QUANTUM THEORY OF A MASSLESS RELATIVISTIC SURFACE AND A TWO-DIMENSIONAL BOUND STATE PROBLEM

by

Jens Hoppe

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c Jens Hoppe 1982

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Submitted to the Department of Physics on January 20, 1982, in partial fulfillment of the requirements for the Degree of Doctor of Philosophy in Physics

ABSTRACT

PART ONE

A massless relativistic surface is defined in a Lorentz invariant way by letting its action be proportional to the volume swept out in Minkowski space. The system is described in light cone coordinates and by going to a Hamiltonian formalism one sees that the dynamics depend only on the transverse coordinates X and Y. The Hamiltonian H is invariant under the group of area preserving reparametrizations whose Lie algebra can be shown to correspond in some sense to the Large N-limit of $\underline{SU(N)}$. Using this one arrives at a $\underline{SU(N)}$ invariant, large N-two-matrix model with a quartic interaction $[X,Y]^2$.

PART TWO

The problem of N particles with nearest neighbors &-function interactions is defined by regularizing the 2 body problem and deriving an eigenvalue integral equation that is equivalent to the Schrödinger equation (for bound states). The 3 body problem is discussed extensively and it is argued to be free of irregularities, in contrast with the known results in 3 dimensions. The crucial role of the dimension is displayed in looking at the limit of a short-range potential.

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PART ONE

QUANTUM THEORY OF A MASSLESS RELATIVISTIC SURFACE

INTRODUCTION

As a natural generalization of the massless string theory,* but also of interest in its own right, as an example in which geometry, classical relativity and quantum mechanics are deeply connected, one can define the dynamics of a massless closed M dimensional surface in a Lorentz- and coordinate invariant way by letting its action be proportional to the M+1 dimensional volume swept out in the D dimensional (generalized) Minkowski space \mathcal{M} . A particular observer with coordinate system $\mathbf{x}^{\mu}=(\mathbf{t},\mathbf{x}^{\ell})$ would describe the shape he sees by $\mathbf{x}^{\ell}(\mathbf{t},\lambda^1\ldots\lambda^M)$, where $\vec{\lambda}$ is a parametrization of the surface and the time like parameter λ^0 of the M+1--dimensional manifold was chosen to be t. Related to the arbitrariness of the choice of parametrization, not all of the \mathbf{x}^{μ} and their conjugate momenta \mathbf{p}^{μ} are independent.

It turns out to be extremely convenient to describe the system in terms of light cone coordinates $\tau(\equiv \frac{1}{2}(t+x^{D-1}))$ $\int (\equiv t-x^{D-1})$ and $\vec{x}(\equiv (x^*, \dots x^{D-2}))$, because the Hamiltonian turns out to be independent of \int and ** one can take \vec{x} and the conjugate momentum \vec{p} as the independent dynamical variables. In the classical theory \int is determined via constraint equations, which are consistent provided $\{\vec{x},\vec{p}\}_{\tau,\vec{k}} = \frac{\partial \vec{x}}{\partial x^*} = \frac{\partial (\vec{p}/\omega(x))}{\partial x^*} = 0$ where $\omega(\lambda)$ is a chosen density. These constraints fortunately do not cause a problem as their poisson bracket (commutator in the quantum theory) with the Hamiltonian is 0. (In the quantum theory they are interpreted as constraints acting on the wave functions ψ .)

^{*}Goddard, Goldstone, Rebbi, Thorn, NP B56 (1973) "Quantum dynamics of a massless relativistic string".

^{**}by picking a particular gauge, called orthonormal gauge.

 $\{\vec{x},\vec{p}\}_{rs}$ are the generators of volume preserving (time independent) $\vec{\lambda}$ -reparametrizations, which form a symmetry group that remains in orthonormal gauge.

After the general theory is described, everything else will be for the case M=2, D=4, with the parameter space $(\lambda^{\perp}, \lambda^{2})$ taken to have the topology of a 2-sphere. (Two examples of solutions to the equations of motion are given to become a little bit more familiar with the geometry of the problem and the parametrizations). The Hamiltonian which becomes

$$H = \int \min d\theta d\phi \left\{ p_{x}^{2} + p_{y}^{2} + \frac{1}{\sin^{2}\theta} \left(\frac{\partial x}{\partial \theta} \frac{\partial y}{\partial \phi} - \frac{\partial y}{\partial \theta} \frac{\partial x}{\partial \phi} \right)^{2} \right\}$$

is invariant under the group G of areapreserving reparametrizations of S^2 (and $x + y + C^{\alpha}(x + y)$). The Lie algebra G consists of all smooth functions* of θ and G, a basis of which one can take to be the usual spherical harmonics (leaving out Y_{00}).

In Part B it will be proved that the structure constants of G in the Y basis are in fact equal to the N+∞ limit of the structure constants of SU(N), in a particular, properly chosen basis. This proof, which from a mathematical point of view turns out to be much more natural than the construction first seems to be, makes use of the fact that the Y are the harmonic polynomials (restricted to the unit sphere S^2) which one writes as S^2 S^2 where S_1 is a N-dimensional representation of SU(N) can then be defined as S^2 S^2 S^2 S^2 S^2 where S_1 is a N-dimensional representation of SU(N) in this basis and *identifying any two differing just by a constant

others differing from $\mathring{T}_{\ell m}$ by N and ℓ dependent normalization factors so to make the structure constants have a finite non-zero totally antisymmetric N+ ∞ limit, can be derived. The SU(N)-invariant Hamiltonian H_N one gets by replacing $\mathbf{x}(\theta,\, \phi)$ by a hermitian N×N matrix $\mathbf{x},\, \{\ ,\ \}$ by $\frac{1}{i}[,], \star \int d\Omega$ by Tr, is a good approximation to H for large N-in the sense that the degrees of freedom corresponding to $\mathbf{Y}_{\ell m}$ with $\ell \leqslant N-1$ are represented correctly up to $O(\frac{1}{N})$, while the higher "frequencies" $(\ell \geqslant N)$ have been cut off.

Note that both H and $H_{\hat{N}}$ are hamiltonians for a gauge theory in 2+1 dimensions with spatial derivatives = 0:

$$H_{(N)} = \sum_{\alpha} \left(\left(p_{\alpha}^{\times} \right)^{2} + \left(p_{\alpha}^{y} \right)^{2} + \left(\sum_{k,c} f_{abc}^{(N)} \times_{k} Y_{c} \right)^{2} \right)$$

$$= \operatorname{Tr} \left(E_{x}^{2} + E_{y}^{2} + B^{2} \right)$$

where $x_b \leftrightarrow A_b^X$, and $B = [A^X, A^Y]$. The conditions fabor $\vec{x}_b \cdot \vec{p}_c = 0$ which are needed as a consistency condition for $H_{(w)}$ to be well defined translates into $[\vec{A}, \vec{E}] = 0$ which is exactly Gauss's law (when the spatial derivatives are 0). Bjorken** has looked at the analogue of this for SU(N=3) in 3 dimensions $(H=Tr(\vec{E}^2+\vec{B}^2))$, with the vectors now having 3 components) and seems to have shown that the lowest

^{*}Please note the misleading notation: this transition has nothing to do with the transition from a classical theory with poisson bracket $\{\}_p$ to a quantum theory with $[x,p]=i\acute{h}$.

^{**&}quot;Elements of quantum chromodynamics", SLAC PUB 2372, Dec. '79.

lying set of energy levels is a rotational band corresponding to 3-dimensional rotations. We have so far been unable to confirm this result. The last chapter contains some work on or related to ${\rm H}_{\rm N}$.

One would hope to be able to find out much about the spectrum of H_{N} by using (or finding new) techniques for large N-matrix models.* The work on this during the past months, however, has provided puzzles rather than insight.

Though the original classical action is manifestly Lorentz invariant, we are quantizing in a particular Lorentz frame and will have to demonstrate the Lorentz-invariance of our theory. A satisfactory method would be to construct the generators of Lorentz transformations, but we have been unable to do this. A weaker method, which would give only a necessary condition, is to show that the spectrum is consistent with Lorentz invariance, i.e., that the states fall into multiplets characterized by mass and spin. We have not carried our study of the dynamics far enough to see if this is true, although there is some indication that $H_N(N\to\infty)$ will have a high degeneracy of its energy levels.

^{*}See e.g., "Planar Diagrams" CMP 59 p.35-51 (1978), by Brezin et al., and the review article about the 1/N expansion by Sidney Coleman: SLAC PUB 2484, 198.

A. THE ACTION AND THE HAMILTONIAN FORMALISM

I. The action S and an example

A massless M-dimensional closed surface moving in D-dimensional Minkowski space can be defined by letting its action be proportional to the M+l dimensional volume swept out in Minkowski space (which is invariant under both Lorentz transformations and general reparametrizations $(\lambda^{\alpha} + \lambda^{\alpha})$ of the surface:

$$S = -T_0 \int_{\lambda' \text{ minimal}}^{\lambda' \text{ final}} \sqrt{G}$$
(A1)

where G is (-) M. the determinant of the metric $G = \frac{\partial x}{\partial \lambda} \frac{\partial x}{\partial \lambda}$ induced on the M+1 dimensional manifold M by Minkowski space; $X'=X'(\lambda^{*'})$ are the space time coordinates of $M: \mu=0, 1...D-1; \alpha=0...M$, $A'b_{\mu}=A^{b_{\mu}}-\sum_{k=1}^{D-1}\alpha^{k}b^{k}$ for two D-vectors; and T_{0} is the surface energy density (tension) of dim $\frac{\text{Energy}}{(\text{length})^{M}}$ which will from now on be put = 1 (one can always put it in on dimensional grounds). Using $A'b = \frac{1}{2}\sqrt{G}C^{*'} A'b = \frac{1}{2}\sqrt{G}C^{*$

$$SS = -\frac{1}{2} \int d^{n+1} \lambda \, \mathcal{C} \, \mathcal{C}^{\mathcal{S}} \, \mathcal{S} \left(\partial_{\alpha} \times r \partial_{\beta} \times_{\mu} \right)$$

$$= \int d^{n+1} \lambda \, \mathcal{C} \, \mathcal{S} \, \mathcal{L}_{\mu} \, \mathcal{L}_$$

Choosing the timelike parameter $\lambda^{\,0}$ of the manifold to be t, one has

$$G_{x} = \begin{pmatrix} 1 - \vec{x}^{2} & -\vec{x} \partial_{x} \vec{x} \\ -\vec{x} \partial_{x} \vec{x} & -g_{rs} \end{pmatrix} \quad \text{where } \vec{x} \equiv \frac{\partial \vec{x}}{\partial t}, \quad \partial_{x} \vec{x} \equiv \frac{\partial \vec{x}}{\partial x^{2}}$$

$$\vec{x} = (x^{1}, \dots, x^{2n-1}) \quad \vec{x} = \sum_{\ell=1}^{2n-1} x^{\ell} x^{\ell}$$
(A3)

$$\Im_{rs} = -\partial_{r} \times r \partial_{s} \times_{r} = + \partial_{r} \overrightarrow{x} \partial_{s} \overrightarrow{x} \left(a_{s} \partial_{r} t = 0 \text{ for } \lambda^{c} t \right)$$

It is convenient to partially fix the parametrization by requiring

(i)
$$G_{00} = 1 - \vec{x}^2 = +g$$
 ($g = det g_{rs} = |g_{rs}|$)
(ii) $G_{0r} = G_{ro} = \vec{x} \partial_r \vec{x} = 0$

This choice is possible provided x^{μ} satisfies the equations of motion: (ii) says that given the parametrization λ^{r} of the surface at time $t=t_{\bullet}$, one chooses the parametrization at a slightly later time t +dt to be such that the intersections of any normal with the two surfaces are at equal λ^{r} . Further one certainly can choose the parametrization such that $0 = \frac{2\lambda^{r}}{2\lambda^{r}} \frac{2\lambda^{r}}{2\lambda^{r}} \qquad \text{is } 1^{-\frac{1}{x^{2}}} \text{ at a given time. But given (ii) (for all times) the } \lambda = 0 \text{ part of Eq. (A2) says that } \lambda = 0 \text{ for all } t$. Note that (A4) is still invariant under volume preserving time independent reparametrizations of the surface (as those are exactly the ones that leave g invariant).

It is not difficult to find a solution of the classical equations of motion for M=2, D=4 (the physical case). The Ansatz

$$X^{n} = \begin{pmatrix} t \\ S(t)_{\vec{m}} \end{pmatrix}, \vec{n} = \begin{pmatrix} \sin\theta \cos\theta, \sin\theta \sin\theta, \cos\theta \end{pmatrix}$$
(A5)

with Θ and ψ being the usual angles of spherical coordinates, and defining $\lambda^1 \Xi - \cos \Theta \Xi \mu$, $\lambda^2 = \psi$ gives

$$G_{x,s} = \begin{pmatrix} 1 - \dot{S}^{2} & 0 & 0 \\ 0 & -S_{sn^{2}0}^{2} & 0 \\ 0 & 0 & -sn^{2}0S^{2} \end{pmatrix} / G = S^{2}\sqrt{1 - \dot{S}^{2}}$$

as $\partial_{\mu}\vec{m} \cdot \partial_{\mu}\vec{m} = 0$ $(\partial_{\mu}\vec{m})^{2} = \frac{1}{5^{-2}\theta}$ $(\partial_{\mu}\vec{m})^{2} = \sin^{2}\theta$ The $\mu=0$ part of (A2) $(1-\dot{S}^{2})^{-1/2}S^{-2}\partial_{\mu}S^{2}(1-\dot{S}^{2})^{-1/2}\partial_{\mu}t$ leads to

$$S^{4} = (const)(1-\dot{S}^{2}) \tag{A6}$$

while the spatial part, which, using (A6) becomes

is trivially satisfied by definition of \vec{n} . The solution of Eq. (A6), which is equivalent to $t/S_0 = \int_{S/S_0}^{\infty} \sqrt{1-\chi^4}$

(S = maximal radius), is a periodic

elliptic function which can easily be expressed in terms of the standard Weierstrass-P-function.

II. General formalism in light cone coordinates We define light cone coordinates by:

$$\begin{cases} \uparrow = \frac{1}{2}(t+2) \end{cases} \iff \begin{cases} t = \gamma + \frac{3}{2} \end{cases}$$

$$\begin{cases} \downarrow = t - \frac{3}{2} \end{cases}$$
(A7)

From now on \vec{x} will always stand for $(x', ..., x^{D-2})$ and no distinction will be made between x_i and x^i (i=1,...N=D-2) $x^{\mu}=(t, \vec{x}, z)$. Choosing $\lambda^{\mu}=1$:

$$G_{ro} = \begin{pmatrix} G_{oo} & G_{or} \\ G_{ro} & -g_{rs} \end{pmatrix} = \begin{pmatrix} 2 \dot{y} - \dot{\vec{x}}^2 & (\partial_r \dot{y} - \dot{\vec{x}} \partial_r \dot{\vec{x}}) \\ \partial_r \dot{y} - \dot{\vec{x}} \partial_r \dot{\vec{x}} & -\frac{\partial \dot{\vec{x}}}{\partial \lambda^r} \frac{\partial \dot{\vec{x}}}{\partial \lambda^s} \end{pmatrix}$$
(Note that $\dot{\vec{x}}^2$ and $\dot{\vec{x}}$

(Note that $\frac{1}{x}^2$ and $\frac{1}{x}$ are differently defined from the $\frac{1}{x}^2$ and $\frac{1}{x}^2$ appearing on p. 0 .) Now x is a D-2-vector and indicates differentiation with respect to op

$$G = (-)^{M} \text{ dit } G_{dp} = \text{ dit } (G_{00} - G_{0r}) = G_{00} g + G_{0r} G_{0s} g^{rs}$$

$$= g (G_{00} + G_{0r} G_{0s} g^{rs}) - (g^{rs} g_{st} = g^{rs})$$
having used the fact that for a completely general square matrix

$$A = \begin{pmatrix} a_0 & \vec{a}^{t} \\ \vec{c} & \vec{B} \end{pmatrix} \text{ with invertible B one has } |A| = |B| \begin{cases} a - \vec{a}^{t} \vec{B}^{-1} \vec{c} \\ \vec{c} & \vec{c} & \vec{c} \\ \vec{c} & \vec{$$

If we define canonical

$$\overrightarrow{P} = \frac{\partial \mathcal{L}}{\partial \dot{x}} - \sqrt{\frac{9}{\Gamma}} \left(\dot{\overrightarrow{x}} - u^{\tau} \partial_{\tau} \overrightarrow{x} \right)$$

$$T = \frac{\partial \mathcal{L}}{\partial \dot{y}} = -\sqrt{\frac{9}{\Gamma}}$$
we find that
$$(A9)$$

$$\vec{p} \cdot \vec{Q}_1 \vec{X} + \vec{T} \vec{Q}_1 \vec{X} = 0$$
(A10)

This constraint is a direct consequence of the invariance of ${\sf S}$ under t-dependent reparametrization,

$$\vec{x} = \vec{x}(\lambda, r) \vec{a}_r \vec{x} \qquad \text{if } = \vec{x}(\lambda, r) \vec{a}_r \vec{x}$$

To go to a Hamiltonian formalism, we express $\overrightarrow{p} \times + \overrightarrow{T} - \mathscr{L} = \mathscr{L}$ as a function of $\overrightarrow{p} \times \overrightarrow{X} = \mathscr{L}$ (This expression is, of course, not unique because of the relation (AlO).)

$$\mathcal{X} = \overset{\rightarrow}{\times} \overset{\rightarrow}{\rho} + \overset{\rightarrow}{5} T + \overset{\rightarrow}{9} \overset{\rightarrow}{2} \overset{\rightarrow}{5} \overset{\rightarrow}{-} \overset{\rightarrow}{2} \overset{\rightarrow}{+} u_{\tau} u^{\tau}$$

$$= \overset{\checkmark}{\sqrt{2}} \overset{\rightarrow}{5} \overset{\rightarrow}{-} \overset{\rightarrow}{\times} u_{\tau} u^{\tau}} \left\{ \overset{\rightarrow}{\times} \overset{\rightarrow}{-} V_{\tau} u^{\tau} - \overset{\rightarrow}{5} + \left(2 \overset{\rightarrow}{5} - \overset{\rightarrow}{\times} \overset{\rightarrow}{+} u_{\tau} u^{\tau} \right) \right\}$$

$$= \overset{\checkmark}{\sqrt{2}} \overset{\rightarrow}{5} \overset{\rightarrow}{-} \overset{\rightarrow}{\times} u_{\tau} u^{\tau}} \left\{ \overset{\rightarrow}{5} - \overset{\rightarrow}{O_{\tau}} \overset{\rightarrow}{5} u^{\tau} \right\} \left(V_{\tau} = \overset{\rightarrow}{\times} \overset{\rightarrow}{O_{\tau}} \overset{\rightarrow}{\times} \right)$$

$$\overset{\rightarrow}{V_{\tau}} \overset{\rightarrow}{=} \overset{\rightarrow}{-} \overset{\rightarrow}{\times} \overset{\rightarrow}{+} u_{\tau} u^{\tau}} \left\{ \overset{\rightarrow}{5} - \overset{\rightarrow}{V_{\tau}} u^{\tau} + \left(\overset{\rightarrow}{O_{\tau}} \overset{\rightarrow}{\times} u^{\tau} \right)^{2} + \left(2 \overset{\rightarrow}{5} - \overset{\rightarrow}{\times} \overset{\rightarrow}{+} u_{\tau} u^{\tau} \right) \right\}$$

$$\overset{\rightarrow}{V_{\tau}} \overset{\rightarrow}{=} \overset{\rightarrow}{-} \overset{\rightarrow}{\times} \overset{\rightarrow}{+} u_{\tau} u^{\tau}} \left\{ \overset{\rightarrow}{5} - \overset{\rightarrow}{V_{\tau}} u^{\tau} + u_{\tau} u^{\tau} \right\} = \overset{\rightarrow}{V_{\tau}} \left(\overset{\rightarrow}{5} \overset{\rightarrow}{\times} \overset{\rightarrow}{+} u_{\tau} u^{\tau} \right)$$

$$\overset{\rightarrow}{V_{\tau}} \overset{\rightarrow}{=} \overset{\rightarrow}{-} \overset{\rightarrow}{\times} \overset{\rightarrow}{+} u_{\tau} u^{\tau}} \left\{ \overset{\rightarrow}{5} - \overset{\rightarrow}{V_{\tau}} u^{\tau} + u_{\tau} u^{\tau} \right\} = \overset{\rightarrow}{V_{\tau}} \left(\overset{\rightarrow}{5} \overset{\rightarrow}{-} \overset{\rightarrow}{\times} \overset{\rightarrow}{+} u_{\tau} u^{\tau} \right)$$

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$$\overset{\rightarrow}{V_{\tau}} \overset{\rightarrow}{=} \overset{\rightarrow}{-} \overset{\rightarrow}{\times} \overset{\rightarrow}{+} u_{\tau} u^{\tau} \right\} = \overset{\rightarrow}{V_{\tau}} \left(\overset{\rightarrow}{5} \overset{\rightarrow}{-} \overset{\rightarrow}{\times} \overset{\rightarrow}{+} u_{\tau} u^{\tau} \right)$$

$$\overset{\rightarrow}{V_{\tau}} \overset{\rightarrow}{=} \overset{\rightarrow}{-} \overset{\rightarrow}{\times} \overset{\rightarrow}{+} u_{\tau} u^{\tau} \right\} = \overset{\rightarrow}{V_{\tau}} \left(\overset{\rightarrow}{\sim} \overset{\rightarrow}{\to} \overset{$$

We can then obtain the equations of motion from the Hamiltonian $X = X + U(\vec{p} \cdot 2\vec{x} + \vec{p} \cdot 2\vec{x})$ treating $\vec{x} \cdot \vec{j} \cdot \vec{p} \cdot \vec{r}$ and U as independent variables:

^{*}For a general discussion of "constrained Hamiltonian systems" one could refer to the long article (with same title) of Hanson, Regge and Teitelboim. Academia Nazionale Dei Lincei, 1976.(Contributi del Centro Linceo Interdisciplinare Di Scienze Matematiche e

$$\frac{\int H}{\int u^{\tau}} = \vec{p} \cdot \vec{\partial}_{x} \vec{x} + \vec{T} \vec{\partial}_{x} \vec{y} = 0 \qquad \left(H' = \int d^{\eta}_{\lambda} \mathcal{X} \right)$$

$$\dot{\vec{S}} = \frac{\int H}{\int T} = \int \frac{2+g}{2\pi^{2}} + u^{\tau} \partial_{\tau} \vec{y}$$

$$\dot{\vec{X}} = -\vec{p}_{\pi} + u^{\tau} \partial_{\tau} \vec{x} \qquad \dot{\vec{T}} = \partial_{\tau} (\vec{T} u^{\tau}) \qquad \text{(All')}$$

$$\dot{\vec{P}} = -\partial_{\tau} + g g^{\tau s} \partial_{s} \vec{x} + \partial_{\tau} (u^{\tau} \vec{p})$$

$$(a \Rightarrow + \int \frac{\partial g}{2\pi} = + \int \frac{\partial g^{\tau s}}{2\pi} \delta g^{\tau s} = -\int \partial_{\tau} (\frac{1}{\pi} g g^{\tau s} \partial_{s} \vec{x}) \delta \vec{x} \qquad \text{and } H = \int \mathcal{X} d^{\eta}_{\lambda}$$
Note that $H' = \int \mathcal{X} d^{\eta}_{\lambda} \vec{y} \sin n \vec{y} \sin n \vec{y} \sin n \vec{y}$

under reparametrization provided that p and || transform as densities. Also as a consequence of Hamilton's equations, \mathcal{U} is equal to \mathcal{U} as defined in (AB) (Just Calculate $\hat{\mathbf{x}} = \hat{\mathbf{x}} = \hat{\mathbf{x}}$) from (AH')

To discuss classical solutions, we can always choose the time variation of the parametrization so that $\mathcal{U}=0$. Since \mathcal{U} is independent of \mathcal{J} , in this gauge $\mathbb{T}=0$. We are still free to make a time-independent reparametrization. Since \mathbb{T} transforms as a density we can make it equal to a constant times a specified λ -dependence, $\pi=-\eta w(\lambda)$. We are then left with the Hamiltonian

$$H = \frac{1}{2\eta} \int \frac{d^n \lambda}{\omega(\lambda)} \left(\vec{p}^2 + g \right)$$
(A12)

to determine the motion of \vec{x} . We call this gauge orthonormal gauge (ONG). The constraint (AlO) becomes $\vec{p} \cdot \partial_1 \vec{x} = \gamma \omega(\lambda) \partial_1 \vec{y}$ which we can solve for \vec{x} provided

$$\partial_{x} \frac{\omega(x)}{1} \overrightarrow{p} \cdot \partial_{x} \overrightarrow{x} - \partial_{x} \frac{\omega(x)}{1} \partial_{x} \overrightarrow{x} = 0$$
Al3)

i.e.
$$ST = \partial_{\tau}(f^{\tau}T)$$
, $S\vec{p} = \partial_{\tau}(f^{\tau}\vec{p})$. (while u^{τ} transforms like a contravariant vector: $Su^{\tau} = (\partial_{S}u^{\tau})f^{S} - (\partial_{S}f^{\tau})u^{S} + f^{\tau}$)

These constraints are consistent with the equations of motion derived from (Al2) because H is still invariant under reparametrizations which leave the measure $w(\lambda)d^M\lambda$ invariant.

