GEOMETRY OF TIME-SPACES.

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Work partly done during the Moduli Theory-year 2006-07 at Institute Mittag-Leffler. Mathematics Subject Classification (2000): 14A22, 14H50, 14R, 16D60, 16G30, 81, 83. Keywords: Associative algebras, modules, simple modules, extensions, deformation theory, moduli spaces, non-commutative algebraic geometry, time, quantum theory.

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§ 0 Introduction

Philosophy. In a first paper on this subject, see [La 6], we sketched a physical "toy model", where the space-time of classical physics became a section of a universal fiber space $\tilde{E}$, defined on the moduli space, $H := \text{Hilb}^{12}(E^3)$, of the physical systems we chose to consider, in this case the systems composed of an observer and an observed, both sitting in Euclidean 3-space, $E^3$. This moduli space is easily computed, and has the form $H = \tilde{H}/\mathbb{Z}_2$, where $H = k[t_1, \ldots, t_6], k = \mathbb{R}$ and $\tilde{H} := \text{Spec}(H)$ is the space of all ordered pairs of points in $E^3$, $\tilde{H}$ is the blow-up of the diagonal, and $\mathbb{Z}_2$ is the obvious group-action. The space $H$, and by extension, $\tilde{H}$ and $\tilde{H}$, was called the time-space of the model.

Measurable time, in this mathematical model, turned out to be a metric $\rho$ on the time-space, measuring all possible infinitesimal changes of the state of the objects.
in the family we are studying. A relative velocity is now an oriented line in the tangent space of a point of $\tilde{H}$. Thus the space of velocities is compact.

This lead to a "physics" where there are no infinite velocities, and where the principle of relativity comes for free. The Galilean group, acts on $\mathbb{E}^3$, and therefore on $\tilde{H}$. The Abelian Lie-algebra of translations defines a 3-dimensional distribution, $\tilde{\Delta}$ in the tangent bundle of $\tilde{H}$, corresponding to 0-velocities. Given a metric on $\tilde{H}$, we define the distribution $\tilde{c}$, corresponding to light-velocities, as the normal space of $\tilde{\Delta}$. We explain how the classical space-time can be thought of as the universal space restricted to a subspace $\tilde{S}(l)$ of $\tilde{H}$, defined by a fixed line $l \subset \mathbb{E}^3$.

Moreover, we observe that the three fundamental "gauge" groups of current quantum theory $U(1)$, $SU(2)$ and $SU(3)$ are part of the structure of the fiber space, $\tilde{E} \twoheadrightarrow \tilde{H}$.

In fact, for any point $t = (o, x)$ in $H$, outside the diagonal $\Delta$, we may consider the line $l$ in $\mathbb{E}^3$ defined by the pair of points $(o, x) \in \mathbb{E}^3 \times \mathbb{E}^3$. We may also consider the action of $U(1)$ on the normal plane $B_o(l)$, of this line, oriented by the normal $(o, x)$, and on the same plane $B_x(l)$, oriented by the normal $(x, o)$. Using parallel transport in $\mathbb{E}^3$, we find an isomorphisms of bundles,

$$P_o : B_o \to B_x, P : B_o \oplus B_x \to B_o \oplus B_x,$$

the partition isomorphism. Using $P$ we may write, $(v, v)$ for $(v, P_o(v)) = P((v, 0))$. We have also seen, in loc.cit., that the line $l$ defines a unique sub scheme $H(l) \subset H$. The corresponding tangent space at $(o, x)$, is called $A_{(o, x)}$. Together this define a decomposition of the tangent space of $H$,

$$T_H = B_o \oplus B_x \oplus A_{(o, x)}.$$

If $t = (o, o) \in \Delta$, and if we consider a point $o'$ in the exceptional fiber $E_o$ of $\tilde{H}$ we find that the tangent bundle decomposes into,

$$T_{H, o'} = C_{o'} \oplus A_{o'} \oplus \tilde{\Delta},$$

where $C_{o'}$ is the tangent space of $E_o$, $A_{o'}$ is the light velocity defining $o'$ and $\tilde{\Delta}$ is the 0-velocities. Both $B_o$ and $B_x$ as well as the bundle $C_{(o, x)} := \{(\psi, -\bar{\psi})\} \in B_o \oplus B_x$, become complex line bundles on $H - \Delta$. $C_{(o, x)}$ extends to all of $\tilde{H}$, and its restriction to $E_o$ coincides with the tangent bundle. Tensorising with $C_{(o, x)}$, we complexify all bundles. In particular we find complex 2-bundles $CB_o$ and $CB_x$, on $H - \Delta$, and we obtain a canonical decomposition of the complexified tangent bundle. Any real metric on $H$ will decompose the tangent space into the light-velocities $\tilde{c}$ and the 0-velocities, $\tilde{\Delta}$, and obviously,

$$T_H = \tilde{c} \oplus \tilde{\Delta}, \quad CT_H = C\tilde{c} \oplus C\tilde{\Delta}.$$

This decomposition can also be extended to the complexified tangent bundle of $\tilde{H}$. Clearly, $U(1)$ acts on $T_{\tilde{H}}$, and $SU(2)$ and $SU(3)$ acts naturally on $CB_o \oplus CB_x$ and $C\tilde{\Delta}$ respectively. Moreover $SU(2)$ acts on $CC_{o'}$, in such a way that their actions should be physically irrelevant. $U(1)$, $SU(2)$, $SU(3)$ are our elementary gauge groups.
The above example should be considered as the most elementary one, seen from the point of view of present day physics. In fact, whenever we try to make sense of something happening in nature, we consider ourselves as observing something else, i.e. we are working with an observer and an observed, in some sort of ambient space, and the most intuitively acceptable such space, today, is obviously the 3-dimensional Euclidean space.

However, the general philosophy behind this should be the following. If we want to study a natural phenomenon, called \( \mathcal{P} \), we would, in the present scientific situation, have to be able to describe \( \mathcal{P} \) in some mathematical terms, say as a mathematical object, \( X \), depending upon some parameters, in such a way that the changing aspects of \( \mathcal{P} \) would correspond to altered parameter-values for \( X \). \( X \) would be a \textit{model} for \( \mathcal{P} \) if, moreover, \( X \) with any choice of parameter-values, would correspond to some, possibly occurring, aspect of \( \mathcal{P} \).

Two mathematical objects \( X(1) \), and \( X(2) \), corresponding to the same aspect of \( \mathcal{P} \), would be called equivalent, and the set, \( \mathcal{M} \), of equivalence classes of these objects should be called the \textit{moduli space} of the models, \( X \). The study of the natural phenomenon \( \mathcal{P} \), would then be equivalent to the study of the \textit{structure} of \( \mathcal{M} \). In particular, the notion of \textit{time} would, in agreement with Aristotle and St. Augustin, see [La 6], be a metric on this space.

With this philosophy, and this "toy"-model in mind we embarked on the study of moduli spaces of representations (modules) of associative algebras in general, see §

\section{Phase Spaces, and the Dirac Derivation.}

For any associative \( k \)-algebra \( A \) we have, in [La 6], and §1, defined a \textit{phase space} \( \text{Ph}(A) \), i.e. a universal pair of a morphism \( \iota : A \rightarrow \text{Ph}(A) \), and an \( \iota \)-derivation, \( d : A \rightarrow \text{Ph}(A) \), such that for any morphism of algebras, \( A \rightarrow R \), any derivation of \( A \) into \( R \) decomposes into \( d \) followed by an \( A \)-homomorphism \( \text{Ph}(A) \rightarrow R \), see [La 6], and [La 7]. These associative \( k \)-algebras are either trivial or non-commutative. They will give us a natural framework for quantization in physics. Iterating this construction we obtain a limit morphism \( \iota^n : \text{Ph}^n(A) \rightarrow \text{Ph}^\infty(A) \) with image \( \text{Ph}^{(n)}(A) \), and a universal derivation \( \delta \in \text{Der}_k(\text{Ph}^\infty(A), \text{Ph}^\infty(A)) \), the \textit{Dirac}-derivation. This Dirac derivation will, as we shall see, create the dynamics in our different geometries, on which we shall build our theory. For details, see §1. Notice that the notion of \textit{superspace} is easily deduced from the the Ph-construction. An affine superspace corresponds to a quotient of some \( \text{Ph}(A) \), where \( A \) is the affine \( k \)-algebra of some scheme.

\section{Non-commutative Algebraic Geometry, and Moduli of Simple Modules.}

The basic notions of affine non-commutative algebraic geometry related to a (not necessarily commutative) associative \( k \)-algebra, for \( k \) an arbitrary field, have been treated in several texts, see [La 2,3,4,5]. Given a finitely generated algebra \( A \), we prove the existence of a non-commutative scheme-structure on the set of isomorphism classes of simple finite dimensional representations, i.e. right modules, \( \text{Simp}_{<\infty}(A) \). We show in [La 4], and [La 5], that any \textit{geometric} \( k \)-algebra (see §2 below) \( A \) may be recovered from the (non-commutative) structure of \( \text{Simp}_{<\infty}(A) \), and that there is an underlying quasi-affine (commutative) scheme-structure on each
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component $\text{Simp}_n(A) \subset \text{Simp}_{<\infty}(A)$, parametrizing the simple representations of dimension $n$, see also [Procesi]. In fact, we have shown that there is a commutative algebra $C(n)$ with an open subvariety $U(n) \subset \text{Simp}_1(C(n))$, an étale covering of $\text{Simp}_n(A)$, over which there exists a "universal representation" $\tilde{V} \simeq C(n) \otimes_k V$, a vector bundle of rank $n$ defined on $\text{Simp}_1(C(n))$, and a versal family, i.e. a morphism of algebras,

$$\tilde{\rho}: A \longrightarrow \text{End}_{C(n)}(\tilde{V}) \rightarrow \text{End}_{U(n)}(\tilde{V}),$$

inducing all isoclasses of simple $n$-dimensional $A$-modules.

