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1. Lukáč, J. Smorodinsky, P. Winternitz

QUANTUM NUMBERS IN THE LITTLE GROUPS OF THE POINCARE GROUP

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I. Lukád, J. Smorodinsky, P. Winternitz

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## I. Introduction

We shall investigate the separation of variables in the equation

$$
\begin{equation*}
\Delta \psi=\lambda \psi \tag{1}
\end{equation*}
$$

where $\Delta$ is the Laplace operator in two-dimensional spaces with constant posttive, zero or negative curvature, and demonstrate how the choice of the coordlnate system is connected with the choice of an operator, diagonalized simultaneously with the Laplace operator, i.e. with the choice of certain integrals of motion and quantum numbers. (We call an operator $L$ diagonal, if the elgenfunctions of (1) satisfy the equation $L \psi=\mu \psi$ ).

Such a problem is interesting from the physical point of view for several reasons.

Firstly, the obtained results can be applled to construct representations of the Poincaré group $/ 1,2 /$ slnce the groups of motions of the considered spaces, i.e. the three-dimensional rotation group $O(3)$, the group of motions of the eucll dean plane $E$, and the three-dimensional Lorentz group $O(2,1)$ are little groups of the Poincare group, corresponding to a time-like, isotroplc or space-like/3/ vector $x /$. It follows that we are actually considering possible parametrizations of the relativistic spin in the physical or unphysical regions (or on the light cone). Such parametrizations should prove useful e.g. in the the analytical continuation of wave functions and amplitudes for particles with spin.

Secondly, the representation theory of compact and speclally noncompact groups finds many applications in physics, e.g. in investigations concerning relativistic generalizations of the SU (6) symmetry; or in the obtaining of complete sets of functions with definite transformation properties, in terms of which physical quantities, like reaction amplitudes, can be expanded. The cholce of the

[^0]diagonal operator $L$ determines the choice of the basls of the corresponding group representation. An analogous problem was first considered in connection with the invariant expansions of relativistic amplitudes in terms of the basls functions of the unltary representations of the Lorentz group $/ 4 /$. It was shown in $5,6 /$, that the diagonalization of operators, which are invariants of the subgroups of the Lorentz group, leads to coordinate systems with one centre. The dlagonelization of other operators, quadratical in the generators of the Lorentz group $/ 6,7 /$ leads to coordinate systems of the elliptical or parabolical type.

Thirdly, in the investigation of higher symmetrles in quantum mechanics, specifically in a search for all potentials, having a symmetry group larger than the evident geometrical symmetry group, it also proved necessary to find all operators, commuting with the Laplace operator and leading to the separation of variables $k, 9$ ?

We shall prove the following assertion:
Theorem. A linear self-adjoint operator $L_{k}$, being a homogeneous quadratical polynomial in the generators of the group of motions of the considered two dimensional space with constant curvature, corresponds to each system of coordinates In which the variables in equation (1) separate. The operator $L_{k}$ is determined by the condition $I_{K} \psi=\mu \psi$, where $\psi$ are solutions of (1), separated in K.As K runs through all types of coordinate systems allowing the separation of variables, $L_{k}$ runs through all types of non-equivalent operators in the considered class. The operators $L$ and $\mathcal{L}$ are considered equivalent, if the relation

$$
\begin{equation*}
\tilde{L}=a L^{\prime}+\beta \Delta \tag{2}
\end{equation*}
$$

hoids, where $a, \beta$ are real constants, and a motion of the space exists, transforming $L$ into $L^{\prime}$.

An analogous theorem for the euclidean space $E_{2}$ has been proved in $/ \sigma /$. Here we shall consider the cases $k>0, k=0$ and $k<0$ separately.