The constants of the motion p^{μ} may be obtained by comparing $P_{x} x^{\mu} = P^{c} x^{c} - P^{2} 2 - \overrightarrow{P} \cdot \overrightarrow{x}$ $= P^{+} + P^{-} \gamma - \overrightarrow{P} \cdot \overrightarrow{x}$

We see that since \overrightarrow{P} generates transverse translations, $-\overrightarrow{P}^{\dagger}$ must generate translations in \overrightarrow{J} and \overrightarrow{P} must be our H which generates the motion in τ . Thus $\overrightarrow{\overrightarrow{P}} = \left(\overrightarrow{P} \stackrel{\mathcal{M}}{\nearrow} \lambda\right)^{1/2}$

$$P^{+} = -\int \Pi d^{n}\lambda = \eta \int \omega(\lambda) d^{n}\lambda$$
and
$$P^{-} = \frac{1}{2\eta} \int (\rho^{2} + g) \frac{d^{n}\lambda}{\omega(\lambda)}$$
(A14)

If for a given choice $w(\lambda)$ (with $\int_{W(\lambda)} d^n \lambda = W$) we choose a complete orthonormal set of functions $\phi_n(\lambda)$, $\int_{M} \int_{M} w(\lambda) d^n \lambda = \int_{Mm} w(\lambda) d^$

 x_m and p_m will be canonically conjugate variables. If we take $\phi_0 = \frac{1}{\sqrt{N}}$, g which depends only on $0, \overline{X}$ will be independent of \overline{x}_0 and we find

$$\overrightarrow{P} = \overrightarrow{\rho}_{0} \, \sqrt{w} \, , \quad \overrightarrow{P}^{+} = \gamma \, w$$

$$\overrightarrow{P}^{-} = \frac{1}{2\gamma} \left(\overrightarrow{\rho}_{0}^{2} + \sum_{m > 0} \overrightarrow{\rho}_{m}^{2} + \int g \, \frac{d^{m} \lambda}{\omega(\lambda)} \right)$$

$$= \frac{1}{2\gamma} \left(\overrightarrow{P}^{2} + w \left\{ \sum_{m > 0} \overrightarrow{\rho}_{m}^{2} + \int g \, \frac{d^{m} \lambda}{\omega(\lambda)} \right\} \right)$$

This relation is of the correct relativistic form,

$$2 P^{+}P^{-} - \vec{P}^{2} \left(= P^{\circ 2} - P^{2} - \vec{P}^{2}\right) = m^{2} \text{ with}$$

$$m^{2} = V \left\{\sum_{n \geq 0} \vec{p}_{n}^{2} + \int g \frac{d^{n}\lambda}{\omega(\lambda)}\right\} \equiv + \text{internal}$$
depending only on the degrees of freedom \vec{x}_{n} , \vec{p}_{n} , $n > 0$.

In order to quantize this theory, we use the Hamiltonian

$$H = - \int \vec{p}^{2} d^{n} \lambda$$

with $\frac{1}{x}(\lambda)$, $\frac{1}{p}(\lambda)$, $\frac{1}{y}(\lambda)$, $\frac{1}{y}(\lambda)$ as canonical variables, obeying e.g.

 $[\times;(\lambda), f;(\lambda')] = i + \int_{0}^{\pi} \int_$

$$\left(\partial_{r}\vec{x}\cdot\vec{p}+\partial_{r}\vec{y}+\partial_{r}\vec{y}+\partial_{r}\vec{y}\right)=0$$
(A15)

These constraints are consistent with each other and with $\psi = E \psi$ since they are the generators of the group of reparametrizations. Since H is independent of ψ we can find eigenstates which are also eigenstates of ψ (λ),

$$T(\lambda)|\psi\rangle = - \eta \omega(\lambda)|\psi\rangle \qquad (A16)$$

(*)
$$\Upsilon \rightarrow \Upsilon' = \sqrt{\frac{1+V}{1-V}} \Upsilon = e^{u} \Upsilon$$
, $3 \rightarrow 3' = \sqrt{\frac{1-V}{1+V}} 3 = e^{-u} 3$

These will not satisfy (Al5). However, (Al5) is equivalent to the condition that the wavefunctions $\psi [\vec{x}(\lambda), T(\lambda)]$ invariant when \overrightarrow{x} , \overrightarrow{x} are transformed by reparametrization. transforms as a density.) We can always construct such a wavefunction from a ψ satisfying (Al6) and invariant under those reparametrizations which leave $w(\lambda)$ invariant. Furthermore we need only consider a single specified form of $w(\lambda)$ since all others may be reached by reparametrization and rescaling of η . This invariance condition is exactly (Al3) interpreted as a The classical discussion is now exactly paralleled by the quantum theory. We must find the eigenstates of H_{int} subject to (Al3). These will also be eigenstates of J_2 . Clearly a necessary condition for Lorentz invariance is that for a given eigenvalue of H_{int} the states can be arranged into SO(3) multiplets (i.e., that the number of states increases as J, decreases). It is possible to see that in a certain sense this is also a sufficient condition, i.e., if it is unitary operators realizing Lorentz invariance can be constructed level by level of H_{int}. However, they would not necessarily be related in any simple way to the canonical variables.

The further discussion will be restricted to the case M=2, D=4, $w(\lambda)d^{\Pi}\lambda$ =sinoded[]. It is convenient to define $\rho = \rho_{din0}$ so that (when $\mu = -(\pi \theta)$) $\left[\times (\theta, \varphi), \rho(\theta, \varphi) \right] = \frac{1}{6\pi \theta} \delta(\theta - \theta) \delta(\varphi - \varphi) = i \delta(\mu - \mu) \delta(\varphi - \varphi)$ Then $\frac{1}{6\pi \theta} = \frac{1}{2} \int \sin \theta \, d\theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \left\{ \rho^2 + \frac{1}{2} \int \sin \theta \, d\theta \, d\varphi \, \right\} \right\} \right\} \right\} \right]$

and
$$\frac{\partial}{\partial x} = \left(\frac{1}{\sin \theta} \left(\frac{\partial x}{\partial \theta} \frac{\partial y}{\partial \psi} - \frac{\partial y}{\partial \theta} \frac{\partial x}{\partial \psi}\right)\right)^2 = \left\{x_i y\right\}^2$$

where we define the Lie bracket of two functions A,B by

$$\left\{ A, B \right\} = \frac{1}{6m0} \frac{\partial \left(A, B \right)}{\partial \left(\theta, \varphi \right)} \left(= \frac{\partial \left(A, B \right)}{\partial \left(\mu, \varphi \right)} \right) - - \cdot \cdot \left(A | B \right)$$

Area preserving transformations are of the form

$$Sx = \frac{\partial x}{\partial \theta} f^{\theta} + \frac{\partial x}{\partial \phi} f^{\phi} = \frac{\partial x}{\partial r} f^{r} + \frac{\partial x}{\partial \phi} f^{\phi}$$
where $\int_{r} f^{r} + \partial_{\phi} f^{\phi} = 0$ so that
$$f^{r} = \partial_{\phi} f, \quad f^{\phi} = -\frac{1}{6} \partial_{\phi} f \text{ and } Sx = \left\{x, f^{\phi}\right\}$$

The constraints (Al3) take the form $\{x, \widetilde{\rho_x}\} + \{y, \widetilde{\rho_y}\}$ on the states. It is seen that the whole theory now depends on the single algebraic structure $\{A,B\}$. Part B will depend essentially on this fact.

III. Another example and a comparison

The Ansatz
$$\vec{X} = \begin{pmatrix} X \\ Y \end{pmatrix} = R(T, \mu) \begin{pmatrix} \cos \theta \\ \sin \theta \end{pmatrix} \equiv R \cdot \vec{m}$$
 (A19)

leads, in orthonormal gauge (*) to
$$\vec{P} = M \cdot \vec{R} \cdot \vec{m}$$
,
$$R = \frac{\partial R}{\partial T} \begin{vmatrix} \cos \theta \\ \cos \theta \end{vmatrix} = \begin{pmatrix} \cos \theta \\ \cos \theta \end{vmatrix}$$

$$R = \frac{\partial R}{\partial T} \begin{vmatrix} \cos \theta \\ \cos \theta \end{vmatrix}$$

and the equation of motion reads (*)

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^{*}See pg. 22.

$$\chi^2 \ddot{R} = R (RR')'$$
(A20)

The constraint $\left\{ \overrightarrow{X}, \overrightarrow{P} \right\} \equiv \left\{ \overrightarrow{\partial x}, \overrightarrow{\partial p} - \frac{\partial p}{\partial x}, \overrightarrow{\partial p} = 0 \right\}$ is satisfied, as $\overrightarrow{m} \partial_{\varphi} \overrightarrow{m} = 0$ (so both terms = 0). Equivalently one can see directly that the equations for \overrightarrow{A} are integrable, for \overrightarrow{X} of the form (Al9):

$$\dot{S} = \frac{\rho^2 + 9}{2\eta^2} = \frac{1}{2} \dot{R}^2 + \frac{1}{2\eta^2} R^2 R^{2}$$

$$\dot{S}' = \frac{1}{2} \dot{P} \cdot \overrightarrow{X}' = \dot{R} R'$$

$$\partial_{\varphi} J = \frac{1}{2} \dot{P} \cdot \partial_{\varphi} \overrightarrow{X} = 0$$
(A21)

The integrability conditions involving derivatives of J with respect to f are trivially satisfied (as J is independent of f), the one involving f gives exactly (A20).

One particular solution of (A20) with μ -cos θ , is $R(T,\mu) \ge R(T)$ sin θ , leading to*

$$\ddot{R} = -R^{3}/2^{2} \left(= R^{4} + 2\chi^{2} \dot{R}^{2} = D\chi^{2}, D = canot \right)$$
(A22)

and (A21) becomes

$$\dot{S} \stackrel{(a)}{=} \frac{1}{2} \dot{R}^{2} \sin^{2}\theta + \frac{1}{2\eta^{2}} R^{4} \cos^{2}\theta
\partial_{\theta} S \stackrel{(a)}{=} R \dot{R} \sin \theta \cos \theta = \frac{1}{2} R \dot{R} \sin^{2}(2\theta)$$
(A23)

This will now be integrated explicitly from the second equation

$$J = -\frac{R\dot{R}}{4}(\omega(20) + f(0) =) \dot{S} = \dot{f} - \frac{1}{4}\cos 20(\dot{R}^2 - R^4/2^2)$$

(using (A22) which has to equal (A23i)

^{*}Please note that 0 is not any geometrical angle, in particular not the angle of the spherical coordinates.

Therefore f has to equal $\frac{1}{4}R^2 + \frac{1}{4}R^2 + \frac{1}{4}R^4$, which—using again (A22)—is $\frac{1}{4}\frac{1}{4R}(R\dot{R})$, so that

$$\vec{S} = \frac{RR}{2} \sin^2 \theta + \mathcal{L}(\tau) \quad ; \quad \dot{\mathcal{L}} = \frac{1}{2\gamma^2} R^4(\tau) \quad (A19)$$

Both because (A22) is exactly the equation found earlier for S(t) (the radius of the breath ing solution in a regular Lorentz frame) and because both (A5) and $\vec{x}=R(\tau)\sin\theta\vec{m}$ are most simple and symmetric solutions, one would think that they are in fact the same solution, just looked at in different frames and with different variables. This appears to be wrong, i.e., the above solution $R=R(T)\cdot\sin\theta$ is not the $R(T,\mu)$ in $\vec{x}=R(T,\mu)e^{-\lambda t}$ that corresponds to the solution x $(t,\nu,\psi)=(t,S(t),\mu)$ nor a simple Lorentz transform of it.

One can, in fact, calculate the parameter μ as a function of t and the geometric angle with the z-axis ${\cal P}$

So far

$$3=3(T,\mu)$$
, $R=R(T,\mu)$; $dR=RdT+Rd\mu$, $dS=3dT+3d\mu$
so that $d\mu=d3-3/3$, dT on the other hand one could extract
 $\mu=\mu(T,3)$ from $3(T,\mu)$, $R(T,\mu)$ and think of R as $R(T,3)$,
so that $dR=RdT+Rd3$, $d\mu=\mu_TdT+\mu_3d3$ where

$$X_{T} = \frac{\partial X}{\partial T} \Big|_{\text{constant}3} \Big|_{X_{3}} = \frac{\partial X}{\partial S} \Big|_{\text{constant}T}$$

By comparing the two expressions for de one finds

$$\mathcal{M}_{\tau} = -\frac{3}{3}, \quad \mathcal{M}_{3} = \frac{1}{3}, \\
\dot{R} = \frac{2R}{2\tau} |_{x} = R_{\tau} + R_{3} \frac{2S}{2\tau}|_{x} = R_{\tau} + R_{3} \frac{S}{2\tau}|_{x} = R_{\tau} + R_{3} \frac{S}{2\tau}$$

$$\dot{\beta} = \frac{1}{2} \left(R_{1} + R_{3} \dot{\beta} \right)^{2} + \frac{1}{2\eta^{2}} R^{2} R_{3}^{2} 3^{2}$$

$$\dot{\beta} = \left(R_{1} + R_{3} \dot{\beta} \right) R_{3} \dot{\beta}^{2}$$

from which one deduces

$$\dot{\zeta} = \frac{1 - R_T R_3}{R_3^2}$$

$$\mathcal{J} = \frac{\gamma}{RR_3} \sqrt{2\vec{J} - (R_T + R_3 \vec{J})^2} = \gamma_{RR_3} \sqrt{1 - 2R_T R_3}$$
Therefore

 $\mu_{T} = \frac{R(R_{T}R_{3}-1)}{2\sqrt{1-2R_{T}R_{3}}}, \mu_{3} = \frac{RR_{3}^{2}}{\sqrt{1-2R_{T}R_{3}}}$ (A25)

This expression is true whenever $X = R(T,\mu)$ (and). Now one specifies: the solution (to A20) $R(T,\mu)$ that corresponds to the solution

(A5) $(x/^{2} = (t, S(t)))$ obeys $R^{2} + 2^{2} = S^{2}, R^{2} + (\tau - 3_{2})^{2} = S^{2}(\tau + 3_{1})$ From this it follows (eg. RVI-2R-R-=5) that

$$\mu_{3} = \frac{1}{4\pi S^{3}} (S\partial_{\xi}S + 2)^{2}, \mu_{7} = \frac{1}{27S^{3}} (S^{2}(S_{\xi}S)^{2} + 2^{2} - 2S^{2})$$
and therefore
$$(2 = 7 - 3L)$$

$$\mu_{z} = -\frac{S+zS}{27S^{2}}, \mu_{t} = \frac{1}{27S^{3}} \left\{ -S^{6} + z^{2} + zSS \right\}^{2}$$

from which one can determine μ as a function of t and z, or t and V; one finds $2\pi M = -\cos 2\theta + \frac{1-54}{2}\sin^2 2\theta + cont.$ (± for collapsing sphere).

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Summary of formulae in orthonormal gauge

For convenience, the important equations (in particular (All')) are written out explicitly for orthonormal gauge:

$$\begin{aligned}
\overline{\Pi} &= - \gamma \omega(\lambda) &, & u_r &= 0 \\
&\text{constraint} : \overrightarrow{P} \partial_t \overrightarrow{X} &= \gamma \omega(\lambda) \partial_t \overrightarrow{J}, & \{\overrightarrow{P}_{\omega(\lambda)}, \overrightarrow{X}\} &= 0 \\
& \overrightarrow{J} &= \frac{P^2 + 9}{2 \gamma^2 \omega(\lambda)} & \overrightarrow{P} &= \gamma \omega(\lambda) \overrightarrow{X} \\
& \overrightarrow{P} &= \frac{1}{2} \partial_r \frac{99^{rs}}{\omega(\lambda)} \partial_s \overrightarrow{X} \\
&=) \overrightarrow{X} &= - \overrightarrow{P} / \overline{\Pi} &= + \frac{1}{2} \frac{1}{2} \omega(\lambda) \partial_r \frac{99^{rs}}{\omega(\lambda)} \partial_s \overrightarrow{X}
\end{aligned}$$

For M=2 a convenient choice is (used in AIII):

$$\lambda' \equiv \mu \equiv -\cos\theta \in [-1,+1] \quad (\theta \in [0,T])$$

$$\lambda^2 = \theta \quad , \quad \omega(\lambda) = 1$$
(if $\mu = +1(\alpha - 1)$, all points in the objective $(q - values)$
have to be solutified $(q - values)$
if $\lambda' = \theta$, then choose $\omega(\lambda) = \sin\theta$

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- B. The surface problem as the limit of a large N matrix problem
- I. The group of area preserving reparametrizations of S² and the structure of its Lie algebra in connection with the surface Hamiltonian

The Hamiltonian found in Section A may be written as:

H is invariant under the group G of area preserving diffeomorphisms of S² (that are connected to the identity)—meaning that the functional dependence of H (on \dot{x} and \dot{p}) will not change under a smooth reparametrization of the parameter space (a 2-sphere): $(\mu, \psi) \rightarrow (\mu', \psi')$, with unit Jacobian. This can be seen by looking at infinitesimal transformations $\mu' = \mu + \delta \mu$, $\psi' = \psi + \delta \psi$, for which the condition

is satisfied (to first order) if $\delta \mu = +\partial_{\mu} f$, $\delta \psi = -\partial_{\mu} f$ with f being any smooth infinitesimal function (defined by these equations up to a constant); it follows that for any function $Z(\mu, \psi)$ one has

$$\delta z = 2(\mu', \psi') - 2(\mu, \psi) = 2\mu 2 + 2\mu 2 - 2\mu 2 + 0(\beta^2)$$

$$= \{2, \beta\} + 0(\beta^2)$$

and (to first order in f):

$$SH = \frac{1}{2} \int d\Omega \left(\left\{ p_x^2 + p_y^2, f \right\} + 2 \left\{ x_i y \right\} \delta \left\{ x_i y \right\} \right) = 0,$$

as $\int O(2 g) \int \int for any g(\mu, \mu)$ (integrate by parts!). Using the Jacobi identity for $\begin{cases} 1 & \text{one has} \end{cases}$

$$\begin{aligned} &\{\{x_iy\} = \{\delta x_iy\} + \{x_i\delta y\} = \{\{x_if\}_iy\} + \{x_i\{y_if\}\}\} \\ &= \{\{x_iy\}_if\} \end{aligned}$$

so that

The equations of motion derived from H are*

$$\ddot{x} = \left\{ \left\{ x_{1} y \right\}_{1} y \right\}$$

$$\ddot{y} = -\left\{ \left\{ x_{1} y \right\}_{1} x \right\}$$
(B2)*

The Lie algebra G of G is the space of all smooth functions $f(\theta, \psi)$ with f and g identified if they differ only by a constant. The Lie bracket on G is

^{*&#}x27;A whole class of solutions is: $x+iy=we^{i(w-y^t)}\sin\theta$; these solutions are, however, not consistent with the light cone description, as constraint(Al3) is not satisfied.

[Note that for more than 2 parameters $\lambda' \dots \lambda^r \dots \lambda^r$, Jacobian=1 would have still given $\partial_r \delta \lambda^r = 0$, which is solved by $f^r \equiv \delta \lambda^r \equiv \partial_s \hat{F}^{rs}$ provided \hat{F}^{rs} is an anti-symmetric tensor; the lie bracket on the space of all divergence-free vector fields $\vec{f}(\vec{\lambda})$ is

$$\{\vec{f}, \vec{g}\}^{T} = f \frac{g_{0}}{g_{\lambda}} - g \frac{g_{0}}{g_{\lambda}} - g \frac{g_{0}}{g_{\lambda}} = g_{0} (f_{0}^{ST} - g_{0}^{S})^{T}$$
(B4)

For M=2 there is only one independent antisymmetric tensor (\in^{rs}) so that $f^r = \partial_s \in^{rs} f$ so that (B4) translates to the Lie bracket (B3) for functions $f \in G$.

As an orthonormal basis of \underline{G} one can take the usual spherical harmonics $\begin{pmatrix} 0 & 4 \\ l_m & 4 \end{pmatrix}$ and define structure constants g by the equation:

(Summation over repeated indices is understood, unless stated otherwise; (lm) will often be abbreviated by a or α or just by ().) For definiteness, the definition* of spherical harmonics is given:

*This and all other conventions concerning angular momentum coupling—coefficients are those of A. Messiah Quantum Mechanics books, referred to as MI and MII (see especially Appendix C of MII).

$$V_{m} = \begin{cases} m & \text{if } m \geq 0 \\ 0 & \text{if } m \leq 0 \end{cases} \quad N_{em} = \sqrt{\frac{2\ell+1}{4\pi}} \frac{(\ell-lm_{1})!}{(\ell+lm_{1})!}$$

$$P_{em}(-\mu) = (-)^{\ell+lm_{1}} \frac{(\ell+\mu^{2})^{\ell+lm_{1}}}{2^{\ell}\ell!} \frac{d\ell+lm_{1}}{d\mu^{2}\ell+lm_{1}} \frac{(\mu^{2}-1)^{\ell}}{(\mu^{2}-1)^{\ell}} = (\ell-\mu^{2})^{\ell+lm_{1}} P_{em}(\mu^{2}-1)^{\ell+lm_{1}} P_{em}(\mu^{2}-$$

is an associated Legendre function of -µEcos0.)

Upon first inspection of (B5), one sees that

$$Q = + \int d\Omega \setminus \{Y_i\} \text{ is real } (\partial_{\varphi} \text{ gives } i_m)$$
and to tally surface the forts!)

$$Y = (-)^m Y_{\ell-m} \text{ (and } g \text{ real }) =)g_{(i_1(i_1)_{\ell}(i_2)_{\ell})} \text{ (B7)}$$

$$Y_{\ell-m} = (-)^{\ell} Y_{\ell-m} \text{ (B7)} = (-)^{\ell} Y_{\ell-m} (G_{i_1}g_{i_2}g_{i_3}g_{i_4}g_{i_5}g_{i_6}g_$$

where in the last step the decomposition of $\cos \Theta Y_{\ell m}$ into the linear combination $\sqrt{\dots}$

The group G itself has been studied in the mathematical literature, and although not relevant for the further discussion of the surface problem contained in this thesis, some properties will be listed.

- G is simple, i.e., has no nontrivial invariant subgroup H $\left(g + g^{-1} = H \quad \forall g \in G \right)$

- The Homotopy classes of G are those of SO(3) [Stephen Smale² proved this for the group of all diffeomorphisms; it then follows from a theorem by Moser³ that the same thing is true for G]
- any ge G has at least two fixed points [N.A. Nikishin⁴
 and C.P. Simon⁵]
- given $P_1...P_R \in S^2$, $Q_1...Q_R \in S^2 \rightarrow \exists g \in G$ with $Q_i = g(P_i)$ and furthermore let $C_1...C_R$ be an arbitrary collection of

¹I would like to thank Augustin Banyaga for telling me this and other things about G; as a reference see: A.B. "Sur la structure du groupe des diffeomorphismes qui préservent une forme symplectique". Comment. Math. Helv. <u>53</u>, 174-227 (1978).

²"Diffeomorphisms of the 2-sphere", Proc. Am. Math. Soc. 10, 1959.

³AMS Transactions 120, 1965, p. 287.

⁴Funct. Anal. Preloz. 8, 84-85 (1974) (in Russian).

⁵Inventiones Math. <u>26</u>, 187-200 (1974)

⁶See "Transformation Groups" by Kobayashi-Nomizu.

disjoint closed curves on s^2 , then there is a 1-parameter group of area preserving transformations with these curves as orbits.