Suppose, in line with our Philosophy that we have uncovered the moduli space of the mathematical models of our subject, and that $A$ is the affine $k$-algebra of this space, assumed to contain all the parameters of our interest, then the above construction furnishes the Geometric landscape on which our Quantum Theory will be based.

Obviously, $\text{End}_{C(n)}(\tilde{V}) \simeq M_n(C(n))$, and we shall use this isomorphism without further warning.

§2. Non-commutative deformations and the structure of the moduli space of simple representations. We have, above, introduced moduli spaces, both for our mathematical objects, modeling the physical realities, and for the dynamical variables of interest to us. Now we have to put these things together to create dynamics in our geometry.

Dynamical Structures. A dynamical structure, see Definition (3.1), defined for a space, or any associative $k$-algebra $A$, is now an ideal $\sigma \subset \text{Ph}^\infty(A)$, stable under the Dirac derivation, and the quotient algebra $A(\sigma) := \text{Ph}^\infty(A)/\langle \sigma \rangle$, will be called a dynamical system.

These associative, but usually highly non-commutative, $k$-algebras are the models for the basic Affine algebras creating the geometric framework of our theory.

As an example, assume that $A$ is generated by the space-coordinate functions, $$\{t_i\}_{i=1}^d$$ of some configuration space, and consider a system of equations,

$$\delta^n t_p := d^n t_p = \Gamma^p(t_i, dt_j, d^2 t_k, ..., d^{n-1} t_l), \ p = 1, 2, ..., d.$$

Let $\sigma := (\delta^n t_p - \Gamma^p)$ be the two-sided $\delta$-stable ideal generated by the equations above, then $\langle \sigma \rangle$ will be called a dynamical structure or a force law, of order $n$, and the $k$-algebra,

$$A(\sigma) := \text{Ph}^\infty(A)/\langle \sigma \rangle,$$

will be referred to as a dynamical system of order $n$.

Producing dynamical systems of interest to physics, is now a major problem. One way is to introduce the notion of Lagrangian, i.e. any element $L \in \text{Ph}^\infty(A)$, and consider the Lagrange equation,

$$\delta(L) = 0.$$

Any $\delta$-stable ideal $\sigma \subset \text{Ph}^\infty(A)$, for which $\delta(L) = 0 \ (\text{mod} \langle \sigma \rangle)$, will be called a solution of the Lagrange equation. This is the non-commutative way of taking care of the parsimony principles of Maupertuis and Fermat in physics.
If we are looking for dynamical systems of order 2, we should impose on the Dirac derivation the form,

$$\delta = \sum_i dt_i \frac{\partial}{\partial t_i} + d^2 t_i \frac{\partial}{\partial dt_i},$$

in the commutative quotient of $Ph(A)$.

Whenever $A$ is commutative and smooth, we may consider classical Lagrangians, like, $L = 1/2 \sum_{i,j} g_{i,j} dt_i dt_j \in Ph(A)$, a non degenerate metric, expressed in some regular coordinate system $\{t_i\}$. Then the Lagrange equations, produces a dynamical structure of order 2,

$$d^2 t_i = -\sum_{j,k} \Gamma^{j}_{i,j,k} dt_j dt_k,$$}

where $\Gamma$ is given by the Levi-Civita connection.

One may also, for a general Lagrangian, $L \in Ph^2(A)$ impose $\delta$ as the time, and use the Euler-Lagrange equations, and obtain force laws, see the discussion later in this introduction, and in the section General Quantum Fields and Lagrangians of §3.

By definition, $\delta$ induces a derivation $\delta_a \in Der_k(A(\sigma), A(\sigma))$, also called the Dirac derivation, and usually just denoted $\delta$.

For different Lagrangians, we may obtain different Dirac derivations on the same $k$-algebra $A(\sigma)$, and therefore, as we shall see, different dynamics of the universal families of the different components of $Simp_n(A(\sigma))$, $n \geq 1$, i.e. for the "particles" of the system.

Quantum Fields and Dynamics.

Any family of components of $Simp(A(\sigma))$, with its versal family $\tilde{V}$, will, in the sequel, be called a family of particles. A section $\phi$ of the bundle $\tilde{V}$, a is now a function on the moduli space $Simp(A)$, not just a function on the configuration space, $Simp_1(A)$, nor $Simp_1(A(\sigma))$. The value $\phi(v) \in \tilde{V}(v)$ of $\phi$, at some point $v \in Simp_n(A)$, will be called a state of the particle, at the event $v$.

$End_{C(n)}(\tilde{V})$ induces also a bundle, of operators, on the étale covering $U(n)$ of $Simp_n(A(\sigma))$. A section, $\psi$ of this bundle will be called a quantum field. In particular, any element $a \in A(\sigma)$ will, via the versal family map, $\tilde{p}$, define a quantum field, and the set of quantum fields form a $k$-algebra.

Physicists will tend to be uncomfortable with this use of their language. A classical quantum field for any traditional physicist is, usually, a function $\psi$, defined on some configuration space, (which is not our $Simp_n(A(\sigma))$), with values in the polynomial algebra generated by certain creation and annihilation-operators in a Fock-space.

As we shall see, this interpretation may be viewed as a special case of our general set-up. But first we have to introduce Planck’s constant(s) and Fock-space. Then in the section Grand picture, Bosons, Fermions, and Supersymmetry, this will be explained. There we shall also focus on the notion of locality of interaction, see [Klein-Spiro], p. 104, where Cohen-Tannoudji gives a very readable explanation of this strange non-quantum phenomenon in the classical theory, see also [Weinberg], the historical introduction.

Notice also, that in physics books, the Greek letter $\psi$ is usually used for states, i.e. sections of $\tilde{V}$, or in singular cases, see below, for elements of the Hilbert space,
on which their observables act, but it is also commonly used for quantum fields. Above we have a situation where we have chosen to call the quantum fields ψ, reserving φ for the states. This is also our language in the section Grand picture, Bosons, Fermions, and Supersymmetry. Other places, we may turn this around, to fit better with the comparable notation used in physics.

Let \( v \in \text{Simp}_n(\mathbf{A}(\sigma)) \) correspond to the right \( \mathbf{A}(\sigma) \)-module \( V \), with structure homomorphism \( \rho_v : \mathbf{A}(\sigma) \to \text{End}_k(V) \), then the Dirac derivation \( \delta \) composed with \( \rho_v \), gives us an element,

\[
\delta_v \in \text{Der}_k(\mathbf{A}(\sigma), \text{End}_k(V)).
\]

Recall now that for any \( k \)-algebra \( B \), and right \( B \)-modules \( V, W \), there is an exact sequence,

\[
0 \to \text{Hom}_B(V, W) \to \text{Hom}_k(V, W) \to \text{Der}_k(B, \text{Hom}_k(V, W)) \to \text{Ext}^1_B(V, W) \to 0,
\]

where the image of,

\[
\eta : \text{Hom}_k(V, W) \to \text{Der}_k(B, \text{Hom}_k(V, W))
\]

is the sub-vectorspace of trivial (or inner) derivations.

Modulo the trivial (inner) derivations, \( \delta_v \) defines a class,

\[
\xi(v) \in \text{Ext}^1_{\mathbf{A}(\sigma)}(V, V),
\]

i.e. a tangent vector to \( \text{Simp}_n(\mathbf{A}(\sigma)) \) at \( v \). The Dirac derivation \( \delta \) therefore defines a unique one-dimensional distribution in \( \Theta_{\text{Simp}_n(\mathbf{A}(\sigma))} \), which, once we have fixed a “universal” family, defines a vector field,

\[
\xi \in \Theta_{\text{Simp}_n(\mathbf{A}(\sigma))},
\]

and, in good cases, a (rational) derivation,

\[
\xi \in \text{Der}_k(\mathbf{C}(n))
\]

inducing a derivation,

\[
[\delta] \in \text{Der}_k(\mathbf{A}(\sigma), \text{End}_{\mathbf{C}(n)}(\tilde{V})),
\]

lifting \( \xi \), and, in the sequel, identified with \( \xi \). By definition of \( [\delta] \), there is now a Hamiltonian operator

\[
Q \in M_n(\mathbf{C}(n)),
\]

satisfying the following fundamental equation, see Theorem (3.2),

\[
\delta = [\delta] + [Q, \tilde{\rho}_V(-)].
\]

This equation means that for an element (an observable) \( a \in \mathbf{A}(\sigma) \) the element \( \delta(a) \) acts on \( \tilde{V} \cong \mathbf{C}(n)^n \) as \( [\delta](a) = \xi(\tilde{\rho}_V(a)) \) plus the Lie-bracket \([Q, \tilde{\rho}_V(a)]\).