It has been shown by Olevsky $10 /$ that two coordinate systems exist in $\mathrm{R}_{2}$ In which the variables in (1) separate. It can be shown that the operator

$$
\begin{equation*}
L_{s}=L_{3}^{2} \tag{3}
\end{equation*}
$$

Is diagonal in the spherical system

$$
x=\sin \rho_{1} \sin \rho_{2} \quad y=\sin \rho_{1} \cos \rho_{2} \quad z=\cos \rho_{1}
$$

and the operator ${ }^{\text {. }}$

$$
\begin{equation*}
L_{E}=L_{z}^{2}+\sin ^{2} 6 L_{2}^{2} \tag{4}
\end{equation*}
$$

is diagonal in the elliptical system

$$
\begin{gathered}
x^{2}=\frac{\left(\rho_{1}-a\right)\left(\rho_{2}-a\right)}{(c-a)(b-a)} y^{2}=\frac{\left.\left(\rho_{1}-b\right) \rho_{2}-b\right)}{(c-b)(a-b)} \quad z^{2}=\frac{\left(\rho_{1}-c\right)\left(\rho_{2}-c\right)}{(a-c)(b-c)} \\
c<\rho_{2}<b<\rho_{2}<a
\end{gathered}
$$

Here $\sin ^{2} f=\frac{a-b}{a-c}, \quad 2 f \quad$ Is the distance between the focusses; $L, i=1,2,3$ are the generators of the group $O(3)$ and the Laplace operator is $\Delta=L_{1}{ }^{2}+L_{2}^{2}+L_{8}^{2}$

To prove the theorem stated in the Introduction we must show that any operator of the type

$$
\begin{equation*}
L=A_{1 k} L_{1} L_{k} \quad A_{1 k}-A_{k 1} \tag{5}
\end{equation*}
$$

(summation from 1 to 3 over repeated indices) is equivalent to $L_{s}$ or $L_{E}$ (or to the zero operator). However, this is obvious since under a motion of $R_{2}$ the generators $L_{1}$ transform like a vector

$$
L_{1}=a_{1 k} L_{k},
$$

where $\mathbf{a}_{1 k}$ are elements of a real orthogonal matrix. The polynomial $L$ is transformed into

$$
L^{\prime}=A_{1 k}^{\prime} L_{1} L_{k}
$$

where $A_{\text {it }}^{\prime}-a_{\ell_{1}} A_{\ell_{m}} a_{m k}$ i.e. in matrix form we have

$$
\begin{equation*}
A^{\prime}=a^{T} A \quad a^{T} \cdot=1 \tag{6}
\end{equation*}
$$

(the letter $T$ means transposition).
It is well known that any real symmetrical matrix $\Lambda_{i k}$ can be diagonalized by an orthogonal transformation (6), so that we have to consider only three possibllities.

1) All three elgenvalues $\lambda_{1}$ of $A$ are different (they are naturally always real). Transformation (6) can be chosen so that

$$
\begin{equation*}
L^{\prime}-\lambda_{1} L_{1}^{2}+\lambda_{2} L_{2}^{2}+\lambda_{8} L_{3}^{2} \tag{7}
\end{equation*}
$$

Let us consider the case $\lambda_{1}>\lambda_{2}>\lambda_{8}$ (any other inequality simply corresponds to an interchange of the axes). We then have

$$
I-\frac{1}{\lambda_{B}-\lambda_{1}}\left(L^{\prime}-\lambda_{1} \Delta\right)=L_{E} \text {. }
$$

2) Two of the elgenvalues coinclde $\left(\lambda_{1}=\lambda_{2} \nleftarrow \lambda_{s}\right)$ : the operator $L^{\prime}$ can be written as

$$
\begin{equation*}
L^{\prime}=\lambda_{1}\left(L_{1}^{2}+L_{2}^{\prime}\right)+\lambda_{3} I_{3}^{2} \tag{8}
\end{equation*}
$$

and

$$
\tilde{L}=\frac{1}{\lambda_{3}-\lambda_{1}}\left(L^{\prime}-\lambda_{1} \Delta\right)=L_{0}
$$

3) All three eigenvalues colncide $\left(\lambda_{1}-\lambda_{2}=\lambda_{3}\right)$.

Thus

$$
\begin{equation*}
L^{\prime}=\lambda_{1}\left(L_{1}^{2}+L_{2}^{2}+L_{g}^{2}\right) \tag{9}
\end{equation*}
$$

and

$$
I=L^{\prime}-\lambda_{1} \Delta=0
$$

This completes the proof of the theorem for $R_{2}$.