Expanding x,y,p, and p, (θ, ψ) in spherical harmonics

one gets

One can think of H=T+V as describing infinitely many particles (labelled by ℓ and m) moving in two dimensions ($\times_{\ell m}$ and $y_{\ell m}$) and interacting through the, not very symmetric, quartic potential V. The unitary transformation

$$\begin{pmatrix}
\widetilde{X}_{l|m|} \\
\widetilde{X}_{l-l|m|}
\end{pmatrix} \equiv \frac{1}{2} \begin{pmatrix}
(-)^{lm|} & 1 \\
+ \cdot (-)^{lm|} & - \cdot \cdot
\end{pmatrix} \begin{pmatrix}
X_{l|m|} \\
X_{l-l|m|}
\end{pmatrix} \text{ the same for } y, p_x, p_y)$$

corresponding to a real basis

will make x real. The structure constants

are still totally antisymmetric (as the $\widetilde{Y}_{\ell,m}$ are orthonormal), but obey fewer selection rules than the $g\alpha\beta$ Some properties of the real \widetilde{Y} -basis:

from
$$Y_{em} = (-)^{m} N_{em} P_{em} e^{-im\varphi}$$
 (see (B6):

$$Y_{elm} = \frac{1}{2} (-)^{lm} Y_{elm} + Y_{e-lm}) = N_{em} P_{em} V_{e} Cos Im I\varphi$$

=) SYem Yem! = See, Sma,
So that
$$\mathfrak{F} = SYSY, YS$$

is still totally antisymmetric.

$$\begin{pmatrix} Y_{\ell | m_1} \\ Y_{\ell - \ell m_1} \end{pmatrix} = \mathcal{U} \begin{pmatrix} Y_{\ell | m_1} \\ Y_{\ell - \ell m_1} \end{pmatrix} \cdot \mathcal{U} = \frac{1}{2^2} \begin{pmatrix} (-1)^{\ell m_1} \\ (-1)^{\ell m_1} \\ (-1)^{\ell m_1} \end{pmatrix} (B9)$$

is unitary:

$$u^{\dagger}u = \frac{1}{2} \left(\frac{(-)^{lml} + i(-)^{lml}}{1 - i} \right) \left(\frac{(-)^{lml}}{1 + i} \right) = 1$$

for
$$X(\theta, \varphi) = \sum_{\ell m} X_{\ell m} = \sum_{\ell m} X_{\ell m} Y_{\ell m}$$

$$\begin{pmatrix} \tilde{X}_{\ell | m |} \\ \tilde{X}_{\ell - (m |} \end{pmatrix}$$

has to transform with the complex conjugate of U:

$$\left(\frac{\chi_{\ell | m|}}{\chi_{\ell - | m|}}\right) = \mathcal{U}^* \left(\frac{\chi_{\ell | m|}}{\chi_{\ell - | m|}}\right).$$

in shorthand notation:

$$\widehat{\overrightarrow{x}} = u^* \overrightarrow{x} , \quad \widehat{\overrightarrow{Y}} = u \overrightarrow{Y}$$

so that

written out

$$\tilde{X}_{e-|m|} = \frac{-(-)^{|m|}}{\mathcal{V}_{e-i}} (x_{e|m|} - x_{e+|m|}^{*}) = \frac{1}{\mathcal{V}_{i}} (-(-)^{|m|} + x_{e-|m|})$$

One has

where

and $\{\}_p$ is now the poisson bracket for functions of the canonical variables $x_{\ell m}^i$ and $p_{\ell m}^i$. The invariance of H under G is now expressed as

$$\left\{ H, \widetilde{g}_{4\beta} \widetilde{X}, \widetilde{P}_{\beta} \right\}_{P} = \left\{ H, K_{\chi} \right\}_{P} = 0$$

(which one can verify explicitly) The constants of the motion

$$K_{\chi} \equiv \int d\Omega \stackrel{\sim}{Y}_{\chi} \left\{ \vec{x}, \vec{p} \right\}$$
(B10)

are the generators of area preserving transformations, and, for the light cone coordinate description to be consistent, one has to have $K_{Lm}=0$ \swarrow_{Lm} . Note that, of course,

(see also below). As mentioned in part A, one can proceed to the quantum theory via the correspondence $\{,\}_{p}(-)-i[,],$ i.e.,

$$\begin{bmatrix} \widetilde{\mathbf{x}}_{\alpha}^{i}, \, \widehat{\mathbf{P}}_{\beta}^{j} \, \end{bmatrix} = i \, \mathcal{S}_{ij} \, \mathcal{S}_{ij} \, \left(\mathbf{t} = 1 \right)$$

for the hermitian operators x_{\star}^{i} and \tilde{p}_{β}^{j} . (From now on drop for the quantum mechanical operators.)

One finds that

as it must (they are a basis of the representation of \underline{G} as operators on Hilbert space),

$$\begin{bmatrix} K_{\alpha}, K_{\alpha} \end{bmatrix} = \begin{bmatrix} \widehat{g}_{\alpha \neq \beta} & \widehat{x}_{\alpha} & \widehat{p}_{\beta} & \widehat{g}_{\alpha \neq \beta} & \widehat{x}_{\beta} & \widehat{p}_{\delta} \end{bmatrix}$$

$$= \widehat{g}_{\alpha \neq \beta} \widehat{g}_{\alpha \neq \beta} \underbrace{\begin{cases} \widehat{x}_{\alpha} & \widehat{p}_{\beta} & \widehat{x}_{\beta} & \widehat{p}_{\beta} & \widehat{x}_{\beta} & \widehat{p}_{\beta} \\ \widehat{x}_{\alpha} & \widehat{y}_{\beta} & \widehat{y}_{\beta} & \widehat{x}_{\alpha} & \widehat{p}_{\beta} & \widehat{y}_{\beta} & \widehat{y}_{\beta} & \widehat{y}_{\beta} \\ \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\beta} & \widehat{y}_{\beta} & \widehat{y}_{\beta} & \widehat{y}_{\beta} & \widehat{y}_{\beta} & \widehat{y}_{\beta} \\ \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\beta} & \widehat{y}_{\beta} \\ \widehat{y}_{\alpha} & \widehat{y}_{\beta} \\ \widehat{y}_{\alpha} & \widehat{y}_{\alpha} \\ \widehat{y}_{\alpha} & \widehat{y}_{\alpha} \\ \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} \\ \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} \\ \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} \\ \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} \\ \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} \\ \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} \\ \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} \\ \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} \\ \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} \\ \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} \\ \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} \\ \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} \\ \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} \\ \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} \\ \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} \\ \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} \\ \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} \\ \widehat{y}_{\alpha} & \widehat{y}_{\alpha} \\ \widehat{y}_{\alpha} & \widehat{y}_{\alpha} \\ \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat{y}_{\alpha} & \widehat$$

Also one can check that $[K_{\alpha},H]=0$. The consistency condition (Al3) requires physical states $|\psi\rangle$ to be singlets under the symmetry group, i.e., $K_{\alpha}|\psi\rangle=0$ $\alpha=(\ell m)$. The change of a wavefunctional $\psi[x]$ under an infinitesimal are preserving reparametrization characterized by a function $f(\mu,\phi)$ is:

$$\int_{\mathcal{A}} \psi = \int_{\mathcal{A}} d\Omega \int_{\mathbf{X}_{i}} \int_{\mathbf{X}_{i}} \frac{\partial \psi}{\partial \mathbf{X}_{i}} \\
= \int_{\mathcal{A}} \int_{\mathcal{A}} \int_{\mathbf{X}_{i}} \int_{\mathbf{X}_{i}} \frac{\partial \psi}{\partial \mathbf{X}_{i}} \\
= + i \int_{\mathcal{A}} \int_{\mathcal{A}} \int_{\mathbf{X}_{i}} \int_{\mathbf{X}_{i}} \int_{\mathbf{Y}_{i}} \psi \\
= - i \int_{\mathcal{A}} \int_{\mathcal{A}} \int_{\mathbf{X}_{i}} \int_{\mathbf{X}_{i}} \int_{\mathbf{Y}_{i}} \psi \\
= - i \int_{\mathcal{A}} \int_{\mathbf{X}_{i}} \int_{\mathbf{X}_{i}} \int_{\mathbf{X}_{i}} \int_{\mathbf{Y}_{i}} \psi \\
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= - i \int_{\mathbf{X}_{i}} \psi \\
= - i \int_{\mathbf{X}_{i}} \psi \\$$

II. Explicit construction and proof that a basis of the fundamental representation of <u>SU(N)</u> can be chosen such that for the structure constants:

$$\lim_{N\to\infty} \hat{f}_{\alpha\beta\gamma}^{(N)} = g_{\alpha\beta\gamma}.$$

$$H_{N} = \frac{1}{2} \operatorname{Tr} \left\{ P_{X}^{2} + P_{y}^{2} - \frac{1}{N} \left[\times_{i} Y \right]^{2} \right\}$$

 (x,y,p_x) and p_y traceless hermitian N×N=matrices). Going from H to H_N is a sort of renormalization as one is cutting off the degrees of freedom corresponding to $Y_{\ell m}$ with $\ell \geqslant N$ ("High frequencies") while representing the low frequencies ($\ell \leqslant N-1$) correctly up to 0(1/N).

(BII) is subdivided into 5 sections as follows.

l. By a correspondence to the solid spherical harmonics $r^{\ell}Y_{\ell m}$ (written as harmonic polynomial) one defines N^2-1 linearly independent real, traceless N×N matrices $\tilde{T}_{\ell m}$ (\$\ell=1...N+1, |m|<\ell\$). They are a basis of the N-dimensional representation of $\underline{SU(N)}$. Also they are, for given \$\ell\$, tensor operators of degree \$\ell\$; so is any $T_{\ell m}$ differing from $\tilde{T}_{\ell m}$ by N and \$\ell\$ dependent (but m-independent!) factor.

^{*}The underlining always denotes the Lie algebra of the corresponding group.

2. Using the Wigner Eckart theorem, the structure constants of $\underline{SU(N)}$, defined by the relation $[T_{\ell m}, T_{\ell'm'}] = f_{\ell \ell' \ell''}^{(N)} T_{\ell''m''}^+$, can mm'm"

be calculated in terms of the reduced matrix elements $R_{N}(\ell)$. The answer also involves Wigner 3j- and 6j-symbols.

3. Instead of actually calculating the structure constants $g_{\alpha\beta\gamma} \text{ of } \underline{G} \text{ (α is a short-hand notation for (ℓm)), a proof is given that}$

$$\frac{1}{T_{\ell m}} \equiv \frac{\frac{1}{T_{\ell m}}}{\left(\frac{N-1}{4}\right)^{\ell-1}}$$

must lead to structure constants $\hat{f}_{\alpha\beta\gamma}$ that in the N+ ∞ limit are equal to the $g_{\alpha\beta\gamma}$ $\biguplus_{\alpha\beta\gamma}$. This proof is the central part of (BII).

4. Knowing this one can deduce the corresponding choice $\hat{R}_N(\boldsymbol{\ell})$, when calculating the N \rightarrow elimit of the structure constants derived in (2). This limit then is the formula for $g_{\alpha\beta\gamma}$. In (5) the correct choice $\hat{R}_N(\boldsymbol{\ell})$ is derived without using 3.

1. Definition of T_{lm}:

Let S_i be an N-dimensional representation of the Lie algebra SO(3), the spin S=N-1/2 representation. Conventionally one chooses a basis S_1 , S_2 , S_3 with

$$\langle Sm' | S_3 | Sm \rangle = m S_{m',m}$$

 $\langle m' | S_1 \pm iS_2 | m \rangle = \sqrt{S(S+1) - m (m\pm 1)} S_{m',m\pm 1}$
(B11)

 S_3 and $S_{\pm} \equiv S_1 \pm i S_2$ are real. One then defines N×N matrices $T_{\ell m}$ as polynomials of degree ℓ in the S_i which correspond in some sense to the $Y_{\ell m}(\theta, \theta)$. One does this by remembering that $r^{\ell}Y_{\ell m}$ are homogeneous, in fact harmonic, polynomials of degree ℓ in the variables $X_1 (\equiv r \cos \sin \theta)$, $X_2 (\equiv r \cos \theta)$ and $X_3 (\equiv r \cos \theta)$

$$\int_{\mathbb{R}^{m}} \int_{\mathbb{R}^{m}} \left(O_{1} \varphi \right) dm \left(O_{1} \varphi \right) dm \left(O_{1} \varphi \right)$$

$$= \sum_{\mathcal{I} = \mathcal{L}} \hat{a}_{1}^{(m)} \times_{1}^{j} X_{1}^{j} X_{2}^{j} \times_{3}^{j} = \sum_{\mathcal{I} = \mathcal{I}} \hat{a}_{1}^{(m)} \times_{1}^{j} \times_{1}^{j} \times_{1}^{j}$$

$$= \sum_{\mathcal{I} = \mathcal{L}} \hat{a}_{1}^{(m)} \times_{1}^{j} X_{2}^{j} \times_{3}^{j} = \sum_{\mathcal{I} = \mathcal{I}} \hat{a}_{1}^{(m)} \times_{1}^{j} \times_{1}^{j} \times_{1}^{j}$$

$$= \sum_{\mathcal{I} = \mathcal{L}} \hat{a}_{1}^{(m)} \times_{1}^{j} X_{2}^{j} \times_{2}^{j} \times_{3}^{j} = \sum_{\mathcal{I} = \mathcal{I} = \mathcal{I}} \hat{a}_{1}^{(m)} \times_{1}^{j} \times_{1}^{j} \times_{1}^{j}$$

$$= \sum_{\mathcal{I} = \mathcal{L}} \hat{a}_{1}^{(m)} \times_{1}^{j} X_{2}^{j} \times_{2}^{j} \times_{3}^{j} = \sum_{\mathcal{I} = \mathcal{I}} \hat{a}_{1}^{(m)} \times_{1}^{j} \times_{1}^{j}$$

$$= \sum_{\mathcal{I} = \mathcal{L}} \hat{a}_{1}^{(m)} \times_{1}^{j} X_{2}^{j} \times_{2}^{j} \times_{3}^{j} = \sum_{\mathcal{I} = \mathcal{I}} \hat{a}_{1}^{(m)} \times_{1}^{j} \times_{1}^{j}$$

$$= \sum_{\mathcal{I} = \mathcal{L}} \hat{a}_{1}^{(m)} \times_{1}^{j} X_{2}^{j} \times_{2}^{j} \times_{3}^{j} = \sum_{\mathcal{I} = \mathcal{I}} \hat{a}_{1}^{(m)} \times_{1}^{j} \times_{1}^{j}$$

$$= \sum_{\mathcal{I} = \mathcal{L}} \hat{a}_{1}^{(m)} \times_{1}^{j} \times_{1}^{j} \times_{1}^{j} \times_{1}^{j}$$

$$= \sum_{\mathcal{I} = \mathcal{L}} \hat{a}_{1}^{(m)} \times_{1}^{j} \times_{1}^{j} \times_{1}^{j} \times_{1}^{j}$$

$$= \sum_{\mathcal{I} = \mathcal{L}} \hat{a}_{1}^{(m)} \times_{1}^{j} \times_{1}^{j} \times_{1}^{j}$$

$$= \sum_{\mathcal{I} = \mathcal{L}} \hat{a}_{1}^{(m)} \times_{1}^{j} \times_{1}^{j} \times_{1}^{j}$$

$$= \sum_{\mathcal{I} = \mathcal{L}} \hat{a}_{1}^{(m)} \times_{1}^{j}$$

$$= \sum_{\mathcal{I} = \mathcal{L}} \hat{a}_{1}^{(m)}$$

$$= \sum_{\mathcal{I}$$

The $\mathring{Q}_{1}^{(m)}$ defined this way are traceless between any two indices ($\Leftrightarrow \nabla^2 \mathcal{V}_{\ell m} = 0$) and totally symmetric. For given ℓ there are $2\ell+1$ independent ones. Then define:

$$\frac{1}{1} = \sum_{i=1}^{3} \hat{\alpha}_{i-i}^{(m)} S_{i} \cdots S_{i}$$

$$\lim_{i \to \infty} \hat{\beta}_{i-i}^{(m)} S_{i} \cdots S_{i}$$
(B13)

The first few ones are:

All $\mathbf{T}_{\ell m}$ are by definition real and traceless, but not hermitian:

$$\left(\stackrel{\circ}{\mathsf{T}}_{\ell m}\right)^{\dagger} = \left(\stackrel{\circ}{\mathsf{T}}_{\ell m}\right)^{\dagger r} = \left(-\right)^{m} \stackrel{\circ}{\mathsf{T}}_{\ell - m}$$

$$(a_{i-1}^{(m)})^* = (-)^m a_{i-1}^{(-m)}$$

For fixed 1, the T_{lm} form a set of tensor operators of rank ${\cal L}$, i.e., for a rotation R:

$$U(R) \stackrel{\circ}{\uparrow}_{\ell m} U(R) = \stackrel{+\ell}{\downarrow}_{\ell m'} \stackrel{\circ}{\uparrow}_{\ell m'} (R) \qquad (B14)$$

where $R_{m'm}^{\ell}$ are the rotation matrices for angular momentum ℓ (see for instance Messiah II, p. 1070) and U(R) is a N-dimensional representation of the rotation R). [If R \in SO(3), N would have to be odd, and later one would take ℓ im $f^{(N)}$; $N \to \infty$

(N odd!)

but are might as

well take RC SU(2) which does not alter anything as the two Lie algebras $\underline{SU(2)}$ and $\underline{SO(3)}$ are the same.] Changing the normalization of the $T_{\ell m}$'s in an m-independent way will not alter the transformation properties. Therefore any $T_{\ell m} = U(\ell,N) T_{\ell m}$ will obey the Wigner Eckart theorem* :

$$\langle Sm_1 | T_{\ell m}^{(N)} | Sm_2 \rangle = (-)^{S-m_1} \begin{pmatrix} S\ell S \\ m_1 m m_2 \end{pmatrix} R(\ell)^{(B15)}$$

where () denotes the 3j-symbol* and $R_N(\ell)$ the reduced matrix element (real for real $U(\ell,N)$). $R_N(\ell) \equiv R_N^0(\ell) \cdot \mathcal{U}(N_\ell \ell)$

has been left general, as different normalizations will be useful in different situations.

$$\begin{bmatrix} T_{\ell m}, T_{\ell'm'} \end{bmatrix} = f_{\ell \ell' \ell''}^{(N)} T_{\ell'' m''}^{\dagger} \qquad (B16)$$

By using (B15) and standard formulae concerning coupling of angular momenta one can proceed to calculate (T,T) and $f^{(N)}$. This is done in the next section.

^{*}See e.g., MII. p. 1056.

2. Calculation of $\text{Tr}(\text{TT}^+)$, Tr(TTT), and choice of $R_N(\ell)$ From (E15) one has

$$T_{T} T_{\ell m} T_{\ell m'}^{\dagger} = \sum_{m_{1}, m_{2}} \langle m_{1} | T_{\ell m} | m_{2} \rangle \langle m_{2} | T_{\ell m'} | m_{1} \rangle$$

$$= R_{N}(\ell) R_{N}(\ell') \sum_{m_{1}, m_{2}, m'} \langle m_{1} | T_{\ell m'} | m_{1} \rangle \langle m_{2} | T_{\ell m'} | m_{1} \rangle$$

$$= R_{N}(\ell) R_{N}(\ell') \sum_{m_{1}, m_{2}, m'} \langle m_{1} | T_{\ell m'} | m_{1} \rangle \langle m_{2} | T_{\ell m'} | m_{1} \rangle$$

As the second 3j-symbol is 0 unless $m_1 = m_2 + m^4$, $2S - m_1 - m_2 + m^4$ has to be even and (-) $2S - m_1 - m_2 + m^4$ therefore =+1. Further

$$\sum ()() = \sum_{m_{1}m_{2}} {SlS \choose m_{1}m_{1}m_{2}} {SlS \choose -m_{2}-m'_{1}-m_{1}}
= \sum {SSl \choose m_{1}m_{2}m} {SlS \choose m_{2}m'_{1}m_{1}} (-)
= (-) l+l' \sum {SSl \choose m_{1}m_{2}m} {SSl \choose m_{1}m_{2}m'_{1}} (-)
= (-) l+l' Sll' Smm'_{2l+1} (-)
= (-) l+l' Smm'_{2l+1} (-)
= (-$$

where in the first step m_1 was changed to $-m_1$; in the second step 2nd and 3rd column of the first 3j-symbol were interchanged (giving factor (-) $^{2S+l}$) and in the second 3j-symbol the sign of the lower row was changed (i.e., $-m_{\alpha}^{\rightarrow}+m_{\alpha}^{\rightarrow}$, giving a factor (-) $^{2S+l}$); in the third step invariance of () under cyclic permutations of the 3 columns was used; the last step is true because of Eq. (Cl5a), p. 1057, MII. Therefore

$$T_{r}\left(T_{\ell m}T_{\ell m'}\right) = \int_{\ell' \ell} \int_{m'm} \frac{R_{N}^{2}(\ell)}{(2\ell+1)}$$
(B17)

i.e., the $T_{\ell m}$'s are orthogonal (with the choice $R_N^{(\ell)} \equiv \sqrt{2\ell+1}$ ' they would be orthonormal.) Note that $T_{\ell m} \equiv 0$ for $\ell > N = 2S+1$, as $\binom{S R S}{\ldots} = 0$ then.

But this means that one has constructed this way exactly N^2-1 ($3+5+\ldots+N-1$) independent traceless real N×N matrices. They, therefore, furnish a basis of the fundamental (i.e., N-dimensional) representation of the Lie algebra SU(N), and the $f_{\alpha\beta\gamma}^{(N)}$ defined via (B16) are the structure constants of SU(N) in this basis. They will now be calculated:

$$T_{\tau} \left(T_{\ell_{1}m_{1}} T_{\ell_{2}m_{2}} T_{\ell_{3}m_{3}} \right) =$$

$$= \sum_{s=1}^{35-m-m'-m''} T_{\ell_{1}m_{2}} \left(S_{\ell_{1}} S_{s} \right) \left(S_{\ell_{2}} S_{s} \right) \left(S_{\ell_{3}} S_{s} \right) \left(S_{\ell_$$

Now change summation variables to $M_2 \equiv -m$ $M_3 \equiv -m$, $M \equiv -m$, in all three 3-j symbols interchange 2nd and 3rd row, picking up a factor of $(-)^{6S+\sum l_i} = (-)^{2S+\sum l_i}$ altogether;

use formula (C33), p. 1064 in MII, with the identification $(l, m_1) \leftarrow (J_1, m_2)$ and $J_1 = J_2 - J_3 \equiv S$ $\prod_i \in \prod_i f_i$ to get: $T_{i} \left(T_{i} \prod_{i \in I_{i}, m_{i}} T_{i} \prod_{i \in I_{i}, m_{i}} T_{i} \right)$

$$= \operatorname{TR}_{N}(\ell_{i}) \left(-\right)^{2s+\sum \ell_{i}} \left(\frac{\ell_{i} \ell_{i} \ell_{2}}{m_{i} m_{2} m_{3}}\right) \left\{ \begin{array}{c} \ell_{i} \ell_{2} \ell_{3} \\ 5 5 5 \end{array} \right\}$$

where {} denotes the Wigner 6j-symbol.

$$= \left(T_{R_{N}} T_{\ell_{1}} m_{1} T_{\ell_{3}} m_{3} \right)$$

$$= \left(T_{R_{N}} (\ell_{i}) \right) (-)^{2S} \begin{pmatrix} \ell_{1} \ell_{2} \ell_{3} \\ m_{1} m_{2} m_{3} \end{pmatrix} \begin{cases} \ell_{1} \ell_{2} \ell_{3} \\ S S S \end{cases}$$

as {} is invariant under interchange of two columns, while () has to be multiplied by $(-)^{l_1+l_2+l_3}$. Therefore,

$$\begin{cases}
R_{N}(l_{1})R_{N}l_{2} \\
R_{N}(l_{3})
\end{cases} (2l_{3}+1) 2 (-)^{N} \begin{pmatrix} l_{1}l_{2}l_{3} \\ m_{1}m_{2}m_{3} \end{pmatrix} \begin{cases} l_{1}l_{1}l_{3} \\ SSS \end{cases} \\
R_{N}(l_{3}) \qquad (2l_{3}+1) 2 (-)^{N} \begin{pmatrix} l_{1}l_{2}l_{3} \\ m_{1}m_{2}m_{3} \end{pmatrix} \begin{cases} l_{1}l_{1}l_{3} \\ SSS \end{cases} (B18)$$

$$\begin{pmatrix} n_{1} & T l_{1} & even \end{pmatrix}$$

Note that f=0 if one ℓ_i >N. Also f=0 unless Σm_i =0 and the ℓ_i satisfy the triangle inequalities. (B18 was obtained from (B16):

(B18) is a formula for the structure constants of $\underline{SU(N)}$ in the basis $T_{\ell m}^{0} = T_{\ell m} \cdot U(N, \ell)$.