Notice that any right \( \mathbf{A}(\sigma) \)-module \( V \) is also a \( \mathbf{Ph}^\infty(\mathbf{A}) \)-module, and therefore corresponds to a family of \( \mathbf{Ph}^n(\mathbf{A}) \)-module-structures on \( V \), for \( n \geq 1 \), i.e. to
an \( A \)-module \( V_0 := V \), an element \( \xi_0 \in \text{Ext}_A^1(V,V) \), i.e. a tangent of the deformation functor of \( V_0 := V \), as \( A \)-module, an element \( \xi_1 \in \text{Ext}_{Ph(A)}^1(V,V) \), i.e. a tangent of the deformation functor of \( V_1 := V \) as \( Ph(A) \)-module, an element \( \xi_2 \in \text{Ext}_{Ph^2(A)}^1(V,V) \), i.e. a tangent of the deformation functor of \( V_2 := V \) as \( Ph^2(A) \)-module, etc. All this is just \( V \), considered as an \( A \)-module, together with a sequence \( \{ \xi_n \} \), \( 0 \leq n \), of a tangent, or a momentum, \( \xi_0 \), an acceleration vector, \( \xi_1 \), and any number of higher order momenta \( \xi_n \). Thus, specifying a point \( v \in \text{Simp}_n(A(\sigma)) \) implies specifying a formal curve through \( v_0 \), the base-point, of the miniversal deformation space of the \( A \)-module \( V \).

Knowing the dynamical structure, \((\sigma)\), and the state of our object \( V \) at a time \( \tau_0 \), i.e. knowing the structure of our representation \( V \) of the algebra \( A(\sigma) \), at that time (which is a problem that we shall return to), the above makes it reasonable to believe that we, from this, may deduce the state of \( V \) at any "later" time \( \tau_1 \). This assumption, on which all of science is based, is taken for granted in most textbooks in modern physics. This paper is, in fact, an attempt to give this basic assumption a reasonable basis. The mystery is, of course, why Nature seems to be parsimonious, in the sense of Fermat and Maupertuis, giving us a chance of guessing dynamical structures.

The dynamics of the system is now given in terms of the Dirac vector-field \([\delta]\), generating the vector field \( \xi \) on \( \text{Simp}_n(A(\sigma)) \). An integral curve \( \gamma \) of \( \xi \) is a solution of the equations of motion. Let \( \gamma \) start at \( v_0 \in \text{Simp}_n(A(\sigma)) \) and end at \( v_1 \in \text{Simp}_n(A(\sigma)) \), with length \( \tau_1 - \tau_0 \). This is only meaningful for ordered fields \( k \), and when we have given a metric (time) on the moduli space \( \text{Simp}_n(A(\sigma)) \). Assume this is the situation. Then, given a state, \( \phi(v_0) \in \hat{V}(v_0) \simeq V_0 \), of the particle \( \hat{V} \), we prove that there is a canonical evolution map, \( U(\tau_0, \tau_1) \) transporting \( \phi(v_0) \) from time \( \tau_0 \) i.e. from the point representing \( V_0 \), to time \( \tau_1 \), i.e. corresponding to some point representing \( V_1 \), along \( \gamma \). It is given as,

\[
U(\tau_0, \tau_1)(\phi(v_0)) = \exp(\int_{\gamma} Q(\phi))
\]

where \( \exp(\int_{\gamma}) \) is the non-commutative version of the classical action integral, related to the Dyson series, to be defined later, see the proof of Theorem (3.3) and the section Grand picture. Bosons, Fermions, and Supersymmetry. In case we work with unitary representations, of some sort, we may also deduce analogies to the \( S \)-matrix, perturbation theory, and so also to Feynman-integrals and diagrams.

**Classical Quantum Theory.** Most of the classical models in physics are either essentially commutative, or singular, i.e. such that either \( Q = 0 \), or \([\delta] = 0 \). General relativity is an example of the first category, classical Yang-Mills theory is of the second kind. In fact, any theory involving connections are singular, and infinite dimensional. But we shall see that imposing singularity on a theory, sometimes recover the classical infinite dimensional (Hilbert-space-based) model as a limit of the finite dimensional simple representations, corresponding to a dynamic system, see Examples 3.6-3.8, where we treat the Harmonic Oscillator.

**Planck’s Constants, and Fock space.** This general model allows us also to define a general notion of a Planck’s constant(s), \( h_i \), as the generator(s) of the "generalized
monoid",
\[ \Lambda(\sigma) := \{ \lambda \in C(n) \mid \exists f_\lambda \in A(\sigma), f_\lambda \neq 0, \]
\[ [Q, \tilde{\rho}(\delta(f_\lambda))] = \tilde{\rho}(\delta(f_\lambda)) - [\delta(\tilde{\rho}(f_\lambda))] = \lambda \tilde{\rho}(f_\lambda) \}\]
which has the property that \( \lambda, \lambda' \in \Lambda(\sigma), f_\lambda f_{\lambda'} \neq 0 \) implies \( \lambda + \lambda' \in \Lambda(\sigma) \). From this definition we may construct a general notion of Fock algebra, or space, and a representation, also named \( F \), on this space. \( F \) is the sub-\( k \)-algebra of \( \text{End}_{C(n)}(\tilde{V}) \) generated by \( \{ a^+_l := f_{l_1}, a^-_l := f_{-l_2} \} \), see (3.7) and (3.9) for a rather complete discussion of the one-dimensional harmonic oscillator in all ranks, and of the quartic anharmonic oscillator in rank 2 and 3. Notice that this is just a natural generalization of standard work on classification of representations of (semi-simple) Lie algebras, see the discussion of fundamental particles in the Example (3.18).

When \( A \) is the coordinate \( k \)-algebra of a moduli space, we should also consider the family of Lie algebras of essential automorphisms of the objects classified by \( \text{Simp}(A(\sigma)) \), and apply invariant theory, like in [La 4], to obtain a general form for Yang-Mills theory, see [Bj-La] and [La-Pf], for the case of plane curve singularities. This would offer us a general model for the notions of gauge particles and gauge fields, coupling with ordinary particles via representations onto corresponding simple modules.

**General Quantum Fields, Lagrangians and Actions.** Perfectly parallel with this theory of simple finite dimensional representations, we might have considered, for given algebras \( A \), and \( B \), the space of algebra homomorphisms, \( \phi : A \to B \).

In the commutative, classical case, when \( A \) is generated by \( t_1, \ldots, t_r \), and \( B \) is the affine algebra of a configuration space generated by \( x_1, \ldots, x_s \), \( \phi \) is determined by the images \( \phi_i := \tilde{\phi}(t_i) \), and \( \phi \) or \( \{ \phi_i \} \) is called a classical field. Any such field, \( \phi \) induces a unique commutative diagram of algebras,

\[
\begin{array}{ccc}
A & \xrightarrow{\phi} & B \\
\downarrow & & \downarrow \\
Ph(A) & \xrightarrow{Ph\phi} & Ph(B).
\end{array}
\]

Given dynamical structures, (say of order two), \( \sigma \) and \( \mu \), defined on \( A \), resp. \( B \), we construct a vector field \( [\delta] \) on the space, \( \mathcal{F}(A(\sigma), B(\mu)) \), of fields, \( \phi : A(\sigma) \to B(\mu) \). The singularities of \( [\delta] \) defines a subset,

\[ \mathcal{M} := \mathcal{M}(A(\sigma), B(\mu)) \subset \mathcal{F}(A(\sigma), B(\mu)) =: \mathcal{F}. \]

There are natural equations of motion, analogous to those we have seen above, see (3.2). Notice that a field \( \phi \in \mathcal{M} \) is said to be on shell, those of \( \mathcal{F} - \mathcal{M} \) are off shell.

We shall explore the structure of \( \mathcal{M} \) in some simple cases.

The actual choice of dynamical structures \( (\sigma), (\mu) \), for the particular physical set-up, is, of course, not obvious. They may be defined in terms of "force laws", but, in general, force laws do not pop up naturally. Instead, physicists are used to
insist on the Lagrangian, an element $L \in Ph(A)$, as a main player in this game. The \textit{Lagrangian density}, $L$ should then be considered an element of the versal family of the iso-classes of $F(A,B)$. In fact, assuming that this space has a local affine algebraic geometric structure, parametrized by some ring $\mathcal{C}$, we may consider the versal family as a homomorphisms of $k$-algebras,

$$\tilde{\phi} : Ph(A) \to \mathcal{C} \otimes_k Ph(B),$$

and put $L := \tilde{\phi}(L)$. Classically one picks a (natural) representation, corresponding to a derivation of $B$,

$$\rho : Ph(B) \to B,$$

and put, $\mathcal{L} := \rho(L)$. One considers the Lagrangian density as a function in $\phi_i, \phi_{i,j} := \frac{\partial \phi_i}{\partial x_j}$, thus as a function on configuration space $\text{Simp}_1(B)$, with coefficients from $\mathcal{C}$. One postulates that there is a functional, or an \textit{action}, which, for every field $\phi$, associates a real or complex value,

$$S := S(\mathcal{L}(\phi_i, \phi_{i,j})),$$

usually given in terms of a trace, or as an integral of $\mathcal{L}$ on part of the configuration space, see below. $S$ should be considered as a function on $\mathcal{F} := F(A,B)$, i.e. as an element of $\mathcal{C}$. The parsimony principles of Fermat and Maupertuis is then applied to this function, and one wants to compute the vector field,

$$\nabla S \in \Theta_\mathcal{F},$$

which mimic our $[\delta]$, derived from the Dirac derivations. The equation of motion, i.e. the equations picking out the subspace $M \subset \mathcal{F}$, is therefore,

$$\nabla S = 0.$$

Here is where classical calculus of variation enters, and where we obtain differential equations for $\phi_i$, the Euler-Lagrange equations of \textit{motion}.

