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# Completeness of multiseparable superintegrability in $E_{2, C}$ 

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#### Abstract

The possibility that Schrödinger's equation with a given potential can separate in more than one coordinate system is intimately connected with the notion of superintegrability. Examples of this type of system are well known. In this paper we demonstrate how to establish a complete list of such potentials using essentially algebraic means. Our approach is to classify all nondegenerate potentials that admit a pair of second-order constants of motion. Here 'nondegenerate' means that the potentials depend on four independent parameters. This is carried out for two-dimensional complex Euclidean space, though the method generalizes to other spaces and dimensions. We show that all these superintegrable systems correspond to quadratic algebras, and we work out the detailed structure relations and their quantum analogues.


## 1. Introduction

It has long been known that Schrödinger's equation with certain special potentials can admit (multiplicative) separation of variables in more than one coordinate system. This is intimately related to the notion of superintegrability [1-3]. This subject has been studied by a number of authors, based on the use of the corresponding differential equations that are implied by the requirement of simultaneous separability [4-17]. Specifically, superintegrability here means that for a Schrödinger equation in dimension $N$ there exist $2 N-1$ functionally independent second-order quantum mechanical observables (i.e., second-order $\dagger$ self-adjoint operators that commute with the Hamiltonian). There is an analogous concept of superintegrability for classical mechanical systems. This relates to the corresponding additive separation of variables of the Hamilton-Jacobi equation. A first step in studying separability in the classical case is to realize that the direct formulation of the simultaneous separability requirement is not obviously tractable. An additional observation is that if we do have simultaneous separability then the resulting constants of motion are observed to close quadratically under repeated application of the Poisson bracket [13]. We also know that, for spaces of constant curvature, separable coordinate systems of the free motion are described by quadratic elements of the corresponding first-order symmetries [18-20].

Although concrete examples of superintegrable systems are easily at hand, a complete classification of all such systems has presented major difficulties. How can one be sure that all systems for free motion have been found? (For example, Rañada's classification [17] omits
$\dagger$ We restrict our definition to second-order symmetries because it is only these that could possibly be related to variable separation. A classification of superintegrable systems involving observables of arbitrary order remains open and it is not clear if our method of integrablitity conditions would be tractable in the more general case.
our system (5a) below.) Once these are determined, how can one be sure that the most general additive potential term has been calculated?

Here we take a new approach to the problem and apply it for the case of two-dimensional complex Euclidean space. In section 2 we classify all nondegenerate potentials that admit a pair of second-order constants of motion. Here 'nondegenerate' means that the potentials depend on four independent parameters. The requirement that a potential admit two constants of motion leads to two second-order partial differential equations obeyed by the potential, and the integrability conditions for these two simultaneous equations permit us to classify all possibilities. (The 'direct approach' would be to classify directly all cases in which there are multiparameter solutions to the two coupled second-order PDEs obeyed by the potential, each of the form (16), a very complicated procedure. We replace this procedure by finding all solutions of the integrability conditions, (23), (24) a system of linear equations. These equations are lengthy to write in detail (we used Maple to handle the calculations) but straightforward to solve. Then the PDEs (16) need to be solved only for the small number of cases in which we know that the integrability conditions are satisfied.) The classification is greatly simplified by the equivalence of two potentials that are related by an action of the complex Euclidean motion group. We then prove that each nondegenerate potential is associated with a pair of constants of motion in the classical case, and a pair of symmetry operators in the quantum case, that generate a quadratic algebra. Furthermore, we verify that there is a one-to-one correspondence between superintegrable systems and freefield symmetry operators that generate quadratic algebras. Finally, we demonstrate explicitly that superintegrability implies multiseparability, i.e., separability in more than one coordinate system.

This systematic classification approach introduces a 'fine structure' into our problem. It is easy to show that potentials admitting two constants of motion cannot depend on more than four parameters. However, potentials that depend on fewer parameters, i.e., that cannot be embedded in a four parameter family, are not associated with a quadratic algebra.

## 2. Completeness in two-dimensional Euclidean space

Due to the close connection between separation of variables and constants of motion [21], a common approach to the classification and study of superintegrable systems is to search for potentials that permit variable separation in more than one coordinate system. The HamiltonJacobi equation $\dagger$ is

$$
\begin{equation*}
H=p_{x}^{2}+p_{y}^{2}+V(x, y)=\left(\frac{\partial S}{\partial x}\right)^{2}+\left(\frac{\partial S}{\partial y}\right)^{2}+V(x, y)=E \tag{1}
\end{equation*}
$$

The additive separation ansatz implies a solution

$$
\begin{equation*}
S=U(u, E, \alpha)+V(v, E, \alpha) \tag{2}
\end{equation*}
$$

in a suitable coordinate system $x=x(u, v), y=y(u, v)$. Here $\alpha$ is the separation constant. In the complex Euclidean plane there are six different separable coordinate systems, listed in the appendix, in what we take to be a standard form. One approach to our problem relates to finding all potentials that permit separation in more than one coordinate system. Here we have to allow for the possibility that the second coordinate system can be subjected to a Euclidean motion consisting of a rotation through the angle $\beta$ and a translation by the vector $(a, b)$. Due to
$\dagger$ Here we omit the usual $\frac{1}{2}$ factor multiplying the kinetic energy term in both the Hamilton-Jacobi and Schrödinger equations. Since the potentials of superintegrable systems are arbitrary up to a multiplicative factor, the $\frac{1}{2}$ is extraneous here.
its complexity this is not a revealing approach. Furthermore, it does not exclude the possibility of superintegrability without variable separation. Instead we adopt a different method, one that does not require variable separation but that leads to it as a consequence. Furthermore, the classification depends on nothing that is intrinsically more complicated than solving a system of linear equations.

Let us assume that, in addition to the classical Hamiltonian, we have two quadratic constants of motion

$$
\begin{equation*}
L_{h}=\sum_{k, j=1}^{2} a_{(h)}^{k j}(x, y) p_{k} p_{j}+W_{(h)}(x, y) \equiv \ell_{h}+W_{(h)} \quad h=1,2 \tag{3}
\end{equation*}
$$

which must satisfy

$$
\left\{H, L_{h}\right\}=0
$$

with $\left\}\right.$ the usual Poisson bracket. We require that the set $\left\{\mathrm{d} H, \mathrm{~d} L_{1}, \mathrm{~d} L_{2}\right\}$ is linearly independent, so that $H, L_{1}, L_{2}$ is a maximal set of functionally independent constants of motion. It is clear that $R=\left\{L_{1}, L_{2}\right\}$ is a constant of motion, so it and $R^{2}$ must be expressible as an analytic function of $H, L_{1}, L_{2}$ :

$$
\begin{equation*}
R^{2}=F\left(L_{0}, L_{1}, L_{2}\right) \quad H \equiv L_{0} \tag{4}
\end{equation*}
$$

Note that $R$ has the form

$$
\begin{equation*}
R=\sum_{k, l, m=1}^{2} c^{k l m} p_{k} p_{l} p_{m}+\sum_{k=1}^{2} d^{k} p_{k} \tag{5}
\end{equation*}
$$

but that it does not follow that $R^{2}$ is necessarily a polynomial as a function of $L_{0}, L_{1}, L_{2}$. We will find conditions that guarantee that $F$ is a third-order polynomial in its arguments.

Using the identity

$$
\begin{equation*}
\{K, G\}=\sum_{h=0}^{2}\left\{K, L_{h}\right\} \frac{\partial G}{\partial L_{h}} \tag{6}
\end{equation*}
$$

for a continuously differentiable function $G\left(L_{h}\right)$, we find the relations

$$
\begin{equation*}
\left\{L_{1}, R\right\}=\frac{1}{2} \frac{\partial F}{\partial L_{2}} \quad\left\{L_{2}, R\right\}=-\frac{1}{2} \frac{\partial F}{\partial L_{1}} \tag{7}
\end{equation*}
$$

Thus, the constants of motion $\left\{L_{1}, R\right\},\left\{L_{2}, R\right\}$ are easily computed once $F$ is known. Further, if $F$ is a polynomial in the invariants, then so are $\left\{L_{1}, R\right\}$, and $\left\{L_{2}, R\right\}$.