Particular choices for $R_{N}(\ell) = R_{N}^{O}(\ell) \cdot U(N, \ell)$ are:

ii)
$$R = \overline{R}_{N}(\ell) \equiv \sqrt{2\ell+1}$$

This choice will make $f^{(N)}$ totally antisymmetric for all N. $(=T_{\ell m}$ -basis orthonormal) (any R differing from \overline{R} by an ℓ -dependent factor will not have this property).

iii)
$$R = \frac{(-)^{N}\sqrt{2\ell+1}}{2}$$
 leads to
$$\begin{cases}
\binom{N}{\ell_1\ell_2\ell_3} = \binom{\ell_1\ell_2\ell_3}{m_1m_2m_3} \begin{cases}
\ell_1\ell_2\ell_3 \\
m_1m_2m_3
\end{cases} \begin{cases}
\ell_1\ell_2$$

totally antisymmetric (\maltese_N) , but in addition, the corresponding $f^{(N)}$ will have a finite, non-zero limit as $N+\infty$, as

$$\left\{\begin{array}{cc} l_1 & l_2 & l_3 \\ s & s & s \end{array}\right\} \times N^{-3/2}$$

for S+ ∞ (see below). Therefore, if there is \underline{any} choice of R (at all) for which

$$\lim_{N\to\infty} \int_{\alpha/\gamma}^{(N)} = g_{\alpha/\gamma}$$
 one can use the above choice $R = N^{3/2} \sqrt{2\ell+1}$ to calculate $g_{\alpha\beta\gamma}$ (up to a constant independent of $\alpha\beta$ and γ —which turns out to be $(16\pi)^{-1/2}$). The "correct" choice is

$$\widehat{R}_{N}(\ell) = \sqrt{\frac{(N+\ell)!}{(N-\ell-1)!}} \frac{\sqrt{N^{2}-1}}{\sqrt{16\pi}} \sqrt{2\ell+1}$$
(B20)

3

0,85° Pigit 115° and is based on the proof given in the next section. (As $N \to \infty$, \hat{R}_N differs from (B19) only by

3. There is a basis $\hat{T}_{\ell m}^{(N)}$ of $\underline{SU(N)}$ with $\lim_{N\to\infty} \hat{f}_{\alpha\beta\gamma}^{(N)} = g_{\alpha\beta\gamma}$ (constructive proof)

First look at the process that determines $g_{\alpha\beta\gamma}$:

$$\iint_{\ell_m} = \int_{\ell_m} \left\{ \sum_{m=1}^{\infty} \hat{a}_{m,m}^{(m)} \times_{n-1} \times_{n} \right\}$$
are harmonic polynomials, \hat{a} traceless (and symmetric of course).

$$\{f(94),g(94)\}=\frac{1}{10}(3f3g-3g3f)$$

is in fact what one gets when one restricts the space of all polynomial functions $f(x_1x_2x_3)$ with the Lie bracket defined as

$$\{f,g\}\equiv \epsilon_{ijk}\times_i \partial_j f \partial_k g$$
 (B21)

to functions on the unit sphere $\left(X_{i}^{2}+X_{i}^{2}+X_{j}^{2}=\overline{X}^{2}=1\right)$ That (B21) really defines a Lie bracket (i.e., { } satisfies the Jacobi identity; {f,f}=0 is trivial) is shown in III; there in fact for C_{ijk} being replaced by any C_{jk}^{i} antisymmetric in j and k and satisfying the Jacobi identity. Therefore:

Ŋ,

$$\begin{cases}
Y_{lm}, Y_{lm} \\
Y_{lm} \\$$

using $\{fg, k\} = f\{g, k\} + \{f, k\}g$ The order is, of course, completely irrelevant as and

The order is, of course, completely irrelevant, as everything commutes, and with $\left\{ \begin{array}{c} \chi_{1,k} \\ \chi_{1,k} \end{array} \right\} = \mathcal{C}_{1,k,l} \mathcal{C}_{k}$ one gets (B21) as one must. There were/are two reasons for having written down (B22). The first is that it makes slightly more apparent the following decomposition of $\left\{ \begin{array}{c} \chi_{l-1} \\ \chi_{l-1} \end{array} \right\} = \mathcal{C}_{1,k,l} \mathcal{C}_{k}$ which is a homogeneous but no longer harmonic polynomial $\left\{ \begin{array}{c} \chi_{l-1} \\ \chi_{l-1} \end{array} \right\} = \mathcal{C}_{1,k,l} \mathcal{C}_{k}$ one gets one gets (B21) as one must. There were/are two reasons for having written down (B22). The first is that it makes slightly more apparent the following decomposition of $\left\{ \begin{array}{c} \chi_{l-1} \\ \chi_{l-1} \end{array} \right\} = \mathcal{C}_{1,k,l} \mathcal{C}_{k}$ which is a homogeneous but no longer harmonic polynomial and lower:

$$P_{\ell+\ell'-1} = \sum_{m'', \ell''=1}^{\ell+\ell'-1} d_{\ell} e_{\ell} e_{\ell} (x^{2})^{\frac{2}{2}} \cdot y^{*}$$

$$(\ell+\ell'-1=\ell'' even)$$

$$(\ell+\ell'-1=\ell'' even)$$

(corresponding to making $\alpha^{(m)}$ $\alpha^{(m')}$ $\alpha^{(m')}$ traceless and totally symmetric). By restriction to the wait

and totally symmetric). By restriction to the unit sphere, one sees that $d_{\ell,\ell,m} = -ig_{\ell,\ell,m}$ (see B5). The second reason is that (B22) stresses the connection between $Y_{\ell,m}$ and $T_{\ell,m}$ as the expression for $[T_{\ell,m},T_{\ell,m}]$ will be exactly like (B22) just with

が、中で、一人の場合に、 佐藤林の様の

$$X_{i} \rightarrow S_{i}$$
 and $\{X_{i} \times_{p}\} \rightarrow [S_{i}, S_{i}]$

$$(= i \in_{i \neq j \neq k} S_{k}).$$

$$\sum_{m'' \ell'' = 1}^{\ell' + \ell' - 1} \int_{-\infty}^{\ell' + \ell' - 1} \int_{-\infty}^{\infty} \int_{-\infty$$

however is more complicated, as the S_i are non commuting objects so that the process of making $a^{(m)}e^{(m')}\in \mathcal{A}_{j,k}$ traceless and symmetric, which involves moving the S_i around,

using
$$S_i S_i = S_i S_i + i \epsilon_{ijk} S_k$$
 (B24)

will give lower order polynomials. Therefore

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with highest order term

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will contain lower powers of N. $X_N = S_1^2 + S_2^2 + S_3^2 = \frac{1}{4} \pm \frac{1}{4}$ of course arises from the trace contributions) But what is important is that all terms in (B23') of degree $\ell + \ell - 1$ (X_N has degree 2, X_N degree ℓ "), i.e., all terms

arose from always picking up the first term in (B24), i.e., treating the $S_{\dot{1}}$ as commuting objects, in effect. Therefore:

$$fee'e'' = ndee'e''$$
mm'm'

mm'm''

(the i,as [,] gives an extra i compared with {}). This means that for $\frac{1}{l_m} = \frac{1}{l_m} / \frac{l_m}{2}$ leading to

one has
$$\int_{N}^{(N)} \left(\frac{1}{X} \right) = \frac{1}{X} \int_{mm'm'}^{2} \left(\frac{1}{X} \right) = i \int_{N}^{2} \left(\frac{1}{X} \right) dt dt dt$$
The second of the second of

$$id = i(-ig) = g \qquad \lim_{N \to \infty} \hat{f}_{(N)}^{(N)} = g \qquad (B28)$$

$$\lim_{N \to \infty} \hat{f}_{(N)}^{(N)} = g \qquad (B28)$$

Thus one has the desired result.

4. Calculation of gαβγ

Because of (B28) one can now use (B18) (for properly chosen $R_N(\mathbb{L})$) to calculate $g_{\alpha\beta\gamma}$. First one has to find out the behavior of $\left\{ \begin{array}{cc} \mathcal{L}_1 \ \mathcal{L}_2 \ \mathcal{L}_3 \end{array} \right\}$ as $N^{+\infty}$.

Racah's formula [see e.g. MII, p. 1065,

for
$$| \leq l_1 \leq l_2 \leq l_3$$
, $l_1 + l_2 \leq 2S = N-1$

is

$$\left\{ \begin{array}{c} \mathcal{L}_{1} \ \mathcal{L}_{2} \ \mathcal{L}_{3} \\ S \ S \ S \end{array} \right\} = \mathcal{L}_{1} \mathcal{L}_{2} \mathcal{L}_{3} \mathcal{L}_{3} \mathcal{L}_{3} \mathcal{L}_{4} \mathcal{L}_{4} \mathcal{L}_{3} \mathcal{L}_{4} \mathcal{L}_$$

•
$$\sqrt{\frac{(2s-\ell_1)!}{(2s+1+\ell_1)!}}$$
 $\sqrt{\frac{(2s-\ell_2)!}{(2s+1+\ell_2)!}}$ $\sqrt{\frac{(2s-\ell_3)!}{(2s+1+\ell_3)!}}$ $(-)^{2s+1}$

$$\sum_{k=0}^{\ell_{1}+\ell_{2}-\ell_{3}} \frac{(2S+\ell_{3}+x+1)!}{(2S+x-\ell_{1}-\ell_{2})!} \frac{(-)^{\times}}{x!(\ell_{1}+\ell_{2}-\ell_{3}-x)!(\ell_{1}-x)!(\ell_{2}-x)!(\ell_{2}+x-\ell_{1})!(\ell_{3}+x-\ell_{1})!$$

$$\equiv \mathcal{J}(\mathcal{L}_{i}) + (\mathcal{L}_{i})(-) \sum_{X=0}^{N} \mathcal{G}_{N}(X_{i},\mathcal{L}_{i}) \frac{(-)^{X}}{F(X_{i},\mathcal{L}_{i})}$$
(B29')

[i.e., J includes all N and x-independent factors; H_N consists of the remaining x-independent factors (apart from $(-)^N$) G_N depends on both x and N, F is independent of N]

$$+ \frac{3}{N} \sqrt{\frac{(N - (l_1 + 1))!}{(N + l_1)!}} = \frac{3}{1!} \left((N + l_1 + 1) - (N - l_1) \right) \xrightarrow{-3} N^2 N^{-l_1 - l_2 - l_3}$$

the leading term in G_N is $N^{\ell_1+\ell_2+\ell_3+\ell}$. However, any independent term in G_N will give 0, as $F(x;\ell_2)$ is invariant under $\times \longrightarrow (\ell_1+\ell_2-\ell_3)-\times$, the # of terms in \sum is even

(as
$$l_1 + l_2 - l_3$$
 is odd), and therefore $\sum_{x} \frac{(-)^x}{F(x)} = 0$

The leading contributing term in

$$G_{N} = \frac{(N+x+l_{3})!}{(N+x-(l_{1}+l_{2}+1))!}$$
 is kerefore
$$N^{l_{1}+l_{2}+l_{3}} (l_{1}+l_{2}+l_{3}+1) \cdot X$$

The leading term in $\begin{cases} l_1 l_2 l_3 \\ s s \end{cases}$ as $N \rightarrow \infty$, is therefore:

$$\frac{\int_{1}^{-3/2} (-)^{N} (l_{1} + l_{2} + l_{3} + 1) \sum_{x=0}^{l_{1} + l_{2} - l_{3}} \frac{\chi(-)^{x}}{F(x)}}{F(x)}$$
(B29")

$$\int_{l_{1}, l_{2}, l_{3}}^{(N)} = \frac{R_{N}(l_{1})R_{N}(l_{2})}{N^{3/2}} \frac{2l_{3}+1}{R_{N}(l_{3})} \frac{2l_{3}+1}{2l_{3}+1} \cdot 2 \cdot \left(\sum l_{1}+1\right) \cdot \underbrace{\sum \frac{\times (-)^{\times}}{F(\kappa_{1} \cdot l_{3})}}_{\chi = 0} \frac{l_{1} + l_{2} - l_{3}}{I_{1} \cdot I_{3}} \frac{(l_{1} l_{1} l_{3})}{I_{2} \cdot I_{3}} \frac{(l_{1} l_{1} l_{3})}{I_{2} \cdot I_{3}} \frac{(l_{1} l_{1} l_{3})}{I_{3} \cdot I_{3}} \frac{(l_{1} l_{1} l_{3})}{I_{3} \cdot I_{3}} \frac{(l_{1} l_{2} l_{3} l_{3})}{I_{3} \cdot I_{3}} \frac{(l_{1} l_{3} l_{3} l_{3} l_{3})}{I_{3} \cdot I_{3}} \frac{(l_{1} l_{3} l_{3} l_{3} l_{3})}{I_{3} \cdot I_{3}} \frac{(l_{1} l_{3} l_{$$

and $F(x) \equiv x! (l_1 + l_2 - l_3 - x)! (l_1 - x)! (l_1 - x)! (l_2 - x)! (x + l_3 - l_1)! (x + l_3 - l_2)! ,$ $| \leq l_1 \leq l_2 \leq l_3, \quad l_1 + l_2 \leq N - 1$ (B30)

(these two conditions are slightly artificial as they are only necessary to write down $\left\{ \begin{array}{c} \mathcal{L}_1 \ \mathcal{L}_2 \ \mathcal{L}_3 \\ \mathcal{L}_3 \ \mathcal{L}_3 \end{array} \right\}$ as explicitly as in (B29)).

Because of (B18), and because \hat{R}_N has to behave like $N^{3/2}$ value const. for large N (so to make (B30) finite and totally antisymmetric as $N\!\!\rightarrow\!\!\infty$), $g_{\alpha\beta\gamma}$ can be calculated as:

$$\lim_{N\to\infty} \int_{\sqrt{3}\chi}^{(N)} \mathbb{C} R_N(\ell) = N^{3/2} \sqrt{2\ell+1} \cdot \text{canot.}$$

where the constant can be determined by comparing g and $\lim_{N\to\infty} f^{(N)} \quad \text{(calculated via (B30)) in just one simple case,}$

e.g., $\begin{bmatrix} 1 & 1 & 1 & 1 \\ 0 & +1 & -1 & 1 \end{bmatrix}$. As not more work is involved one calculates $f^{(N)}$ for the case $\begin{bmatrix} 1 & 2 & 2 \\ 0 & m & -m \end{bmatrix}$: with

(B30) gives, for $R=N^{3/2}\sqrt{2(1-1)}$ (const.):

$$f_{1RR} = (carst.)(2l+1)V_3(2l+2)(l!)^2 \cdot (2l-1)! \frac{1}{(2l+2)! \cdot l!(l-1)!}$$

$$= 2V_3 m(-1)^m carst. \frac{(-1)^m m}{\sqrt{l(l+1)(2l+1)}}$$

^{*}See e.g. MII, p. 1060.

which agrees with

$$\int_{0m-m}^{\infty} |e| = m(-)^m \sqrt{\frac{3}{4\pi}} \quad (\text{from B8})$$

provided const = \sqrt{kT} , and provides already a first check, as ℓ and m are general in ℓ ℓ so ℓ o ℓ o ℓ so

$$R_{N}(\ell) = \frac{1}{\sqrt{16\pi}} N^{3/2} \sqrt{2\ell+1}$$
(B31)

can be used to calculate

$$=\frac{1}{\sqrt{4\pi}}\left(\ell_{1}+\ell_{2}+1\right)\sqrt{\sqrt{2}\ell_{1}+1}\int_{-1}^{1}\left(\frac{\ell_{1}\ell_{2}\ell_{3}}{m_{1}m_{1}m_{3}}\right)\frac{\ell_{1}+\ell_{2}-\ell_{3}}{F(x_{1}^{*},\ell_{1}^{*})}$$
(B32)

Since g is totally antisymmetric, it can be calculated via (B32) for all ($\alpha\beta\gamma$). We check (B32) for $\begin{bmatrix} 2 & \mathcal{L} + 1 \\ 0 & m - m \end{bmatrix}$: Using

$${2 l l + 1 \choose 0 m - m} = (-)^{l - m} 2_m \sqrt{\frac{6(l + m + 1)(l - m + 1)}{(2l + 4)(2l + 3)(2l + 2)(2l + 1)2l}}$$

one finds that $\begin{cases} 2 & \text{let} \end{cases}$ calculated via (B32), agrees with

$$\frac{9}{32\ell \ell + 1} = 3\sqrt{\frac{5}{4\pi}} m(-)^{m} \sqrt{\frac{(\ell + m + 1)(\ell + 1 - m)}{(2\ell + 1)(2\ell + 3)}}$$
from (B8')

Finally it is useful to calculate f^N for $\ell_3 = \ell_1 + \ell_2 - 1$: with the notation as in (B29¹) one has for this case

$$\sum_{X=0}^{l+l_1-l_3} G(X, l_1) = \sum_{Y=0}^{l} \frac{(N+x+l_3)!}{(N+x-(l_1+l_1+l_1)!)!} G(X, l_1)$$

$$= \frac{1}{F(0)} \left(\frac{(N+l_3)!}{N-(l_1+l_2+l)!} - \frac{(N+l_3+l)!}{(N+l-(l_1+l_2+l))!} \right)$$

$$= \frac{1}{F(0)} \left(\frac{(N+l_3)!}{N-(l_1+l_2+l)!} - \frac{(N+l_3+l_1)!}{(N+l-(l_1+l_2+l_1))!} \right)$$

$$= \frac{1}{F(0)} \left(\frac{(N+l_3)!}{(N-l_1-l_2)!} - \frac{(N-l_3-l_3-l_3)!}{(N-l_1-l_2)!} \right)$$

$$= \frac{-2}{F(0)} \left(\frac{(N+l_3+l_2-l_3)!}{(N-l_1-l_2)!} - \frac{(l_1+l_2)!}{(l_1+l_2)!} \right)$$

Therefore, using (B18) and (B29):

$$\int_{\ell_{1}\ell_{2}}^{(N)} \ell_{3} = \frac{R_{N}(\ell_{1})R_{N}(\ell_{2})(2\ell_{1}+2\ell_{2}-1)}{R_{N}(\ell_{1}+\ell_{2}-1)} 2 \left(\frac{\ell_{1}\ell_{1}\ell_{1}+\ell_{1}-1}{m_{1}m_{2}m_{3}} \right) \\
- \frac{2(\ell_{1}+\ell_{2})}{F(0)} \int_{-1}^{\infty} \left(\frac{\ell_{1}\ell_{2}-1}{k_{1}+\ell_{2}-1} \right) \left(\frac{1}{N-\ell_{2}-1} \right) \left(\frac{1}{N-\ell_{2}$$

The proof in (BII3) showed in particular that R corresponding

to

leads to $f^{(N)}$ that is independent of N for $\mathcal{L}_3 = \mathcal{L}_1 + \mathcal{L}_2 - 1$ A factor

must therefore be contained in $\overset{\wedge}{R_N}$. This factor $\overset{\wedge}{R_N}$ as $N^{+\infty}$. To have $\overset{\wedge}{R_N}$ $N^{3/2}$ (as $N^{+\infty}$), which is needed so that \hat{f} has a finite non-zero limit, one must include another factor that is $\approx N^{1-\mathcal{L}}$ (as $N^{+\infty}$). Because of

one is led to choose this factor to be $\sqrt{N^2/l}$. Putting all together

$$\stackrel{\wedge}{R} = \frac{\sqrt{2\ell+1}}{\sqrt{|GT|}} \sqrt{\frac{(N+\ell)!}{(N-\ell-1)!}} \sqrt{\frac{N^2-1}{N^2-1}} = \ell$$
(B20)

5. Direct calculation of \hat{R}_N

Rather than by making heavy use of the proof BII3 and deducing \hat{R}_N by the above arguments, \hat{R}_N can be derived directly from the correspondence to the $Y_{\ell m}$ and the properties of that \hat{R} will then provide a check on the proof, rather than relying on it!

$$\frac{\partial}{\partial t} = \frac{(-1)^{\ell}}{\ell!} \frac{(2\ell+1)!}{4\pi} \left(S_{\chi} + i S_{\chi} \right)^{\ell}$$
(from (B12))

$$\frac{\overset{\circ}{R}_{N}^{2}(\ell)}{2\ell+1} = T_{T} \left(\overset{\circ}{T}_{\ell} + \overset{\circ}{T}_{\ell} + \frac{(2\ell+1)!}{(4\pi)(\ell!)^{2}} \frac{(2\ell+1)!}{2^{2\ell}} \right) (B17)$$

$$T_{r}\left(S_{+}^{2}S_{-}^{2}\right)$$

$$= \sum_{m} \langle m|S_{+}|m-1\rangle\langle m-1|S_{+}|m-2\rangle \dots \langle m+1-\ell|S_{+}|m-\ell\rangle} \langle m-\ell|S_{-}|m-\ell\rangle \cdot \langle m-\ell|S_{-}|m-\ell$$

(B34) can be proved in the following way: Since one has

$$(|-x|)^{-m-1} = \sum_{\tau} {m+\tau \choose \tau} x^{\tau}$$
one has
$$(|-x|)^{-m-1} (|-x|)^{-m-1} = \sum_{\tau \mid S} {m+\tau \choose \tau} x^{\tau+S} {m+s \choose S}$$

but also

$$= (1-x)^{-m-m-2} = \sum_{t} {m+m+1+t \choose t} x^{t}$$

$$= \sum_{so \text{ that } r+s=t} {m+r \choose r} {m+s \choose s} = {m+m+1+t \choose t}$$

$$= \left(\frac{m+m+1+t}{m+m+1} \right)$$

Since

$$\begin{pmatrix} \chi + \ell \\ \ell \end{pmatrix} \begin{pmatrix} N - 1 - \chi \\ \ell \end{pmatrix} \equiv \begin{pmatrix} \chi + \ell \\ \chi \end{pmatrix} \begin{pmatrix} N - 1 - \chi \\ N - 1 - \ell - \chi \end{pmatrix}$$

one obtains (B34) by identifying

Using (B30) one has

$$\mathring{R}_{N}(\ell) = \sqrt{N+\ell!} \frac{1}{\sqrt{4\pi}} \sqrt{2\ell+1}$$

and
$$\hat{R}_{N} = \frac{\hat{R}_{N}}{\sqrt{N}} e^{-1} = \frac{\sqrt{2\ell+1}}{\sqrt{KT}} \sqrt{\frac{N+\ell+1}{N-\ell-1}!} \sqrt{N^{2}-1} - \ell$$
(which agrees with (\$\mathbb{2}\mathcal{2}\mathcal{2}\mathcal{2}\mathcal{2}\mathcal{2}\mathcal{2}\mathcal{2}\mathcal{2}\mathcal{2}\mathcal{2}\mathcal{2}\mathcal{2}\mathcal{2}\mathcal{2}\mathcal{2}}

As was done in the previous section, one can explicitly see, that for this choice, the structure constants of the stretched position $(\ell_3=\ell_1+\ell_2-1)$ are independent of N, a fact whose significance appears in the next section.

III. AN UNDERLYING MATHEMATICAL REASON FOR THE ABOVE CONSTRUCTION

Having found explicitly the correspondence between the Lie algebra of area preserving transformations and an N-dimensional representation of SU(N) (N+ ∞) by constructing a basis (the $T_{\ell m}$) as polynomials in the S_i (a basis of the N-dimensional representation of SO(3)) one might wonder whether there is not an underlying mathematical reason for this construction to work. This would provide some additional understanding and also possibly lead to generalizations. In particular, most statements would be independent of a particular representation.

It turns out that it is the space of the $Y_{\ell m}$'s with $\{,\}$ and the role of the abstract Lie algebra $\underline{SO(3)}$ which have a natural generalization, while $\underline{SU(N)}$ arises as the space in which N-dimensional unitary representation of $\underline{SO(3)}$ lie. (In this sense $\underline{SO(3)}$ is special, as for a general Lie algebra there will not be exactly one irreducible inequivalent representation for each N. Also there will be in general more than one Casimir operator that when going to an N-dimensional representation will carry the N-dependence.)

Let \underline{G} be a Lie algebra over the complex numbers, whose adjoint representation is completely reducible, and G be the adjoint group.* Let $x_1 ldots x_n$ be a basis of G. The enveloping

^{*}Note: G will not correspond to the group of area preserving reparametrizations of S^2 (which was called G in BI), but rather to SO(3).

algebra U(G), which is defined in rather abstract terms *, can be taken** to be the tensor algebra $\tau(G)$ (i.e., the space of all polynomials $a_1 \cdots a_m a_1 \cdots a_$

is therefore a basis of $U(\underline{G})$, and $u \in U$ will be written as

$$\sum a_{j_1\cdots j_m}^{(u)} \times_{i_1}^{j_1} \cdots \times_{i_m}^{j_m}$$

Define U as the space of all u & U with degree

and the $U_{\mathcal{L}}$ are called a filtration of U. There is a natural Poisson bracket defined on U: [U,U']=uu'-u'u; then $[U_k,U_k]\subseteq U_{k+\ell-1}$.

The symmetric algebra S(G) is defined as the space of all polynomials in n commuting objects $x_1 \cdots x_n$. This space also has

$$\left\{ \times_{i}^{j_{1}} : \ldots \times_{m}^{j_{m}} \mid j_{\gamma} \geqslant 0, \Sigma j_{\gamma} > 0 \right\}$$

^{*}See e.g. "Lie Algebras" by Jacobson (Interscience, 1962).

^{**}Poincaré-Birkhoff-Witt theorem, see *.

as a basis, but xx'-x'x=0 in S(G). S(G) can (and will from now on) be regarded as the space of pol nomial functions f on the dual space $G' = \mathbb{R}^n$, (then s $\mathfrak E$ S is a polynomial in n real variables, with complex coefficients). Let S_kCS be the set of all homogeneous polynomials of degree=k. One can define a Poisson bracket $\{\}$ on S, with $\{S_k, S_k\} \in S_{k+\ell-1}$, by defining the following surjective homomorphism $\{S_k, S_k\} \in S_{k+\ell-1}$, by (which has $S_k \in S_k$) as kernel):

$$U_{k} = \sum_{\sum i_{s} \leq k} a_{i_{1} \dots i_{m}}^{u} \times_{i_{1} \dots i_{m}}^{i_{1} \dots i_{m}} \times_{i_{1} \dots i_{m}}^{i_{1} \dots i_{m}}^{i_{1} \dots i_{m}} \times_{i_{1} \dots i_{m}}^{i_{1} \dots i_{m}}^{i_{1} \dots i$$

where the u_k are some elements of U_k with k_i (u_k) = k_i . {} is well defined, as u_k is ambiguous only in the terms of degree k_i , so that $[u_{k_1}, u_{k_2}]$ is some $u_{k_1} + k_2 - 1$, with an ambiguity only in the terms of degree $k_1 + k_2 - 1$, which makes $k_1 + k_2 - 1$ ($u_{k_1} + k_2 - 1$) $k_1 + k_2 - 1$ unambiguous (uniquely defined). The so defined Poisson bracket $\{f,g\}$ of two polynomial functions $f,g \in S(G)$ is in fact equal to

where c_{jk}^i are the (not necessarily totally antisymmetric) structure constants * of \underline{G} . One can verify explicitly that (C9) defines a Poisson bracket, i.e.,

and

$$\sum_{r} \left\{ \left\{ f_{r}, f_{r2} \right\}, f_{r3} \right\}$$

$$(r cyclic permulation of $(l_{12,3})$) = $C_{jk} \cap C_{jk} \cap C_{st} \times C_{st} \cap C_{st} \times C_{st} \cap C_{st} \times C_{st} \cap C_{s$$$

the first term = $\chi(C_{st}^{j}C_{jk}^{j}+C_{ks}^{j}C_{jt}^{j}+C_{ks}^{j}C_{js}^{j})$ of $\chi(C_{st}^{j}C_{jk}^{j}+C_{ks}^{j}C_{jt}^{j}+C_{ks}^{j}C_{js}^{j})$ one then sees that the second and third term cancel, $\chi(C_{st}^{j}C_{jk}^{j})$ using only $c_{jk}^{i}=-c_{kj}^{i}$.

^{*}in the basis $x_1...x_n$

Following B. Kostant*, one can characterize the structure of U and S and the relation between them in the following way:

1. S=J BH (every element of S can be written as

where J is defined as the space of all polynomials invariant under the group action [which is induced by the adjoint action of G on G, for matrices:

XE E -> g'xg E E]

and H = the set of all G-harmonic polynomials, i.e., all

fe S such that 7 f=0 for every homogeneous differential operator

with constant coefficients, that commutes with the group

action. In our case:

the only such ∂ is ∇^2 (and functions of ∇^2), $H \equiv \text{Space of harmonic polynomials in the usual sense}(\nabla^2 h=0)$.