Notice now that in an infinite dimensional representation, the \textit{Trace} is an integral on the spectrum. The equation of motion defining $M \subset \mathcal{F}$, now corresponds to,

$$\delta S := \delta \int \tilde{\rho}(\mathcal{L}) = 0.$$

The calculus of variation produces Euler-Lagrange equations, and so picks out the singularities of $\nabla S$, the replacement for $[\delta]$, without referring to a dynamical structure, or to (uni)versal families. See the Examples (3.7) and (3.8), where we treat the harmonic oscillator, and where we show that the classical infinite dimensional representation is a limit of finite dimensional simple representations. We also show that the Lagrangian of the harmonic oscillator produces a vector field $\nabla S$ on $\text{Simp}_2(A(\sigma))$ which is different from the one generated by the Dirac derivation for the dynamical system deduced from the Euler-Lagrange equations for the same Lagrangian. However, the sets of singularities for the two vector fields coincide.
This should never the less be cause for worries, since the world we can test is finite. The infinite dimensional mathematical machinery is obviously just a computational trick.

Another problem with this reliance on the Parsimony Principle via Lagrangians, and the (commutative) Euler-Lagrange equations, is that, unless we may prove that $\nabla S = [\delta]$, for some dynamical structure $\sigma$, the philosophically satisfying realization, that a preparation in $A(\sigma)$ actually implies a deterministic future for our objects, disappear, see above.

Otherwise, it is clear that the theory becomes more flexible. It is easy to cook up Lagrangians.

In QFT, quantizing fields, physicists are, however, usually strangely vague; suddenly they consider functions, $\{\phi_i, \phi_{i,j}\}$, on configuration space, as elements in a k-algebra, introduce commutation relations and start working as if these functions on configuration space were operators. This is, maybe, due to the fact that they do not see the difference between the role of $B$ in the classical case, and the role of $Ph(B)$, in quantum theory.

**Grand Picture. Bosons, Fermions, and Supersymmetry.** With this done, we sketch the big picture of QFT that emerges from the above ideas. This is then used as philosophical basis for the treatment of the harmonic oscillator, general relativity, electromagnetism, Dirac's equation, spin and quarks, which are the subjects of the examples, (3.6) to (3.18).

In particular, we sketch, here and in §4, how we may treat the problems of **Bosons, Fermions, Anyons, and Super-symmetry.**

**Connections and the Generic Dynamical Structure.** Moreover we shall see that, on a space with a non-degenerate metric, there is a unique **generic dynamical structure**, $(\sigma)$, which produces the most interesting physical models. In fact, any connection on a bundle, induces a representation of $A(\sigma)$. We shall use this method to quantize the **Electromagnetic Field**, as well as the **Gravitational Field**, obtaining generalized Maxwell, Dirac and "Einstein" equations, with interesting properties, see (3.5), (3.17) and (3.18). The Levi-Civita connection turns out to be a very particular singular representation for which the Hamiltonian is identified with the Laplace-Beltrami operator.

**Clocks and Classical Dynamics.** At this point we need to be more interested in how to measure time. We therefore discuss the notion of clocks in this picture, and we propose two rather different models, one called The Western clock, modeled on a free particle in dimension 1, i.e. one with $d^2\tau = 0$, and another, called the Eastern Clock, modeled on the harmonic oscillator in dimension 1, i.e. one with $d^2\tau = \tau$.

**Time-Space and Space-Time.** The application to the case of point-like particles in the $\tilde{H}$-model is treated in Example (3.5), mainly as an introduction to the study of the Levi-Civita connection, in our tapping. Coupled with the non-trivial geometry of $\tilde{H}$ we see a promising possibility of defining notions like mass and charge, of different colors, related to the structure of $\tilde{H}$ along the diagonal $\Delta$. A catchy way of expressing this would be that every point in our real world is a "black hole", outfitted with a density of, at least, mass and charge. Notice that the dimension of $\Delta$ is 5, which brings about ideas like those of Kaluza and Klein.

In particular, the definition of mass, and the deduction of Newton's law of gravitation, from the assumption that mass is a property of the geometry of $\tilde{H}$, related
to the blow up along the diagonal, seems promising. A simple example in this
direction leads to a Schwarzschild-type geometry. The corresponding equations of
motion reduces to Kepler’s laws, see the Example 3.16. As another example, we
shall again go back to our ”toy” model, where the standard Gauge groups, \( U(1), \)
\( SU(2), \) and \( SU(3) \) pop up canonically and show that the results above can be used
to construct a general geometric theory closely related to general relativity and to
quantum theory, generalizing both. See the examples (3.17),(3.18), where the action
of the natural gauge group, on the canonically decomposed tangent bundle of \( H, \)
as described above, sets up a nice theory for elementary particles, spin, isospin,
hypercharge, including quarks. Here the notion of non-commutative invariant space,
plays a fundamental role. In particular, notice the possible models for ”light” and
”dark” matter, or energy, hinted upon in the examples (3.17), and (3.18). Notice
also that, in this ”toy” model light cannot be described as point-particles. There are
no ”radars” available for point-particles, like in current general relativity. However,
the quantized E-M works well to explain communication with light. Moreover, as
one might have expected, a reasonable model of the process creating the universe
as we see it, will provide a better understanding of what we are modeling. This is
the subject of the next section.

\textit{Cosmology, Big Bang and all that.} Our ”toy-model”, i.e. the moduli space, \( H, \)
of two points in the Euclidean 3-space, or its étale covering, \( \tilde{H}, \) turns out to be \textit{created}
by the deformations of the obvious (non-commutative) singularity in 3-dimensions,
\( U := k < x_1, x_2, x_3 > / (x_1, x_2, x_3)^2. \) In fact, it is easy to see that the versal
space of the deformation functor of the \( k \)-algebra \( U \) contains a flat component (a
room in the modular suite, see [La-Pf]) isomorphic to \( \tilde{H}. \) The modular stratum
(the inner room) is reduced to the base point. This furnishes a nice model for The Universe
with easy relations to classical cosmological models, like those of Friedman-Robertson-Walker, and Einstein-de Sitter.

\textit{Interaction and Non-commutative Algebraic Geometry.} In §4, we shall introduce
interactions, lifetime, decay and creation of particles. The possibility of treating
interaction between fields in a perfectly geometric way, with the usual metrics
and connections replaced with a ”non-commutative metric” is, maybe, the most
interesting aspect of the model presented in this paper.

The essential point is that, in non-commutative algebraic geometry, say in the
”space” of representations of an algebra \( B, \) there is a tangent space, \( T(V,W) := \text{Ext}^1_B(V,W), \)
between any two points, \( V, W. \) In particular, if \( B = \text{Ph}(k[x_1, ..., x_n]), \)
then any 1-dimensional representation of \( B \) is represented as a pair \( (q, \xi), \) of a
closed point \( q \) of \( \text{Spec}(k[z]), \) and a tangent \( \xi \) at that point. Given two such points,
\( (q_i, \xi_i), i = 1, 2, \) an easy calculation proves that \( T((q_1, \xi_1), (q_2, \xi_2)) \) is of dimension
1 if \( q_1 \neq q_2, \) of dimension \( n \) if \( q_1 = q_2, \xi_1 \neq \xi_2, \) and of dimension \( 2n \) when \( (q_1, \xi_1) =
(q_2, \xi_2), \) see Example(1.1),(ii).

Now, just as we may talk about vector-fields, as the assignment of a tangent
vector to any point in space, and consider metrics as functions that associate a
length to any tangent-vector, we may consider fields of tangents between any two
points, and extend the notion of metric to measure the length of such.

If we do, we find very nice models for treating the notion of identical particles,
and interaction between fields, see the Examples (4.3), (4.4).

Finally, we shall not resist the temptation to attempt a formalization, in our
language, of the notion of Alternative Histories, see [Gell-Mann], p.140, and the
paper [Gell-Mann and Hartle]. The result is another application of noncommutative

deformation theory which seems to be a promising tool in mathematical physics.

Apology. Referring to the historical introduction of Weinberg’s, The Quantum

Theory of Fields, see [Weinberg], where he quotes Heisenberg’s 1925-Manifesto, I

must confess that the present paper is based on the same positivistic philosophy as

the one Weinberg rules out.

But then, I am not a physicist, and this paper is a paper on geometry of cer-

tain finitely generated non-commutative algebraic schemes, where I have taken the

liberty of using my version of the physicist jargon to make the results more palpable.

