\alpha}{2}+\frac{1}{2} & \frac{\beta}{2}+\frac{1}{2} & \frac{\mu}{2}-\frac{1}{2} \\
-\frac{\alpha}{2}+\frac{1}{2} & -\frac{\beta}{2}+\frac{1}{2} & -\frac{\mu}{2}-\frac{1}{2}
\end{array}\right) \\
= & -\partial_{z}^{2}+\left(\alpha^{2}-1\right) \frac{1}{4 z^{2}(1-z)}+\left(\beta^{2}-1\right) \frac{1}{4 z(1-z)^{2}}-\left(\mu^{2}-1\right) \frac{1}{4 z(1-z)} . \tag{3.34}
\end{align*}
$$

It will be natural to introduce the parameters

$$
\begin{equation*}
\kappa:=\frac{1}{2}\left(\alpha^{2}-\beta^{2}\right), \quad \delta:=\frac{1}{2}\left(\alpha^{2}+\beta^{2}\right) \tag{3.35}
\end{equation*}
$$

In some cases the parameter $\kappa$ will be replaced by

$$
\begin{equation*}
\tau:=\frac{\mathrm{i}}{2}\left(\alpha^{2}-\beta^{2}\right)=\mathrm{i} \kappa \tag{3.36}
\end{equation*}
$$

We will describe two complex classes of exactly solvable potentials depending on two complex parameters. Within each complex class there will be three real classes of exactly solvable potentials depending on two real parameters.

### 3.3 Trigonometric Pöschl-Teller potential

In this subsection we consider Ansatz (1) of (3.30). We set

$$
\begin{equation*}
z=\sin ^{2} \frac{r}{2}, \quad \text { which solves } z^{\prime}=z^{\frac{1}{2}}(1-z)^{\frac{1}{2}} \tag{3.37}
\end{equation*}
$$

This leads to the Schrödinger equation

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V_{\delta, \kappa}^{\mathrm{tPT}}(r)-\frac{\mu^{2}}{4}\right) \phi(r)=0 \tag{3.38}
\end{equation*}
$$

where

$$
\begin{align*}
V_{\delta, \kappa}^{\mathrm{tPT}}(r) & :=\left(\alpha^{2}-\frac{1}{4}\right) \frac{1}{4 \sin ^{2} \frac{r}{2}}+\left(\beta^{2}-\frac{1}{4}\right) \frac{1}{4 \cos ^{2} \frac{r}{2}}  \tag{3.39}\\
& =\left(\delta-\frac{1}{4}\right) \frac{1}{\sin ^{2} r}+\kappa \frac{\cos r}{\sin ^{2} r}
\end{align*}
$$

This potential was proposed and solved by G. Pöschl and E. Teller [PT]. It is usually called the PöschlTeller potential, sometimes also the Pöschl-Teller potential of the first kind or the trigonometric Scarf potential.

A natural real domain for this potential is $] 0, \pi[$. If $\kappa, \delta$ are real, then the potential is real on this domain.

Explicitly, the reduction of (3.38) to the hypergeometric equation is derived as follows:

$$
\left.\left.\begin{array}{rl} 
& -z^{\frac{\alpha}{2}+\frac{1}{4}}(1-z)^{\frac{\beta}{2}+\frac{1}{4}} \mathcal{F}\left(\frac{\alpha+\beta+\mu+1}{2}, \frac{\alpha+\beta-\mu+1}{2} ; 1+\alpha ; z, \partial_{z}\right) z^{-\frac{\alpha}{2}-\frac{1}{4}}(1-z)^{-\frac{\beta}{2}-\frac{1}{4}} \\
= & -z(1-z) \mathcal{P}\left(\begin{array}{ccc}
0 & 1 & \infty \\
\frac{\alpha}{2}+\frac{1}{4} & \frac{\beta}{2}+\frac{1}{4} & \frac{\mu}{2} \\
-\frac{\alpha}{2}+\frac{1}{4} & -\frac{\beta}{2}+\frac{1}{4} & -\frac{\mu}{2}
\end{array}\right) \\
= & -z(1-z)\left(\partial_{z}\right.
\end{array}\right) .\left(\frac{1}{2 z}-\frac{1}{2(1-z)}\right) \partial_{z}\right) .
$$

### 3.4 Hyperbolic Pöschl-Teller potential

We continue with (1) of (3.30). We set

$$
\begin{equation*}
z=-\sinh ^{2} \frac{r}{2}, \quad \text { which solves } z^{\prime}=-(-z)^{\frac{1}{2}}(1-z)^{\frac{1}{2}} \tag{3.41}
\end{equation*}
$$