Only nonconstant $f(x_1x_2x_3)$ can be written as $\sum_{k=1}^{\infty} j_k(\tau) \left(\sum_{m=1}^{\infty} Q_{km} \tau^k Y_{km}(Q_1 Y_1)\right)$, separation of variables.

2. Let O_X denote the G-orbit in G of $x \in G$, and let $S(O_X)$ be the ring of all functions on O_X defined by restricting S to O_X ; let \underline{r} be the rank of \underline{G} ; then $\widehat{O_X} \subseteq M - \underline{f}$ and

^{*&}quot;Lie group representations on polynomial rings", Am. J.M. 85, 1963, p. 327-404. I would like to thank Prof. Kostant, Alex Uribe and Robin Ticciati very much for several discussions and much patience. This Section (BIII) would not exist without their ideas and help.

for every $x \in G$ such that $\dim O_{\mathbf{x}} = \mathbf{n} - \mathbf{r}$, H and $S(O_{\mathbf{x}})$ are isomorphic as G-modules [a G-module is a vector space V together with a map $G \to GL(V)$, $g_1(g_2V) = (g_1g_2)V$]. For G = SO(3), $\dim O_{\mathbf{x}} = 3 - 1 = 2V_{\mathbf{x} \in G}$.

3. U=Z @ E where Z \equiv Center of U (i.e., all Z with $[\mathbf{Z},\mathbf{u}]=0$ $\forall_{\mathbf{u}\in U}$ and E \equiv space spanned by all powers $\mathbf{x}^{\mathbf{k}}$, for all nilpotent elements $\mathbf{x}\in G$ ($\mathbf{x}\in G$ is called nilpotent if $(ad\mathbf{x})^M=0$ for some M, where adx is the adjoint representation of x, which is a nxn matrix) [for G=SO(3):

$$2 \equiv \text{all polarounials in } X \equiv X_1^2 + X_2^2 + X_3^2$$

$$ad X_1 \equiv S_1 \quad \text{with } (S_1)_{jk} = -i \in_{ijk},$$

$$A = a_1 S_1 = -i \begin{pmatrix} 0 & a_3 & -a_2 \\ -a_3 & 0 & a_4 \\ a_2 & -a_1 & 0 \end{pmatrix}$$

satisfies

$$A^{3} + A \left(a_{1}^{2} + a_{1}^{2} + a_{3}^{2} \right) = 0$$

$$= \overline{a}^{2}$$

as

$$\det (A - \lambda \mathbf{1}) = -\left(\lambda^3 + \lambda \vec{a}^2\right)$$

has to vanish for $\lambda = A$. Therefore $A^3 = 0$ for $\frac{\lambda^2}{a^2} = 0$ and one has:

x nilpotent \rightleftharpoons $\overset{1}{a}^2=0$ \rightleftharpoons $\mathscr X$ totally traceless. As $\mathscr X$ is by definition a symmetric tensor, one has: $E(\underline{SO(3)}) \subset U(\underline{SO(3)})$ is the space of all u with $\mathscr X$ traceless (and symmetric).

 $E_k \leq U_k \cap E$ is (2k+1) dimensional.]

- 4. J @ H and Z @ E are isomorphic as G-modules.
- 5. Now look at the Poisson structures of S and U; using 1 and 3 one finds

$$\begin{bmatrix}
e_{km_k}, e_{lm_l} \\
= \sum_{i=1}^{k+l-1} \frac{d_{im}(E_i)}{d_{im_i}(X_i)} \\
= \sum_{i=1}^{k+l-1} \frac{d_{im}(E_i)}{d_{im_i}(X_i)} \\
= \sum_{i=1}^{l} \frac{d_{im_i}(E_i)}{d_{im_i}(X_i)} \\
= \sum_{i=1}^{l} \frac{d_{im_i}(E_i)}{d_{im_i}(E_i)} \\$$

where the

 $C_{jm_{j}} \left(| \leq m_{j} \leq \dim (E_{j}) \right)$ denote a basis of E_{j} and $d_{km_{k}\ell m_{k}} \left(\times_{\alpha} \right) \in Z$ is a polynomial in the independent Casimir operators X_{α}

$$\left\{ h_{k} m_{k}, h_{\ell m_{\ell}} \right\} = S_{k+\ell-1} = \sum_{i=1}^{k+\ell-1} \widetilde{\mathcal{A}}_{k}^{im_{i}} h_{k} m_{\ell} m_{\ell}$$

is a basis of H_j which one chooses to be the one given by the isomorphism between E and H and \widetilde{d} \widetilde

is just a set of complex numbers when restricting S(G) to $S(O_X)$. [The Y 's are a basis of $S(O_X)$ for G=SO(3) and |X|=1.] Because of the way $\{\}$ was defined via $\{\}$ in terms of [,] one has $A_{kmk}^{n=k+k-1}$, A_{kmk}^{n} = $A_{kmk}^{n=k+k-1}$, A_{kmk}^{n} [if one has chosen A_{kmk}^{n}] which is, of course, independent of A_{kmk}^{n}], A_{kk}^{n} anyway (just counting

is, of course, independent of \swarrow anyway (just counting powers) it is therefore also the same for all representations of U. Let the mapping \Vdash from U into the set of all complex NxN matrices be such a N-dimensional representation of U.

$$T_N(X_{\alpha})$$
 and $d_{km_{\alpha},l_{m_{\alpha}}}^{im_{\alpha}}(T_N(X_{\alpha}))$

then just become a set of numbers and $T_N(E)$ is a Lie algebra with structure constants d, that depend on N is $X_N^{\alpha} = T_N(X_N)$. The earlier proof that $\int_{-\infty}^{\infty} f(N) = G$ relied on the

The earlier proof that $\lim_{N\to\infty} \hat{f}^{(N)} = g$ relied on the

fact that for $\underline{SO(3)}$ there is only one independent Casimir operator \mathcal{K} ($\equiv S_1^2 + S_2^2 + S_3^2$), (exactly) one irreducible representation for each N, and $\mathbb{T}_{N}(\mathcal{K}) = \frac{N^2-1}{4}$ $\longrightarrow \infty$ as $N \to \infty$

C. THE NATURE OF THE SPECTRUM OF $H_{_{ m N}}$

I. Some general remarks

In Section B it was shown that the structure constants $\textbf{g}_{\alpha\beta\gamma}$ appearing in

$$\begin{aligned} & + \frac{1}{2} \int d\Omega \left(\rho_{x}^{2} + \rho_{y}^{2} + \left\{ x_{i} y \right\}^{2} \right) \\ & = \frac{1}{2} \sum_{\alpha = (\ell_{m})} \left(\vec{\rho}_{x} \vec{\rho}_{\alpha}^{*} + g_{\alpha s_{i}} g_{\alpha s_{i}} x_{i} x_{i} x_{i} x_{i} \right) \end{aligned}$$

are equal to the N+∞ limit of the SU(N) structure constants $f_{\alpha\beta\gamma}^{\,\,(N)}$. The Hamiltonian

involving only a finite number of degrees of freedom is, therefore, a good approximation to H as $N+\infty$. It is invariant under the finite group SU(N). Defining traceless hermitian $N\times N$ matrices $X=X_{\alpha}\hat{T}_{\alpha},Y=\ldots$, the above Hamiltonian becomes

$$\hat{C} \cdot \left\{ \frac{1}{2} \operatorname{Tr} \left(P_x^2 + P_y^2 - \left[X, Y \right]^2 \right) \right\}$$
where
$$\operatorname{Tr} \left(\hat{T}_x \hat{T}_x^{\dagger} \right) = \hat{C}^{-1} \int_{\mathcal{A}_x'} \left(\frac{N^3}{16\pi} + O(N^2) \right)$$

One is, of course, always free to change the relative strength of potential to kinetic energy by rescaling X and Y. (See B17 and B20).

One could have gone <u>directly</u> from the surface Hamiltonian H to the above matrix hamiltonian noticing that H depends only on the algebraic structure $\{,\}$ which is preserved when

replacing
$$\{X(\theta, Y), Y(\theta, Y)\}$$
 by $\frac{1}{2}[x, Y]$ (and $\int d\Omega \rightarrow \hat{C} T_{r}$).

Note that, as already mentioned in the introduction, this transition has nothing to do with the transition from the classical surface Hamiltonian to the quantum theory, although the 1/i formally comes from the extra i in $[S_i, S_j] = i \in S_k$ compared to $\{X_i, X_j\}_p = C_{ijk} \times_k$ (compare fax $\{X_j, X_j\}_p = C_{ijk} \times_k$) In order to obtain a sensible $\{X_j, X_j\}_p = C_{ijk} \times_k$

In order to obtain a sensible N-plimit one rescales X and Y by $N^{1/6}$, absorbs the overall factor $CN^{1/3}$ in the surface tension T_0 and defines the SU(N) invariant Hamiltonian

$$H_{N} = \frac{1}{2} \operatorname{Tr} \left(P_{x}^{2} + P_{y}^{2} - \frac{1}{N} \left[X_{i} Y_{j}^{2} \right]^{2} \right)$$
(c1)

From what is known about large N-matrix models in general,* H_N will have a groundstate with energy of $O(N^2)$ (which one subtracts) and the level spacing of the excited states will be of O(1).

From now on the matrices X,Y ... are most conveniently expanded in hermitian orthonormal generators $\neg a$ $(\neg e \times x_a \neg a) = (\neg T_a \neg T_c) = (\neg a e)$ with real coefficients. With $[T_a,T_b] = if_{abc}T_c$ one then has, e.g., for the potential

$$V = \frac{1}{2N} \text{ fall fade } \times_{e} Y_{c} \times_{d} Y_{e} , \quad (C2)$$
for SU(N=2) this is $V = \frac{1}{4} (\overrightarrow{x} \times \overrightarrow{y})^{2}$

^{*}Following the work of Brezin, Itzykson, Parisi, Zuber, Communications in mathematical physics 59, p. 35-51 (1978).

The generators of SU(N) symmetry transformations are $K_{\alpha} = \frac{1}{2} T_{\alpha} \left(T_{\alpha} \left(T_{\alpha} P_{\alpha} T_{\alpha} + T_{\alpha} P_{\alpha} T_{\alpha} \right) \right)$

and one is interested in $K_a=0$ (classically), $K_a \mid \psi \rangle = 0$ (for the quantum theory). [These constraints, unfortunately, exclude the class of solutions $X+iY=e^{i\omega t} \omega(S_X+iS_Y) \sqrt{N}$ which solve the classical equations of motion derived from (C1):

$$\ddot{X} = \frac{1}{N} [Y [X, Y]], \ddot{Y} = \frac{1}{N} [X, [Y, X]]$$
 (C3)

The S_i (i=1,2,3) denote 3N×N matrices satisfying $[S_i,S_j]=iE_{ijk}S_k$ $K_a = -2\omega^3 N T_r (T_a S_2) \neq 0$

One can further see that, at least for SU(N=2) that these solutions are unstable against small perturbations. Note that one can rewrite (C3) in the slightly more compact form:

Although $V \geqslant O(as A=[x,y]=-A^+, V=+\frac{g^2}{4N}Tr(A^+A)\geqslant 0)$ one might wonder whether the potential C2 confines* or not, as V=0 for a rather large subspace of configuration space (for fixed X, all matrices Y that commute with X). The simplest case, SU(N=2)

$$V_2 = \frac{1}{4} (\vec{x} \times \vec{y})^2$$

which is 0 for $x \parallel y$ (the classical partition function diverges as a result) The simplest quartic potential of type (C2) one could possibly think of is $V=x^2y^2$ (in fact, one is lead to something very similar for $O(2) \times O(3)$ singlet states of V_2) which will be looked at in the next section. As the answer there is that V confines, one is led to believe that the spectrum of H_N is discrete.

^{*}i.e.has a purely discrete spectrum

The x^2y^2 -problem and the "B.O." approximation

We consider the spectrum of $H = \int_{x}^{2} + \int_{y}^{2} + \chi^{2}y^{2}$ Although there is a short mathematical proof* that the spectrum of H is discrete**

[H>H' = $\frac{1}{2}(p_x^2+p_y^2+|x|+|y|)$; spectrum of H' discrete \Longrightarrow spectrum of H discrete] it might be worth looking at the problem in the following way. As the question of binding should not have much to do with the shape of the potential in a finite region, assume $V=\infty$ for $x \leq \bigwedge$, $\bigwedge \gg 1$ and try to solve the problem

$$\left(-\left(\partial_{x}^{2}+\partial_{y}^{2}\right)+\chi^{2}\right)\psi(x|y)=E\psi f. \times \lambda \qquad (C5)$$
Changing variables to \(\frac{1}{2}\) and \(\frac{1}{2}\)

Changing variables to $\{ \}$ 0 and $\{ \}$ by writing $x = \{ \} + \{ \}$, $\{ \} = \{ \}$

$$\begin{aligned}
&H \widetilde{\psi}(\xi, \gamma) = E \widetilde{\psi}, \ \widetilde{\psi} = 0 \quad \text{an} \quad \xi = 0 \\
&H = \Lambda \left\{ -\frac{1}{\Lambda} \partial_{\xi}^{2} + \left(-\partial_{\gamma}^{2} + \widetilde{V}(\xi, \gamma) \right) \right\} \\
&\sim 0
\end{aligned}$$

$$\widehat{V}(\xi, \gamma) = (1 + \xi_{\Lambda})^2 \gamma^2 \equiv \omega^2(\xi) \cdot \gamma^2$$
Now one first solves the α redoperate.

$$(-\partial_{\chi}^{2} + \omega^{2}(\xi) \chi^{2}) \psi(\chi) = E \psi$$
gives

$$E = E_m(\zeta) = 2 (m - \frac{1}{2}) \omega(\zeta), m = 1, 2, ...$$
the same sense as Born and C.

In the same sense as Born and Oppenheimer treated the electron energy (calculated as a function of the nuclei distance) as a potential for the two nuclei, $E_m(\)$ will now be treated as a potential for the } -coordinate, i.e., for given m solve for the eigenvalues and eigenstates of

*pointed out by Barry Simon in private communication **as 2 = (e-x2) dxdy ~ 8 dx with R.Jackiw. still diverges (logarithmically) one could say that the discreteness is due to the uncertainty

$$H(m) = (2m-1) \wedge + (2m-1)^{2/3} (-\partial_u^2 + u)$$

$$u = (2m-1)^{1/3} \xi$$
Calling the eigenvalues of $(-\partial_u^2 + u)$, as before,
$$H = -(\partial_x^2 + \partial_y^2) + \chi^2 y^2$$
will therefore (within the Born Oppenheimer approximation) have the eigenvalues

$$E_{3}^{m} = (2m-1)\Lambda + (2m-1)^{43}E_{3}$$

One can show quite generally that the Born-Oppenheimer approximation gives a lower bound for the true ground state energy [so that $E_{B.O.} \leq \text{true } E_O \leq E_{\text{var}}^O$]. Proof: Consider a general Hamiltonian $H = H(p,q;p',q') = p^2 + H(q;q',p')$ where p' and q' are abbreviating all degrees of freedom different from q and p. Define $H_{B.O.}$ to be the Hamiltonian obtained from H by replacing H(q;q',p') for fixed q by its eigenvalues $E_m(q)$, i.e., $H_{B.O.} = p^2 + E_m(q)$. Using $(\bigvee_{(q)}, \bigvee_{(q)})$ as an abbreviation for integrating $\bigvee_{(q)} (q,q')$ only over q'-coordinates, one has $E_O(q) \leq (\bigvee_{(q)}, \bigvee_{(q)})$ by the variational principle and, therefore, for all I_{I_I} :

$$\langle \psi | H | \psi \rangle = \int dq (\psi(q), H \psi(q))$$

 $= \int dq \left\{ \left(\frac{3\psi}{5q}, \frac{3\psi}{5q} \right) + \left(\frac{\psi(q), H(q)}{(\psi(q), \psi(q))} (\psi(q), \psi(q)) \right\} \right\}$
 $\geq \int dq \left(\psi(q), (\rho^2 + E_0(q)) \psi(q) \right) = \mathcal{L} \psi | H_{8.0.} | \psi \rangle_{q.e.d.}$

For our case one can do an Explicit calculation and comparison of $E_{B.O.}$ and E_{var} :

a) taking $e^{-1/2\omega(x^2+y^2)}$ as trial wave function and minimizing with respect to ω gives $E_{Ver} = 2\left(3\sqrt{\frac{1}{16}} + \sqrt[3]{\frac{1}{128}}\right) \approx 1.2$ b) $E_{3o} = \text{lowest eigenvalue of } (-3x + |x|). \text{ One therefore has to find the smallest E for which}$ $\left\{f(t) - t f(t) = 0 : t = (|x| - E) \in [-E, +\infty), f(+\infty) = 0, f(-E) = 0\right\}$ has a solution. For z > 0 one takes $f(z) = \sqrt[3]{3\sqrt{2}} \cdot \left(\frac{1}{3} \cdot \frac{1}{3} \cdot \frac{1}{3} \cdot \frac{1}{3}\right) \in \mathbb{R}$ and by analytical continuation (H⁽¹⁾) and J defined as in

Jahnke Emde.) $\frac{e^{-1/2\omega(x^2+y^2)}}{e^{-1/2}} = \frac{2}{3\sin\frac{2\pi}{3}} = \int_{-2\sqrt{3}}^{2} \left(\frac{2}{3} \cdot \frac{3}{2}\right) - \int_{-2\sqrt{3}}^{2} \left(\frac{2}{3} \cdot \frac{3}{2}\right) = 0$ $= \int_{-2\sqrt{3}}^{2} \left(\frac{2}{3} \cdot \frac{3}{2}\right) - \int_{-2\sqrt{3}}^{2} \left(\frac{2}{3} \cdot \frac{3}{2}\right) = 0$ $= \int_{-2\sqrt{3}}^{2} \left(\frac{2}{3} \cdot \frac{3}{2}\right) - \int_{-2\sqrt{3}}^{2} \left(\frac{2}{3} \cdot \frac{3}{2}\right) = 0$ $= \int_{-2\sqrt{3}}^{2} \left(\frac{2}{3} \cdot \frac{3}{2}\right) - \int_{-2\sqrt{3}}^{2} \left(\frac{2}{3} \cdot \frac{3}{2}\right) = 0$ $= \int_{-2\sqrt{3}}^{2} \left(\frac{2}{3} \cdot \frac{3}{2}\right) - \int_{-2\sqrt{3}}^{2} \left(\frac{2}{3} \cdot \frac{3}{2}\right) = 0$ $= \int_{-2\sqrt{3}}^{2} \left(\frac{2}{3} \cdot \frac{3}{2}\right) - \int_{-2\sqrt{3}}^{2} \left(\frac{2}{3} \cdot \frac{3}{2}\right) = 0$ $= \int_{-2\sqrt{3}}^{2} \left(\frac{2}{3} \cdot \frac{3}{2}\right) - \int_{-2\sqrt{3}}^{2} \left(\frac{2}{3} \cdot \frac{3}{2}\right) = 0$

III. calculating
$$2(g^2) = \int dX dY e^{-\frac{1}{2}Tr(X^2+Y^2-\frac{g^2}{2N}(X_1Y_1)^2)}$$

Although it does not provide any information about the spectrum of H_N , the integral $Z(g^2) = (dX dYe - Hg)$ will be calculated below, where

This integral is interesting in its own right as, at least to the best of my knoledge, integrals of this type (i.e a two-matrix-model with coupled quartic interaction) have not been calculated so far in the literature, while the one-matrix-model with quartic selfinteraction, and the multi-matrix-problem with quartic self- but only quadratic nearest neighbour interactions have been solved*.

In the case at hand, one first integrates over all but N of the original 2N2 (real) variables explicitely (arriving at (C7)). The

where W is of $O(N^{\frac{1}{2}})$. Therefore, as $N \rightarrow \infty$, the integral will be $\approx e^{-WC\{X\}} \qquad \text{where } \{\overline{\lambda}_i\} \qquad \text{minimizes W.By defining the density } \mathcal{U}(\lambda) \equiv \frac{1}{N} \sum_{i=1}^{N} \delta(\lambda - \lambda_i) \qquad \text{the problem of minimizing W becomes}$ ~ -WC(I) that of solving a singular integral equation for $u(\lambda)$ (see C 40) One can do so, but instead of calculating (e.g.) the first moment of u (i.e. $\int \lambda^2 u(\lambda) d\lambda$) as a function of g , we are only able to explicitely calculate it as a function of a parameter b, where b is given as a function of g via an implicit equation involving complete elliptic integrals(see C | 8iii). The formula for $\langle q^{(1)} | Y | 2 \rangle$, (the expectation value of the potential) is given in terms of $\int \lambda^2 u(\lambda) d\lambda$ (see C ||).

** Brezin et al | and Mehta et al ,J.Phys.A.Math.Gen. |4(1981)

$$\widetilde{Z} = \int dX dY e^{-\frac{1}{4}g} \\
= \int dX dY e^{-\frac{1}{4}\chi^{2} + \frac{1}{2}Y^{2} - \frac{9}{4}\chi [X_{1}Y_{2}]^{2}} \\
= C \int T dX_{1} dY e^{-\frac{1}{2}\sum_{i,j} |Y_{ij}|^{2} (1 + \frac{9}{2}\chi [X_{1} - X_{j}]^{2})}$$

where x_i are the eigenvalues of X and then, using X diagonal and $y^+=y$: $Tr[x,y] = 2Tr(xyxy-x^2y^2)$

$$=2\left(\sum_{i,j}x_{i}x_{j}|y_{ij}|^{2}-x_{i}^{2}|y_{ij}|^{2}\right)=-\sum_{i,j}(x_{i}-x_{j})^{2}|y_{ij}|^{2}$$

and writing the exponent as

 $-\frac{1}{2}\sum_{x_{i}}^{2}-\frac{1}{2}\sum_{x_{i}}^{N}\sum_{x_{i}}^{2}-\sum_{x_{i}}^{2}\left(\left(\Re\left(x_{i}\right)\right)^{2}+\left(\Im\left(x_{i}\right)^{2}\right)^{2}\left(1+\left(\Im\left(x_{i}\right)^{2}-X_{i}\right)^{2}\right)$ the integral $\begin{cases} dy \text{ is simply a product of gaussian integrals so} \\ \text{that (with } \lambda_{i}\equiv \frac{x_{i}}{\sqrt{2N}} \end{cases}$

$$\frac{2}{2} = C \int_{-\infty}^{+\infty} d\lambda_i e^{-N \sum_{j=1}^{N} \lambda_j^2} \frac{(\lambda_i - \lambda_j)^2}{1 + g^2(\lambda_i - \lambda_j)^2}$$

One can also calculate the integral in a more symmetrical way by introducing an auxiliary matrix ϕ —to get rid of the quartic interaction, then integrating over x and the λ i appearing in the above formula are then the eigenvalues of

9/6 N

$$\frac{(C6)}{2} = \frac{1}{2(16)} \int dX dY e^{-\frac{1}{2} \left\{ b^{-1}(x^{2} + Y^{2}) - \frac{g^{2}}{2N} b^{-1}(x^{2} + Y^{2}) \right\}} dX dY e^{-\frac{1}{2} \left\{ b^{-1}(x^{2} + Y^{2}) - \frac{g^{2}}{2N} b^{-1}(x^{2} + Y^{2}) \right\}} dX dY e^{-\frac{1}{2} \left\{ b^{-1}(x^{2} + Y^{2}) - \frac{g^{2}}{2N} b^{-1}(x^{2} + Y^{2}) \right\}} dX dX dY e^{-\frac{1}{2} \left\{ b^{-1}(x^{2} + Y^{2}) - \frac{g^{2}}{2N} b^{-1}(x^{2} + Y^{2}) - \frac{g^{2}}{2N} b^{-1}(x^{2} + Y^{2}) \right\}} dX dX dY e^{-\frac{1}{2} \left\{ b^{-1}(x^{2} + Y^{2}) - \frac{g^{2}}{2N} b^{-1}(x^{2} + Y^{2}) - \frac$$

W is of $O(N^2)$, so that Z can be computed, in the large N-limit, by minimizing W with respect to the λ ;

$$0 = \frac{\partial W}{\partial \lambda^{2}} = 2N\lambda_{t} - 2\sum_{s+t}^{2} \frac{1}{1+g^{2}(\lambda_{t} - \lambda_{s})^{2}} + 2\sum_{s+t}^{2} \frac{g^{2}(\lambda_{t} - \lambda_{s})}{1+g^{2}(\lambda_{t} - \lambda_{s})^{2}}$$
Introducing the size of t

Introducing the eigenvalue density $u(\lambda) = \frac{1}{\lambda_1} \sum_{k=1}^{N} \delta(\lambda - \lambda_{r})$

the above equation can be written as

$$\lambda = \int_{-a}^{+a} \frac{u(\mu)}{\lambda^{-\mu}} d\mu - \int_{-a}^{+a} \frac{u(\mu)(\lambda^{-\mu})}{\sqrt{2}^{2} + (\lambda^{-\mu})^{2}} d\mu$$
 (C10)

which is a singular integral equation for $u(\lambda)$, which has to be solved subject to the constraint $\int_{1}^{\infty} u(\lambda) d\lambda = \int_{1}^{\infty} Before$ outlining how to solve equation (Cl0) note how one can, e.g., determine $\langle v \rangle$ once $u(\lambda)$ is known:

$$\langle V \rangle = \langle -g^{2}_{4N} tr[X,Y]^{2} \rangle = -g^{2} \frac{\partial z}{\partial g^{2}} = \langle g^{2} \frac{\partial W}{\partial g^{2}} \rangle$$

$$= + \sum_{T \in S} g^{2} \frac{(\lambda_{T} - \lambda_{S})^{2}}{1 + g^{2}(\lambda_{T} - \lambda_{S})^{2}} \qquad (\lambda_{T} Satisfying C9)$$

$$= N^{2} \left\{ \frac{1}{2} - \int_{-a}^{b} u(\lambda) \lambda^{2} d\lambda - \int_{2N}^{b} duicard \right\} - \cdots (C11)$$