Even though I see a lot of difficulties in the interpretation of the mathematical

notions of my models, in a physics context, I hope that the model I propose may help

other mathematicians to gain faith in their jugend-traums; sometime, somehow, to

be able to understand some physics.

An attentive reader will also see that, if my modelist philosophy about Nature,

see above, should be taken seriously, it would reduce the physicists work to define,

in a mathematical language, the model of the "object" she is studying, then with

the help of a mathematician work out the moduli space of all such models, define

the infinite phase space of this moduli space, guess on a metric to define time, and

a corresponding dynamical structure, and give the result to the computer algebra

group in Kaiserslautern, and hope for the best.

§1 Phase spaces and the Dirac derivation

Given a $k$-algebra $A$, denote by $A/k - \text{alg}$ the category where the objects are

homomorphisms of $k$-algebras $\kappa : A \to R$, and the morphisms, $\psi : \kappa \to \kappa'$ are

commutative diagrams,

\[
\begin{array}{ccc}
A & \xrightarrow{\kappa} & \kappa' \\
\downarrow & & \downarrow \\
R & \xleftarrow{\psi} & R'
\end{array}
\]

and consider the functor,

\[
\text{Der}_k(A, -) : A/k - \text{alg} \to \text{Sets}.
\]

It is representable by a $k$-algebra-morphism,

\[
i : A \to Ph(A),
\]

with a universal family given by a universal derivation,

\[
d : A \to Ph(A).
\]

Ph (A) is relatively easy to compute. It can be constructed as the non-commutative

versal base of the deformation functor of the morphisme $\rho : A \to k[\epsilon]$, see [La 6]

and [La 7].

Clearly we have the identities,

\[
d_* : \text{Der}_k(A, A) = \text{Mor}_A(Ph(A), A),
\]
and,
\[ d^* : \text{Der}_k(A, Ph(A)) = \text{End}_A(Ph(A)), \]
the last one associating \( d \) to the identity endomorphism of \( Ph \). Let now \( V \) be a right \( A \)-module, with structure morphism \( \rho : A \to \text{End}_k(V) \). We obtain a universal derivation,
\[ c : A \longrightarrow \text{Hom}_k(V, V \otimes_A Ph(A)), \]
defined by, \( c(a)(v) = v \otimes d(a) \). Using the long exact sequence, see the introduction,
\[
0 \to \text{Hom}_A(V, V \otimes_A Ph(A)) \to \text{Hom}_k(V, V \otimes_A Ph(A)) \to^\iota \text{Der}_k(A, \text{Hom}_A(V, V \otimes_A Ph(A))) \to^\kappa \text{Ext}^1_A(V, V \otimes_A Ph(A)) \to 0,
\]
we obtain the non-commutative Kodaira-Spencer class,
\[ c(V) : = \kappa(c) \in \text{Ext}^1_A(V, V \otimes_A Ph(A)), \]
inducing the Kodaira-Spencer morphism,
\[ g : \Theta_A : = \text{Der}_k(A, A) \longrightarrow \text{Ext}^1_A(V, V), \]
via the identity, \( \delta_\ast \). If \( c(V) = 0 \), then the exact sequence above proves that there exist a \( \nabla \in \text{Hom}_k(V, V \otimes_A Ph(A)) \) such that \( \delta = \iota(\nabla) \). This is just another way of proving that \( \delta \) is given by a connection,
\[ \nabla : \text{Der}_k(A, A) \longrightarrow \text{Hom}_k(V, V). \]
As is well known, in the commutative case, the Kodaira-Spencer class gives rise to a Chern character by putting,
\[ \text{ch}^i(V) : = 1/i! \ c^i(V) \in \text{Ext}^i_A(V, V \otimes_A Ph(A)), \]
and if \( c(V) = 0 \), the curvature \( R(V) \) induces a curvature class,
\[ R_\Sigma \in H^2(k, A; \Theta_A, \text{End}_A(V)). \]
Any \( Ph(A) \)-module \( W \), given by its structure map,
\[ \rho_W : Ph(A) \longrightarrow \text{End}_k(W) \]
corresponds bijectively to an induced \( A \)-module structure on \( W \), and a derivation \( \delta_\rho \in \text{Der}_k(A, \text{End}_k(W)) \), defining an element,
\[ [\delta_\rho] \in \text{Ext}^1_A(W, W), \]
see the introduction. Fixing this element we find that the set of \( Ph(A) \)-module structures on the \( A \)-module \( W \) is in one to one correspondence with,
\[ \text{End}_k(W) / \text{End}_A(W). \]
Conversely, starting with an $A$-module $V$ and an element $\delta \in \text{Der}_k(A, \text{End}_k(V))$, we obtain a $\text{Ph}(A)$-module $V_\delta$. It is then easy to see that the kernel of the natural map,

$$\text{Ext}^1_{\text{Ph}(A)}(V_\delta, V_\delta) \to \text{Ext}^1_A(V, V),$$

induced by the linear map,

$$\text{Der}_k(\text{Ph}(A), \text{End}_k(V_\delta)) \to \text{Der}_k(A, \text{End}_k(V))$$

is the quotient,

$$\text{Der}_A(\text{Ph}(A), \text{End}_k(V_\delta))/\text{End}_k(V).$$

Remark. Since $\text{Ext}^1_A(V, V)$ is the tangent space of the miniversal deformation space of $V$ as an $A$-module, see e.g. [La 4], or the next section §2, we see that the non-commutative space $\text{Ph}(A)$ also parametrizes the set of generalized momenta, i.e. the set of pairs of a simple module $V \in \text{Simp}(A)$, and a tangent vector of $\text{Simp}(A)$ at that point.

Example 1.1. (i) Let $A = k[t]$, then obviously, $\text{Ph}(A) = k < t, dt >$ and $d$ is given by $d(t) = dt$, such that for $f \in k[t]$, we find $d(f) = J_1(f)$ with the notations of [La 5], i.e. the non-commutative derivation of $f$ with respect to $t$. One should also compare this with the non-commutative Taylor formula of loc.cit. If $V \simeq k^2$ is an $A$-module, defined by the matrix $X \in M_2(k)$, and $\delta \in \text{Der}_k(A, \text{End}_k(V))$, is defined in terms of the matrix $Y \in M_2(k)$, then the $\text{Ph}(A)$-module $V_\delta$ is the $k < t, dt >$-module defined by the action of the two matrices $X, Y \in M_2(k)$, and we find

$$e^1_{V_\delta} : = \dim_k \text{Ext}^1_A(V, V) = \dim_k \text{End}_A(V) = \dim_k \{ Z \in M_2(k) \mid [X, Z] = 0 \}$$

$$e^1_{V_\delta} : = \dim_k \text{Ext}^1_{\text{Ph}(A)}(V_\delta, V_\delta) = 8 - 4 + \dim \{ Z \in M_2(k) \mid [X, Z] = [Y, Z] = 0 \}.$$

We have the following inequalities,

$$2 \leq e^1_{V_\delta} \leq 4 \leq e^1_{V_\delta} \leq 8.$$

(ii) Let $A = k^2 \simeq k[x]/(x^2 - r^2)$, $r \in k$, $r \neq 0$, then,

$$\text{Ph}(A) = k < x, dx > /((x^2 - r^2), x \cdot dx + dx \cdot x).$$

Notice that $\text{Ph}(A)$ just has 2 points, i.e. simple representations, given by,

$$k(r) : x = r, dx = 0, \ k(-r) : x = -r, \ dx = 0.$$

An easy computation shows that,

$$\text{Ext}^1_{\text{Ph}(A)}(k(\alpha), k(\alpha)) = 0, \ \alpha = r, -r, \ \text{Ext}^1_{\text{Ph}(A)}(k(\alpha), k(-\alpha)) = k \cdot \omega,$$

where $\omega$ is represented by the derivation given by $\omega(x) = 2r$, $\omega(dx) = t \in k$ where $t$ is the tension of this string of dimension $-1$, see end of §2, and end of §3. Notice also that this is an example of the existence of tangents between different points, in non-commutative algebraic geometry.
(iii) Now, let \( A = k[x] := k[x_1, x_2, x_3] \) and consider,
\[
Ph(A) = k < x_1, x_2, x_3, dx_1, dx_2, dx_3 > /([x_i, x_j], d([x_i, x_j])).
\]

Any rank 1 representation of \( A \) is represented by a pair of a closed point \( q \) of \( Spec(k[x]) \), and a tangent \( p \) at that point. Given two such points, \((q_i, p_i), i = 1, 2\), an easy calculation proves,
\[
dim_k \text{Ext}^1_{PhA}(k(q_1, p_1), k(q_2, p_2)) = 1, \text{for } q_1 \neq q_2
\]
\[
dim_k \text{Ext}^1_{PhA}(k(q_1, p_1), k(q_2, p_2)) = 3, \text{for } q_1 = q_2, p_1 \neq p_2
\]
\[
dim_k \text{Ext}^1_{PhA}(k(q_1, p_1), k(q_2, p_2)) = 6, \text{for } (q_1, p_1) = (q_2, p_2)
\]

Put \( x_j(q_i, p_i) := q_{i,j}, \ dx_j((q_i, p_i)) := p_{i,j}, \ \alpha_j = q_{1,j} - q_{2,j}, \ \beta_j = p_{1,j} - p_{2,j}. \) See that for any element \( \alpha \in \text{Hom}_k(k((q_1, p_1)), k((q_2, p_2))) \) we have,
\[
x_j \alpha = q_{1,j} \alpha, \ \alpha x_j = q_{2,j} \alpha, \ dx_j \alpha = p_{1,j} \alpha, \ adx_j = p_{2,j} \alpha,
\]
with the obvious identification. Any derivation
\[
\delta \in \text{Der}_k(PhA, \text{Hom}_k(k((q_1, p_1)), k((q_2, p_2))))
\]
must satisfy the relations,
\[
\delta([x_i, x_j]) = [\delta(x_i), x_j] + [x_i, \delta(x_j)] = 0
\]
\[
\delta([dx_i, x_j] + [x_i, dx_j]) = [\delta(dx_i), x_j] + [dx_i, \delta(x_j)] + [\delta(x_i), dx_j] + [x_i, \delta(dx_j)] = 0.
\]

Using the above left-right action-rules, the result follows from the long exact sequence computing \( \text{Ext}^1_{PhA} \). The two families of relations above give us two systems of linear equations.