## 3. Integrals of Motion on the Euclldean Plane $E_{2}$ (Case of an Isotropical

Vector)

The varlables in equation (1) can be separated in four coordinate systems in the space $E_{2} / 11 /$. Let us denote the generators of the group of mottons $P_{1}, P_{2}$ (Infinitesimal translations) and $M$ (an Infinitesimal rotation). The Laplace operator is $\Delta=P_{1}^{2}+P_{2}^{2}$. The separable coordinate systems and corresponding diagonal operators are:

1) Cartesian system $x, y$ with the operator

$$
\begin{equation*}
\mathrm{L}_{\mathrm{D}}=\mathrm{P}_{1}^{2}-\mathrm{P}_{2}^{2} \tag{10}
\end{equation*}
$$

2) The polar system $x=r \cos \phi, y=r \sin \phi \quad$ with the operator
3) The elliptical system $x=\ell L_{s}-M^{2}$
4) $y=\ell \sqrt{\left(\xi^{2}-1\right)\left(1-\eta^{2}\right)} \quad(R>0$ is the focal distance) with the operator

$$
\begin{equation*}
L_{E}=\Lambda^{2}+\frac{\ell^{2}}{2}\left(P_{1}^{2}-P_{2}^{2}\right) \tag{i2}
\end{equation*}
$$

4) The parabolical system $x=2 / 2\left(\xi^{2}-\eta^{2}\right)$; $y=\xi$ with the operator

$$
\begin{equation*}
L_{D}=M P_{2}+P_{2} M \tag{13}
\end{equation*}
$$

It has been proved in $/ 6 /$ that the most general operator of the consldered type

$$
\begin{equation*}
L=a M^{2}+b_{1}\left(M P_{1}+P_{1} M\right)+b_{2}\left(M P_{2}+P_{2} M\right)+c_{1} P_{1}^{2}+2 c_{2} P_{1} P_{2}+c_{s} P_{2}^{2} \tag{14}
\end{equation*}
$$

is equivalent (in the sense of equation (2)) to one of these four types.

## 4. Integrals of Motion on the Hyperbolold $L_{2}$ (Case of a SpaceLike Vector)

Nine types of separable coordinate systems exist in the two-dimensional space $L_{2} / 10 /$. Let us denote the generators of the group $0(2,1)-K_{1}, K_{2}$ (infinitesimal hyperbolic rotations) and $M_{3}$ (an infinitesimal space rotation), the Laplace operator is $\Delta-K_{1}^{2}+K_{2}^{2}-M_{3}^{2}$.

Let us first enumerate these coordinate systems and the corresponding diagonal operators. As usual we denote the Welerstrass coordinates $x, y$, $t$ (for simplicity we set $k=-1$ ).

1) Horocyclic system

$$
\begin{gather*}
\left.x=1 / 2\left[e^{\rho_{1}}+e^{-\rho_{1}}\left(\rho_{2}^{2}-1\right)\right] y=\rho_{2} e^{-\rho_{1}} t=1 / 2 e^{\rho_{1}}+e^{-\rho_{1}}\left(\rho_{2}^{2}+1\right)\right] \\
\left.L_{0}-\left(K_{1}+M_{3}\right)^{2}-K_{1}^{2}+M_{3}^{2}+K_{1} M_{3}+M_{3} K_{1}\right] \tag{15}
\end{gather*}
$$

2) Equidistant system

$$
\begin{array}{ll}
x=\operatorname{ch} \rho_{1} \operatorname{sh} \rho_{2} & y=\operatorname{sh}_{1} \quad t=\operatorname{ch} \rho_{1} \operatorname{ch} \rho_{2} \\
L_{E q}=K_{2}^{2} \tag{16}
\end{array}
$$

3) Spherical system

$$
\begin{gather*}
x=\operatorname{ch} \rho_{1} \cos \rho_{2} \quad y=\operatorname{sh} \rho_{1} \sin \rho_{2} \quad t-\operatorname{ch} \rho_{1} \\
L_{B}=M_{8}^{2}, \tag{17}
\end{gather*}
$$