We first determine the conditions that the function

$$
\begin{equation*}
L=\sum_{j, k=1}^{2} a^{j k}(x, y) p_{k} p_{j}+W(x, y) \quad a^{j k}=a^{k j} \tag{8}
\end{equation*}
$$

must satisfy to be a constant of motion. The requirement is $\{H, L\}=0$ where

$$
\begin{equation*}
\{f, g\}=\sum_{j=1}^{2}\left(\frac{\partial f}{\partial p_{j}} \frac{\partial g}{\partial x_{j}}-\frac{\partial f}{\partial x_{j}} \frac{\partial g}{\partial p_{j}}\right) \quad\left(x_{1}, x_{2}\right)=(x, y) \tag{9}
\end{equation*}
$$

and

$$
\begin{equation*}
H=p_{1}^{2}+p_{2}^{2}+V(x, y) \tag{10}
\end{equation*}
$$

The conditions are thus

$$
\begin{array}{ll}
\frac{\partial a^{11}}{\partial x}=0 & 2 \frac{\partial a^{12}}{\partial x}+\frac{\partial a^{11}}{\partial y}=0 \\
\frac{\partial a^{22}}{\partial y}=0 & \frac{\partial a^{22}}{\partial x}+2 \frac{\partial a^{12}}{\partial y}=0 \tag{11}
\end{array}
$$

and

$$
\begin{equation*}
\frac{\partial W}{\partial x}-a^{11} \frac{\partial V}{\partial x}-a^{12} \frac{\partial V}{\partial y}=0 \quad \frac{\partial W}{\partial y}-a^{12} \frac{\partial V}{\partial x}-a^{22} \frac{\partial V}{\partial y}=0 \tag{12}
\end{equation*}
$$

The solution for the terms quadratic in the $p_{j}$ is

$$
\begin{align*}
& a^{11}=\alpha_{1} y^{2}+\alpha_{2} y+\alpha_{3}^{\prime}  \tag{13}\\
& a^{12}=-\alpha_{1} x y-\frac{1}{2} \alpha_{2} x-\frac{1}{2} \alpha_{4} y+\frac{1}{2} \alpha_{5}  \tag{14}\\
& a^{22}=\alpha_{1} x^{2}+\alpha_{4} x+\alpha_{3}^{\prime \prime} \tag{15}
\end{align*}
$$

where the $\alpha_{k}$ are constants. The requirement that $\partial_{x} W_{y}=\partial_{y} W_{x}$ leads from (12) to the second-order partial differential equation for the potential [22-26]

$$
\begin{gather*}
\frac{1}{2}\left(2 \alpha_{1} x y+\alpha_{2} x+\alpha_{4} y-\alpha_{5}\right)\left(V_{x x}-V_{y y}\right)+\left(\alpha_{1}\left[y^{2}-x^{2}\right]+\alpha_{2} y-\alpha_{4} x+\alpha_{3}\right) V_{x y} \\
=\left(-3 \alpha_{1} y-\frac{3}{2} \alpha_{2}\right) V_{x}+\left(3 \alpha_{1} x+\frac{3}{2} \alpha_{4}\right) V_{y} \tag{16}
\end{gather*}
$$

where $\alpha_{3}=\alpha_{3}^{\prime}-\alpha_{3}^{\prime \prime}$. We denote the solution space of this equation by

$$
\begin{equation*}
\left[\alpha_{1}, \ldots, \alpha_{5}\right] \tag{17}
\end{equation*}
$$

Let us now return to our assumption that the Hamilton-Jacobi equation admits two constants of motion:

$$
L_{h}=\sum_{j, k=1}^{2} a_{(h)}^{j k} p_{k} p_{j}+W_{(h)} \quad h=1,2 .
$$

These two operators together with $H$ are assumed functionally independent. The constant of motion $L_{1}$ leads to the condition that the potential $V$ belong to the solution space (17), whereas $L_{2}$ leads to the solution space

$$
\begin{equation*}
\left[\beta_{1}, \ldots, \beta_{5}\right] \tag{18}
\end{equation*}
$$

Thus the potential must lie in the intersection of the solution spaces (17) and (18). It follows that the equations

$$
\begin{equation*}
V_{x x}-V_{y y}=A V_{x}+B V_{y} \quad V_{x y}=C V_{x}+D V_{y} \tag{19}
\end{equation*}
$$

must hold, where

$$
\begin{align*}
& A \mathcal{E}=\frac{3}{2} H_{12}\left(x^{2}+y^{2}\right)-3 H_{14} x y+3 H_{13} y-\frac{3}{2} H_{24} x+\frac{3}{2} H_{23} \\
& B \mathcal{E}=\frac{3}{2} H_{14}\left(x^{2}+y^{2}\right)-3 H_{12} x y-3 H_{13} x+\frac{3}{2} H_{24} y+\frac{3}{2} H_{34} \\
& 2 C \mathcal{E}=-3 H_{14} y^{2}+\left(-\frac{3}{2} H_{24}+3 H_{15}\right) y+\frac{3}{2} H_{25}  \tag{20}\\
& 2 D \mathcal{E}=3 H_{12} x^{2}+\left(-\frac{3}{2} H_{24}-3 H_{15}\right) x-\frac{3}{2} H_{45} \\
& 2 \mathcal{E}=-H_{12} x y^{2}+H_{14} x^{2} y-H_{12} x^{3}+H_{14} y^{3}-2 H_{13} x y+H_{24}\left(x^{2}+y^{2}\right) \\
& \quad+H_{15}\left(x^{2}-y^{2}\right)+\left(H_{34}-H_{25}\right) y+\left(H_{45}-H_{23}\right) x-H_{35}
\end{align*}
$$

and $H_{k \ell}=-H_{\ell k}=\alpha_{k} \beta_{\ell}-\alpha_{\ell} \beta_{k}$.
From the fundamental equations (19) we can compute all of the third partial derivatives of $V$. Indeed
$V_{x x x}=\left(A_{x}+B C+C_{y}+C^{2}+A^{2}\right) V_{x}+\left(B_{x}+D B+D_{y}+C D+A B\right) V_{y}$ $+(A+D) V_{y y}$
$V_{x x y}=\left(C_{x}+D C+A C\right) V_{x}+\left(D_{x}+D^{2}+B C\right) V_{y}+C V_{y y}$
$V_{x y y}=\left(C_{y}+C^{2}\right) V_{x}+\left(D_{y}+C D\right) V_{y}+D V_{y y}$
$V_{y y y}=\left(-A_{y}+C_{x}+D C\right) V_{x}+\left(-B_{y}-A D+D_{x}+D^{2}+B C\right) V_{y}+(C-B) V_{y y}$.

Thus if the potential $V$ belongs to the solution spaces (17) and (18), then $V$ can depend on at most three parameters, in addition to a trivial additive constant. We can choose these parameters to be $V_{x}\left(x_{0}, y_{0}\right), V_{y}\left(x_{0}, y_{0}\right), V_{y y}\left(x_{0}, y_{0}\right)$ for any fixed regular point $\left(x_{0}, y_{0}\right)$. Then $V_{x x}\left(x_{0}, y_{0}\right)$ and all higher derivatives can be computed by successive differentiation of relations (21). We require that our potential be nondegenerate, i.e., that it depend on three arbitrary parameters.