This leads to the Schrödinger equation

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V_{\delta, \kappa}^{\mathrm{hPT}}(r)+\frac{\mu^{2}}{4}\right) \phi(r)=0 \tag{3.42}
\end{equation*}
$$

where

$$
\begin{align*}
V_{\delta, \kappa}^{\mathrm{hPT}}(r) & :=\left(\alpha^{2}-\frac{1}{4}\right) \frac{1}{4 \sinh ^{2} \frac{r}{2}}-\left(\beta^{2}-\frac{1}{4}\right) \frac{1}{4 \cosh ^{2} \frac{r}{2}}  \tag{3.43}\\
& =\left(\delta-\frac{1}{4}\right) \frac{1}{\sinh ^{2} r}+\kappa \frac{\cosh r}{\sinh ^{2} r}
\end{align*}
$$

This potential was also proposed and solved by G. Pöschl and E. Teller PT. In the literature it is known as hyperbolic, generalized Pöschl-Teller potential, or the Pöschl-Teller potential of the second kind.

A natural real domain for this potential is $] 0, \infty[$. If $\kappa, \delta$ are real, then the potential is real on this domain.

To see that (3.42) can be solved in terms of the hypergeometric equation, we first repeat the computations leading to (3.40), and then set $z=-\sinh ^{2} \frac{r}{2}$.

### 3.5 Scarf potential

We continue with (1) of (3.30). We set

$$
\begin{equation*}
z=\frac{1}{2}-\mathrm{i} \cosh \frac{r}{2} \sinh \frac{r}{2}, \quad \text { which solves } z^{\prime}=(-z)^{\frac{1}{2}}(1-z)^{\frac{1}{2}} \tag{3.44}
\end{equation*}
$$

This leads to the Schrödinger equation

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V_{\delta, \tau}^{\mathrm{S}}(r)+\frac{\mu^{2}}{4}\right) \phi(r)=0 \tag{3.45}
\end{equation*}
$$

where

$$
V_{\delta, \tau}^{\mathrm{S}}(r):=-\left(\delta-\frac{1}{4}\right) \frac{1}{\cosh ^{2} r}-\tau \frac{\sinh r}{\cosh ^{2} r}
$$

This potential was proposed and solved by F. Scarf [Sca. In the literature it is often called the hyperbolic Scarf potential.

A natural real domain for this potential is $]-\infty, \infty[$. If $\delta, \tau$ are real, then the potential is real on this domain.

To see that (3.45) can be solved in terms of the hypergeometric equation, we first repeat the computations leading to (3.40), and then set $z=\frac{1}{2}-\mathrm{i} \cosh \frac{r}{2} \sinh \frac{r}{2}$.

### 3.6 More about Pöschl-Teller and Scarf potentials

Both kinds of Pöschl-Teller potentials and the Scarf potential are real cases of the same complex case. To see this, consider eg. the hyperbolic Pöschl-Teller potential as a function of the complex parameter $r$. It is holomorphic away from singularities at $\mathrm{i} \pi n, n \in \mathbb{Z}$.

For real $\delta, \kappa, V_{\delta, \kappa}^{\mathrm{hPT}}(r)$ is real on $\mathrm{i} \mathbb{R}$ and $\mathbb{R}+\mathrm{i} \pi n$. On each halfline $] 0, \infty[+\mathrm{i} \pi n$ and $]-\infty, 0[+\mathrm{i} \pi n$ we obtain the hyperbolic Pöschl-Teller potential. On each interval ]i $\pi n, \mathrm{i} \pi(n+1)$ [ we obtain the trigonometric Pöschl-Teller potential.

For real $\delta, \mathrm{i} \kappa,, V_{\delta, \kappa}^{\mathrm{hPT}}(r)$ is real on $\mathbb{R}+\mathrm{i} \pi\left(n+\frac{1}{2}\right)$. On each of these lines we obtain the Scarf potential.
Above, we used Ansatz (1) to derive Scarf and both kinds of Pöschl-Teller potentials. Alternatively, one can use Ansatzes (2) or (3). To see this it is enough to consider Ansatz (3). In fact, we first repeat computations analogous to (3.40):

$$
\begin{aligned}
& -z(1-z)^{2} \mathcal{P}\left(\begin{array}{ccc}
0 & 1 & \infty \\
\frac{\alpha}{2}+\frac{1}{4} & \frac{\beta}{2} & \frac{\mu}{2}+\frac{1}{4} \\
-\frac{\alpha}{2}+\frac{1}{4} & -\frac{\beta}{2} & -\frac{\mu}{2}+\frac{1}{4}
\end{array}\right) \\
=- & -z(1-z)^{2}\left(\partial_{z}^{2}+\left(\frac{1}{2 z}-\frac{1}{1-z}\right) \partial_{z}\right) \\
& +\left(\alpha^{2}-\frac{1}{4}\right) \frac{1-z}{4 z}+\frac{\beta^{2}}{4}-\left(\mu^{2}-\frac{1}{4}\right) \frac{(1-z)}{4} .
\end{aligned}
$$