4) Ellptical system

$$
\begin{gather*}
\mathbf{x}^{2}=\frac{\left(\rho_{1}-b\right)\left(\rho_{2}-b\right)}{(a-b)(b-c)} \quad y^{2}=\frac{(\rho-a)\left(a-\rho_{2}\right)}{(a-b)(a-c)} c^{2}=\frac{\left(\rho_{1}-c\right)\left(a-\rho_{2}\right)}{(a-c)(b-c)} \\
c<b<\rho_{2}<a<\rho_{1} \\
L_{E}=M_{B}^{2}+\operatorname{sh}^{2} f K_{2}^{2} \tag{18}
\end{gather*}
$$

where $\operatorname{sh}^{2} t=\frac{a-b}{5-c}$ and $2 f$ is the focal distance
5) Hyperbolic system

$$
\begin{gather*}
\mathrm{x}=\frac{\left(\rho_{1}-c\right)\left(c-\rho_{2}\right)}{(a-c)(b-c)} \quad y^{2}-\frac{\left(\rho_{1}-a\right)\left(a-\rho_{2}\right)}{(a-b)(a-c)} t^{2}=\frac{\left(\rho_{1}-b\right)\left(b-\rho_{2}\right)}{(a-b)(b-c)} \\
\rho_{a}<c<b<a<\rho_{1} \\
L_{B}=K_{2}^{2}-\sin ^{2} a L_{8}^{2} \tag{19}
\end{gather*}
$$

where $\sin ^{2} \alpha^{2} \frac{b-c}{a-c}$ and $2 a$ Is the angle between two focal lines.
6) Semihyperbolic system

$$
\begin{align*}
& x^{2}=-\frac{\left(\rho_{1}-a\right)\left(a-\rho_{2}\right)}{2\left[(a-\gamma)^{2}+\delta^{2}\right]}-\gamma+\frac{1}{2 \delta} \sqrt{\frac{\left[\left(\rho_{1}-\gamma\right)^{2}+\delta^{2}\right]\left[\left(\rho_{2}-\gamma\right)^{2}+\delta^{2}\right]}{(a-\gamma)^{2}+\delta^{2}}} \\
& \bar{y}^{2}=\frac{\left(\rho_{1}-a\right)\left(a-\rho_{2}\right)}{(a-\gamma)^{2}+\delta^{2}} \quad \rho_{2}<a<\rho_{1} \\
& t^{2}=\frac{\left(\rho_{1}-a\right)\left(a-\rho_{2}\right)}{2\left[(a-\gamma)^{2}+\delta^{2}\right]}+1 /+\frac{1}{2 \delta} \sqrt{\frac{\left[\left(\rho_{1}-\gamma\right)^{2}+\delta^{2}\right]\left[\left(\rho_{2}-\gamma\right)^{2}+\delta^{2}\right]}{(a-\gamma)^{2}+\delta^{2}}} \\
& \mathrm{~L}_{\mathrm{BH}}=\mathrm{M}_{\mathrm{s}} \mathrm{~K}_{1}+\mathrm{K}_{1} \mathrm{M}_{3}+\operatorname{sh} 2 \mathrm{fK}_{2}^{2} \text {, } \tag{20}
\end{align*}
$$

where $\operatorname{sh} 2 f=\frac{a-\gamma}{\delta}$ and $2 f$ is the distance between the focus of the semihyperbolas and the basis line of their equidistant curve:
7) Elliptic-parabolical system

$$
\begin{align*}
& 2 x^{\prime}=\frac{\left(\rho_{1}-a\right)\left(a-\rho_{2}\right)}{(a-b)^{3 / 2}\left[\left(\rho_{1}-b\right)\left(\rho_{2}-b\right)\right]^{1 / 2}}+\left[\frac{a-b}{\left(\rho_{1}-b\right)\left(\rho_{2}-b\right)}\right]^{b / 2}-\left[\frac{\left(\rho_{1}-b\right)\left(\rho_{2}-b\right)}{a-b}\right]^{1 / 2} \\
& 2 t=\frac{\left(\rho_{1}-a\right)\left(a-\rho_{2}\right)}{(a-b)^{3 / 2}\left[\left(\rho_{1}-b\right)\left(\rho_{2}-b\right)\right]^{1 / 2}}+\left[\frac{\left(\rho_{1}-a\right)\left(a-\rho_{2}\right)}{\left(\rho_{1}-b\right)\left(\rho_{2}-b\right)}\right]^{1 / 2}+\left[\frac{\left(\rho_{1}-b\right)\left(\rho_{2}-b\right)}{}\right]^{1 / 2} \\
& L_{E P}=(a-b) K_{2}^{2}+K_{1}^{2}+M_{3}^{2}+K_{1} M_{B}+M_{3} K_{1}
\end{align*}
$$