The last step could be made because $\frac{\partial W}{\partial \lambda} = 0$

$$0 = \sum_{k} \frac{2k}{2k}$$

$$= 2N \sum_{k} \frac{2}{k} - (N^{2}N) + 2 \sum_{k \in S} \frac{2(\lambda_{k} - \lambda_{k})^{2}}{1 + q^{2}(\lambda_{k} - \lambda_{k})^{2}} + a$$
solution of (C10), first for g=0: Defining $F(z) = \int_{z-\lambda}^{u(\lambda) d\lambda}$

which is real for real $Z \notin [-a,+a]$, behaves like 1/Z for $|z| \rightarrow \infty$ is analytic in the complex Z-plane except for a cut along

[-a,+a], and--approaching the cut from above -

$$\lim_{\epsilon \to 0} \operatorname{Re} F(\lambda \pm n\epsilon) = \frac{1}{2} \lim_{n \to \infty} \int_{-n}^{+n\epsilon} u(n) dn = \frac{1}{\lambda - n - n\epsilon} = \int_{-n}^{+n\epsilon} u(n) dn = \lambda$$

while
$$\lim_{E\to 0} \operatorname{Jm} F(\lambda \pm n E) = \frac{1}{2n} \lim_{E\to 0} \operatorname{Jm} U(\mu) \left(\frac{1}{\lambda_{jn} \pm n E} - \frac{1}{\lambda_{jn} \mp n E}\right)$$

$$= \frac{1}{2n} \operatorname{Jm} U(\mu) \left(\mp 2n \operatorname{T} S(\lambda_{jn})\right) = \mp \operatorname{Tr} U(\lambda)$$

the (unique) function having these properties is

$$F(z) = 2 - \sqrt{z^2 - 2}$$

as is easy to see and even simpler to check:

$$F(2) = \frac{2}{a^2} \left(2 - \sqrt{2^2 - a^2}\right)$$
satisfies the first 3 criteria, while $\lim_{\epsilon \to 0} \Re(\lambda \pm i\epsilon) = \frac{2\lambda}{a^2} = \lambda$
gives $a = \sqrt{2}$. Then one calculates $u(\lambda)$ as

$$U(\lambda) = \frac{1}{\pi} \lim_{\epsilon \to 0} \int_{0}^{\infty} F(\lambda + i\epsilon) = \frac{1}{2\pi i} \lim_{\epsilon \to 0} (-2\sqrt{(\lambda + i\epsilon)^{2}-2'})$$

$$= + \sqrt{2-\lambda^{2}}$$
(C12)

As a check one can calculate

$$\int_{-a}^{+a} u(\lambda) d\lambda = \frac{1}{\pi} \int_{-R}^{+\infty} \sqrt{2-\lambda^2} d\lambda = \frac{4}{\pi} \int_{0}^{T/L} \cos^2\theta d\theta = 1$$

as it must be. Also, according to the general formula (C11) for $\langle v \rangle$, $(\lambda^2 \mu i \lambda) \delta \lambda$ has to be +1/2 for g=0, so that $\langle v \rangle = 0$;

indeed:
$$\frac{1}{T} \int_{-12}^{12} \sqrt{2-\lambda^2} \lambda^2 d\lambda = \frac{(2)^2 2 \cdot 2}{T} \int_{0}^{11/2} \sin^2 \theta d\theta = \frac{8}{T} \frac{11}{2} \frac{1}{3} = \frac{1}{2}$$

For $g\neq 0$, define $G(Z) \equiv -i\left(F(2+i/2g)-F(2-i/2g)\right)$ $= -\frac{1}{g}\int_{-a}^{a} \frac{u(t)dt}{(2-t)^2+i/4g^2} \quad \text{which}$ assuming $u(\lambda)=u(-\lambda)$ behaves like $\frac{1}{g^{2^2}}-\frac{(3(\lambda^2u(\lambda)d\lambda-\frac{1}{4g^2})+O(\frac{1}{2}6))}{g^{2^2}}$

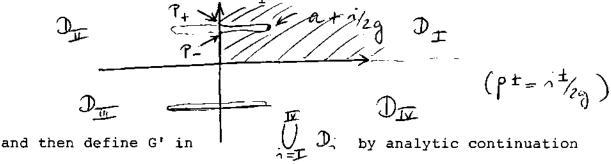
and has, because of (ClO), the property:

$$\int_{m} G(\lambda \pm \sqrt{2}g) = \pm \lambda \int_{a} A \in Ga, +a$$
(irrespective of approach from above or below).

 $G'_{=}-gz^{2}+G$ this translates to:

 $\int_{-\infty}^{\infty} G' = 0 \quad \text{for } 2 = \lambda \pm \frac{1}{2}g \quad \lambda \in [-\alpha, +\alpha]$ (from above or below) abor: G(z) = G(-z), G(z) = G(z) (and real) (C13) and at ∞ : $G(z) \approx -g z^2 + 8/z^2 + 6/z^4 + \cdots$ (where γ necessarily = -1/g and $(1^2 n(\lambda) d\lambda = \frac{1}{3} (\frac{1}{49})$ so that the knowledge of δ will yield $\langle V \rangle$ via (C10').

In order to find such a function G', analytic everywhere except at the two cuts [-a,+a]±i/2g, think of G' as being first only defined in the domain D, shown in the figure below:



(Cl3ii) and (iii) gives which, using

for 2eDIV: G'(2) = G'*(2*)

for $Z \in \mathcal{D}_{\pi}$: G'(z) = G'(-z)

for 2 G DI G'(2) = G/* (-2*)

This shows that, in fact, ImG' vanishes on the entire boundary of Dr. Therefore G'(z) can in fact be taken to be, up to real constants, the conformal transformation $(z \rightarrow \ \)$ mapping D_T onto the upper half plane. This transformation ((z)), mapping $P \rightarrow -1$, $a+\frac{1}{2}$ -C, P \rightarrow -b<-C<0, real z into real $\frac{1}{2}$, is given implicitly by the equation(s)*

 $2 = A \cdot \int_0^3 \frac{(t+c)dt}{\sqrt{f(t+l)(t+e)}}$

^{*}See, e.g., Fuchs and Shabat "Functions of a complex variable", Vol. 1, Problem 9 in Ch. 8, but note the mistakes in the last two lines before Problem 10.

$$\frac{1}{2c_{3}} \stackrel{(a)}{=} A \int \frac{(c-g)dg}{\sqrt{(l-g)(l-g)g}}, \quad \alpha = A \int \frac{(c-g)dg}{\sqrt{(l-g)(g-l)g}}$$

$$\alpha = \int \frac{(g-c)dg}{\sqrt{(l-g)(g-l)g}}$$
Although it would be nice (and Simplify to know 3 (z), which

Although it would be nice (and Simply to know $\int (z)$, which is an elliptic function in closed form, e.g., expressed in terms of the Weierstraß function P(z) of the same periods (and, possibly, $\int (z)$, one can calculate $\int (z)$, the coefficient of 1/z in $\int (z)$ as $z \to \infty$ also directly and therefore give a formula for (z) (which, however, will be very complicated and not much less implicit than (Cl4), as the g-dependence of z

and b can only be given implicitly). From C14:

$$\frac{2}{A} = \int_{0}^{3} \left(\frac{t+c}{\sqrt{t(t+1)(t+b)}} - \frac{1}{\sqrt{t}}\right) dt + 2\sqrt{3}$$

$$= 2\sqrt{3} + \int_{0}^{\infty} \left(\frac{t+c}{\sqrt{t(t+1)(t+b)}} - \frac{1}{\sqrt{t}}\right) dt - \int_{0}^{\infty} \left(\frac{t+c}{\sqrt{-1}}, -\frac{1}{\sqrt{t}}\right) (c15)$$

The third term will be expanded, the second term $(\equiv \mathcal{I})$ can be shown to be =0, using (Cl4) (i)-(iii):

$$\mathcal{Z} = 2 \int_{(x = \sqrt{\epsilon})}^{\infty} dx \left(\frac{x^2 + c}{\sqrt{(x^2 + 1)(x^2 + b)}} - 1 \right) = \int_{-\infty}^{+\infty} dx$$

As the integrand \int behaves like $1/x^2$ at ∞ , one can close the contour (at ∞) without altering \mathcal{L} , and as \int is analytic in \mathcal{L}_2 upper half-plane except for a cut between i and i.b., alter the contour to the closed path \int shown below.

erengi.

The dx=0, and therefore, with
$$g = -x^2$$

$$Z = 2 \int \frac{c-g}{\sqrt{g(g-1)(b-g)}} dg$$

which is 0 because of (Cl4ii) + (iii) (added together). The constant term in the expansion of z/2A for large ζ is

therefore 0, and
$$\frac{2}{2A} = \sqrt{3} + \sqrt{3} + \frac{E}{3^{2}} + \frac{F}{5^{2}} + \cdots$$

(higher order terms will not be needed to calculate δ). From (C15) the coefficients D,E and F can be calculated :

$$D = \left(\frac{b+1}{2} - C\right) \ge 0 \qquad E = -\left(\frac{b^2+1}{8} - \frac{C(b+1)}{6} + \frac{b}{12}\right) \le 0$$

$$F = +\frac{1}{5} \left\{ \frac{5}{16} (k^3+1) - \frac{bc}{4} - \frac{3}{8} k^2 (c - \frac{b}{2}) - \frac{3}{8} (c - \frac{b}{2}) \right\}$$
from which

$$\frac{2^{2}}{4A^{2}} = \frac{3}{7} + \frac{3}{3} + \frac{3}{12} + \cdots$$

$$(\sqrt{2}D^{2} + 2E < 0, \beta = 2F + 6DE + 2D^{2}, 3 = 3 + 2D)^{(C16')}$$

and therefore

$$S = \frac{2^{2}}{4A^{2}} - \frac{4 \times A^{2}}{2^{2}} - \frac{16(x^{2}+3)A^{4}}{2^{4}} + \cdots$$
 (C1611)

so that
$$G'(3(t)) = -49A^{2}(3(t)+2D)$$

$$= -9t^{2} + \frac{16A^{4}\chi g}{2^{2}} + \frac{64gA^{4}(\chi^{2}+\beta)}{24} + \cdots$$
(C17)

will have the required behavior at ∞ . Also it must be that

$$\chi = 16A^{4} \times 3 = -\frac{1}{3} \tag{c17'}$$

and, extracting $\begin{cases} \text{as the coefficient of } 1/2^4 \text{ in (Cl7), one} \end{cases}$ has, using (Cl1):

$$\lim_{N\to\infty} \left\langle \frac{V}{N^2} \right\rangle = \frac{1}{2} - \frac{1}{3} \left(\frac{1}{4g^2} - 64g^2 A^6 (\alpha^2 + \beta) \right) (c1711)$$

The problem with the above formula(e) is that they are rather useless unless one can determine bc and A (a is not needed) as functions of g via (Cl4i)-(ii)--which seems to be very difficult. What one can do without much work, however, is to derive equations for C, A and g as functions of b: (Cl4ii)-

(iii) gives
$$C = C(k) = k \frac{\int_{0}^{\pi/2} \frac{1}{|k|^{2} + k^{2} + k^{2}$$

b)c-) ∞ (c<b+1), using the expansions of the complete elliptic integrals E(x) and K(x) for x+0 and x+1

$$K(x) = \begin{cases} T_{2}(1+x_{4}^{2}+\frac{9}{64}x_{4}^{4}+\cdots) & f(x) = 0 \\ \int_{0}^{\infty} \frac{4}{x_{4}^{2}} + \frac{x^{12}}{4} \int_{0}^{\infty} \frac{4}{x_{4}^{2}}(1+x_{4}^{2}+x_{4}^{2}) + \cdots & f(x) = 0 \end{cases}$$

$$E(x) = \begin{cases} T_{2}(1-x_{4}^{2}-\frac{3}{64}x_{4}^{4}-\cdots) & f(x) = 0 \\ 1+\frac{x^{12}}{2} \int_{0}^{\infty} \frac{4}{x_{4}^{2}}(1+x_{4}^{2}-\frac{1}{2}) + \cdots & f(x) = 0 \end{cases}$$

$$E(x) = \begin{cases} T_{2}(1+x_{4}^{2}-\frac{3}{64}x_{4}^{4}-\cdots) & f(x) = 0 \\ 1+\frac{x^{12}}{2} \int_{0}^{\infty} \frac{4}{x_{4}^{2}}(1+x_{4}^{2}-\frac{1}{2}) + \cdots & f(x) = 0 \end{cases}$$

SUMMARY

The Lorentz-invariant action and the transition to a Hamiltonian formalism are given for a closed M-dimensional surface moving in D-dimensional Minkowski space. The definition of the system, the use of light cone coordinates and much more is in close analogy to the theory of a massless relativistic string, although the important role which the group of volume preserving reparametrizations plays is new. For the case M=2, D=4 this group and in particular its Lie algebra are studied, and the latter can be shown to correspond in some sense to the limit of SU(N) (as $N \rightarrow \infty$). This fact is used to transform the surface Hamiltonian into a large N two-matrix hamiltonian with the quartic interaction [X,Y]², a problem formulated in a much more familiar language. However, we have been so far unable to find out much about the spectrum of this Hamiltonian, apart from being almost certainly purely discrete, and some hints that its levels are highly degenerate which is needed for the theory to be Lorentz invariant. We hope that the states of each energy level of $\mathbf{H}_{\mathbf{N}}$ could be arranged into multiplets of total spin S. As the "energy" is really the square of the restmass, the states would then be characterized by spin and mass, as they should in a relativistic theory.

PART TWO A TWO DIMENSIONAL BOUND STATE PROBLEM

INTRODUCTION

Attempts to relate field theories of the strong interactions, in particular QCD, to string models of hadrons lead one* to study the nonrelativistic system of N distinguishable particles of equal mass (labelled 1 to N) moving in two dimensions with an attractive $\{-$ function potential between particles r and r+1:

$$+ \prod_{N=1}^{W} = \frac{1}{2} \sum_{T=1}^{N} \overrightarrow{T}_{T}^{2} - (2\pi\lambda) \sum_{T} S^{(2)} (\overrightarrow{\xi} - \overrightarrow{\xi}_{T})$$

where the second sum runs from either 1 to N-1 ("open case") or 1 to N ("closed case", (N+1) \equiv (1)). Solving the two-body problem one encounters divergences which are regularized by introducing a cut-off Λ to the divergent integral(s) and choosing the coupling constant λ in a cut-off dependent way so to make the two-body binding energy Δ_2 of the bound state independent of Λ :

$$\Lambda^2 e^{-2/\lambda(\Lambda)} = \Lambda_2$$

(which one then sets =1). The question then is what happens to the N(>2)-body problem, with λ given by the above equation? While in 3 dimensions the spectrum of the 3-body problem will

not be

^{*}See e.g. C.B. Thorn, Phys. Rev. D19 (1979) 639 Thorn J.L. Gervais, A. Neveu, N.P. B163 (1980) 189.

bounded from below (when regularizing the 2-body problem in an analogous way) , the answer for D=2 seems to be that the open (closed) 3-body system has only one (two) bound state(s) at energy -2.5* (-16 and -1.5), and is free of any irregularities. One can conclude this by deriving an eigenvalue-integral equation that is equivalent to the Schrödinger equation for bound states (but no longer contains \ nor \(\lambda\)

How delicate a border case D=2 is (note that for D€2 no regularization is necessary at all) can be illustrated by looking at the δ -function as a limit of a short-range potential

$$V=\frac{S}{a^2}f(T_a), f(g)=0f.g\geq 1, a\rightarrow 0$$

One finds out how the choice of S=S(a), that will give one bound state at finite energy (-1, say) depends crucially on the dimension:

$$S = O(a^{2-2})f \cdot D < 2 \quad (\cong a f \cdot D = 1)$$

$$S \cong \frac{2}{|ma|} f \cdot D = 2$$

$$S \cong 2e \quad f \cdot D = 2 + e (e \ll 1)$$

$$S \cong T^{2} f \cdot D = 3$$

^{*}See also Bruch/Tjon, Phys. Rev. A19, No. 2 (79)p.425-432, I.V. Simenog (1980) "Regularisation of the zero range interaction limit in a one- and two-dimensional manyparticle problem."

D=2, looked at it this way, is more like D(2 as $\lim S(a)=0$ / $a \to 0$

$$f. D \le 2$$
, while $\lim_{a \to 0} S(a) = (aust \neq 0) f. D > 2$

For D>2 both kinetic and potential energy diverge (logarithmically for D=2 but with the kinetic energy contained in the classically allowed region r < a finite; as a negative power for D>2) and the total energy (-1) arises from a delicate cancellation between them.

For the general N-body problem one can, using the consistency relation for λ , again derive an integral equation that does not contain λ nor Λ and is equivalent to the Schrodinger equation for bound states.

In an earlier work* the following results were derived for the open case (they will only be stated here in the introduction)

The N-body system binds

$$\Delta_{\text{M+N}} \gg \Delta_{\text{M}} + \Delta_{\text{N}} + 1$$

and in a random phase approximation is found to have phonon like excitations that come arbitrarily close to the ground state energy as $N\rightarrow\infty$:

^{*}J.H. Master Thesis, MIT, 1980.

When this result is used in the hadron models, one obtains a relation between the slope of the Regge trajectories and the QCD perturbation theory scale parameter \bigwedge . E(N) will be the same for any short-range potential, while for an arbitrary interaction

$$\sum_{r=1}^{N-1} V\left(\frac{1}{2} - \frac{1}{2} \right), \quad \sqrt{3}$$

has to be replaced by $(-g_{XX}(0)^{-1/2})$ where $g_{XX}(w)$ is a response function for the corresponding two-body problem. A (diagrammatic) random phase approximation is used to obtain

$$E_o(N) \approx -1.4 N + (2.06) - \frac{TV_3}{12} \frac{1}{N} + O(\frac{1}{N^2})$$

as an approximation to the ground state energy, which should be compared with a second order perturbation theory result: $E_{0}\left(N\right)\approx-(1.3)\,N+1.6.$

See Thorn, which also contains some of the above mentioned results.

A. The two-body problem (exact solution)

In two dimensions the Hamiltoni n is

$$H_{2}^{ur} = \frac{1}{2} \left(\vec{T}_{1}^{2} + \vec{T}_{2}^{2} \right) - (2\pi\lambda) \delta^{(2)} \left(\vec{\xi}_{1} - \vec{\xi}_{2} \right)$$

As the potential depends only on the relative coordinate, the problem separates in the center of mass system:

$$H_{z}^{Ur} = \frac{1}{4} \overrightarrow{P}^{2} + (\overrightarrow{p}^{2} - (ET\lambda)) S^{(2)}(\overrightarrow{x})$$

where
$$\vec{X} = \vec{\xi}_1 - \vec{\xi}_2$$
, $\vec{p} = \frac{1}{2} (\vec{T}_1 - \vec{T}_2)$

and \overrightarrow{P} is the total momentum $\overrightarrow{T_l} + \overrightarrow{T_z}$. The problem is therefore reduced to finding the spectrum of

$$h = \vec{p}^2 - (2\pi\lambda) \delta^{(2)}(\vec{x})$$

The equation for a bound state is $h \mid B \rangle = -\Delta \mid B \rangle$.

Multiply by $\langle p \mid$ to get

$$\vec{p}^{2} \langle \vec{p} | \mathcal{B} \rangle - (2\pi\lambda) \langle \vec{p} | \delta^{(2)}(\vec{x}) | \mathcal{B} \rangle = -\Delta \langle \vec{p} | \mathcal{B} \rangle$$

insert a complete set of states, use $\langle \vec{p} | \vec{k}^{(2)} | \vec{p}' \rangle = 1$ and rearrange terms to get

$$(\vec{p}^2 + \Delta) \langle \vec{p} | \mathcal{B} \rangle = (2\pi\lambda) \int \frac{d\vec{p}'}{(2\pi)^2} \langle \vec{p}' | \mathcal{B} \rangle = const.$$

Therefore there is only one bound state \backslash B> of the two-body system (with binding energy $\Delta \equiv \Delta_2$)

$$\widetilde{V}_{\mathcal{B}}(\vec{p}) = \langle \vec{p} | \mathcal{B} \rangle = \frac{(\text{const.})}{(\vec{p}^2 + \Delta_2)} = \frac{\sqrt{4T\Delta_2}}{p^2 + \Delta_2}$$

$$\left(\text{demanding } \langle \mathcal{B} | \mathcal{B} \rangle = 1 \right)$$

Putting $\langle p | B \rangle$ back into the original equation gives the consistency relation for λ :

$$(2\pi\lambda)\int \frac{d^{2}\rho'}{(2\pi)^{2}} \frac{1}{\rho^{2}+\Lambda_{2}} = 1$$

The integral $\left(=\frac{1}{4\pi}\int_{0}^{\infty}\frac{dE}{E+\Delta_{2}}\right)$ diverges; introducing

a cutoff \bigwedge^2 it becomes equal to $\frac{1}{4\pi} \ln \left(\bigwedge^2 / \Delta_2 \right)$ and therefore

$$\Delta_{2} = \Lambda^{2} e^{-2/\lambda}$$

In order to have Δ_2 finite, λ has to go to 0 as $\Lambda \to \infty$. The parameter of this model problem is therefore not λ , but the two-body binding energy Δ_2 .* From now on all energies will be measured in units of Δ_2 , i.e., Δ_2 =1.

^{*}In a slightly more mathematical treatment Λ , would appear as the one real parameter of the class of self adjoint extensions of $h_0=p^2$. For a mathematically precise treatment of point interactions in general see: Albeverio, Fenstadt (cont.)

and the self consistency relation for λ is

$$(2\pi\lambda)\int_{0}^{\infty}\frac{d^{2}p}{(2\pi)^{2}}\frac{1}{p^{2}+1}=1$$
(A1)

Because $v(x) = \int_{0}^{2} (x)$ is a (special case of a) separable potential, the scattering problem $h(y) = C_{y}(y)$ can be solved exactly by using the Lippman Schwinger equation. One finds

$$\langle \vec{p} | \chi^{\pm} \rangle = 2\pi^2 \delta^{(2)} (\vec{p} - \vec{p}_{\delta}) - \frac{4\pi}{(p_{\delta}^2 - p_{\delta}^2 + i\epsilon)(\ln p_{\delta}^2 + i\pi)}$$

and Hoegh Krohn, "Singular Perturbation and Nonstandard Analysis", Transac. AMS, Vol. 252, August 1979, and for 3 dimensions, the good review article by G. Flamand, "Mathematical Theory of Nonrelativistic 2- and 3-particle Systems with Point Interactions" in Cargese lectures in Theor. Phys., Gordon and Breach, N.Y., 1967, Lurcat, ed.. Also I would like to thank Prof. T.T. Wu for interesting discussions.

from which

$$\langle \mathcal{B}|\vec{p}|\gamma\rangle = \langle 4\pi \frac{\vec{p}_r}{\rho_r^2 + 1}$$

$$\langle \mathcal{S}|\vec{p}|\gamma\rangle \quad \text{will not be used, } \langle \mathcal{B}|\vec{p}|\mathcal{B}\rangle \text{ is =0.}$$

$$1 = |\mathcal{B}\rangle\langle \mathcal{B}| + \int \frac{d^2p_r}{\partial r} 2|\gamma^{\pm}\rangle\langle \gamma^{\pm}|$$

The normalisations of position and momentum eigenstates and the definition of Fourier transformation are listed below:

$$\langle \vec{x} | \vec{p} \rangle = e^{\gamma \vec{p} \cdot \vec{x}}$$

$$\langle \vec{x} | \vec{x}' \rangle = \delta^{(2)}(\vec{x} - \vec{x}') , \langle \vec{p} | \vec{p}' \rangle = (2\pi)^2 \delta^{(2)}(\vec{p} - \vec{p}')$$

$$\int d^2x |\vec{x}\rangle \langle \vec{x}| = 1 = \int \frac{d^2p}{(2\pi)^2} |\vec{p}\rangle \langle \vec{p}|$$

$$\int (\vec{p}) = \langle \vec{p} | f \rangle = \int d^2x e^{-\gamma \vec{p} \cdot \vec{x}} f(\vec{x})$$

$$\int (\vec{x}) = \langle \vec{x} | f \rangle = \int \frac{d^2p}{(2\pi)^2} e^{+\gamma \vec{p} \cdot \vec{x}} f(\vec{p})$$

The δ -function as the limit of a short-range potential

$$V = \frac{S}{a^2} f(\sqrt{a}) \equiv \sqrt{a^2}$$

with S and f dimensionless, and f normalized to $\int_{V}^{\infty} |dg = -1$; i.e., f determines the shape of \overline{V} , S its strength. By defining a rescaled variable $g \equiv r/a$ one writes the two-body hamiltonian

$$l_{12} = -\nabla^{2} + V = -\nabla^{2} + \frac{S}{a^{2}} f(\sqrt[4]{a}) \text{ as}$$

$$k_z = \frac{1}{a^2} \left(-\nabla_{\!\!g}^2 + S f(g) \right) \equiv \frac{\overline{k_z}}{a^2}$$
 (A3)

For h_2 to have exactly one bound state at a given finite energy (- $\{ say \}$ as $a \rightarrow 0$, S has to be chosen appropriately as a function of a (and $\{ say \}$) so that were \bar{V} just binds (\bar{h}_2 with bound state at energy - $\{ sa_2 \rightarrow 0 \}$, as $a \rightarrow 0 \}$.