We set

$$
\begin{equation*}
z=\operatorname{tgh}^{2} \frac{r}{2}, \quad \text { which solves } z^{\prime}=z^{\frac{1}{2}}(1-z) \tag{3.46}
\end{equation*}
$$

obtaining the hyperbolic Pöschl-Teller potential,

$$
\begin{equation*}
z=-\operatorname{tg}^{2} \frac{r}{2}, \quad \text { which solves } \quad z^{\prime}=-(-z)^{\frac{1}{2}}(1-z) \tag{3.47}
\end{equation*}
$$

obtaining the trigonometric Pöschl-Teller potential, or

$$
\begin{equation*}
z=\operatorname{ctgh}^{2} \frac{r}{2}, \quad \text { which solves } z^{\prime}=z^{\frac{1}{2}}(1-z) \tag{3.48}
\end{equation*}
$$

obtaining the Scarf potential.

### 3.7 Manning-Rosen potential

Let us consider Ansatz (3) of (3.30). We set

$$
\begin{equation*}
z=\frac{1}{1+\mathrm{e}^{2 r}}, \quad \text { which solves } z^{\prime}=2 z(z-1) \tag{3.49}
\end{equation*}
$$

This leads to the Schrödinger equation

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V_{\kappa, \mu}^{\mathrm{MR}}(r)+\delta\right) \phi(r)=0 \tag{3.50}
\end{equation*}
$$

where

$$
V_{\kappa, \mu}^{\mathrm{MR}}(r):=-\kappa \frac{\sinh r}{\cosh r}-\left(\frac{\mu^{2}}{4}-\frac{1}{4}\right) \frac{1}{\cosh ^{2} r}
$$

This potential was proposed and solved by M. F. Manning and N. Rosen MR. In the literature it is also called the Woods-Saxon potential SW] (for instance in [Flü]), also the hyperbolic Rosen-Morse potential.

A natural real domain for this potential is $]-\infty, \infty[$. The potential is real if $r \in]-\infty, \infty\left[\right.$ and $\kappa, \beta^{2}$ are real.

Here is a derivation of (3.50) from the hypergeometric equation:

$$
\begin{align*}
& -4 z^{1+\frac{\alpha}{2}}(1-z)^{1+\frac{\beta}{2}} \mathcal{F}\left(\frac{\alpha+\beta+\mu+1}{2}, \frac{\alpha+\beta-\mu+1}{2} ; 1+\alpha ; z, \partial_{z}\right) z^{-\frac{\alpha}{2}}(1-z)^{-\frac{\beta}{2}} \\
= & -4 z^{2}(z-1)^{2} \mathcal{P}\left(\begin{array}{ccc}
0 & 1 & \infty \\
\frac{\alpha}{2} & \frac{\beta}{2} & \mu+\frac{1}{2} \\
-\frac{\alpha}{2} & -\frac{\beta}{2} & -\mu+\frac{1}{2}
\end{array}\right) \\
= & -4 z^{2}(1-z)^{2}\left(\partial_{z}^{2}+\left(\frac{1}{z}-\frac{1}{1-z}\right) \partial_{z}\right)+\alpha^{2}(1-z)+\beta^{2} z-\left(\mu^{2}-1\right) z(1-z)  \tag{3.51}\\
= & -\partial_{r}^{2}+\delta+\kappa \frac{\mathrm{e}^{2 r}-1}{1+\mathrm{e}^{2 r}}-\left(\mu^{2}-1\right) \frac{\mathrm{e}^{2 r}}{\left(1+\mathrm{e}^{2 r}\right)^{2}} . \tag{3.52}
\end{align*}
$$