Then, the conditions $\partial_{x} V_{x x y}=\partial_{y} V_{x x x}, \partial_{y} V_{x x y}=\partial_{x} V_{x y y}, \partial_{y} V_{x y y}=\partial_{x} V_{y y y}$ for the fourth partial derivatives lead to the integrability conditions
$\partial_{x}(2 C-B)=\partial_{y}(2 D+A) \quad$ (satisfied identically)
$C_{x x}-C_{y y}-A_{x y}=2 C C_{y}-D A_{y}-2 C D_{x}+A A_{y}-A C_{x}+C B_{y}+B C_{y}$
$D_{x x}-D_{y y}-B_{x y}=-2 D D_{x}-C B_{x}+2 D C_{y}-B B_{x}-B D_{y}+D A_{x}+A D_{x}$.
Note that if we have another constant of motion $L_{3}$ associated with a nondegenerate potential, then $L_{3}$ must be a linear combination of $H, L_{1}, L_{2}$. Indeed, if $L_{3}$ is not a linear combination of the basis functions, then the potential $V$ must satisfy an equation (16) that is linearly independent of the equations associated with $L_{1}, L_{2}$. This means an additional constraint on the solution space and that $V$ can depend on at most two parameters, which is a contradiction.

We will use the conditions (23) and (24) to classify the possible potentials $V$ and the corresponding constants of motion $L_{1}, L_{2}$. For this we note that it is only the three-dimensional subspace spanned by $H, L_{1}, L_{2}$ that matters; we can choose any basis for this subspace. Hence we can replace the subspace bases (17) and (18) by linear combinations of themselves without changing the potential. Moreover, to simplify the results we note that we can always subject the coordinates $(x, y)$, and $L_{1}, L_{2}$ to a simultaneous Euclidean motion, i.e., we regard all translated and rotated potentials as members of the same equivalence class.

Multiplying both sides of (23) and (24) by $\mathcal{E}^{3}$ we obtain polynomial identities in $x$ and $y$. Equating the coefficients of the various powers $x^{n} y^{m}$ we obtain conditions on the parameters $H_{j k}$. The simplest nontrivial condition, which is associated with the coefficient of a fifth-order power in either of the equations, is

$$
\begin{equation*}
2 H_{15}\left(H_{14}^{2}-H_{12}^{2}\right)+H_{24}\left(H_{14}^{2}+H_{12}^{2}\right)-4 H_{14} H_{12} H_{13}=0 . \tag{25}
\end{equation*}
$$

We exploit these and the remaining conditions, and Euclidean motions to classify the possibilities for the $L_{j}$. The full conditions (23) and (24), expressed in terms of the parameter $H_{i j}$, take several pages to list and are complicated to solve in general. (Indeed a symbol manipulation program was an important aid to our computations.) However, by dividing the problem up into special cases and using Euclidean motions, we can simplify the conditions and obtain a full solution. In the listing that follows we use the fact that the constants of motion can each be expressed as a quadratic element in the enveloping algebra of the Euclidean group in the plane with basis elements

$$
p_{x}, p_{y}, M=x p_{y}-y p_{x}
$$

to which a potential term $W(x, y)$ is added. (Strictly speaking, conditions (23) and (24) are only necessary conditions for existence of nondegenerate potentials. However, in our case-bycase study we have found that they are also sufficient: all solutions of these equations lead to nondegenerate potentials.)

Suppose $I \equiv H_{12}^{2}+H_{14}^{2} \neq 0$. Via an appropriate coordinate rotation through complex angle $\theta$ we obtain a new set of equations (19) in the rotated coordinates where the new parameters $H_{12}^{\prime}, H_{14}^{\prime}$ are related to the original ones by

$$
\begin{equation*}
H_{12}^{\prime}=H_{12} \cos \theta+H_{14} \sin \theta \quad H_{14}^{\prime}=H_{14} \cos \theta-H_{12} \sin \theta \tag{26}
\end{equation*}
$$

and for which $I=I^{\prime}$. Thus, by an appropriate choice of $\theta$, we can assume $H_{12}=0$. (Similarly, the translation $x=x^{\prime}+a, y=y^{\prime}+b$ induces new parameters $H_{i j}^{\prime}$
$H_{12}^{\prime}=H_{12} \quad H_{14}^{\prime}=H_{14} \quad H_{13}^{\prime}=H_{13}+b H_{12}-a H_{14}$
$H_{15}^{\prime}=H_{15}-a H_{12}-b H_{14} \quad H_{24}^{\prime}=H_{24}-2 a H_{12}+2 b H_{14}$
$H_{45}^{\prime}=H_{45}+a\left(H_{24}+2 H_{15}\right)-2 a^{2} H_{12} \quad H_{25}^{\prime}=H_{25}+b\left(2 H_{15}-H_{24}\right)-2 b^{2} H_{14}$
$H_{23}^{\prime}=H_{23}-a H_{24}+2 b H_{13}+a^{2} H_{12}-2 a b H_{14}+b^{2} H_{12}$
$H_{34}^{\prime}=H_{34}-2 a H_{13}+b H_{24}+a^{2} H_{14}-2 a b H_{12}+b^{2} H_{14}$
$H_{35}^{\prime}=H_{35}+a\left(H_{23}-H_{45}\right)+b\left(H_{25}-H_{34}\right)-a^{2}\left(H_{15}+H_{24}\right)+2 a b H_{13}+b^{2}\left(H_{15}-H_{24}\right)$

$$
+a^{3} H_{12}-a^{2} b H_{14}+a b^{2} H_{12}-b^{3} H_{14}
$$

expressed in terms of the original $H_{i j}$ parameters.) Then, by an appropriate Euclidean translation that leaves $H_{12}, H_{14}$ unchanged, it follows from (20) that we can assume $H_{13}=H_{24}=0$. Then (25) implies $H_{15}=0$ so, since $H_{i j}=\alpha_{i} \beta_{j}-\alpha_{j} \beta_{i}$, we can assume $H_{35}=H_{23}=0$. Further, the fourth-order integrability conditions give $H_{34}=H_{45}=0$.

Case (1): $\quad H_{12}^{2}+H_{14}^{2} \neq 0$

$$
\begin{equation*}
[1,0,0,0,0] \quad[0,0,0,1,0] . \tag{27}
\end{equation*}
$$

Here,

$$
\begin{align*}
& L_{1}=4 M^{2}+W^{(1)} \quad L_{2}=-2 M p_{y}+W^{(2)}  \tag{28}\\
& V(x)=\frac{\alpha}{\sqrt{x^{2}+y^{2}}}+\frac{1}{\sqrt{x^{2}+y^{2}}}\left[\frac{\beta}{\sqrt{x^{2}+y^{2}}+x}+\frac{\gamma}{\sqrt{x^{2}+y^{2}}-x}\right] . \tag{29}
\end{align*}
$$

This potential allows separation in parabolic or polar coordinates. (Note: due to expressions (12), it is always a straightforward integration to compute the terms $W^{(j)}$, and we will not list these explicitly. Furthermore, once the constants of motion for a superintegrable system are known, it is relatively straightforward to determine the possible separable coordinate systems associated with this system. This is due to the fact that the constants of motion for the separable coordinate systems are already known (see the appendix). The only complication is that we may have to apply a Euclidean transformation to the standardized constant of motion for the separable system to obtain the corresponding symmetry of the superintegrable system. Here, we simply list the separable systems associated with each superintegrable system, and provide details only in the cases where shifted coordinates occur.)

If, on the other hand, $H_{12}= \pm \mathrm{i} H_{14} \neq 0$, then, via translation, we can assume $H_{24}=0$. In this case (25) implies $H_{15}=\mathrm{i} H_{13}$.

Case (2): $\quad H_{12}= \pm \mathrm{i} H_{14} \neq 0$

$$
\begin{equation*}
\left[1,0, \alpha_{3}, 0,-\mathrm{i} \alpha_{3}-\mathrm{i} \beta_{3}^{2}\right] \quad\left[0,-1, \beta_{3}, \mathrm{i}, \mathrm{i} \beta_{3}\right] . \tag{30}
\end{equation*}
$$

Here
$L_{1}=M^{2}-\frac{\beta_{3}^{2}}{4} p_{+}^{2}+\left(\frac{\alpha_{3}}{2}+\frac{\beta_{3}^{2}}{4}\right) p_{-}^{2}+W^{(1)} \quad L_{2}=M p_{+}+\frac{\beta_{3}}{2} p_{+}^{2}+W^{(2)}$
where $p_{ \pm}=p_{x} \pm \mathrm{i} p_{y}, z=x+\mathrm{i} y$, and $\bar{z}=x-\mathrm{i} y$. There are two subcases to consider. If $\mu=2 \alpha_{3}+\beta_{3}^{2} \neq 0$, then via a rotation about the origin we achieve the following.