The first, in the variables \( \delta(x_1), \delta(x_2), \delta(x_3), \delta(dx_1), \delta(dx_2), \delta(dx_3) \), with matrix,
\[
\begin{pmatrix}
-\beta_2 & \beta_1 & 0 & -\alpha_2 & \alpha_1 & 0 \\
-\beta_3 & 0 & \beta_1 & -\alpha_3 & 0 & \alpha_1 \\
0 & -\beta_3 & \beta_2 & 0 & -\alpha_3 & \alpha_2
\end{pmatrix},
\]
and the second, in the variables \( \delta(x_1), \delta(x_2), \delta(x_3) \), with matrix,
\[
\begin{pmatrix}
-\alpha_2 & \alpha_1 & 0 \\
-\alpha_3 & 0 & \alpha_1 \\
0 & -\alpha_3 & \alpha_2
\end{pmatrix}.
\]

In particular we see that the trivial derivation given by,
\[
\delta(x_i) = \alpha_i, \ \delta(dx_j) = \beta_j,
\]
of course satisfies the relations, and the generator of \( \text{Ext}^1_{PhA}(k(q_1, p_1), k(q_2, p_2)) \) is represented by,
\[
\delta(x_i) = 0, \ \delta(dx_j) = \alpha_i.
\]
This is, in an obvious sense, the vector \(-q_1, q_2\), and we notice that this generator is of the type \(\delta(d-1\hat{\epsilon})\), so it is an acceleration in Simp_1(k|\mathcal{X}|), see the interpretation of this as an interaction in §4.

It is not difficult to extend this result from dimension 3 to any dimension \(n\), see the introduction, and §4.

(iv) Consider now the space of 2-dimensional representation of \(\text{Ph}(A)\). It is an easy computation that any such is given by the actions,

\[
x_1 = \begin{pmatrix} a_1 & 0 \\ 0 & a_2 \end{pmatrix}, \quad x_2 = \begin{pmatrix} b_1 & 0 \\ 0 & b_2 \end{pmatrix}, \quad x_3 = \begin{pmatrix} c_1 & 0 \\ 0 & c_2 \end{pmatrix},
\]

and,

\[
dx_1 = \begin{pmatrix} \alpha_{1,1} & (a_1 - a_2) \\ (a_2 - a_1) & \alpha_{2,2} \end{pmatrix},
\]

\[
dx_2 = \begin{pmatrix} \beta_{1,1} & (b_1 - b_2) \\ (b_2 - b_1) & \beta_{2,2} \end{pmatrix},
\]

\[
dx_3 = \begin{pmatrix} \gamma_{1,1} & (c_1 - c_2) \\ (c_2 - c_1) & \gamma_{2,2} \end{pmatrix}
\]

The angular momentum is now given by,

\[
L_{1,2} := x_1 dx_2 - x_2 dx_1 = \begin{pmatrix} (a_1\beta_{1,1} - b_1\alpha_{1,1}) & (a_2b_1 - a_1b_2) \\ (a_1b_2 - a_2b_1) & (a_2\beta_{2,2} - b_2\alpha_{2,2}) \end{pmatrix},
\]

etc. And the isospin, see (3.18) and (3.19), has the form,

\[
I_1 := [x_1, dx_1] = \begin{pmatrix} 0 & (a_1 - a_2)^2 \\ (a_2 - a_1)^2 & 0 \end{pmatrix},
\]

etc.

(v) Let \(A = M_2(k)\), and let us compute \(\text{Ph}(A)\). Clearly the existence of the canonical homomorphism, \(i : M_2(k) \to \text{Ph}(M_2(k))\) shows that \(\text{Ph}(M_2(k))\) must be a matrix ring, generated, as an algebra, over \(M_2(k)\) by \(d\epsilon_{i,j}\), \(i, j = 1, 2\), where \(\epsilon_{i,j}\) is the elementary matrix. A little computation shows that we have the following relations,

\[
de_{1,1} = \begin{pmatrix} 0 & (de_{1,1})_{1,2} = -(de_{2,2})_{1,2} \\ (de_{1,1})_{2,1} = -(de_{2,2})_{2,1} & 0 \end{pmatrix}
\]

\[
de_{2,2} = \begin{pmatrix} 0 & (de_{2,2})_{1,2} = -(de_{1,1})_{1,2} \\ (de_{2,2})_{2,1} = -(de_{1,1})_{2,1} & 0 \end{pmatrix}
\]

\[
de_{1,2} = \begin{pmatrix} (de_{2,2})_{1,2} & (de_{1,2})_{1,2} = -(de_{2,1})_{1,2} \\ 0 & -(de_{2,2})_{2,1} \epsilon_{1,2} \end{pmatrix}
\]

\[
de_{2,1} = \begin{pmatrix} (de_{2,1})_{1,2} & 0 \\ (de_{1,2})_{1,2} = -(de_{1,1})_{1,2} & \epsilon_{2,1}(de_{1,1})_{1,2} \end{pmatrix}
\]

From this follows that any section, \(\rho : \text{Ph}(M_2(k)) \to M_2(k)\), of \(i : M_2(k) \to \text{Ph}(M_2(k))\), is given in terms of an element \(\phi \in M_2(k)\), such that \(\rho(da) = [\phi, a]\).
The phase-space construction may, of course, be iterated. Given the \( k \)-algebra \( A \) we may form the sequence, \( \{ \text{Ph}^n(A) \}_{1 \leq n} \), defined inductively by

\[
\text{Ph}^0(A) = A, \quad \text{Ph}^1(A) = \text{Ph}(A), \quad \ldots, \quad \text{Ph}^{n+1}(A) := \text{Ph}(\text{Ph}^n(A)).
\]

Let \( i^n_0 : \text{Ph}^n(A) \to \text{Ph}^{n+1}(A) \) be the canonical imbedding, and let \( d_n : \text{Ph}^n(A) \to \text{Ph}^{n+1}(A) \) be the corresponding derivation. Since the composition of \( i^n_0 \) and the derivation \( d_{n+1} \) is a derivation \( \text{Ph}^n(A) \to \text{Ph}^{n+2}(A) \), there exist by universality a homomorphism \( i^{n+1}_1 : \text{Ph}^{n+1}(A) \to \text{Ph}^{n+2}(A) \), such that,

\[
d_n \circ i^{n+1}_1 = i^n_0 \circ d_{n+1}.
\]

Notice that we here compose functions and functors from left to right. Clearly we may continue this process constructing new homomorphisms,

\[
\{ i^j_n : \text{Ph}^n(A) \to \text{Ph}^{n+1}(A) \}_{0 \leq j \leq n},
\]

with the property,

\[
d_n \circ i^{j+1}_{n+1} = i^j_n \circ d_{n+1}.
\]

It is easy to see, [La 7], that,

\[
i^n_{p,q} = i^{n-1}_{q-1} i^1_p, \quad p < q
\]

\[
i^n_{q,p} = i^n_{p+1}
\]

\[
i^n_{p,q} = i^{n+1}_{q+1}, \quad q < p,
\]

i.e. the \( \text{Ph}^*A \) is a semi-cosimplicial algebra.