8) Hyperbolic-parabolical system.

$$
\begin{align*}
& 2 x=\frac{\left(\rho_{1}-a\right)\left(a-\rho_{2}\right)}{(a-b)^{g / 2}\left[\left(\rho_{1}-b\right)\left(b-\rho_{2}\right)\right]^{k}}+\left[\frac{a-b}{\left(\rho_{1}-b\right)\left(b-\rho_{2}\right)}\right]^{1 /}-\left[\frac{\left(\rho_{1}-b\right)\left(b-\rho_{2}\right)}{a-b}\right]^{1 / 4} \\
& y^{2}=\frac{\left(\rho_{1}-a\right)\left(a-\rho_{2}\right)}{(a-b)^{2}} \quad \rho_{2}<b<a<\rho_{1} \\
& 2 t=\frac{\left(\rho_{1}-a\right)\left(a-\rho_{2}\right)}{(a-b)^{8 / 2}\left[\left(\rho_{1}-b\right)\left(b-\rho_{2}\right)\right]^{1 / 2}}+\left[\frac{a-b}{\left(\rho_{1}-b\right)\left(b-\rho_{2}\right)}\right]^{1 /}+\left[\frac{\left(\rho_{1}-b\right)\left(b-\rho_{2}\right)}{a-b}\right]^{1 / 2} \\
& L_{H P}=-(\mathrm{a}-\mathrm{b}) \mathrm{K}_{\mathrm{i}}^{2}+\mathrm{K}_{1}^{2}+M_{9}^{2}+\left(\mathrm{K}_{1} M_{3}+M_{3} K_{1}\right) \text {, } \tag{22}
\end{align*}
$$

9) Semicircular-parabolical system

$$
\begin{align*}
& z=\frac{\left(\rho_{1}-\rho_{2}\right)^{2}}{8\left[\left(\rho_{1}-a\right)\left(a-\rho_{2}\right)\right)^{\beta / 2}}-1 /\left[\left(\rho_{1}-a\right)\left(a-\rho_{2}\right)\right]^{3} \\
& 2 g=\left[\frac{\rho_{1}-a}{a-\rho_{2}}\right]^{1 /}-\left[\frac{a-\rho_{2}}{\rho_{1}-a}\right]^{1 / 2} \rho_{2}<a<\rho_{1} \\
& =\frac{\left(\rho_{1}-\rho_{2}\right)^{2}}{\left.8\left[\left(\rho_{1}-a\right)\left(a-\rho_{2}\right)\right]^{8 / 2}+1 / A\left(\rho_{1}-a\right)\left(a-\rho_{2}\right)\right]^{K}} \\
& L_{C P}=K_{1} K_{2}+K_{2} K_{1}+K_{2} M_{3}+M_{8} K_{2} \tag{23}
\end{align*}
$$

We have Introduced names for the coordinate systems, corresponding to the form of the coordinate lines. The definitions of the corresponding curves can be found e.g. in $12 /$.

As in the cases considered above we must prove that any operator of the type

$$
\begin{equation*}
L=a K_{1}^{2}+b\left(K_{1} K_{2}+K_{2} K_{1}\right)+c K_{2}^{2}+d\left(K_{1} M_{2}+M_{3} K_{1}\right)+c\left(K_{2} M_{3}+M_{3} K_{2}\right)+f M_{s}^{2} \tag{24}
\end{equation*}
$$

Is equivalent (cf. equation (2)) to one of the operators (15)-(23).
Proof.
Let us consider the symmetrical matrix

$$
x=\left(\begin{array}{lll}
a & b & d  \tag{25}\\
b & c & c \\
d & e & f
\end{array}\right)
$$

and the characteristic equation

$$
\operatorname{det} X(\mu)=\left(\begin{array}{ccc}
a-\mu & b & d  \tag{26}\\
b & c-\mu & c \\
d & e & f+\mu
\end{array}\right)=0
$$