Case (2a):

$$
\begin{equation*}
\left[1,0,2 c^{2}, 0,0\right] \quad[0,-1, \mathrm{i} c, \mathrm{i},-c] \tag{32}
\end{equation*}
$$

where

$$
\begin{align*}
& L_{1}=M^{2}+c^{2} p_{x}^{2}+W^{(1)} \quad L_{2}=M p_{+}+\frac{\mathrm{i} c}{2} p_{+}^{2}+W^{(2)} \\
& V(x)=\frac{\alpha z}{\left(c^{2}-z^{2}\right)^{\frac{1}{2}}}+\frac{\beta}{\sqrt{(c-z)(c+\bar{z})}}+\frac{\gamma}{\sqrt{(c+z)(c+\bar{z})}} \tag{33}
\end{align*}
$$

The corresponding Hamilton-Jacobi and Schrödinger equations for this system separates in elliptical coordinates (see the appendix), as well as shifted elliptical coordinates. In terms of the cartesian coordinates $x_{E}, y_{E}$ associated with the elliptic coordinates, the shifted coordinates are $X=x_{E}-c, Y=y_{E}-\mathrm{i} c$. The corresponding operator in this case is $\left(M+\mathrm{i} c\left(p_{x}+\mathrm{i} p_{y}\right)\right)^{2}+c^{2} p_{x}^{2}$. If $\mu=0$ we have the following.

Case (2b):

$$
\begin{equation*}
[1,0,2,0,2 \mathrm{i}] \quad[0,-1,2 \mathrm{i}, \mathrm{i},-2] \tag{34}
\end{equation*}
$$

where

$$
\begin{align*}
& L_{1}=M^{2}+p_{+}^{2}+W^{(1)} \quad L_{2}=\left(M+2 \mathrm{i} p_{+}\right)^{2}+p_{+}^{2}+W^{(2)} \\
& V(x)=\frac{\alpha}{z^{2}}+\frac{\beta}{\sqrt{z^{3}(\bar{z}+2)}}+\frac{\gamma}{\sqrt{z(\bar{z}+2)}} \tag{35}
\end{align*}
$$

This system separates in terms of hyperbolic coordinates (see the appendix) and displaced hyperbolic coordinates. In terms of the cartesian coordinates $x_{H}, y_{H}$ associated with the hyperbolic coordinates, the displaced coordinates are $X=x_{H}-2, Y=y_{H}-2 \mathrm{i}$.

Now suppose $H_{12}=H_{14}=0$.

Case (3): $\quad H_{12}=H_{14}=0, \alpha_{1} \neq 0$

$$
\begin{equation*}
\left[1,0,0,0, \alpha_{5}\right] \quad\left[0,0, \beta_{3}, 0, \beta_{5}\right] \tag{36}
\end{equation*}
$$

This breaks up into three subcases. A rotation through complex angle $\theta$ has the effect
$H_{15}^{\prime}=H_{15} \cos 2 \theta+H_{13} \sin 2 \theta \quad H_{13}^{\prime}=H_{13} \cos 2 \theta-H_{15} \sin 2 \theta \quad H_{24}^{\prime}=H_{24}$.
Thus, if $H_{15}^{2}+H_{13}^{2} \neq 0$ we can achieve $H_{15}=0$. Then, we can assume via translation that $\alpha_{2}=\alpha_{4}=0$, so $H_{24}=0$. An integrability condition gives $H_{35}=0$. Thus,

Case ( $3 a$ ): $\quad H_{15}^{2}+H_{13}^{2} \neq 0$

$$
\begin{equation*}
[1,0,0,0,0] \quad[0,0,1,0,0] . \tag{38}
\end{equation*}
$$

Here

$$
\begin{align*}
& L_{1}=M^{2}+W^{(1)} \quad L_{2}=p_{x}^{2}+W^{(2)}  \tag{39}\\
& V(x)=\alpha\left(x^{2}+y^{2}\right)+\frac{\beta}{x^{2}}+\frac{\gamma}{y^{2}} \tag{40}
\end{align*}
$$

This potential permits separation in polar, elliptic and cartesian coordinates.
If, however, $H_{15}= \pm i H_{13} \neq 0$, we can again translate to get $H_{24}=0$, and find two possibilities, depending on whether $H_{35}=0$. (Here case ( $3 b$ ) can be considered as a limit of case ( $3 c$ ) as the parameter $c \rightarrow 0$.)

Case (3b): $\quad H_{15}^{2}+H_{13}^{2}=0$

$$
\begin{equation*}
[1,0,0,0,0] \quad[0,0,2,0, \pm 2 \mathrm{i}] . \tag{41}
\end{equation*}
$$

Here

$$
\begin{align*}
& L_{1}=M^{2}+W^{(1)} \quad L_{2}=p_{+}^{2}+W^{(2)} \\
& V(x)=\alpha \frac{x^{2}+y^{2}}{(x+\mathrm{i} y)^{4}}+\frac{\beta}{(x+\mathrm{i} y)^{2}}+\gamma\left(x^{2}+y^{2}\right) \tag{42}
\end{align*}
$$

(There is a similar solution where the term $p_{+}^{2}$ in $L_{2}$ is replaced by $p_{-}^{2}$.) The potential permits separation in hyperbolic and polar coordinates.

Case (3c): $\quad H_{15}^{2}+H_{13}^{2}=0$

$$
\begin{equation*}
\left[1,0, c^{2}, 0,0\right] \quad[0,0,2,0, \pm 2 \mathrm{i}] . \tag{43}
\end{equation*}
$$

Here

$$
\begin{align*}
& L_{1}=M^{2}+c^{2} p_{x}^{2}+W^{(1)} \quad L_{2}=p_{+}^{2}+W^{(2)} \\
& V(x)=\frac{\alpha z}{\sqrt{z^{2}-c^{2}}}+\frac{\beta \bar{z}}{\sqrt{z^{2}-c^{2}}\left(z+\sqrt{z^{2}-c^{2}}\right)^{2}}+\gamma z \bar{z} \tag{44}
\end{align*}
$$

The potential permits separation in hyperbolic and elliptic coordinates. Indeed, instead of the basis $L_{1}, L_{2}$ let us consider the basis $M^{2}+c^{2} p_{x}^{2}$ and $M^{2}+\frac{1}{4} c^{2}\left(p_{x}-\mathrm{i} p_{y}\right)^{2}$. Corresponding to these operators are (1) the normal choice of elliptic coordinates and (2) the choice of hyperbolic coordinates $x=\frac{1}{2} c x_{H}$ and $y=-\frac{i}{2} c y_{H}$ (see the appendix).

Case (4): $\quad H_{12}=H_{13}=H_{14}=H_{15}=0, \alpha_{2} \neq 0, H_{24} \neq 0$

$$
\begin{equation*}
[0,1,0,0,0] \quad[0,0,0,1,0] . \tag{45}
\end{equation*}
$$

Here

$$
\begin{align*}
& L_{1}=-2 M p_{x}+W^{(1)} \quad L_{2}=-2 M p_{y}+W^{(2)}  \tag{46}\\
& V(x)=\frac{\alpha}{\sqrt{x^{2}+y^{2}}}+\beta \frac{\left(\sqrt{x^{2}+y^{2}}+x\right)^{\frac{1}{2}}}{\sqrt{x^{2}+y^{2}}}+\gamma \frac{\left(\sqrt{x^{2}+y^{2}}-x\right)^{\frac{1}{2}}}{\sqrt{x^{2}+y^{2}}} . \tag{47}
\end{align*}
$$

Separation of variables is possible in two types of parabolic coordinates, the usual parabolic coordinates and the interchanged parabolic coordinates $x=\mu \nu, y=\frac{1}{2}\left(\mu^{2}-\nu^{2}\right)$.