### 3.8 Eckart potential

We still consider Ansatz (3) of (3.30).
We set

$$
\begin{equation*}
z=\frac{1}{1-\mathrm{e}^{-2 r}}, \quad \text { which solves } \quad z^{\prime}=2 z(1-z) \tag{3.53}
\end{equation*}
$$

and use (3.35). This leads to the Schrödinger equation

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V_{\kappa, \mu}^{\mathrm{E}}(r)+\delta\right) \phi(r)=0 \tag{3.54}
\end{equation*}
$$

where

$$
V_{\kappa, \mu}^{\mathrm{E}}(y):=-\kappa \frac{\cosh r}{\sinh r}+\left(\frac{\mu^{2}}{4}-\frac{1}{4}\right) \frac{1}{\sinh ^{2} r}
$$

This potential was proposed and solved by C. Eckart [Eck]. In the literature (for instance [Flü]) it is also called the Hulthen potential Hul, sometimes also the generalized Morse potential, because of its similarity to the Morse potential, see Subsect. 4.4.

A natural real domain for this potential is $] 0, \infty\left[\right.$. If $\kappa, \beta^{2}$ are real, then the potential is real on this domain.

To derive the Eckart potential, we first repeat the computations leading to (3.51), and then set $z=\frac{1}{1-\mathrm{e}^{-2 r}}$.

### 3.9 Rosen-Morse potential

Once again, we consider Ansatz (3) of (3.30).
We set

$$
\begin{equation*}
z=\frac{1}{1-\mathrm{e}^{2 \mathrm{i} r}}, \quad \text { which solves } z^{\prime}=2 \mathrm{i} z(1-z) \tag{3.55}
\end{equation*}
$$

This leads to the Schrödinger equation

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V_{\tau, \mu}^{\mathrm{RM}}(r)-\delta\right) \phi(r)=0 \tag{3.56}
\end{equation*}
$$

where

$$
V_{\tau, \mu}^{\mathrm{RM}}(y):=\tau \frac{\cos r}{\sin r}+\left(\frac{\mu^{2}}{4}-\frac{1}{4}\right) \frac{1}{\sin ^{2} r}
$$

This potential is known as the Rosen-Morse potential, also the trigonometric Rosen-Morse potential (altough this name is widely used in the literature, we were unable to explain decisively its origin).

A natural real domain for this potential is $] 0, \pi\left[\right.$. If $\tau, \mu^{2}$ are real, then the potential is real on this domain.

To derive the Rosen-Morse potential, we first repeat the computations leading to (3.51), then set $z=\frac{1}{1-\mathrm{e}^{2 \mathrm{i} r}}$.

### 3.10 More about Manning-Rosen, Eckart and Rosen-Morse potentials

The Manning-Rosen, Eckart and Rosen-Morse potentials are all real cases of the same complex case. To see this, consider eg. the Eckart potential as a function of the complex parameter $r$. It is holomorphic away from singularities at $\mathrm{i} \pi n, n \in \mathbb{Z}$.

For real $\kappa, \mu^{2}, V_{\kappa, \mu}^{\mathrm{E}}(r)$ is real on the lines $\mathbb{R}+\frac{\mathrm{i} \pi n}{2}$. On each line $\mathbb{R}+\mathrm{i}\left(m+\frac{1}{2}\right) \pi, m \in \mathbb{Z}$, we obtain the Manning-Rosen potential. On each halfline $] 0, \infty[+\mathrm{i} \pi m$ and $]-\infty, 0[+\mathrm{i} \pi m$, we obtain the Eckart potential.

For real $\mathrm{i} \kappa, \mu^{2}, V_{\kappa, \mu}^{\mathrm{E}}(r)$ is real on $\mathbb{i} \mathbb{R}$. On each interval $] \mathrm{i} \pi m, \mathrm{i} \pi(m+1)$ [ we obtain the Rosen-Morse potential.

Ansatzes (5) and (6) lead to the same classes of exactly solvable potentials as Ansatz (4). To see this it is enough to consider Ansatz (5). In fact, we first repeat computations analogous to (3.51):

$$
\begin{aligned}
& -4 z^{-\frac{\alpha+1}{2}}(1-z)^{2+\frac{\beta}{2}} \mathcal{F}\left(\frac{\alpha+\beta+\mu+1}{2}, \frac{\alpha+\beta+-\mu+1}{2} ; 1+\alpha ; z, \partial_{z}\right) z^{-\frac{\alpha+1}{2}}(1-z)^{-\frac{\beta}{2}} \\
= & -4(1-z)^{2} \mathcal{P}\left(\begin{array}{cccc}
0 & 1 & \infty \\
\alpha+\frac{1}{2} & \frac{\beta}{2} & \frac{\mu}{2} & z, \partial_{z} \\
-\alpha+\frac{1}{2} & -\frac{\beta}{2} & -\frac{\mu}{2}
\end{array}\right) \\
= & -4(1-z)^{2}\left(\partial_{z}^{2}-\frac{1}{1-z} \partial_{z}\right) \\
& +\left(\alpha^{2}-1\right) \frac{1-z}{z^{2}}+\beta^{2} \frac{1}{z}-\mu^{2} \frac{1-z}{z} .
\end{aligned}
$$