Under a motion of the space $L_{2}$ (hyperbolic motion) the operator $L$ transforms into $L^{\prime}$, determined by the matrix $X^{\prime}$, where

$$
\begin{equation*}
X^{\prime}=a^{T} X a \tag{27}
\end{equation*}
$$

and $a$ is the matrix of a hyperbolic motion i.e.

$$
a^{T} I a-I \quad I=\left(\begin{array}{ccc}
1 & 0 & 0  \tag{28}\\
0 & 1 & 0 \\
0 & 0 & -1
\end{array}\right)
$$

Hyperbolic motions are investigated e.g. in $/ 12,13 /$ in connection with the classification of second order curves (conics) on the Lobachevsky plane. It is shown that the hyperbolic motions leave the following set of quantities invariant

1) Roots of the characteristic equation (26): $\mu_{1}, \mu_{2}, \mu_{3}$; or the quantities

$$
\begin{aligned}
& S=a+c-f=\mu_{1}+\mu_{2}+\mu_{8} \\
& T=A+C-F=-\left(\mu_{1} \mu_{2}+\mu_{2} \mu_{8}+\mu_{\mathrm{g}} \mu_{1}\right) \\
& R=\operatorname{det} X--\mu_{1} \mu_{2} \mu_{\mathrm{s}},
\end{aligned}
$$

where $A, C$ and $F$ are the minors of $X$ corresponding to the elements $a, c$ and $f$.
2) The quantitles
$r_{1}=$ rank of $X\left(\mu_{1}\right)$
$\epsilon_{1}=\operatorname{sgn}\left[A\left(\mu_{1}\right)+C\left(\mu_{1}\right)+F\left(\mu_{1}\right)\right]$
$K_{i}=\epsilon_{1}\left[A\left(\mu_{1}\right)+C\left(\mu_{1}\right)-F\left(\mu_{i}\right)\right]$.
3) For $r_{1} \leq 1$ an additional imariant exists

$$
\eta_{1}=\operatorname{sgn}\left[\mu_{1}\left(a+c+f-\mu_{1}\right)\right]
$$

It is proved in $/ 12 /$ that the values of these invariants completely determine the class, to which a considered conic belongs. Here we shall briefly recapitulate this classification and establish a correspondence with the classification of operators (15)-(23).

Let us first enumerate all "canonical" forms to which the operator (24) can be reduced by hyperbolic motions in dependence on the invariants of the matrix X .

1. All characteristic numbers $\quad \mu_{1}, \mu_{2}, \mu_{3} \quad \therefore$ are different and real:

$$
\begin{equation*}
L_{1}=\mu_{1} K_{1}^{2}+\mu_{2} K_{2}^{2}-\mu_{8} K_{8}^{2} \tag{29}
\end{equation*}
$$

2. One real characteristic number $\mu_{1}$ and two complex ones $\mu_{2,3}=a \pm i \beta$ :

$$
\begin{gather*}
L_{2}=\mu_{1}\left(K_{1}^{2}+K_{2}^{2}\right)-\left(\mu_{1}+2 q\right) M_{8}^{2}+p\left(K_{1} M_{8}+M_{3} K_{1}\right)  \tag{30}\\
\alpha=\mu_{1}+q \quad \beta^{2}{ }^{2} p^{2}-q^{2}>0
\end{gather*}
$$

3. A simple characteristic number $\mu_{1}$ and a double one $\mu_{2}=\mu_{8}$ :
a) $c_{1}=-1, \quad c_{2}=1$

$$
\begin{equation*}
L_{1}-\mu_{1}\left(K_{1}^{2}+K_{2}^{2}\right)-\left(2 \mu_{2}-\mu_{1}\right) M_{8}^{2}+\left(\mu_{2}-\mu_{1}\right)\left(K_{1} M_{3}+M_{3} K_{1}\right), \tag{31}
\end{equation*}
$$

b) $c 1=-1, c a=-1$
$L_{4}-\left(2 \mu_{2}-\mu_{1}\right) K_{1}^{2}+\mu_{1}\left(K_{2}^{2}-M_{8}^{2}\right)+\left(\mu_{1}-\mu_{2}\right)\left(K_{1} M_{3}+M_{3} K_{1}\right)$,
c) $\epsilon_{1}=1, c_{2}=0$

$$
\begin{equation*}
L_{5}=\mu_{2}\left(K_{1}^{2}+K_{2}^{2}\right)-\mu_{1} M_{3}^{2}, \tag{33}
\end{equation*}
$$

d) $\epsilon_{1}-1, c_{2}=0$

$$
\begin{equation*}
L_{B}=\mu_{2}\left(K_{1}^{2}-M_{3}^{2}\right)+\mu_{1} K_{2}^{2} \tag{34}
\end{equation*}
$$