Case (5): $\quad H_{12}=H_{13}=H_{14}=H_{15}=0, \alpha_{2} \neq 0, H_{24}=0$

$$
\begin{equation*}
\left[0,1, \alpha_{3}, \alpha_{4}, \alpha_{5}\right] \quad\left[0,0, \beta_{3}, 0, \beta_{5}\right] . \tag{48}
\end{equation*}
$$

If $\left(H_{34}-H_{25}\right)^{2}+\left(H_{45}-H_{23}\right)^{2} \neq 0$ we can make a complex rotation to achieve $H_{45}=H_{23}$. Then we can take $H_{25}=1$ and the constant term integrability conditions yield $H_{34}=-H_{23}^{2}=1$.

Case (5a): $\quad\left(H_{34}-H_{25}\right)^{2}+\left(H_{45}-H_{23}\right)^{2} \neq 0$

$$
\begin{equation*}
\left[0,1, \alpha_{3}, \pm \mathrm{i}, 0\right] \quad[0,0, \pm \mathrm{i}, 0,1] . \tag{49}
\end{equation*}
$$

Here we choose the typical case

$$
\begin{align*}
& L_{1}=4 \mathrm{i} M p_{-}+p_{+}^{2}+W^{(1)} \quad L_{2}=p_{-}^{2}+W^{(2)}  \tag{50}\\
& V(x)=\alpha(x-\mathrm{i} y)+\beta\left(x+\mathrm{i} y-\frac{3}{2}(x-\mathrm{i} y)^{2}\right)+\gamma\left(x^{2}+y^{2}-\frac{1}{2}(x-\mathrm{i} y)^{3}\right) \tag{51}
\end{align*}
$$

The possible separable coordinates are semihyperbolic coordinates corresponding to operator $M p_{-}+p_{+}^{2}$ and shifted semihyperbolic coordinates with operator $M p_{-}+\delta p_{-}^{2}+p_{+}^{2}$. This corresponds to the standard coordinates shifted via the transformation $x \rightarrow x+\delta, y \rightarrow y+\mathrm{i} \delta$.

Case (5b): $\quad\left(H_{34}-H_{25}\right)^{2}+\left(H_{45}-H_{23}\right)^{2}=0$.
Here $\left(H_{34}-H_{25}\right)= \pm \mathrm{i}\left(H_{45}-H_{23}\right) \neq 0$. The constant term integrability conditions, and a translation in $y$, yield the solutions

$$
\begin{equation*}
\left[0,1,0, \alpha_{4}, 0\right] \quad\left[0,0, \beta_{3}, 0,1\right] \quad \beta_{3}= \pm \mathrm{i} \quad \alpha_{4} \neq \beta_{3} \tag{52}
\end{equation*}
$$

Then, an appropriate rotation about the origin takes this to

$$
[0,0,0,1,0] \quad[0,0,1,0, \pm \mathrm{i}]
$$

This system, for which now $\alpha_{2}=0$, corresponds to case ( $6 b$ ) below.

Case (6): $\quad H_{12}=H_{13}=H_{14}=H_{15}=0, \alpha_{2}=0, \alpha_{3} \neq 0$.
The constant term integrabilty condition is $H_{45}\left(H_{45}^{2}+H_{34}^{2}\right)=0$. There are two cases.
Case (6a): $\quad H_{45}=0$

$$
\begin{equation*}
[0,0,1,0,0] \quad[0,0,0,1,0] . \tag{53}
\end{equation*}
$$

Here

$$
\begin{align*}
& L_{1}=p_{x}^{2}+W^{(1)} \quad L_{2}=-2 M p_{y}+W^{(2)}  \tag{54}\\
& V(x)=\alpha\left(4 x^{2}+y^{2}\right)+\beta x+\frac{\gamma}{y^{2}} . \tag{55}
\end{align*}
$$

The possible separable coordinates are cartesian and parabolic.

Case (6b): $\quad H_{45} \neq 0$

$$
\begin{equation*}
[0,0,1,0, \pm \mathrm{i}] \quad[0,0,0,1,0] \tag{56}
\end{equation*}
$$

Here we take

$$
\begin{align*}
& L_{1}=2 p_{y} p_{+}+W^{(1)} \quad L_{2}=M p_{y}+W^{(2)}  \tag{57}\\
& V(x)=\frac{\alpha}{\sqrt{x+\mathrm{i} y}}+\beta x+\gamma \frac{2 x+\mathrm{i} y}{\sqrt{x+\mathrm{i} y}} \tag{58}
\end{align*}
$$

There is the possibility of separability in parabolic coordinates $\left\{M p_{y}\right\}$ or displaced parabolic coordinates $\left\{\left(M+\delta\left(p_{x} \pm \mathrm{i} p_{y}\right)\right) p_{y}\right\}$ for suitable $\delta$.

Now we demonstrate that there is a quadratic algebra associated with each nondegenerate potential. Because we are working in two dimensions there can only be three functionally independant constants at most. Consequently, all Poisson brackets must be functionally dependent on $H=L_{0}, L_{1}$ and $L_{2}$. We want to show that in fact $R^{2}=\left\{L_{1}, L_{2}\right\}^{2}=$ $F\left(L_{0}, L_{1}, L_{2}\right)$ is a polynomial in these variables.

Note that for arbitrary $L_{1}, L_{2}$, the $F$ is in general not a polynomial. Consider the example

$$
L_{0}=p_{x}^{2}+p_{y}^{2} \quad L_{1}=M^{2}+p_{x} p_{y} \quad L_{2}=p_{x}^{2}
$$

Then we have $R=\left\{L_{1}, L_{2}\right\}=4 M p_{x} p_{y}$ and

$$
R^{2}=F\left(L_{0}, L_{1}, L_{2}\right)=16 L_{1} L_{2}\left(L_{0}-L_{2}\right)-16 L_{2}^{\frac{3}{2}}\left(L_{0}-L_{2}\right)^{\frac{3}{2}}
$$

Here, although $F$ is defined and bounded at the point $\left(L_{0}, L_{1}, L_{2}\right)=(0,0,0)$, it is not analytic at this point. Thus it has no power series expansion about the origin. We conjecture that this is an illustration of the general problem: if $F$ is not a polynomial, then there are branch points or cuts at $(0,0,0)$.

We will show, however, that for nondegenerate potentials the associated $F$ is a polynomial. First, we can verify that this is true when the potential is turned off, i.e., if we consider only the functions

$$
\ell_{h}=\sum_{j, k=1}^{2} a_{(h)}^{j k} p_{k} p_{j} \quad i=h, 2 \quad \ell_{0}=p_{x}^{2}+p_{y}^{2}
$$

where $L_{h}=\ell_{h}+W^{(h)}$. Let $\mathcal{R}=\left\{\ell_{1}, \ell_{2}\right\}$. Then for each of the cases listed above it is straightforward to check that $\mathcal{R}^{2}=\mathcal{P}_{3}\left(\ell_{0}, \ell_{1}, \ell_{2}\right)$ where $\mathcal{P}_{3}$ is a homogeneous third-order polynomial in its arguments $\dagger$. It follows that

$$
\begin{equation*}
R^{2}=F\left(L_{0}, L_{1}, L_{2}\right)=\mathcal{P}_{3}\left(L_{0}, L_{1}, L_{2}\right)+F_{4}\left(s, L_{0}, L_{1}, L_{2}\right) \tag{59}
\end{equation*}
$$

where $F_{4}$ is a fourth-, second- and zeroth-order polynomial in the momenta $p_{x}, p_{y}$, and $F_{4}\left(\mathbf{0}, L_{0}, L_{1}, L_{2}\right)=0$. Here, the parameters in the potential are denoted by $s=\left(V_{x}^{0}, V_{y}^{0}, V_{y y}^{0}\right)$, evaluated at some fixed point $\left(x_{0}, y_{0}\right)$ and $F_{4}$ is a polynomial function of these parameters.