The system of \( k \)-algebras and homomorphisms of \( k \)-algebras \( \{ \text{Ph}^n(A), i^n_j \}_{n,0 \leq j \leq n} \) has an inductive (direct) limit, \( \text{Ph}^\infty A \), together with homomorphisms,

\[
i_n : \text{Ph}^n(A) \to \text{Ph}^\infty(A)
\]

satisfying,

\[
i^n_j \circ i_{n+1} = i_n, \quad j = 0, 1, \ldots, n.
\]

Moreover, the family of derivations, \( \{ d_n \}_{0 \leq n} \) define a unique derivation,

\[
\delta : \text{Ph}^\infty(A) \to \text{Ph}^\infty(A),
\]

such that,

\[
i_n \circ \delta = d_n \circ i_{n+1}.
\]

Put

\[
\text{Ph}^{(n)}(A) := \text{im} \ i_n \subseteq \text{Ph}^\infty(A)
\]

Let for any associative algebra \( B \), \( \text{Rep}(B) \) denote the category of \( B \)-modules. The set of isomorphism-classes of \( B \)-modules is just a set, and the map induced by the obvious forgetful functor,

\[
\omega : \text{Rep}(\text{Ph}^\infty(A)) \to \text{Rep}(A),
\]

is just a set-theoretical map, although having a well defined tangent map,

\[
T\omega : \text{Ext}^1(\text{Ph}^\infty(A))(V,V) \to \text{Ext}^1_A(V,V).
\]

As we shall see, assuming the algebra \( A \) of finite type, the set of simple finite dimensional \( A \)-modules form an algebraic scheme, \( \text{Simp}_{\leq \infty}(A) \). Moreover,
Preparation 1.2. The canonical morphism \( i_0 : A \rightarrow Ph^\infty(A) \) parametrizes simple representations of \( A \) with fixed momentum, acceleration, and any number of higher order momenta.

This should be understood in the following way. Consider, for any simple \( A \)-module \( V \), the exact sequence,

\[
0 \rightarrow \text{End}_A(V) \rightarrow \text{End}_k(V) \rightarrow \text{Der}_k(A, \text{End}_k(V)) \rightarrow \text{Ext}^1_A(V,V) \rightarrow 0.
\]

Let \( \rho : A \rightarrow \text{End}_k(V) \) define the structure of \( V \), then any morphism \( \rho_1 : Ph(A) \rightarrow \text{End}_k(V) \) extending \( \rho \), corresponds to a derivation, \( \xi_1 : A \rightarrow \text{End}_k(V) \), and therefore, via the maps in the exact sequence above, to a tangent vector, also called \( \text{End}_k(V) \), to an element \( \text{End}_k(V) \), and there-

fore, via the maps in the exact sequence above, to a tangent vector, also called \( \xi_1 \), in the tangent space of the point \( v \in \text{Simp}(A) \) corresponding to \( V \). So any such \( \xi_1 \) corresponds to a couple \((v,\xi)\) of a point in \( \text{Simp}(A) \), and an infinitesimal deformation of that point, i.e. a momentum. Any morphism \( \rho_2 : Ph^2(A) \rightarrow \text{End}_k(V) \) extending \( \rho_1 \) corresponds therefore to the triple \((v,\xi_1,\xi_2)\), corresponding to a point and a momentum, and to an infinitesimal deformation of this, etc. Since we have canonical morphisms \( Ph^r(A) \rightarrow Ph^\infty(A) \), it is clear that any morphism \( \xi : Ph^\infty(A) \rightarrow \text{End}_k(V) \) extending \( \rho \), produces a sequence, of any order, of such tuples. A simple consequence of the definition of \( Ph^\infty(A) \), that we identify all \( s^q,q = 1,\ldots,n \), shows that the set of such morphisms \( \xi \), extending a given structure-

morphism, parametrizes the set of formal curves in \( \text{Simp}(A) \) through \( v \).

A fundamental problem of (our model of) physics, see the Introduction, can now be stated as follows: If we \text{prepare} an object so that we know its momentum, and its higher order momenta up to a certain order, what can we infer on its behavior in the future?

In our mathematical model, a preparation made on the \( A \)-module, the object, \( V \), by fixing its structure as a \( Ph^\infty(A) \)-module, now forces it to change in the following way: The Dirac derivation \( \delta \in \text{Der}_k(Ph^\infty(A)) \) maps via the structure homomorphism of the module \( V \),

\[
\rho_V : Ph^\infty(A) \longrightarrow \text{End}_k(V)
\]

to an element \( \delta_V \in \text{Der}_k(Ph^\infty(A), \text{End}_k(V)) \) and via the composition of the canonical linear maps,

\[
\text{Der}_k(Ph^\infty(A), \text{End}_k(V)) \longrightarrow \text{Ext}^1_{Ph^\infty(A)}(V,V) \longrightarrow \text{Ext}^1_A(V,V)
\]

to the element \( \delta(V) \in \text{Ext}^1_A(V,V) \), i.e. to a tangent vector of \( \text{Simp}_n(A) \) at the point, \( v \), corresponding to \( V \), see [La 4].

Suppose first that, \( \delta(V) = 0 \). This means that \( \delta_V \) in \( \text{Der}_k(Ph^\infty(A), \text{End}_k(V)) \) is an inner derivation given by an endomorphism \( Q \in \text{End}_k(V) \), such that for every \( f \in Ph^\infty(A) \), we find \( \delta(f)(v) = (Qf - fQ)(v) \). This \( Q \) is the corresponding Hamiltonian, (or Dirac operator in the terminology of Connes, see [Schücker]), and we have a situation that is very much like classical quantum mechanics, i.e. a set-

up where the objects are represented by a fixed Hilbert space \( V \) and an algebra of observables \( Ph^\infty(A) \) acting on it, with time, and therefore also energy, represented by a special Hamiltonian operator \( Q \).

A system characterized by a \( Ph^\infty(A) \)-module \( V \), for which \( \delta(V) = 0 \), will be called \text{stable} or \text{singular}. It is said to be in the \text{state} \( \psi \) if we have chosen an
element $\psi \in V$. The Dirac derivation $\delta$ defines a Hamiltonian operator $Q$, (a Dirac operator), and time, i.e. $\delta$, now push the state $\psi$ into the state,

$$\exp(\tau Q)(\psi) \in V,$$

corresponding to the isomorphism of the module $V$ defined by the inner isomorphism of the algebra of observables, $\Phi^\infty(A)$ defined by $U := \exp(\tau \delta)$, whenever this is well defined. This is a well known situation in classical quantum mechanics, corresponding to the equivalence between the set-ups of Schrödinger and Heisenberg.

To treat the situation when $[\delta] \neq 0$, we first have to take a new look at non-commutative algebraic geometry, as developed in [La 3,4,5].

§2. Non-commutative deformations and the structure of the moduli space of simple representations

In [La 2], [La 3] and [La 4,5], we introduced non-commutative deformations of families of modules of non-commutative $k$-algebras, and the notion of swarm of right modules (or more generally of objects in a $k$-linear abelian category). Let for any associative $k$-algebra $S$, $a_S$ be the category of $S$-valued associative $k$-algebras, the objects of which are the diagrams of $k$-algebras,

$$S \xrightarrow{\iota} R \xrightarrow{\pi} S$$

such that the composition of $\iota$ and $\pi$ is the identity. In particular, $a_r$ denotes the category of $r$-pointed $k$-algebras, i.e. $a_{k^r}$, with $S = k^r$. Any such $r$-pointed $k$-algebra $R$ is isomorphic to a $k$-algebra of $r \times r$-matrices $(R_{i,j})$.

The radical of $R$ is the bilateral ideal $\text{Rad}(R) := \ker\pi$, such that $R/\text{Rad}(R) \cong k^r$. The dual $k$-vector space of $\text{Rad}(R)/\text{Rad}(R)^2$ is called the tangent space of $R$.

For $r = 1$, there is an obvious inclusion of categories

$$\mathfrak{l} \subseteq a_1$$

where $\mathfrak{l}$, as usual, denotes the category of commutative local Artinian $k$-algebras with residue field $k$.

Fix a not necessarily commutative $k$-algebra $A$ and consider a right $A$-module $M$. The ordinary deformation functor

$$\text{Def}_M : \mathfrak{a} \rightarrow \text{Sets}$$

is then defined. Assuming $\text{Ext}_A^i(M, M)$ has finite $k$-dimension for $i = 1, 2$, it is well known, see [Sch], or [La 2], that $\text{Def}_M$ has a pro-representing hull $H$, the formal moduli of $M$. Moreover, the tangent space of $H$ is isomorphic to $\text{Ext}_A^1(M, M)$, and $H$ can be computed in terms of $\text{Ext}_A^i(M, M)$, $i = 1, 2$ and their matrix Massey products, see [La 1, 2], [La 7].

In the general case, consider a finite family $\mathcal{V} = \{V_i\}_{i=1}^r$ of right $A$-modules. Assume that,

$$\dim_k \text{Ext}_A^1(V_i, V_j) < \infty.$$

Any such family of $A$-modules will be called a swarm. We shall define a deformation functor,

$$\text{Def}_\mathcal{V} : \mathfrak{a} \rightarrow \text{Sets}$$
generalizing the functor \( Def_M \) above. Given an object \( \pi : R = (R_{i,j}) \to k^r \) of \( \mathcal{A}_r \), consider the \( k \)-vector space and left \( R \)-module \((R_{i,j} \otimes_k V_j)\). It is easy to see that \( \text{End}_R((R_{i,j} \otimes_k V_j)) \cong (R_{i,j} \otimes_k \text{Hom}_k(V_i, V_j)) \). Clearly \( \pi \) defines a \( k \)-linear and left \( R \)-linear map, \( \pi : (R_{i,j} \otimes_k V_j) \to \bigoplus_{i=1}^r V_i \), inducing a homomorphism of \( R \)-endomorphism rings, \( \tilde{\pi}(R) : (R_{i,j} \otimes_k \text{Hom}_k(V_i, V_j)) \to \bigoplus_{i=1}^r \text{End}_k(V_i) \).