4. A triple characteristic number $\mu_{10} \mu_{2} \mu_{1} \mu_{8}$
a) $\quad c_{1}=-1$,

$$
\begin{equation*}
L_{T}=K_{1}^{2}+K_{2}^{2}-M_{3}^{2}+K_{1} K_{2}+K_{2} K_{1}-K_{2} M_{s}-M_{3} K_{2} \tag{35}
\end{equation*}
$$

b) $\quad c_{1}=0, \quad \eta_{1}=1$

$$
\begin{equation*}
L_{3}=2 K_{1}^{2}+K_{2}^{2}-K_{1} M_{3}-M_{3} K_{1}, \tag{36}
\end{equation*}
$$

c) $\epsilon_{1}=0, \quad \eta_{1}=-1$

$$
\begin{equation*}
L_{9}=K_{2}^{2}-2 M_{3}^{2}+K_{1} M_{3}+M_{3} K_{1}, \tag{37}
\end{equation*}
$$

d) $f_{1}=\eta_{1}=0$

$$
\begin{equation*}
L_{10}=K_{1}^{2}+K_{2}^{2}-M_{3}^{2} \tag{38}
\end{equation*}
$$

Let us discuss the forms (29)-(38).
Geometrically five different types of curves correspond to (29), depending on the possible values of $\mu_{1}$. However, with the help of the transformation

$$
\begin{equation*}
\overline{\mathrm{L}}=\alpha \mathrm{L}+\beta \Delta \tag{39}
\end{equation*}
$$

and an interchange of the coordinate axes they can be reduced to the cases

1. $\quad \mu_{1}>\mu_{2}>\mu_{3} \quad L_{1}$ reduces to $L_{E}$
2. $\quad \therefore \mu_{1}>\mu_{3}>\mu_{2} \quad L_{1}$ reduces to $L_{H}$
3. A single curve corresponds to (30). It is easy to see that all the in-
 $L_{2}$ can be reduced to $L_{s H}$.
4. Two curves correspond to (31) (different for $\frac{\mu_{1}}{\mu_{2}}>1$ and for $\frac{\mu_{1}}{\mu_{2}}<1$ ), however, the corresponding operators are connected by (39). The operators $L_{s}$ and $L_{E P}=\frac{\mu_{1}-\mu_{2}}{a-b} L_{E P}+\mu_{3} \Delta$ have the same imariants, so that $L_{s}$ reduces to $L_{E P}$ 5. Three curves correspond to (32) but in all cases $L_{4}$ can be reduced to
$\mathbf{L}_{H P}$ in complete analogy with the previous case.
5. Formula (33) describes three types of curves, but can always be reduced to $L_{s}$.
6. Formula (34) describes two types of curves, but can always be reduced to $L_{E q}$.
7. Formula (35) describes a single curve. The operators $L_{V}$ and $\mathrm{L}_{\mathrm{CP}}=\mathrm{L}_{\mathrm{CP}}+\Delta$ have the same invariants, so that $L_{T}$ can be reduced to $L_{C P}$.
8. Formulae (36) and (37) describe different curves, but $L_{8}$ and $L_{8}$ can both be transformed using (39) into $\left(K_{1}-M_{8}\right)^{2}$ connected with $L_{0}$ by a tryperbollc motion, so that they both reduce to $L_{0}$.
9. Formula (38) describes a single curve and, $L_{10}$ is equivalent to the zero operator.

Thus although there are 20 types of nondegenerate curves (det $x \neq 0$ )
in the Lobachevsky plane ( 12 real, 2 imaginary, 5 ideal ones and the absolute), we have shown that they correspond to our 9 classes of operators. The degenerate case $\operatorname{det} X=0$ need not be considered, since degenerate curves can always be transformed into nondegenerate ones by the transformation (39). This completes the proof.