From (7) we have

$$
\begin{aligned}
& \left\{\ell_{1}, \mathcal{R}\right\}=\frac{1}{2} \frac{\partial \mathcal{P}_{3}}{\partial \ell_{2}}\left(\ell_{0}, \ell_{1}, \ell_{2}\right) \\
& \left\{\ell_{2}, \mathcal{R}\right\}=-\frac{1}{2} \frac{\partial \mathcal{P}_{3}}{\partial \ell_{1}}\left(\ell_{0}, \ell_{1}, \ell_{2}\right)
\end{aligned}
$$

hence

$$
\begin{aligned}
& \left\{L_{1}, R\right\}=\frac{1}{2} \frac{\partial \mathcal{P}_{3}}{\partial L_{2}}\left(L_{0}, L_{1}, L_{2}\right)+\frac{1}{2} \frac{\partial F_{4}}{\partial L_{2}}(s) \\
& \left\{L_{2}, R\right\}=-\frac{1}{2} \frac{\partial \mathcal{P}_{3}}{\partial L_{1}}\left(L_{0}, L_{1}, L_{2}\right)-\frac{1}{2} \frac{\partial F_{4}}{\partial L_{1}}(s)
\end{aligned}
$$

where the $\partial F_{4} / \partial L_{h}(s)$ have only terms of orders two and zero in the momenta. It follows that the $\partial F_{4} / \partial L_{h}(s)$ must be expressible as linear combinations of the $L_{h}$. This shows that the commutators $\left\{L_{h}, R\right\}$ can be expressed as polynomials in $L_{0}, L_{1}, L_{2}$. It is then a simple matter to verify that $F$ itself is a polynomial in $L_{0}, L_{1}, L_{2}$.

We now list the quadratic algebra relations for each of the cases studied above. In view of relations (7) it is sufficient to give the relation $R^{2}=F\left(L_{0}, L_{1}, L_{2}\right)$ for each case.

Case (1): $\quad[1,0,0,0,0],[0,0,0,1,0]$

$$
\begin{aligned}
R^{2}=16 L_{1}^{2} H & -16 L_{2}^{2} L_{1}-32(\beta+\gamma) L_{2}^{2} \\
& +64 \alpha(\beta-\gamma) L_{2}+16 \alpha^{2} L_{1}-256 \beta \gamma H-32 \alpha^{2}(\beta+\gamma)
\end{aligned}
$$

Case (2a): $\quad\left[1,0,2 c^{2}, 0,0\right],[0,-1, \mathrm{i} c, \mathrm{i},-c]$

$$
\begin{aligned}
R^{2}=\frac{1}{2} c^{4} H^{3} & -4 \mathrm{i} c L_{2}^{3}+2 c^{2} L_{2}^{2} H-4 L_{2}^{2} L_{1}-c^{2} H^{2} L_{1}-\mathrm{i} c^{3} H^{2} L_{2}+\frac{\mathrm{i}}{2} c^{4} \alpha H^{2} \\
& +2 \mathrm{i} \alpha c^{2} L_{2}^{2}-c^{2} \alpha^{2} L_{1}+\mathrm{i} c\left(\beta^{2}+\gamma^{2}-c^{2} \alpha^{2}\right) L_{2} \\
& +\frac{1}{2}\left(-c^{2} \beta^{2}+c^{4} \alpha^{2}+c^{2} \gamma^{2}\right) H+\frac{1}{2}\left(2 \beta \gamma+\mathrm{i} c^{2} \alpha^{2}\right) c^{2} \alpha .
\end{aligned}
$$

$\dagger$ Moreover, it is straightforward to verify that the cases corresponding to nondegenerate potentials are the only cases where $\mathcal{P}_{3}$ is a homogeneous third-order polynomial in its arguments. Thus the possible quadratic algebras generated by second-order elements in the Euclidean Lie algebra correspond one-to-one with nondegenerate potentials.

Case (2b): $\quad[1,0,2,0,2 \mathrm{i}],[0,-1,2 \mathrm{i}, \mathrm{i},-2]$

$$
\begin{aligned}
R^{2}=-2 L_{1}^{3}- & 2 L_{2}^{3}+2 L_{2}^{2} L_{1}+2 L_{1}^{2} L_{2}+32 \alpha H L_{1} \\
& +32 \alpha H L_{2}-8 \beta \gamma L_{1}+8 \beta \gamma L_{2}+16 \beta^{2} H+16 \alpha \gamma^{2} .
\end{aligned}
$$

Case (3a): $\quad[1,0,0,0,0],[0,0,1,0,0]$
$R^{2}=-16 L_{2}^{2} L_{1}+16 L_{2} L_{1} H-16 \beta H^{2}-16(\beta+\gamma) L_{2}^{2}-16 \alpha L_{1}^{2}+32 H L_{2}+64 \alpha \beta \gamma$.

Case (3b): $\quad[1,0,0,0,0],[0,0,2,0, \pm 2 \mathrm{i}]$

$$
\begin{aligned}
R^{2}=-16 L_{1}^{3} & +32 L_{1}^{2} L_{2}-16 L_{2}^{2} L_{1}+16 \alpha H^{2} \\
& +16 \beta H L_{2}-16 \beta H L_{1}-64 \alpha \gamma L_{1}-16 \beta^{2} \gamma
\end{aligned}
$$

Case (3c): $\quad\left[1,0, c^{2}, 0,0\right],[0,0,2,0, \pm 2 \mathrm{i}]$

$$
\begin{aligned}
R^{2}=4 c^{2} L_{2}^{3}+ & 8 c^{2} L_{2}^{2} H-16 L_{2}^{2} L_{1}+4 c^{2} H^{2} L_{2}+\left(4 \mathrm{i} \beta-2 c^{4} \gamma\right) H^{2}+8 \mathrm{i} c^{2} \alpha L_{2}^{2} \\
& +\left(4 c^{4} \gamma^{2}-16 \mathrm{i} \beta \gamma\right) L_{1}+\left(4 \alpha^{2} c^{2}-c^{6} \gamma^{2}+4 \mathrm{i} c^{2} \gamma \beta\right) L_{2} \\
& +\left(-2 c^{6} \gamma^{2}+8 \mathrm{i} c^{2} \gamma \beta+8 \beta \alpha\right) H+2 c^{4} \alpha^{2} \gamma-2 \mathrm{i} c^{6} \gamma^{2} \alpha-8 c^{2} \alpha \beta \gamma-4 \mathrm{i} \alpha^{2} \beta .
\end{aligned}
$$

Case (4): $\quad[0,1,0,0,0],[0,0,0,1,0]$
$R^{2}=4 H L_{1}^{2}+4 H L_{2}^{2}+4\left(\beta^{2}-\gamma^{2}\right) L_{2}-8 \beta \gamma L_{1}-4 \alpha^{2} H-4 \alpha\left(\beta^{2}+\gamma^{2}\right)$.
Case (5a): $\quad[0,-4 \mathrm{i}, 2,1,2 \mathrm{i}],[0,0,2,0,-2 \mathrm{i}]$
$R^{2}=64 L_{2}^{3}-64 \gamma H^{2}-128 \alpha L_{2}^{2}+128 \beta H L_{2}$

$$
+64 \gamma L_{2} L_{1}+64 \alpha^{2} L_{2}+64 \beta^{2} L_{1}-128 \beta \alpha H
$$

Case (6a): $\quad[0,0,1,0,0],[0,0,0,1,0]$
$R^{2}=16 L_{1}^{3}-32 L_{1}^{2} H+16 H^{2} L_{1}-16 \alpha L_{2}^{2}-8 \beta H L_{2}+8 \beta L_{2} L_{1}-64 \alpha \gamma L_{1}-4 \beta^{2} \gamma$.