As the dimensionality of the problem turns out to be an interesting point, one defines the problem in 2+ ϵ dimensions $(-1 \le \epsilon \le +\epsilon)$ by writing down the Schrödinger equation for radially symmetric bound state wavefunctions $\forall (\epsilon)$ in 2+ ϵ dimensions:

$$\frac{1}{a^2} \left(\dot{\varphi}'' + \frac{1+\epsilon}{s} \dot{\varphi}' - S f(s) \dot{\varphi}(s) \right) = + S \dot{\varphi}(s) \qquad (A4)$$

from now on f(g) will be taken to be simply $-\theta(1-g)$. (A4) is, of course, solved by solving for the regions g < 1 and g > 1 (from now on referred to just as \langle and \rangle) and then matching function and logarithmic derivatives at g = 1 (giving a condition on g)

For r < a (A4) becomes

$$\psi'' + \frac{1+\epsilon}{s}\psi' + (s-sa^2)\psi(s) = 0$$
(A51)

with solution*

$$S^{-\epsilon_{l_2}} \overline{J_{\epsilon_{l_2}}}(\mathcal{H}_S)$$

assuming

for r>a (A4) becomes

$$\psi''(\tau) + \frac{1+\epsilon}{\tau} \psi'(\tau) - \delta \psi(\tau) = 0$$
(A52)

with solution*

^{*}Conventions used, in particular for Bissel functions J and K, are those of "Table of integrals series and products", GradSteyn and Ryzhik.

matching

$$\frac{\psi'}{\psi} \text{ at } \tau = \alpha : \frac{\mathcal{E}}{\alpha} \frac{J_{1+\epsilon_{12}}(\mathcal{E})}{J_{\epsilon_{12}}(\mathcal{E})} = \frac{\sqrt{\delta} K_{1+\epsilon_{12}}(\alpha \sqrt{\delta})}{K_{\epsilon_{12}}(\alpha \sqrt{\delta})}$$
(A6)

Requiring that (A6) has only $\delta = 1$ as a solution independent of $a \leftrightarrow 0$, which is equivalent to h having exactly one bound state with energy $-a^{\frac{1}{2}}$, one finds:

D=2:
$$S = O(a^{2-D})$$
 ($\cong a \mid D=1$)

 $D=2: S \cong \frac{2}{1 \ln a \cdot 1}$ (A7)

 $D = 2 + \varepsilon (o \leqslant \alpha \mid): S \cong 2\varepsilon$
 $D = 3: S \cong T_4^2 + 2a \cong T_4^2$

(where x=y for two functions of a means that x=y (1+h(a)) with $\lim_{\alpha \to 0} L(\alpha) = 0$) ---- (A8)

(A7) shows that for $D \le 2$, $\lim_{a \to 0} S(a) = 0$ while $\lim_{a \to 0} S(a) > 0$

for D>2; this is of interest as it suggests that--despite the
fact that

$$\int \frac{d^{3}p}{p^{2}+\Delta} diverges for D \geqslant 2$$

while for D(2 everything is finite--the binding in 2 dimensions is more like D(2 rather than D>2, and, therefore, a more regular phenomenon than for D>2.4n particular for D=3 where the spectrum of the corresponding 3-body problem is not bounded from below, both "Thomas"- and Efimov-effect are known to occur.*

It is interesting to calculate the expectation values of the potential, the kinetic energy contained in the inside region r(a and the outside region r)a, and where the wavefunction is concentrated. With \cong defined by (A8): __for D=2 take $\left(S\cong \frac{2}{|\mathcal{L}_{AB}|}\right)$:

$$\psi(\tau) = \frac{1}{2\pi} \begin{cases} K_{\bullet}(\tau) & \text{outside} \\ |\text{lua}| J_{\bullet}(\mathcal{E}_{\pi}) & \text{wiside} \end{cases} (49)$$

^{*}See e.g., L.H. Thomas "The Interaction between a neutron and a proton and the structure of H³", Phys. Rev. 47, 1935. Minlos and Faddeev [M F] "Comment on the problem of three particles with point interactions", Soviet Physics JETP Vol. 14, No. 6, 1962. S. Albevenio, Hoegh-Krohn and Tsai Tsun Wu "A class of exactly solvable three-body quantum mechanical problems and the universal low energy behavior", Phys. Lett. 83A, No. 3, 1981, and [Flamand].

then
$$\int |\psi|^2 r dr d\psi = \overline{\psi}_{\psi}$$
, $\int_{C} |\psi|^2 r dr d\psi = \frac{1}{2} a^2 \ln a$
 $\langle V \rangle = \int_{V} |\psi|^2 r dr d\psi = \frac{-2}{a^2 |\ln a|} \int_{C} |\psi|^2 r dr d\psi = -|\ln a| + cant.$
 $(\overline{\nabla}\psi)^2 = \frac{1}{2\pi} \int_{C} (-K_1(\pi))^2 \text{ wiside}$
and so $\int_{C} (-K_1(\pi))^2 \text{ distide}$
and so $\int_{C} (-K_1(\pi))^2 dr + r d\psi = \int_{a^2} \int_{a^2} (\frac{1}{2}V_s^2 - \frac{1}{a})^2 r dr$
 $\int_{C} = \int_{C} (\overline{\nabla}\psi)^2 dr + r d\psi = \int_{a^2} \int_{a^2} (\frac{1}{2}V_s^2 - \frac{1}{a})^2 r dr$
 $\int_{C} = \int_{C} (\overline{\nabla}\psi)^2 r dr d\psi = \int_{a} \int_{C} K_1(\pi) r dr = \int_{a} \frac{dr}{r}$
 $\int_{C} (\overline{\nabla}\psi)^2 r dr d\psi = \int_{a} K_1(\pi) r dr = \int_{a} \frac{dr}{r}$
 $\int_{C} (-K_1(x)) = \int_{C} (-K_1(x$

One sees that both $T_>$ and <V> diverge logarithmically as a+0 and the finite binding energy (-1) arises as a delicate cancellation between them.

On the other hand for D=3 one has

$$\Psi = \frac{1}{\sqrt{2\pi}} \left\{ \frac{e^{-T/\tau}}{\sqrt{\tau}} \text{ outside } \right\}$$

$$(A11)$$

$$\frac{1}{\sqrt{\tau}} \left\{ \frac{1}{\sqrt{\tau}} \left(\frac{\pi}{\tau} \right) \right\} \text{ in side }$$

Thus (the approximation lies in taking T_{ℓ} instead of \sqrt{s} in the expression for ℓ ; $\langle \psi_i \rangle$ and $\langle T_i \rangle$ are exact however):

$$\int |\psi|^2 r^2 dr \int \frac{\partial \partial \partial \phi}{\partial r} = | \int |\psi|^2 r^2 dr d\Omega \cong \alpha$$

$$\int V |\psi|^2 r^2 dr d\Omega = -\frac{\pi^2}{4a^2} \int_{\Gamma} |\psi|^2 r^2 dr d\Omega = -\frac{\pi^2}{4a}$$

$$\int \frac{1}{2a} \int_{\Gamma} \left(-\frac{1}{2a} + \frac{1}{2a} + \frac{1}{2a$$

One sees that in 3 dimensions not only V and T but also

T diverge, all like 1/a, and one can check that, again, the

divergent terms cancel in the expression for the total energy

B. The 3-body problem

As in the two-body case, one can separate the center of mass motion also in the open 3-body problem by going to relative coordinates

$$\vec{X}_1 \equiv \vec{\xi}_1 - \vec{\xi}_2$$
 and $\vec{X}_2 \equiv \vec{\xi}_2 - \vec{\xi}_3$

The Hamiltonian becomes

$$H_{3} = \vec{p}_{1}^{2} + \vec{p}_{2}^{2} - \vec{p}_{1} \cdot \vec{p}_{2} - (2\pi\lambda) \left(S^{(2)}(\vec{x}_{1}) + S^{(2)}(\vec{x}_{2}) \right)$$

Multiplying the equation for a bound state

gives

$$\begin{split} \left(\vec{p}_{1}^{2} + \vec{p}_{2}^{2} - \vec{p}_{1} \cdot \vec{p}_{2} + \Delta\right) \vec{\psi} (\vec{p}_{1}, \vec{p}_{2}) \\ &= (2\pi\lambda) \int \frac{d^{2}p'}{(2\pi)^{2}} \left(\vec{\psi} (\vec{p}_{1}, \vec{p}_{2}) + \vec{\psi} (\vec{p}_{1}, \vec{p}_{2}')\right) \\ &= g_{2} (\vec{p}_{2}) + g_{1} (\vec{p}_{1}) \end{split}$$

Because H₃ is invariant under interchange of 1 and 2, one can use

$$\widetilde{\psi}(\vec{p},\vec{q}) = \pm \widetilde{\psi}(\vec{q},\vec{p})$$

i.e., $g_1 = \pm g_2 \equiv g$ so that

$$\widetilde{\psi}(\vec{p},\vec{q}) = \frac{g(\vec{p}) \pm g(\vec{q})}{p^2 + 2^2 - \vec{p}\vec{q} + \Delta},$$

and from above

$$\begin{split}
\mathcal{G}(\vec{\mathbf{p}}_{1}) &= 2\pi\lambda \int \frac{d^{2}\rho_{2}}{(2\pi)^{2}} \widehat{\psi}(\vec{p}_{1},\vec{p}_{2}) = 2\pi\lambda \int \frac{d^{2}\rho_{2}}{(2\pi)^{2}} \frac{\mathcal{G}(\vec{p}_{1}) \pm \mathcal{G}(\vec{p}_{2})}{\rho_{1}^{2} + \rho_{2}^{2} - \vec{p}_{1}\vec{p}_{2} + \Delta} \\
&= (2\pi\lambda) \mathcal{G}(\vec{p}_{1}) \int \frac{d^{2}\rho_{2}}{(2\pi)^{2}} \frac{\mathcal{G}(\vec{p}_{1}) \pm \mathcal{G}(\vec{p}_{2})}{(\rho_{1}^{2} + \rho_{2}^{2} - \vec{p}_{1}\vec{p}_{2} + \Delta)} \\
&\pm (2\pi\lambda) \int \frac{d^{2}\rho_{2}}{(2\pi)^{2}} \frac{\mathcal{G}(\vec{p}_{2})}{\rho_{1}^{2} + \rho_{2}^{2} - \vec{p}_{1}\vec{p}_{2} + \Delta}
\end{split}$$

Dividing by $2\pi\lambda$, using the consistency relation (Al) for λ , and subtracting the first term on the right hand side

$$g(\vec{p}_{1}) \left\{ \int \frac{d^{2}p}{(2\pi)^{2}} \frac{1}{p^{2}+1} - \int \frac{d^{2}p_{2}}{(2\pi)^{2}} \frac{1}{p_{1}^{2}+p_{2}^{2}-\vec{p}_{1}} \vec{p}_{2}^{2} + \Delta \right\}$$

$$= \pm \left\{ \int \frac{d^{2}p_{2}}{(2\pi)^{2}} \frac{g(\vec{p}_{2})}{p_{1}^{2}+p_{2}^{2}-\vec{p}_{1}} \vec{p}_{2}^{2} + \Delta \right\}$$

Changing variables from \vec{p}_2 to $\vec{p}_{\pm}\vec{p}_2$ — $(1/2)\vec{p}_1$ on the left hand side and then from \vec{p}^2 to E, the curly bracket becomes

$$\lim_{\Lambda \to \infty} \left\{ \frac{1}{4\pi} \int_{0}^{\Lambda^{2}} \frac{dE}{E+1} - \frac{1}{4\pi} \int_{0}^{\Lambda^{2}} \frac{dE}{E+(\frac{3}{4}\rho_{1}^{2} + \Delta)} \right\}$$

$$= \frac{1}{4\pi} \ln \left(\frac{3}{4} \rho_{1}^{2} + \Delta \right) = \frac{1}{4\pi} \left(\ln \Delta + \ln \left(1 + \frac{3}{4} \rho_{1}^{2} \right) \right)$$

Defining rescaled variables $\vec{p} = \vec{p} / \sqrt{\Delta}$, $\vec{q} = \vec{p}^2 / \sqrt{\Delta}$ and $f(\vec{p}) \equiv g(\vec{p} \cdot (\vec{\Delta}))$, the resulting equation is:

$$-\ln\Delta f(\vec{p}) = \ln(1+\frac{3}{4}p^{2})f(\vec{p}) + \frac{1}{\pi}\int_{p^{2}+2^{2}-\vec{p}}\frac{f(\vec{q})d\vec{q}}{p^{2}+2^{2}-\vec{p}}\frac{f(\vec{q})d\vec{q}}{p^{2}+2^{2}-\vec{p}}\frac{f(\vec{q})d\vec{q}}{p^{2}+2^{2}-\vec{p}}\frac{f(\vec{q})d\vec{q}}{p^{2}+2^{2}-\vec{p}}\frac{f(\vec{q})d\vec{q}}{p^{2}+2^{2}-\vec{p}}\frac{f(\vec{q})d\vec{q}}{p^{2}+2^{2}-\vec{p}}\frac{f(\vec{q})d\vec{q}}{p^{2}+2^{2}-\vec{p}}\frac{f(\vec{q})d\vec{q}}{p^{2}+2^{2}-\vec{p}}\frac{f(\vec{q})d\vec{q}}{p^{2}+2^{2}-\vec{p}}\frac{f(\vec{q})d\vec{q}}{p^{2}+2^{2}-\vec{p}}\frac{f(\vec{q})d\vec{q}}{p^{2}+2^{2}-\vec{p}}\frac{f(\vec{q})d\vec{q}}{p^{2}+2^{2}-\vec{p}}\frac{f(\vec{q})d\vec{q}}{p^{2}+2^{2}-\vec{p}}\frac{f(\vec{q})d\vec{q}}{p^{2}+2^{2}-\vec{p}}\frac{f(\vec{q})d\vec{q}}{p^{2}-\vec{p}}\frac{$$

which can be rewritten as

$$f(\vec{p}) = \pm \frac{1}{\pi} \int \frac{f(\vec{q}) d\vec{q}}{\ln(\Delta(1+\vec{q})p^2)(p^2+q^2-\vec{p}\vec{q}+1)}$$

$$= \int K(\vec{p},\vec{q}) f(\vec{q}) d\vec{q} = (Kf)(\vec{p})$$
(B1')

These equations are equivalent to the Schrödinger equation for bound states

$$H_3 |\psi\rangle = -\Delta |\psi\rangle$$
(B2)

in the sense that

if
$$f(\vec{p})$$
 satisfies (B1') then
$$\begin{cases} |\psi\rangle \text{ with } \\ \widehat{\psi}(\vec{p},\vec{q}) \equiv \underbrace{f(\vec{p},\vec{\Delta}) \pm f(\vec{q},\vec{\Delta})}_{p^2+q^2-\vec{p}\cdot\vec{q}+\Delta} \\ \text{Satisfies (B2)} \end{cases}$$

Although λ and Λ and δ -functions do not appear in the equation(s) (Bl⁽¹⁾), which on a naive level might suggest that with the two-body system also the 3-body (and hopefully N-body) problem has been successfully regularized, one really still has to show that (Bl⁽¹⁾) is free of irregularities,-preferably that there is only a finite number of bound states, i.e., that: "He values of Δ for which (Bl⁽¹⁾) (and be Solved, form a finite discrete Set."

Neither the question per se nor the task of actually proving (B4) are of academic nature, as the following discussion—which is an uncompleted attempt to Tigorously answer (34) positively for D=2 and the fact that (B4) is in fact wrong for D=3 (although the corresponding equation is also free of the naive divergencies) show.

For D=3 the equation corresponding to (B1) is

$$\left(\sqrt{1+\frac{3}{4}\rho^{2}}-\frac{1}{\sqrt{\Delta}}\right)f(\vec{p})=+\frac{1}{2\pi^{2}}\int\frac{d^{2}q}{\rho^{2}+q^{2}+\vec{p}\cdot\vec{q}+1}$$
(B5)

which at least for S-waves (f=f()) has been studied extensively in the literature. Even after a continuum of solutions is removed by orthogonality conditions, (B5) still admits solutions for an infinite set of values for A, that extends to + , so that there is no ground state. These results sharpen the difficulty pointed out as early as 1935 by L.H. Thomas, who—in the formulation of the problem as the limit of particles interacting by short—range potentials—constructed a complicated trial wavefunction (whose derivatives are not everywhere continuous e.g.) for

$$\mathcal{L}_3 = -\frac{1}{a^2} \left(-\nabla_1^2 - \overrightarrow{\nabla}_1 \cdot \overrightarrow{\nabla}_2 - \nabla_2^2 + S f(g_1) + S f(g_2) \right)$$

which has infinite Binding energy as a 0. (The attempt to find the analogous trial wavefunction for D=2 leads to one containing Bessel functions and complete elliptic integrals; however, Ever, turns out to go to +∞ (rather than -∞) as a→0)

Then equations (S1) p. 259 in [F], with $\angle +1$, $\lambda^2 + \Delta$, an extra 1/2 in front of the integral (open case!) and cussion of (B5). $\beta = C$ [F] contains a long discussion of (B5).

²First derived by Skornyakov and Ter-Martirosjan, JETP 4,648(1957)

³Danilov, J.E.T.P. 13, 349 (1961).

⁴Minlos and Faddeev, J.E.T.P. Vol. 14, No. 6, 1962.

 $^{^{5}}$ L.H. Thomas, Phys. Rev. 47, 903 (1935).

This article is often quoted but never cursed at for its misprints at crucial places.*

After this brief discussion of the 3-body problem in 3-dimensions, (Bl') will be discussed (trying to prove (B4)): It is not too difficult to prove that Eq. (Bl') has no solution $f \in L^2$

if $\Delta > e^{4/3}$. One does this by noting that,

| larefore:

$$f = Kf =$$
) $||f|| = ||Kf|| \le |K| ||f||$,
where $|K| = ||R||$ (and $x \in P$) of $dx \in P$ of $dx \in P$ $dx \in P$

$$I_{(ni)} = \bigoplus_{s=0}^{\infty} \int_{4\pi}^{2\pi} K_{o}^{2}(\mu_{s}) \left\{ \frac{T^{2}}{4s_{a}^{2}} + \frac{T^{2}}{4s^{2}} \lambda + \frac{2\frac{T}{6}(\frac{2T}{3^{2}\nu_{c}-1})}{(S^{2})} + O(\frac{\alpha}{54}) dV_{2} \right\}$$
and eq. (27):
$$I_{ni} = \bigoplus_{s=0}^{\infty} \int_{-\infty}^{\infty} dv dv dv$$

^{*}In particular, Eq. (28) should read:

f=Kf cannot have a solution $f\neq 0$ ($\in L^2$) if |K| < 1. As K is clearly a monotonically decreasing function of Δ for the kernel of (Bl'), one in fact needs only to show that |K| is finite (then for some big enough $\Delta = \widetilde{\Delta}$, |K| < | and there cannot be a bound state with binding energy $\Delta > \widetilde{\Delta}$) However, accidentally |K| can be computed exactly (as a function of Δ) for

$$|K(\vec{p}, \vec{q})| = \pm \frac{1}{\pi} \frac{1}{\ln(\Delta(1+\frac{3}{4}p^2))(p^2+q^2-\vec{p}\vec{q}+1)}$$

$$|K|^2 = \frac{1}{\pi^2} (2\pi)(\frac{1}{2})^2 \int_{-\infty}^{\infty} \frac{dx dy}{\Delta(1+\frac{3}{4}x)} \int_{0}^{2\pi} \frac{d\varphi}{(x+y+1-i\eta)(p\varphi)^2}$$

$$= \int \frac{dx}{\ln^2 \Delta (1 + \frac{3}{4}x)} \int \frac{(x+y+1) dy}{((x+y+1)^2 - xy)^{\frac{3}{2}}}$$

using
$$\int \frac{dy}{(y^2 + by + c)^{3/2}} = \frac{2(2y + b)}{(4c - b^2)\sqrt{y^2 + by + c}}$$

ż

$$\int \frac{y \, dy}{(y^2 + by + c)^{3/2}} = -2 \frac{2c + by}{(4c - b^2)\sqrt{y^2 + by + c}}$$

(with b=(x+2), c=(x+1)², $(4c-b^2)=4x(1+3x/4) \ge 0$) one gets

$$|K|^{2} = \int_{0}^{\infty} \frac{dx}{\ln^{2} \Delta(H^{\frac{3}{4}}X)} \frac{1}{(H^{\frac{3x}{4}})} = \frac{4}{3} \frac{1}{\ln \Delta}$$
(B7)

So
$$|K| < 1 \quad \text{for} \quad \Delta > e^{4/3} \approx 3.79$$

Unfortunately one has to allow for a larger class of functions than L²--because

$$\|\psi\|^{2} = \int |\widetilde{\psi}(\vec{p},\vec{q})|^{2} \frac{d^{2}p}{(2\pi)^{4}}$$

$$= \frac{1}{2\pi} \int \frac{|f(\vec{p})|^{2}}{1+\frac{3}{4}p^{2}} \frac{d^{2}p}{(2\pi)^{2}} + 2Re \int_{[p^{2}+p^{2}+1](2\pi)^{4}}^{(p)} \frac{d^{2}p}{(2\pi)^{4}}$$
(using B3)

is finite for a larger class of functions L. L includes, e.g., ${\tt L}^{1+p}^2$, defined as the space of functions f with

$$\|f\|_{H^{p^{2}}} \equiv \left(\int \frac{|f|^{2}}{|+|f|^{2}} \frac{d^{2} f}{(2\pi)^{2}}\right)^{\frac{1}{2}} < \infty$$

For this space one would write (Bl') as

$$f(\vec{p}) = \int \frac{d^2q}{(2\pi)^2(1+q^2)} \widetilde{K}(\vec{p}, \vec{q}) f(\vec{q})$$

$$= \pm \frac{1}{\pi} \int \frac{(1+q^2)f(\vec{q})}{(2\pi)^2(1+q^2)} \frac{d^2q}{(2\pi)^2(1+q^2)}$$

$$= (\widetilde{K}f)(\vec{p})$$

and $|K|^2 = \int |K|^2 \frac{d^2\rho d^2q}{(2\pi)^4(\rho^2+1)(q^2+1)}$

no longer converges, so that the proof based on (B6) ceases to hold. (However, the fact that $|\widetilde{K}|_{(r)^2}$ is infinite, does not necessarily mean that (B4) is wrong.)

Looking at(Bl()) for rotationally symmetric functions $f(\vec{\rho}) \equiv h(\rho^2) \text{ simplifies the formulata little bit, but does}$ not help much:

$$\mathcal{L}(x) = \pm \int_{0}^{\infty} \frac{\mathcal{L}(y) dy}{\ln \Lambda(1 + \frac{3}{4}x) \sqrt{(x + y + 1)^{2} - xy}} \equiv (\mathbf{K}_{0} h)(x)$$
(B9)

The bound $\Delta < e^{\frac{4}{3}}$ (B7) for L²-functions is not much improved: instead of getting

$$|K|^2 = \frac{4}{3} \int_0^\infty \frac{dt}{(\ln \Delta + t)^2} = \frac{4}{3} \frac{1}{\ln \Delta}$$

(compare B7, $1 + \frac{3}{4}x = e^{\frac{t}{2}}$)

one gets

$$|K_0|^2 = \frac{4}{3} \int_{0}^{\infty} \frac{dt}{(\ln \Delta + t)^2} \left\{ \frac{\cos^{-1} \left(\frac{2 + e^{-t}}{4 - e^{-t}} \right)}{\sqrt{\frac{4}{3} \left(1 - e^{-t} \right)^2}} \right\}$$
(B10)

With
$$\frac{2+e^{-t}}{4-e^{-t}} = Co\theta$$
 the curly bracket becomes
$$\left\{ \frac{\theta}{2+a_1\theta_{12}} \right\}$$

which, instead of being =1 (in the calculation for $|K|^2$), varies slightly, but not much: Its minimal value in the range of integration is $\frac{1}{2G} \approx 0.907$.

Rewriting (B9) as

$$-\ln \Delta h(x) = \ln \left(1 + \frac{3}{4}x\right) h(x)$$

$$-\int_{0}^{\infty} \frac{h(y) dy}{\sqrt{(x+y+1)^{2}-xy}}$$
(B**9**)

(now restricting oneself also to symmetric wavefunctions for every antisymmetric $|\psi\rangle$ there is always a symmetric $|\psi\rangle$ with lower energy) one could naively apply the variational principle by thinking of the right hand side as a Hamiltonian H acting on h: (Hh)(x) with eigenvalue $-\ln\Delta$. It is not difficult to find normalized trial wavefunctions $h \in L^{1+x}$ with arbitrarily large binding energy: take

then
$$L \in L^{1+x}$$
, and in fact

$$\|L\|_{1+x} = \left(\int_{0}^{x} \frac{L(x) dx}{1+x}\right)^{\frac{1}{2}} = \left(2 \in \int_{0}^{x} \frac{dx}{(1+x)^{1+2\varepsilon}}\right)^{\frac{1}{2}} = 1$$

independent of ε . then

$$\left(\widetilde{H}\right)_{k} = 2 \in \int_{0}^{x} \frac{L(1+\frac{3}{4}x) dx}{(1+x)^{1+2\varepsilon}} = 2 \in \int_{0}^{x} \frac{L(1+x)^{2}}{(1+x)^{1+2\varepsilon}} = 2 \in \int_{$$

3

However, \widetilde{H} acting on L^{1+x} is not a self-adjoint operator, so that the "variational principle" (i.e., the statement that the true ground state energy $E_0 < \langle \widetilde{H} \rangle_h \quad \forall_h \in L^{I+x}$ does not hold.