We set

$$
\begin{equation*}
z=1-\mathrm{e}^{2 r}, \quad \text { which solves } \quad z^{\prime}=-2(1-z) \tag{3.57}
\end{equation*}
$$

obtaining the Eckart potential,

$$
\begin{equation*}
z=1+\mathrm{e}^{2 r}, \quad \text { which solves } z^{\prime}=2(1-z) \tag{3.58}
\end{equation*}
$$

obtaining the Manning-Rosen potential,

$$
\begin{equation*}
z=1-\mathrm{e}^{\mathrm{i} 2 r}, \quad \text { which solves } \quad z^{\prime}=-2 \mathrm{i}(1-z) \tag{3.59}
\end{equation*}
$$

obtaining the Rosen-Morse potential,

## 4 Schrödinger operators reducible to the rescaled confluent equation

One of basic exactly solvable equations is the confluent equation, given by the operator

$$
\begin{equation*}
\mathcal{F}\left(a ; c ; z, \partial_{z}\right):=z \partial_{z}^{2}+(c-z) \partial_{z}-a \tag{4.60}
\end{equation*}
$$

It is convenient to consider (4.60) in parallel with the equation given by the operator

$$
\begin{equation*}
\mathcal{F}\left(a, b ;-; z, \partial_{z}\right):=z^{2} \partial_{z}^{2}+(-1+(1+a+b) z) \partial_{z}+a b \tag{4.61}
\end{equation*}
$$

The equation given by (4.61) is sometimes called the ${ }_{2} F_{0}$ equation. Note that

$$
z^{a} \mathcal{F}\left(a, b ;-; z, \partial_{z}\right) z^{-a}=w \mathcal{F}\left(a ; 1+a-b ; w, \partial_{w}\right), \quad w=-z^{-1}, \quad z=-w^{-1}
$$

Hence the ${ }_{2} F_{0}$ equation is equivalent to the confluent equation. The relationship between the parameters is

$$
c=1+a-b, \quad b=1+a-c
$$

Another exactly solvable equation that we will consider in this section is sometimes called the ${ }_{0} F_{1}$ equation. It is given by

$$
\begin{equation*}
\mathcal{F}\left(c ; z, \partial_{z}\right):=z \partial_{z}^{2}+(c-z) \partial_{z}-a \tag{4.62}
\end{equation*}
$$

and is equivalent to the Bessel equation.
Clearly, the confluent, ${ }_{2} F_{0}$ and ${ }_{0} F_{1}$ equations belong to the class of hypergeometric type equations.
The basic equation of this section will be (1.5), given by

$$
\begin{equation*}
\mathcal{F}\left(a ; c ; \gamma ; z ; \partial_{z}\right):=z \partial_{z}^{2}+(c-\gamma z) \partial_{z}-a \tag{4.63}
\end{equation*}
$$

It will be called the rescaled confluent equation. Note that in the case $\gamma=1$ (4.63) coincides with the confluent equation. If $\gamma \neq 0$, 4.63) can be reduced to the confluent equation by rescaling (and hence also to the ${ }_{2} F_{0}$ equation). If $\gamma=0$, (4.63) coincides with the ${ }_{0} F_{1}$ equation.

The Bose invariant of the rescaled confluent equation equals

$$
{ }_{1} I_{1}(a, b ; \gamma ; z)=\frac{-\gamma^{2} z^{2}+2(c \gamma-2 a) z+2 c-c^{2}}{4 z^{2}}
$$

### 4.1 Symmetries of the rescaled confluent equation

Let us first describe symmetries of the rescaled confluent equation.

$$
\begin{align*}
z^{c-1} \mathcal{F}\left(a ; c ; \gamma ; z, \partial_{z}\right) z^{1-c} & =\mathcal{F}\left(a-c \gamma+\gamma ; 2-c ; \gamma ; z, \partial_{z}\right)  \tag{4.64}\\
\mathrm{e}^{-\gamma z} \mathcal{F}\left(a ; c ; \gamma ; z, \partial_{z}\right) \mathrm{e}^{\gamma z} & =-\mathcal{F}\left(c-a ; c ; \gamma ; w, \partial_{w}\right), \quad z=-w \tag{4.65}
\end{align*}
$$