Case (6b): $\quad[0,0,-2 \mathrm{i}, 0,2],[0,0,0,1,0]$
$R^{2}=2 \mathrm{i} L_{1}^{3}+L_{1}^{2} H-\beta H L_{2}-2 \mathrm{i} \beta L_{2} L_{1}-\gamma^{2} L_{2}-\mathrm{i} \alpha \gamma L_{1}+\frac{1}{4} \beta \alpha^{2}$.

## 3. Quantum superintegrability in two-dimensional Euclidean space

Here we give the analogous quantum algebras for superintegrable systems arising from the potentials we have already computed. The only difference is that the Poisson bracket is now replaced by the commutator bracket $[A, B]=A B-B A$ and the operators $H, L_{1}$ and $L_{2}$ are the obvious (formally self-adjoint) symmetry partial differential operators:
$H=\partial_{x}^{2}+\partial_{y}^{2}+V(x, y) \quad L_{h}=\sum_{k, j=1}^{2} \partial_{k}\left(a_{(h)}^{k j}\right) \partial_{j}+W_{(h)}(x, y) \quad h=1,2$.
Just as for the Hamilton-Jacobi case, if we have another constant of motion $L$ associated with a maximal potential, then $L$ must be a linear combination of $H, L_{1}, L_{2}$. Indeed, if $L$ is in self-adjoint form, then the conditions that $[H, L]=0$ are identical with (11) and (12). Thus, if $L$ is not a linear combination of the basis functions, then the potential $V$ must satisfy an equation (16) that is linearly independent of the equations associated with $L_{1}, L_{2}$. This means
an additional constraint on the solution space and that $V$ can depend on at most two parameters, which is a contradiction.

Furthermore the proof of the existence of quadratic algebra relations at the end of section 2 goes through almost unchanged for the operator case: $\left[L_{1}, L_{2}\right]^{2}=R^{2}$ and $\left[L_{1}, R\right],\left[L_{2}, R\right]$ can be expressed as (symmetric) polynomials in the operators $H, L_{1}, L_{2}$. To make the prior construction go through, one need only note that since $R^{2}$ is a formally self-adjoint sixth-order differential symmetry operator, the fifth-order terms are fixed linear functions of the sixthorder terms. The expressions $\{A, B\}=A B+B A$ and $\{A, B, C\}=A B C+C A B+B C A$ are operator symmetrizers. The explicit relations are as follows.

## Case (1):

$$
\begin{aligned}
& {\left[L_{2}, R\right]=8 L_{2}^{2}+8 H L_{2}+8 \alpha^{2}} \\
& {\left[L_{1}, R\right]=8\left\{L_{2}, L_{1}\right\}+16(1+2 \beta+2 \gamma) L_{2}+32 \alpha(\gamma-\beta)} \\
& R^{2}=16 L_{1}^{2} H-\frac{8}{3}\left\{L_{2}, L_{2}, L_{1}\right\}-16\left(2 \beta+2 \gamma+\frac{11}{3}\right) L_{2}^{2}-\frac{176}{3} H L_{1}+64 \alpha(\beta-\gamma) L_{2} \\
& \quad+16 \alpha^{2} L_{1}+\left(-\frac{32}{3}+96 \gamma+96 \beta+256 \beta \gamma\right) H-\frac{32}{3} \alpha^{2}(3 \beta+3 \gamma-1) .
\end{aligned}
$$

Case (2a):

$$
\begin{aligned}
& {\left[L_{1}, R\right]=-\frac{1}{2} \mathrm{i} c^{3} H^{2}-6 \mathrm{i} c L_{2}^{2}+2 c^{2} H L_{2}-2\left\{L_{1}, L_{2}\right\}} \\
& \quad+\left(2 \mathrm{i} \alpha c^{2}-1\right) L_{2}-\frac{1}{2} \mathrm{i} c^{3} \alpha^{2}+\frac{1}{2} \mathrm{i} c \beta^{2}+\frac{1}{2} \mathrm{i} c \gamma^{2} \\
& {\left[L_{2}, R\right]=\frac{1}{2} c^{2} H^{2}+2 L_{2}^{2}+\frac{1}{2} c^{2} \alpha^{2}} \\
& R^{2}=\frac{1}{2} c^{4} H^{3} \\
& -4 \mathrm{i} c L_{2}^{3}+2 c^{2} L_{2}^{2} H-\frac{2}{3}\left\{L_{1}, L_{1}, L_{2}\right\}-c^{2} H^{2} L_{1}-\mathrm{i} c^{3} H^{2} L_{2} \\
& \\
& \quad+\left(\frac{\mathrm{i}}{2} c^{2} \alpha+\frac{1}{12}\right) H^{2}+\left(2 \mathrm{i} \alpha c^{2}-\frac{11}{3}\right) L_{2}^{2}-c^{2} \alpha^{2} L_{1} \\
& \\
& \quad+\mathrm{i} c\left(\beta^{2}+\gamma^{2}-c^{2} \alpha^{2}\right) L_{2}+\frac{1}{2}\left(-c^{2} \beta^{2}+c^{4} \alpha^{2}+c^{2} \gamma^{2}\right) \\
& \\
& \quad+\frac{1}{2}\left(2 \beta \gamma+\mathrm{i} c^{2} \alpha^{2}+\frac{1}{12} \alpha\right) c^{2} \alpha .
\end{aligned}
$$

Case (2b):

$$
\begin{aligned}
& {\left[L_{1}, R\right]=L_{1}^{2}-3 L_{2}^{2}+\left\{L_{1}, L_{2}\right\}+16 \alpha H+L_{1}-L_{2}+4 \beta, \gamma} \\
& {\left[L_{2}, R\right]=3 L_{1}^{2}-L_{2}^{2}-\left\{L_{1}, L_{2}\right\}-16 \alpha H+L_{1}-L_{2}+4 \beta, \gamma} \\
& R^{2}=-2 L_{1}^{3}-2 L_{2}^{3}+\frac{1}{3}\left\{L_{2}, L_{2}, L_{1}\right\}+\frac{1}{3}\left\{L_{1}, L_{1}, L_{2}\right\}-\frac{11}{3} L_{1}^{2}-\frac{11}{3} L_{2}^{2} \\
& \\
& \quad+32 \alpha H L_{1}+32 \alpha H L_{2}+\frac{11}{3}\left\{L_{1}, L_{2}\right\}-8 \gamma \beta L_{1}+8 \gamma \beta L_{2} \\
& \quad+\left(-\frac{16}{3} \alpha+16 \beta^{2}\right) H+16 \alpha \gamma^{2} .
\end{aligned}
$$

Case (3a):

$$
\begin{aligned}
& {\left[L_{2}, R\right]=-8 L_{2}^{2}+8 H L_{2}-16 \alpha L_{1}+8 \alpha} \\
& {\left[L_{1}, R\right]=-8 H L_{1}+8\left\{L_{2}, L_{1}\right\}-8(1+2 \beta) H+16(1+\beta+\gamma) L_{2}} \\
& R^{2}=-\frac{8}{3}\left\{L_{2}, L_{2}, L_{1}\right\}+8 H\left\{L_{2}, L_{1}\right\}-4(3+4 \beta) H^{2}-16\left(\beta+\gamma-\frac{11}{3}\right) L_{2}^{2}-16 \alpha L_{1}^{2} \\
& \quad+16\left(2 \beta+\frac{11}{3}\right) H L_{2}+\frac{176}{3} \alpha L_{2}+16 \alpha\left(3 \beta+3 \gamma+4 \beta \gamma+\frac{2 \alpha}{3}\right) .
\end{aligned}
$$

## Case (3b):

$$
\begin{aligned}
& {\left[L_{2}, R\right]=-8 L_{2}^{2}-32 \alpha \gamma} \\
& {\left[L_{1}, R\right]=16\left\{L_{2}, L_{1}\right\}-8 \beta H+16 L_{2}} \\
& R^{2}=-\frac{8}{3}\left\{L_{2}, L_{2}, L_{1}\right\}+16 \alpha H^{2}-\frac{176}{3} L_{2}^{2}+16 \beta H L_{2}-64 \alpha \gamma L_{1}+\left(\frac{64}{3} \alpha \gamma-16 \beta^{2} \gamma\right) .
\end{aligned}
$$