## 5. Conclusions

We have investigated the group theoretical origin of coordinate systems, allowing the separation of variables in equation (1) for the spaces $R_{2}, E_{2}$ and $\mathrm{L}_{2}$

Let us recapitulate the main results.

1. The diagonallzation of a homogeneous quadratical polynomial in the generators of the group of motions corresponds to each separating coordinate system.
2. All different types of coordinate systems can be obtained by classifying all operators of the considered type into classes of equivalency, where two operators $L$ and $L$ belong to the same class, if they are related by transformation (2). A coordinate system of a certain type corresponds to each class of operators and systems of the same type, but shifted or rotated with respect to each other, correspond to different pperators in the same class,
3. Operators that are invariants of subgroups of the group of motions, correspond to coordinate systems with one centre (i.e. systems that are not elliptical or parabolical). Different types of coordinate systems correspond to different and unequivalent subgroups (spherical coordinates in $\mathrm{R}_{2}$, cartesian and polar coordinates in $E_{2}$, horocyclic, equidistant and spherical coordinates in $L_{2}$ ).

According to our opinion, the fact that operators, commuting with the Laplace operators of a given group can be distributed into a finite and quite definite number of equivalency classes, is of great physical interest. Mathematically it linits the choice of the basis for the group representations. In non-relativistic quantum mechanics this fact severely limits the number of potentials allowing a dynamical invariance group (this has been consldered for the two-dimensional case in $/ 8,9 /$ and an investigation of three-dimensional potentials is now under way).

From the point of vew of relativistic spin theory our results mean that it is only possible to measure quite definite quantities, characterizing the "projection" of the spin, if the total spin is fixed. For physical particles with non-zero
mass it is elther possible to measure the usual component (e.g. $L_{s}$ ) or the eigervalue of the operator $L_{E}=L_{3}^{2}+\sin ^{2} f L^{2}$, vut no other independent quantity. The situation is similar, but more sophisticated in the case of zero mass particles, or unphysical particles with $\mathrm{m}^{2}<0$.

The group $O(2,1)$ 'describes particles with imaginary mass, but is physically important in the study of scattering problems, in which the symmetry properties of amplitudes, not of single-particle wave functions, are considered. The group $\mathrm{O}(2,1)$ is the symmetry group related to a reaction in a given channel, considered from the cross-channel (it describes the amplitude in the non-physical region).

It would be interesting to conslder possible experiments determining the quantity $L_{3}^{2}+\sin ^{2} \mathrm{fL}_{2}^{2}$ for physical particles with $\mathrm{m}^{2}>0$, to clarify the role of the focus distance $f$ in the corresponding experimental device and to consider the connection of the found operators for $m^{2}=0$ with the Stokes polarization parameters. A discussion of these problems, as well as an investigation of the properties of cross-channel amplitudes connected with the "components" enumerated in §:4 will be presented separately. Although the representation theory of the groups $O(3), O(2,1)$ and $E_{2}$ is well known, explicit representations have only been constructed in the polar coordinate, system, thus corresponding to the diagonalization of a space rotation generator: It would be of interest to construct representations, corresponding to the diagonalization of all other independent operators, found in this paper, i.e. to give expressions for the basls functions, matrix elements of infinitesimal and finite transformations etc, and to find the operators, connecting these quantities in various systems. Such an investigation, specifically concerning the relation between the results of this paper and the representation theory of the group $O(2,1)$, developped by Bargmann, will be published. elsewhere.

It should be noted that the results of this paper can be generalized to spaces with higher dimensions, however in such cases it is necessary to classi-' fy sets of commuting operators, quadratical and symmetrical in the generators of the corresponding group, instead of a single operator, as in the case of the groups of rank I, considered in the present paper.

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[^0]:    $x /$ The little group, corresponding to a fixed null-vector is simply the homogeneous Lorentz group. The separation of variables on the corresponding threedimensional hyperboloid $x_{0}^{2}-x_{1}^{2}-x_{2}^{2}-x_{8}^{2}=I$ has been considered in 10 and its connection with diagonal operators, quadratical in the generators of the Lorentz group in 6,7 ).