One final argument will be given, strongly suggesting that the 3-body spectrum is bounded from below: leaving the cutoff parameter \bigwedge in the integral equation, instead of taking $\bigwedge - \rangle \infty$ — once \bigwedge has disappeared and the appearing expressions are finite as $\bigwedge -) \infty$ — one has, for S-waves:

$$g(x) = \int_{0}^{2} \frac{g(y) dy}{F(x,\Lambda)\sqrt{(x+y+\Delta)^{2}-xy}} = (K_{\Lambda}g)(x)$$
(B12)

where

$$\begin{aligned}
&F(x,\Lambda) = F(\rho^{2},\Lambda) \\
&= \frac{1}{\pi} \int_{0}^{2} d^{2}q \left\{ \frac{1}{q^{2}+1} - \frac{1}{\rho^{2}+2^{2}-\rho^{2}q^{2}+\Delta} \right\} \\
&= \int_{0}^{2} dx - \dots = \ln(\frac{3}{4}x + \Delta) + \ln(1 + \frac{1}{\Lambda^{2}}) - \ln\left(\frac{1}{2}\sqrt{1 + \frac{x+2\Delta}{\Lambda^{2}} + \frac{(x+\Delta)^{2}}{\Lambda^{2}}} + \frac{1}{2} + \frac{x+2\Delta}{\Lambda^{2}} \right)
\end{aligned}$$

and g(x) is assumed to be Lebesque-integrable on $[0, \Lambda]$.

with
$$||g||_{\Lambda}^{2} \equiv \int_{0}^{\Lambda^{2}} dx |g|^{2}$$
and
$$||K_{\Lambda}|^{2} \equiv \int_{0}^{\Lambda^{2}} |K_{\Lambda}(x,y)|^{2} dx dy$$

one has, as before (compare Eq. (B6)):

$$g = \frac{1}{\sqrt{3}} =$$

(see (B7) and (B10)), which is independent of \bigwedge for large \bigwedge , so that (B12) cannot have a solution $g\neq 0$ for any large \bigwedge , if $\bigtriangleup \geqslant e^{4/3}$.— From now on (B4) will be assumed to be true with \bigtriangleup max $\leq e^{4/3}$.

Strengthened by the above argument, one performs a variational calculation for H (defined in Bl), with

$$f(\vec{p}) = \frac{\sqrt{4\pi a}}{p^2 + a} \left(\|f\|_{L^2} = 1 \right)$$

as trial wavefunctions (a as parameter). One finds:

* rie.
$$F(x,\Lambda) = \left(\ln\left(\frac{3}{4}x + \Delta\right)\right)\left(1 + G(x,\Lambda)\right)$$
, where $\lim_{\Lambda \to \infty} G(x,\Lambda) = 0$ even if one allows $x(\Lambda^2)$ to be a diverging further of Λ .

In order to arrive at the above form of W, Feynman's trick of combining denominators was used first. The results of a numerical calculation* for different values of a, which are listed below, gave $a \approx 3/4$ to be the value which leads to a maximal lower bound, on Δ_3 , giving ≈ 2.4 .

α	\bigvee	T	WT	$\Delta_3 \geqslant$	2W-T	Δ_3
1/2	1.443	0.588	0.855	2.350	2298	
3/4	1.60	0.740	0.871	2.389	2.482	
ĺ	1.726	0.863	0.863	2.370	2.589	
4/2	1.836	ł			2672	
_	1.318	1.116			2.720	
2		1.216			2.746	
5/2	2.054	FP C.1			2.762	15.848
11/4	2.084	1.405			2.763	13.010
3	2.10	1.460	•		2.762	
4	2.193	1.648			2.738	

(2W-T has been listed, as it turns out to be the lower bound

^{*}I would like to thank Slobodan Tepic for having done this computation

Finally it will be shown that the closed 3-body problem

(i.e., all 3 particles mutuarly interacting) is exactly the same as the open case, apart from a factor of 2 in front of the integral in the integral equation(s) (B10):

$$H_{3} = \frac{1}{2} \left(\overrightarrow{T}_{1}^{2} + \overrightarrow{T}_{2}^{2} + \overrightarrow{T}_{3}^{2} \right) - (2\pi\lambda) \left(S^{(2)}(\overrightarrow{\xi}_{1} - \overrightarrow{\xi}_{2}) + S^{(2)}(\overrightarrow{\xi}_{2} - \overrightarrow{\xi}_{3}) + S^{(2)}(\overrightarrow{\xi}_{1} - \overrightarrow{\xi}_{3}) \right)$$

Multiplying
$$H_3 |\psi\rangle = -\Delta |\psi\rangle \log \langle \vec{T}_1 \vec{T}_2 \vec{T}_3 |$$

gives

$$\left(\frac{1}{2}(\vec{\Pi}_{1}^{2}+\vec{\Pi}_{2}^{2}+\vec{\Pi}_{3}^{2})+\Delta'\right)(\vec{\vec{\Pi}}_{1}\vec{\vec{\Pi}}_{2}\vec{\vec{\Pi}}_{3})$$

$$=9_{3}(\vec{\vec{\Pi}}_{3})+9_{2}(\vec{\vec{\Pi}}_{2})+9_{1}(\vec{\vec{\Pi}}_{1}) \quad (B14)$$

where

can be shown to be a function of
$$\frac{1}{11}$$
 only $\left(\text{for } \frac{3}{11} = 0 \right)$.

$$9 = (2\pi\lambda) \int d\vec{\xi}_1 d\vec{\xi}_2 d\vec{\xi}_3 e^{-i(\vec{\eta}_1 \vec{\xi}_1 + \vec{\eta}_2 \vec{\xi}_2 + \vec{\eta}_3 \vec{\xi}_3)} \cdot \int (\vec{\xi}_2 - \vec{\xi}_3) \psi(\vec{\xi}_1 \vec{\xi}_2 \vec{\xi}_3)$$

with

$$\vec{T}_{2}\vec{\xi}_{2} + \vec{T}_{3}\vec{\xi}_{3} = (\vec{T}_{2} + \vec{T}_{3})(\vec{\xi}_{2} + \vec{\xi}_{3}) + (\vec{T}_{2} - \vec{T}_{3})(\vec{\xi}_{2} - \vec{\xi}_{3}) \\
= (\vec{T}_{2} + \vec{T}_{3})(\vec{\xi}_{2} + \vec{\xi}_{2}),$$
and
$$S^{2}(\vec{\xi}_{3} - \vec{\xi}_{3}) = \int_{(2\pi)^{\frac{1}{2}}} e^{i\vec{\xi}_{2}}(\vec{\xi}_{2} - \vec{\xi}_{3})$$

one gets:

$$g_{1} = (2\pi\lambda) \int_{2\pi}^{d_{2}} \widetilde{\varphi}_{1}^{2} \widetilde{\varphi}_{1}^{2}$$

 g_2 and g_3 are given by the same expression with the arguments of $\widetilde{\Psi}$ being cyclicly permuted. Restricting oneself to totally symmetric solutions $|\psi\rangle$, $g_1=g_2=g_3=g$ therefore, and—using (B14)—one thus has:

$$g(\vec{\tau}_{i}) = (2\pi\lambda) \int d\vec{q} \int g(\vec{z} - \vec{\tau}_{i/2}) + g(\vec{z} - \vec{\tau}_{i/2}) + g(\vec{\tau}_{i}) + g(\vec{\tau}_{$$

$$Q(\vec{\pi}_{1}) = 2\pi \lambda q(\vec{\pi}_{1}) \int \frac{d^{2}_{2}}{2^{2}} \frac{1}{2^{2} + (\sqrt{4}\pi_{1}^{2} + \Delta')} + (2\pi \lambda) \cdot 2 \cdot \int \frac{d^{2}_{2}}{2\pi} \frac{Q(\vec{q} - \vec{\pi}_{1}/2)}{2^{2} + (\sqrt{4}\pi_{1}^{2} + \Delta')} \frac{Q(\vec{q} - \vec{\pi}_{1}/2)}{2^{2} + (\sqrt{4}\pi_{1}^{2} + \Delta')}$$

and

$$\widehat{V} = \frac{g(\overline{T_1}) + g(\overline{T_2}) + g(-\overline{T_1} - \overline{T_2})}{T_1^2 + T_2^2 + \overline{T_1} + \overline{T_2} + \Delta'}$$

Changing q to -q, assuming g to be an even function* and with the identification

this is, apart from a factor of 2 in front of the second term,

*
$$g(-\overline{\pi}_{1}) = 2\pi\lambda \int \frac{d^{2}g}{(2\pi)^{2}} \mathcal{V}(-\overline{\pi}_{1}) \frac{1}{2} + \overline{\pi}_{1/2} - \overline{g} + \overline{\pi}_{1/2})$$

$$= 2\pi\lambda \int \frac{d^{2}g}{(2\pi)^{2}} \mathcal{V}(-\overline{\pi}_{1}) - (\overline{g}^{2} - \overline{\pi}_{1/2}) - (-\overline{g}^{2} - \overline{\pi}_{1/2})$$

$$= g(+\overline{\pi}_{1})$$

$$= g(+\overline{\pi}_{1})$$
14) to have positive parity, i.e., $\mathcal{V}(---) = + \mathcal{V}(+++)$; Note that H_{3} is invariant under $\overline{g} \rightarrow -\overline{g}$. \forall

the same as the equation considered in the open case and the lower bounds on the ground state binding energy for the closed system (Δ_3 '), corresponding to trial wave functions

of the form
$$\sqrt{4\pi a}/\rho^2 + a$$
 are now given as e^{2W-T}

instead of e^{W-T} . $a \approx 11/4$ led to a maximal bound on Δ_3' : $\Delta_3' > 15.8$. That the binding energy of the closed three-body system comes out so large might be explained by noting that $\Delta_2' = \infty$ in a sense, because the coupling strength had been adjusted to make Δ_2 come out finite.

Because of the additional factor of 2 multiplying the kernel of the integral equation, one has

$$|K|^2 = 4 \cdot \left(\frac{4}{3} \frac{1}{\ln \Delta}\right)$$

for the L^2 -case, so that one knows that for $\Delta > e^{16/3}$ there is no square integrable solution of

$$f(\vec{p}) = \pm \frac{2}{\pi} \int \frac{f(\vec{q})d\vec{q}}{k(\Delta(H^{\frac{3}{4}}p^{\frac{2}{4}})(\rho^{2}+q^{2}-\vec{p}\vec{q}+1)}$$
(B15)

Bruch and Tjon* have in fact calculated numerically the eigenvalues of (B15) as Λ_3^* =16.1 (±0.2) (so that the above variational calculation gave in fact an astonishingly good bound) and a second eigenvalue at Δ '=1.25 (±0.05). For the open case it follows from their numerical calculation that $\Delta \approx 2.5$, which is in very good agreement with the above variational calculation.

^{*}Phys. Rev. 19, No. 2; Only after having done the work presented in this thesis did I find this article. I believe the numerical calculation, although in the theoretical treatment they take the calculation corresponding to (B7) for the closed case (and for S-waves) as proof of (B4) without worrying about functions not in L² (which does not convince me) Maybe they assume even in the numerical calculation, that the eigenfunctions are square integrable.

C. The N-body problem

Changing variables to relative coordinates x = 3-3

and setting the total momentum $\overrightarrow{P} = \overset{\nearrow}{\sum} \overrightarrow{T}$ (conjugate to $\overset{?}{\sum} = \overset{?}{\sum} \overset{?}{\sum} \overset{?}{\sum}$) equal to $\overset{?}{0}$ gives

$$+ \prod_{N=1}^{N-1} \left(\vec{p}_{r}^{2} - (2\pi\lambda) \vec{S}(\vec{x}_{r}) \right) - \sum_{T=1}^{N-2} \vec{p}_{r} \cdot \vec{p}_{r+1}$$
 (C2)

For the closed case (i.e., particles 1 and N also interacting) one can show that

$$\begin{aligned}
& + \frac{N}{N} = \sum_{i=1}^{N} (\vec{p}_{i}^{2} - (2\pi\lambda) S^{(2)}(\vec{x}_{i})) - \sum_{i=1}^{N} \vec{p}_{i} \cdot \vec{p}_{i+1} \\
& = \sum_{i=1}^{N} \frac{1}{2} (\vec{p}_{i} - \vec{p}_{i+1})^{2} - (2\pi\lambda) \sum_{i=1}^{N} S^{(2)}(\vec{x}_{i}) \\
& = T(\vec{p}_{i} ... \vec{p}_{i}) + V \qquad (\vec{p}_{N+1} = \vec{p}_{i})
\end{aligned}$$

is equivalent to "H_N" (closed case) with $\sum \overrightarrow{T}_{r} = \overrightarrow{O}$ "

provided that one restricts oneself to states with $\sum \overrightarrow{X}_{r} = \overrightarrow{O}$ (Note: $\begin{bmatrix} + \\ N \end{bmatrix} = 0$, $\begin{bmatrix} + \\ N \end{bmatrix} = 0$)

As was done for the 3-body system, one can eliminate λ and derive an integral equation from the Schrödinger equation for bound states: $H_N^{(\prime)} \mid \psi \rangle = -\Delta \mid \psi \rangle$. Multiply by a momentum Eigen-bra $\langle p_1 \dots p_N |$ to get

vector notation $\stackrel{\bullet}{\longrightarrow}$ will from now on be dropped. Defining the right hand side of (C3) to be $\sum_{r} g_r(p_1 \dots p_r \dots p_N)$ where p_r indicates that this variable does not occur, one has

$$g_{T} = (2\pi\lambda) \int \frac{d^{2}r}{(2\pi)^{2}} \left(\begin{array}{c} \rho_{1} \dots \rho_{r} \\ \rho_{r} \dots \rho_{r} \end{array} \right)$$

$$= 2\pi\lambda \sum_{S} \int \frac{d^{2}r}{(2\pi)^{2}} \frac{g_{S}(\rho_{1} \dots \rho_{r}, \rho_{s} \dots \rho_{n})}{T(\rho_{1} \dots \rho_{r} \dots \rho_{n}) + \Delta}$$

$$T(\rho_{1} \dots \rho_{r} \dots \rho_{n}) + \Delta$$

$$G_{T} = (2\pi\lambda) \int d^{2}x \dots d^{2}x \dots d^{2}x \quad e^{-i\sum_{S} \rho_{S} \times S} \psi_{(x_{1} \dots x_{s} = \rho_{s} \dots x_{n})}$$

is $(2\pi\lambda)$ times the Fourier transform of the wave function in position space, with the r-% coordinate \mathbf{x}_r fixed at the origin. Separating out the diagonal term in the above equation for \mathbf{g}_r one has:

$$\begin{aligned}
g_{\tau}(\rho_{1}...\rho_{r}...\rho_{N}) \left(\frac{1}{2\pi\lambda} - \int \frac{d^{\frac{1}{2}}_{\sigma}}{(2\pi)^{2}} \frac{1}{T(\rho_{1}...\rho_{r}...\rho_{N}) + \Delta}\right) \\
&= \sum_{S \neq \tau} \int \frac{d^{\frac{1}{2}}_{\sigma}}{(2\pi)^{2}} \frac{g_{S}(\rho_{1}...\rho_{r},\rho_{S}...\rho_{N})}{T(...) + \Delta}
\end{aligned}$$

which, using the consistency relation for $\pmb{\lambda}$ (Eq. Al), leads to an integral equation for the $g_{\bf r}$, not containing $\pmb{\lambda}$:

$$\frac{1}{2\pi\lambda} - \int \frac{d\xi}{(2\pi)^2} \frac{1}{T+\Delta}$$

$$= \frac{1}{4\pi} \int \frac{dE}{E+1} - \int \frac{d\xi_{\sigma}}{(2\pi)^2} \frac{1}{T_{1}^{\sigma-2} + \frac{1}{2}(\rho-\xi_{1}) + \frac{1}{2}(\xi_{r} - \rho_{r+1}) + T_{r+1}^{N} + \Delta}$$
where
$$\int_{-1}^{2\pi} \frac{1}{(\rho-\xi_{1})^{2}} \frac{1}{(\rho-\xi_{1})^{2} + \frac{1}{2}(\rho-\xi_{1})^{2} + T_{1}^{N} + \Delta}$$
oftensic

in the second term change integration variable from \vec{q}_r to $\vec{q}_r = 1/2 (\vec{p}_{r-1} + \vec{p}_{r+1})$ (note that the integral is only logarythmically diverging) and then to $\mathbf{E} \in \vec{q}^2$, so that the denominator becomes

and, by combining the two integrals, one gets

so that

$$g_{\tau} \cdot \ln \left(\Delta + T_{1}^{\tau-2} + T_{\tau+1}^{N} + \frac{1}{4} (\rho_{1+1} - \rho_{\tau-1})^{2} \right)$$

$$= \int_{0}^{2} g_{s} (\rho_{1} \cdot \rho_{\tau} + \rho_{\tau-1})^{2} d\rho_{s} \cdot (\rho_{1} \cdot \rho_{\tau-1} + \rho_{\tau-1})^{2} d\rho_{s} \cdot (\rho_{\tau-1} + \rho_{\tau-1} + \rho_{\tau-1})^{2} d\rho_{s} \cdot (\rho_{\tau-1} + \rho_{\tau-1} + \rho_{\tau-1})^{2} d\rho_{s} \cdot (\rho_{\tau-1} + \rho_{\tau-1} + \rho_{\tau-1} + \rho_{\tau-1})^{2} d\rho_{s} \cdot (\rho_{\tau-1} + \rho_{\tau-1} + \rho_{\tau-1} + \rho_{\tau-1})^{2} d\rho_{s} \cdot (\rho_{\tau-1} + \rho_{\tau-1} + \rho_{$$

$$=\frac{1}{\pi}\sum_{S\neq T}\int dq^{2} \frac{g_{s}(P_{1}...q_{r}P_{s}...P_{N})}{T(P_{1}...q_{r}...P_{N})+\Delta}$$

Scaling all momenta by \bigwedge and with $f_r(...p_s...) = g_r(...p_s \bigwedge$ one finally arrives at

$$- lm \Delta \cdot f_{\tau}(p_{1}...p_{\tau}...p_{N})$$

$$= lm \left(1 + T_{1}^{\tau-2} + \frac{1}{4}(p_{\tau+1} - p_{\tau-1})^{2} + T_{\tau+1}^{N}\right) \cdot f_{\tau}$$

$$- \frac{1}{T} \sum_{S \neq n} \int d^{2}q_{\tau} \frac{f_{S}(p_{1}...q_{\tau}p_{S}...p_{N})}{1 + T(p_{1}...q_{\tau}...p_{N})} \stackrel{(C5)}{=} H_{\tau S} f_{S} = (H\vec{f})_{\tau}(p_{1}...p_{\tau}...p_{N}) ; \tau = 1/2,...N.$$

Also, by definition of g_r :

Although the above derivation is written out for the closed case, all corresponding equations for the open case can be obtained by simply setting $p_N \equiv 0$ everywhere (p_N non-existing).

For the closed case one can simplify (C5) considerably by making use of the fact that H_N^{\bullet} is invariant under cyclic permutations (r \rightarrow r+1) and also reflections (N \leftrightarrow 1, N-1 \leftrightarrow 2,...). Restricting oneself to states that are singlets under these transformation i.e.,

$$\widetilde{\psi}(\rho_1...\rho_N) = \widetilde{\psi}(\rho_N \rho_1...\rho_{N-1}) = \widetilde{\psi}(\rho_N \rho_{N-1}...\rho_1)$$

one has

$$\begin{aligned}
\Im_{+}(P_{1}...P_{T}...P_{N}) &= 2\pi\lambda \int \widehat{\psi}(P_{1}...P_{1}...P_{N}) \, \frac{d\widehat{q}_{-}}{(2\pi)^{2}} \\
&= 2\pi\lambda \int \widehat{\psi}(P_{N}P_{1}...P_{T}...P_{N}) \, \frac{d\widehat{q}_{-}}{(2\pi)^{2}} \\
&= 9_{T+1}(P_{N}P_{1}...P_{T}...P_{N-1})
\end{aligned}$$
(C7)

(analogously)

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$$= g_{\tau-1}(\rho_2 \cdots \rho_{\tau} \cdots \rho_N \rho_1)$$

$$= g_{N+1-\tau}(\rho_N \cdots \rho_{\tau} \cdots \rho_1)$$

Using (C7) (in fact only (i)), (C5) becomes

$$- \ln \Delta \cdot f(\vec{p}_{z}\vec{p}_{3}...\vec{p}_{N})$$

$$= \ln \left(1 + \sum_{t=2}^{N-1} \frac{1}{2} (p_{r+t} - p_{r})^{2} + \frac{1}{4} (p_{N} - p_{z})^{2} \right) \cdot f_{(C8)}$$

- 1 Solg f (P3 P4... PN 9,)+f (P4P5... PN 9, P2)+...+ f (9, P2... PN 9)

1+ \frac{1}{2} (P_1 \frac{P_1}{P_1})^2 + \frac{1}{2} (q_1 \frac{P_2}{P_2})^2 + \frac{1}{2} \frac{1}{2} (P_1 - P_1)^2

($f = f_1$, all other f_r are obtained from f via C7). (C8) is a single Schrodinger-like equation for a function f of N-1 variables \vec{p}_r . It is important, however, to remember that (C8) (and also (C5), for the closed case) is subject to the constraint $\vec{r} = \vec{r}$ which translates to

$$f(p_2...p_N) = f(p_2+k,...,p_N+k)$$
 (C9)

(in general $f_r(p_s...)=f_r(...p_s+k...)$). Also one must not forget the condition (C6), which e.g., for N=4 says that $\int d^2pf(p,q,q')$ is invariant under all permutations of the arguments of f. For the case N=3 can (C9) be used to further reduce the number of variables explicitly: for N=3:

$$-\ln \Delta f(\rho_{2}\rho_{3}) = \ln \left(1 + \frac{3}{4} (\rho_{3} - \rho_{2})^{2}\right) f$$

$$-\frac{1}{4} \int d\hat{q}, \quad \frac{f(\rho_{3}, \rho_{1}) + f(\rho_{1}, \rho_{2})}{1 + \frac{1}{2} (\rho_{3} - \rho_{2})^{2} + \frac{1}{2} (\rho_{3} - \rho_{2})^{2} + \frac{1}{2} (\rho_{3} - \rho_{2})^{2}}$$

(C9) \Rightarrow $f(p_2,p_3) = f(p_3-p_2)$; by shifting the integration variable in the first term to q_1-p_3 , in the second to q_1-p_2 (and using f(x)=f(-x), from parity invariance of $H_N^{1,*}$) one sees that both terms are, in fact, equal to

$$-\frac{1}{\pi}\int d^{2}_{q_{1}}\frac{f^{(\vec{2}_{1})}}{1+(\vec{p}_{3}-\vec{p}_{2})^{2}+\vec{q}_{1}^{2}-\vec{q}_{1}\cdot(\vec{p}_{3}-\vec{p}_{2})}$$

which agrees with Eq. (Bl5) $(\vec{p} \neq \vec{p}_3 - \vec{p}_2)$.

N=3 is a special case ;

As

for a function of two variables, reflection invariance is equivalent to invariance under cyclic permutations, (C8) is the correct equation also for the open case (which has only reflection symmetry) putting $\frac{\Delta}{P_{N=3}} = 0$ (which up to Eq. (C5) was the simple and correct procedure of getting the corresponding equation for the open case). (C8) then is

$$-\ln \Delta f(\vec{p}_2) = \ln (1 + \frac{3}{4} \rho_2^2) f$$

$$-\frac{1}{4} \int d^2q \frac{f(q_1)}{\rho_2^2 + q_1^2 - \vec{p}_2 \cdot \vec{q}_1 + 1}$$

which is exactly Bl. The important new feature of (C8) for N>3 is that it cannot be brought into the form f=Kf with K nonsingular. As the f in the different terms in the integral of (C8) contains all the variables $p_2 \dots p_N$ and the integration variable q_1 (in cyclic permutations), K necessarily involves many δ -functions.

^{*}Or use (C7iii) for N=3, r=2 $g_2(p)=g_2(-p) \Rightarrow g_1(p)=g_2(-p)$.

SUMMARY

Two particles attracting each other by a δ -function will have infinite binding energy in 2 (or more) dimensions, unless one chooses the coupling constant to be infinitesimal and regularizes the δ -function by introducing a cutoff to the divergent integrals. Equivalently, one can define the δ -function as a limit of a short-range potential. It turns out that then 2 dimensions are more similar to lower dimensions (D \leq 2), where there is no regularization needed in the first place.

For the N-body problem, one can derive an integral equation for the Schrodinger equation for bound states, which is free of any naive divergencies. However, one has to make sure that this equation cannot be solved for arbitrarily large binding energy.

For the 3-body case this is argued not to happen (in contrast to the analogous equation in 3 dimensions, where there are eigenfunctions explicitly known for any large binding energy). The major problem is that one has to allow for a rather large class of functions f in the integral equation, as the physical wavefunction will be square integrable even if falls off much slower at ∞ (in momentum space).

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without mentioning Yoshiko I would leave this page incomplete.

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