Besides, the scaling acts as follows:

$$
\mathcal{F}\left(a ; c ; \gamma ; z, \partial_{z}\right)=\gamma \mathcal{F}\left(a / \gamma ; c ; 1 ; w, \partial_{w}\right), \quad w=\gamma z
$$

It is convenient to introduce new parameters $\alpha, \nu$ :

$$
\begin{array}{ll}
a=\frac{\gamma+\alpha \gamma-\nu}{2}, & c=1+\alpha \\
\alpha=c-1=a-b, & \nu=c \gamma-2 a=1-a-b .
\end{array}
$$

In the new parameters the Bose invariant of the rescaled confluent equation has a more symmetric form:

$$
{ }_{1} I_{1}(\alpha, \nu, \gamma ; z)=-\frac{\gamma^{2}}{4}+\frac{\nu}{2 z}+\frac{1-\alpha^{2}}{4 z^{2}}
$$

Thus the starting point for the further analysis will be the equation

$$
\begin{equation*}
\left(-\partial_{z}^{2}+\frac{\gamma^{2}}{4}-\frac{\nu}{2 z}+\left(\frac{\alpha^{2}}{4}-\frac{1}{4}\right) \frac{1}{z^{2}}\right) \phi(z)=0 \tag{4.66}
\end{equation*}
$$

We will describe 3 classes of Schrödinger operators solved using the confluent equation corresponding to three obvious choices for $\left(r^{\prime}(z)\right)^{2}$ :

$$
\begin{equation*}
\text { (1) } \frac{1}{4}, \quad \text { (2) } \frac{1}{4 z}, \quad \text { (3) } \frac{1}{4 z^{2}} \tag{4.67}
\end{equation*}
$$

### 4.2 Hydrogen atom

We consider (1) of (4.67). We set

$$
\begin{equation*}
z=2 r \tag{4.68}
\end{equation*}
$$

This leads to the Schrödinger equation

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V_{\alpha, \nu}(r)+\gamma^{2}\right) \phi(r)=0 \tag{4.69}
\end{equation*}
$$

with the potential

$$
V_{\alpha, \nu}(r):=-\frac{\nu}{r}+\left(\frac{\alpha^{2}}{4}-\frac{1}{4}\right) \frac{1}{r^{2}}
$$

(4.69) is the radial part of the Schrödinger equation for the hydrogen atom.
$] 0, \infty\left[\right.$ is a natural real domain for this equation. If $\nu, \alpha^{2}$ are real, then so is $V_{\alpha, \nu}$ on $] 0, \infty[$.
The derivation of (4.69) from (4.66) is immediate:

$$
\begin{aligned}
& 4\left(-\partial_{z}^{2}+\frac{\gamma^{2}}{4}-\frac{\nu}{2 z}+\left(\frac{\alpha^{2}}{4}-\frac{1}{4}\right) \frac{1}{z^{2}}\right) \\
= & -\partial_{r}^{2}+\gamma^{2}-\frac{\nu}{r}+\left(\frac{\alpha^{2}}{4}-\frac{1}{4}\right) \frac{1}{r^{2}} .
\end{aligned}
$$

### 4.3 Rotationally symmetric harmonic oscillator

We consider (2) of (4.67). We set

$$
\begin{equation*}
z=r^{2}, \quad \text { which solves } \quad z^{\prime}=2 \sqrt{z} \tag{4.70}
\end{equation*}
$$

This leads to the Schrödinger equation

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V_{\alpha, \gamma}(r)-2 \nu\right) \phi(r)=0 \tag{4.71}
\end{equation*}
$$

with the potential

$$
V_{\alpha, \gamma}(r):=\gamma^{2} r^{2}+\left(\alpha^{2}-\frac{1}{4}\right) \frac{1}{r^{2}}
$$

$] 0, \infty\left[\right.$ is a natural real domain. For real $\gamma^{2}, \alpha^{2}$, the potential is real on $] 0, \infty[$. By scaling, the only different real cases are $\gamma^{2}=1,0,-1$.

For $\gamma^{2}=1$ the equation is the radial part of the rotationally symmetric harmonic oscillator. For $\gamma^{2}=0$ the equation is the radial part of the Helmholtz equation.