Case (3c):

$$
\begin{aligned}
& {\left[L_{2}, R\right]=L_{2}^{2}+8 \mathrm{i} \beta \gamma-2 c^{4} \gamma^{2}} \\
& {\left[L_{1}, R\right]=2 c^{2} H^{2}+6 c^{2} L_{2}^{2}+8 c^{2} H L_{2}-8\left\{L_{1}, L_{2}\right\}} \\
& +\left(-16+8 \mathrm{i}^{2} \alpha\right) L_{2}-\frac{1}{2} c^{6} \gamma^{2}+2 \alpha^{2} c^{2}+2 \mathrm{i} c^{2} \beta \gamma \\
& R^{2}=4 c^{2} L_{2}^{3}+8 c^{2} L_{2}^{2} H-\frac{8}{3}\left\{L_{2}, L_{2}, L_{1}\right\} \\
& +4 c^{2} H^{2} L_{2}+\left(4 \mathrm{i} \beta-2 c^{4} \gamma\right) H^{2}+\left(8 \mathrm{i}^{2} \alpha-\frac{176}{3}\right) L_{2}^{2} \\
& +\left(4 c^{4} \gamma^{2}-16 \mathrm{i} \beta \gamma\right) L_{1}+\left(4 \alpha^{2} c^{2}-c^{6} \gamma^{2}+4 \mathrm{i} c^{2} \gamma \beta\right) L_{2} \\
& +\left(-2 c^{6} \gamma^{2}+8 \mathrm{i} c^{2} \gamma \beta+8 \beta \alpha\right) H+2 c^{4} \alpha^{2} \gamma-2 \mathrm{i} c^{6} \gamma^{2} \alpha \\
& -8 c^{2} \alpha \beta \gamma-4 \mathrm{i} \alpha^{2} \beta-\frac{4}{3} \gamma^{2} c^{4}+\frac{16 \mathrm{i}}{3} \beta \gamma .
\end{aligned}
$$

Case (4):
$\left[L_{2}, R\right]=4 H L_{1}-4 \beta \gamma$
$\left[L_{1}, R\right]=-4 H L_{2}+2\left(\gamma^{2}-\beta^{2}\right)$
$R^{2}=4 L_{1}^{2} H+4 L_{2}^{2} H+4 H^{2}+4\left(\beta^{2}-\gamma^{2}\right) L_{2}-8 \beta \gamma L_{1}-4 \alpha^{2} H-4 \alpha\left(\gamma^{2}+\beta^{2}\right)$.
Case (5a):

$$
\begin{aligned}
& {\left[L_{2}, R\right]=32 \gamma L_{2}+32 \beta^{2}} \\
& {\left[L_{1}, R\right]=-96 L_{2}^{2}-64 \beta H+128 \alpha L_{2}-32 \gamma L_{1}-32 \alpha^{2}} \\
& R^{2}=64 L_{2}^{3}-64 \gamma H^{2}-128 \alpha L_{2}^{2}+128 \beta H L_{2}+32 \gamma\left\{L_{2}, L_{1}\right\}+64 \alpha^{2} L_{2} \\
& \quad+64 \beta^{2} L_{1}-128 \beta \alpha H-256 \gamma^{2} .
\end{aligned}
$$

Case (6a):

$$
\begin{aligned}
& {\left[L_{2}, R\right]=8 H^{2}+24 L_{1}^{2}-32 L_{1} H+4 \beta L_{2}-\alpha(24+32 \gamma)} \\
& {\left[L_{1}, R\right]=4 \beta H+16 \alpha L_{2}-4 \beta L_{1}} \\
& R^{2}=16 L_{1}^{3}-32 H L_{1}^{2}+16 H^{2} L_{1}-16 \alpha L_{2}^{2}-8 \beta H L_{2}+4 \beta\left\{L_{2}, L_{1}\right\} \\
& \quad-(64 \alpha \gamma+176 \alpha) L_{1}+128 \alpha H-\left(4 \gamma \beta^{2}+3 \beta^{2}\right) .
\end{aligned}
$$

Case (6b):
$\left[L_{2}, R\right]=3 \mathrm{i} L_{1}^{2}+H L_{1}-\mathrm{i} \beta L_{2}-\frac{\mathrm{i}}{2} \alpha \gamma$
$\left[L_{1}, R\right]=\frac{1}{2} \beta H+\mathrm{i} \beta L_{1}+\frac{1}{2} \gamma^{2}$
$R^{2}=2 \mathrm{i} L_{1}^{3}+L_{1}^{2} H-\beta H L_{2}-\mathrm{i} \beta\left\{L_{2}, L_{1}\right\}-\gamma^{2} L_{2}-\mathrm{i} \alpha \gamma L_{1}+\frac{1}{4}\left(\beta^{2}+\alpha^{2} \beta\right)$.
We note that the quadratic relations in the quantum case provide useful information relating the special functions that occur as (separable) eigenfunctions for each superintegrable case [16].

## 4. Conclusions

In this paper we have used the concept of a 'nondegenerate potential' to add structure to the study of superintegrable classical and quantum mechanical systems in $E(2, C)$. We have shown how to classify all such systems in a straightforward manner, so that gaps can be avoided. Furthermore, we have shown the following:
(1) Each system is associated with a pair of constants of motion in the classical case, and a pair of symmetry operators in the quantum case, that generate a quadratic algebra.
(2) There is a one-to-one correspondence between superintegrable systems and free-field symmetry operators that generate quadratic algebras.
(3) Second-order superintegrability implies multiseparability, i.e., separability in more than one coordinate system.

In a forthcoming paper we will prove the analogous results for superintegrable systems on the complex 2 -sphere.

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## Appendix

As is well known $[4,18,20]$ there are essentially six coordinate systems on the complex Euclidean plane in which the free-particle Hamilton-Jacobi equation separates: cartesian, polar, parabolic, elliptic, hyperbolic and semi-hyperbolic. We describe these coordinate systems and their corresponding free-particle constants of motion $L$. (We adopt the basis $p_{x}, p_{y}, M=x p_{y}-y p_{x}$ for the Lie algebra $e(2, C)$ and define $p_{ \pm}=p_{x} \pm \mathrm{i} p_{y}$.) There is one orbit of constants of motion, with representative $M p_{+}$, that is not associated with variable separation [21]. The systems are as follows.

## Cartesian coordinates

$$
\begin{equation*}
x, y \quad L=p_{x}^{2} . \tag{61}
\end{equation*}
$$

## Polar coordinates

$$
\begin{equation*}
x=r \cos \theta \quad y=r \sin \theta \quad L=M^{2} \tag{62}
\end{equation*}
$$

Parabolic coordinates

$$
\begin{equation*}
x_{P}=\frac{1}{2}\left(\xi^{2}-\eta^{2}\right) \quad y_{P}=\xi \eta \quad L=M p_{y} \tag{63}
\end{equation*}
$$

## Elliptic coordinates

(in algebraic form)

$$
\begin{align*}
& x_{E}^{2}=c^{2}(u-1)(v-1) \quad y_{E}^{2}=-c^{2} u v  \tag{64}\\
& L=M^{2}+c^{2} p_{x}^{2} .
\end{align*}
$$

## Hyperbolic coordinates

$$
\begin{align*}
& x_{H}=\frac{r^{2}+r^{2} s^{2}+s^{2}}{2 r s} \quad y_{H}=\mathrm{i} \frac{r^{2}-r^{2} s^{2}+s^{2}}{2 r s}  \tag{65}\\
& L=M^{2}+p_{+}^{2} .
\end{align*}
$$

## Semi-hyperbolic coordinates

$$
\begin{array}{ll}
x_{S H}=-\frac{1}{4}(w-u)^{2}+\frac{1}{2}(w+u)  \tag{66}\\
L=2 M p_{+}+p_{-}^{2}
\end{array}
$$

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