Here is an explicit derivation of (4.71) from (4.66):

$$
\begin{aligned}
& 4 z^{1-\frac{1}{4}}\left(-\partial_{z}^{2}+\frac{\gamma^{2}}{4}-\frac{\nu}{2 z}+\left(\frac{\alpha^{2}}{4}-\frac{1}{4}\right) \frac{1}{z^{2}}\right) z^{\frac{1}{4}} \\
= & -4 z \partial_{z}^{2}-2 \partial_{z}+\gamma^{2} z-2 \nu+\left(\alpha^{2}-\frac{1}{4}\right) \frac{1}{z} \\
= & -\partial_{r}^{2}+\gamma^{2} r^{2}-2 \nu+\left(\alpha^{2}-\frac{1}{4}\right) \frac{1}{r^{2}} .
\end{aligned}
$$

### 4.4 Morse potential

We consider (3) of (4.67). We set

$$
\begin{equation*}
z=\mathrm{e}^{-r}, \quad \text { which solves } \quad z^{\prime}=-z \tag{4.72}
\end{equation*}
$$

This leads to the Schrödinger equation

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V_{\nu, \gamma}(r)+\frac{\alpha^{2}}{4}\right) \phi(r)=0 \tag{4.73}
\end{equation*}
$$

where

$$
V_{\nu, \gamma}(r)=\frac{\gamma^{2}}{4} \mathrm{e}^{-2 r}-\frac{\nu}{2} \mathrm{e}^{-r}
$$

$]-\infty, \infty\left[\right.$ is a natural real domain. For real $\nu, \gamma^{2}$, the potential is real on $]-\infty, \infty[$. By translation, the only different real cases are $\gamma^{2}=1,0,-1$.

Here is an explicit derivation of (4.73) from (4.66):

$$
\begin{aligned}
& z^{2-\frac{1}{2}}\left(-\partial_{z}^{2}+\frac{\gamma^{2}}{4}-\frac{\nu}{2 z}+\left(\frac{\alpha^{2}}{4}-\frac{1}{4}\right) \frac{1}{z^{2}}\right) z^{\frac{1}{2}} \\
= & -z^{2} \partial_{z}^{2}-z \partial_{z}+\frac{\gamma^{2}}{4} z^{2}-\frac{\nu}{2} z+\frac{\alpha^{2}}{4} \\
= & -\partial_{r}^{2}+\frac{\gamma^{2}}{4} \mathrm{e}^{-2 r}-\frac{\nu}{2} \mathrm{e}^{-r}+\frac{\alpha^{2}}{4} .
\end{aligned}
$$

## 5 Schrödinger equations reducible to the translated harmonic oscillator

The last family of exactly solvable Schrödinger equations that we describe is (1.6), given by

$$
\begin{equation*}
-\partial_{z}^{2}+\theta^{2} z^{2}+\rho z+\lambda \tag{5.74}
\end{equation*}
$$

We will call (5.74) the translated harmonic oscillator equation. It is already in the canonical form.
If $\theta^{2} \neq 0$, then (5.74) is just the translation of the Schrödinger equation for usual harmonic oscillator. It is then equivalent to the Hermite equation, given by

$$
\mathcal{G}\left(a, y, \partial_{y}\right):=\left(\partial_{y}^{2}-2 y \partial_{y}-2 a\right)
$$

In fact, we have

$$
\begin{aligned}
& \theta \mathrm{e}^{-\frac{y^{2}}{2}} \mathcal{G}\left(a, y, \partial_{y}\right) \mathrm{e}^{\frac{y^{2}}{2}} \\
= & \theta\left(-\partial_{y}^{2}+y^{2}+2 a-1\right) \\
= & -\partial_{z}^{2}+\theta^{2} z^{2}+\rho z+\frac{\rho^{2}}{2 \theta^{2}}+2 a-1
\end{aligned}
$$

where $y=\sqrt{\theta}\left(z+\frac{\rho}{2 \theta^{2}}\right)$.
For $\theta=0, \rho \neq 0,(5.74)$ is equivalent to the Airy equation

$$
\left(\partial_{y}^{2}+y\right) \psi(y)=0
$$

which in turn is equivalent to a special case of the ${ }_{0} F_{1}$ equation:

$$
\begin{aligned}
& -\partial_{z}^{2}+\rho z+\lambda \\
= & \rho^{\frac{2}{3}}\left(-\partial_{y}^{2}+y\right) \\
= & (\rho 3)^{\frac{2}{3}} w^{\frac{1}{3}} \mathcal{F}\left(\frac{2}{3} ; w, \partial_{w}\right),
\end{aligned}
$$

where we set $y=\rho^{\frac{1}{3}} z+\rho^{-\frac{2}{3}} \lambda$, and $w=3^{-2} y^{3}$.
We will describe 3 classes of Schrödinger operators that can be reduced to (5.74) corresponding to three obvious choices for $\left(r^{\prime}(z)\right)^{2}$ :

$$
\begin{array}{lll}
(1) & (2)  \tag{5.75}\\
z & (3) z^{2} .
\end{array}
$$

### 5.1 Translated harmonic oscillator

We consider Ansatz (1) of (5.75), which corresponds to the most obvious choice of the energy, that is $-\lambda$. We rename the variable $z=r$. This leads to the Schrödinger equation

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V_{\theta, \rho}(r)+\lambda\right) \phi(r)=0 \tag{5.76}
\end{equation*}
$$

where

$$
V_{\theta, \rho}(r)=\theta^{2} r^{2}+\rho r .
$$

### 5.2 Special potential I

We can choose the energy to be $-\rho$, that is Ansatz (2) of (5.75). This corresponds to the substitution

$$
z=\left(\frac{3 r}{2}\right)^{\frac{2}{3}}, \quad \text { which solves } \quad z^{\prime}=z^{-\frac{1}{2}}
$$

and leads to the Schrödinger equation

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V_{\theta, \lambda}(r)+\rho\right) \phi(r)=0 \tag{5.77}
\end{equation*}
$$

where

$$
V_{\theta, \lambda}(r)=\theta^{2}\left(\frac{3 r}{2}\right)^{\frac{2}{3}}+\lambda\left(\frac{2}{3 r}\right)^{\frac{2}{3}}-\frac{5}{36} \frac{1}{r^{2}}
$$

In fact,

$$
\begin{aligned}
& z^{\frac{1}{4}-1}\left(-\partial_{z}^{2}+\theta^{2} z^{2}+\rho z+\lambda\right) z^{-\frac{1}{4}} \\
= & -\frac{1}{z} \partial_{z}^{2}+\frac{1}{2 z^{2}} \partial_{z}+\theta^{2} z+\rho+\frac{\lambda}{z}-\left(\frac{1}{4}+\frac{1}{4^{2}}\right) \frac{1}{z^{3}} \\
= & -\partial_{r}^{2}+\theta^{2}\left(\frac{3 r}{2}\right)^{\frac{2}{3}}+\rho+\lambda\left(\frac{2}{3 r}\right)^{\frac{2}{3}}-\left(\frac{1}{4}+\frac{1}{4^{2}}\right)\left(\frac{2}{3 r}\right)^{2} .
\end{aligned}
$$

Note the following intriguing feature of the above potential: the coefficient at $r^{-2}$ is fixed and one cannot change it by rescaling the variable $r$.

### 5.3 Special potential II

We can choose the energy to be $-\theta^{2}$, that is Ansatz (3) of (5.75). This corresponds to the substitution

$$
z=(2 r)^{\frac{1}{2}}, \quad \text { which solves } \quad z^{\prime}=z^{-1}
$$

and leads to the Schrödinger equation

$$
\begin{equation*}
\left(-\partial_{r}^{2}+V_{\rho, \lambda}(r)+\theta^{2}\right) \phi(r)=0 \tag{5.78}
\end{equation*}
$$

where

$$
V_{\rho, \lambda}(r)=\frac{\rho}{(2 r)^{\frac{1}{2}}}+\frac{\lambda}{2 r}-\frac{3}{16} \frac{1}{r^{2}}
$$

In fact,

$$
\begin{aligned}
& z^{\frac{1}{2}-2}\left(-\partial_{z}^{2}+\theta^{2} z^{2}+\rho z+\lambda\right) z^{-\frac{1}{2}} \\
= & -\frac{1}{z^{2}} \partial_{z}^{2}+\frac{1}{z^{3}} \partial_{z}+\theta^{2}+\frac{\rho}{z}+\frac{\lambda}{z^{2}}-\left(\frac{1}{2}+\frac{1}{2^{2}}\right) \frac{1}{z^{4}} \\
= & -\partial_{r}^{2}+\theta^{2}+\frac{\rho}{(2 r)^{\frac{1}{2}}}+\frac{\lambda}{2 r}-\left(\frac{1}{2}+\frac{1}{2^{2}}\right) \frac{1}{2^{2} r^{2}} .
\end{aligned}
$$

Again, the coefficient at $r^{-2}$ is fixed and one cannot change it by rescaling the variable $r$.

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