The right \( A \)-module structure on the \( V_i \)'s is defined by a homomorphism of \( k \)-algebras, \( \eta_0 : A \to \bigoplus_{i=1}^r \text{End}_k(V_i) \). Let

\[
\text{Def}_V(R) \in \text{Sets}
\]

be the set of isoclasses of homomorphisms of \( k \)-algebras,
\[
\eta' : A \to (R_{i,j} \otimes_k \text{Hom}_k(V_i, V_j))
\]

such that,
\[
\tilde{\pi}(R) \circ \eta' = \eta_0,
\]

where the equivalence relation is defined by inner automorphisms in the \( k \)-algebra \((R_{i,j} \otimes_k \text{Hom}_k(V_i, V_j))\) inducing the identity on \( \bigoplus_{i=1}^r \text{End}_k(V_i) \). One easily proves that \( \text{Def}_V \) has the same properties as the ordinary deformation functor and we prove the following, see [La 1, 2], [La 7], Theorem 2.1.

**Theorem 2.1.** The functor \( \text{Def}_V \) has a pro-representable hull, i.e. an object \( H \) of the category of pro-objects \( \mathcal{A}_r \), together with a versal family,

\[
\tilde{V} = \{(H_{i,j} \otimes V_j)\}_{i=1}^r \in \varprojlim_{n \geq 1} \text{Def}_V(H/m^n),
\]

where \( m = \text{Rad}(H) \), such that the corresponding morphism of functors on \( \mathcal{A}_r \),

\[
\kappa : \text{Mor}(H, -) \to \text{Def}_V
\]

defined for \( \phi \in \text{Mor}(H, R) \) by \( \kappa(\phi) = R \otimes_\phi \tilde{V} \), is smooth, and an isomorphism on the tangent level. Moreover, \( H \) is uniquely determined by a set of matric Massey products defined on subspaces,

\[
D(n) \subseteq \text{Ext}^1(V_i, V_{j_1}) \otimes \cdots \otimes \text{Ext}^1(V_{j_{n-1}}, V_k),
\]

with values in \( \text{Ext}^2(V_i, V_k) \). The right action of \( A \) on \( \tilde{V} \) defines a homomorphism of \( k \)-algebras, the versal family,

\[
\eta : A \to O(V) := \text{End}_H(\tilde{V}) = (H_{i,j} \otimes \text{Hom}_k(V_i, V_j)),
\]

and the \( k \)-algebra \( O(V) \) acts on the family of \( A \)-modules \( \mathcal{V} = \{V_i\} \), extending the action of \( A \). If \( \dim_k V_i < \infty \), for all \( i = 1, \ldots, r \), the operation of associating \( (O(V), \mathcal{V}) \) to \((A, \mathcal{V})\) is a closure operation.

Moreover, we prove the crucial result,
A generalized Burnside theorem 2.2. Let $A$ be a finite dimensional $k$-algebra, $k$ an algebraically closed field. Consider the family $V = \{V_i\}_{i=1}^r$ of all simple $A$-modules, then

$$\eta: A \longrightarrow O(V) = (H_{i,j} \otimes \text{Hom}_k(V_i, V_j))$$

is an isomorphism.

Using this theorem, we may generalize the above Theorem 2.1. In fact, let $S$ be any finite dimensional $k$-algebra, and consider the category $\mathfrak{a}_S$. An object is a diagram,

$$S \rightarrow R \rightarrow_{\pi} S$$

where $R$ is a finite dimensional $k$-algebra, and the composition $\iota \pi = \text{id}$. Morphisms are the obvious commutative diagrams. Suppose we are given a homomorphism of $k$-algebras,

$$\rho_0: A \rightarrow \text{End}_S(V),$$

where $V$ is finite dimensional $k$-vector space. For any object, $S \rightarrow R \rightarrow_{\pi} S \in \mathfrak{a}_S$, we may consider the left $R$-module $R \otimes_S V$, and the homomorphism of $k$-algebras,

$$\pi_*: \text{End}_R(R \otimes_S V) \rightarrow \text{End}_S(V).$$

The same arguments as above, then proves the more general result,

**The O-construction Theorem.** Let $V$ be any left $S$-and right $A$-module, then the functor,

$$\text{Def}_{(SVA)}: \mathfrak{a}_S \rightarrow \text{Sets},$$

defined by,

$$\text{Def}_{(SVA)}: (R) = \{\rho \in \text{Mor}_k(A, \text{End}_R(R \otimes_S V)|\phi \pi_* = \rho_0\}.$$

This functor has a pro-representable hull, i.e. an object $S \rightarrow H \rightarrow_{\pi} S$, of the category of pro-objects $\hat{\mathfrak{a}}_S$ of $\mathfrak{a}_S$, together with a versal family,

$$\tilde{V} \in \lim_{n \geq 1} \text{Def}_{SVA}(H/m^n),$$

where $m = \text{Rad}(H) := \ker \hat{\pi}$, such that the corresponding morphism of functors on $\mathfrak{a}_S$,

$$\kappa: \text{Mor}(H, -) \rightarrow \text{Def}_V$$

defined for $\phi \in \text{Mor}(H, R)$ by $\kappa(\phi) = R \otimes_{\hat{\mathfrak{a}}} \tilde{V}$, is "smooth", and an isomorphism on the "tangent level". Moreover, $H$ is uniquely determined by a set of matric Massey products with values in $\text{Ext}^2(V_i, V_k)$. The right action of $A$ on $\tilde{V}$ defines a homomorphism of $k$-algebras, the versal family,

$$\eta: A \longrightarrow O(SVA) := \text{End}_H(\tilde{V}) = (H_{i,j} \otimes \text{Hom}_k(V_i, V_j)).$$

The $k$-algebra $O := O(SVA)$ acts on the module $V$, extending the right-action of $A$, commuting with the $S$-action. If $\dim_k V_i < \infty$, for all $i = 1, \ldots, r$, the operation of associating $(O(SVA), SV_O)$ to $(A, SVA)$, is a closure operation.
Notice that the proof of the closure property of the $O$-construction, as proposed in [La 4], is incomplete, and should be replaced with the above. Details will occur in a projected book.

We also prove that there exists, in the non-commutative deformation theory, an obvious analogy to the notion of pro-representing (modular) substratum $H_0$, see [La 0] and [La-Pf]. The tangent space of $H_0$ is determined by a family of subspaces

$$\text{Ext}_A^1(V_i, V_j) \subseteq \text{Ext}_A^1(V_i, V_j), \quad i \neq j$$

the elements of which should be called the almost split extensions (sequences) relative to the family $V$, and by a subspace,

$$T_0(\Delta) \subseteq \prod_i \text{Ext}_A^1(V_i, V_i)$$

which is the tangent space of the deformation functor of the full subcategory of the category of $A$-modules generated by the family $V = \{V_i\}_{i=1}^r$, see [La 1]. If $V = \{V_i\}_{i=1}^r$ is the set of all indecomposables of some Artinian $k$-algebra $A$, we show that the above notion of almost split sequence coincides with that of Auslander, see [R].

In [La 2], we consider the general problem of classification of iterated extensions of a family of modules $V = \{V_i\}_{i=1}^r$, and the corresponding classification of filtered modules with graded components in the family $V$, and extension type given by a directed representation graph $\Gamma$, see §3. The main result is the following result, see [La 4], Proposition 2.3.

**Proposition 2.3.** Let $A$ be any $k$-algebra, $V = \{V_i\}_{i=1}^r$ any swarm of $A$-modules, i.e. such that,

$$\dim_k \text{Ext}_A^1(V_i, V_j) < \infty \quad \text{for all } i, j = 1, \ldots, r.$$

(i): Consider an iterated extension $E$ of $V$, with representation graph $\Gamma$. Then there exists a morphism of $k$-algebras

$$\phi : H(V) \to k[\Gamma]$$

such that

$$E \simeq k[\Gamma] \otimes_{\phi} \tilde{V}$$

as right $A$-modules.

(ii): The set of equivalence classes of iterated extensions of $V$ with representation graph $\Gamma$, is a quotient of the set of closed points of the affine algebraic variety

$$\mathbb{A}[\Gamma] = \text{Mor}(H(V), k[\Gamma])$$

(iii): There is a versal family $\tilde{V}[\Gamma]$ of $A$-modules defined on $\mathbb{A}[\Gamma]$, containing as fibers all the isomorphism classes of iterated extensions of $V$ with representation graph $\Gamma$.

Let $\text{Mod}_V^\Gamma$ denote the full subcategory of $\text{Mod}_A$ generated by the iterated extensions of $V$. As usual denote by $H := H(V)$ the formal moduli of $V$. Then we have the following structure theorem, generalizing a result, of Beilinson,
Theorem. Let $A$ be any $k$-algebra, and fix a swarm, $V = \{V_i\}_{i=1}^r$, of $A$-modules, then there exists a functor,
\[
\kappa : \text{Mod}_{H(V)} \to \text{Mod}^V_A
\]
which is an isomorphism on equivalence-classes of objects, and monomorphic on morphisms. If $V$ consists of simple modules, then $\kappa$ is an equivalence.

Proof. Any right $H(V)$-module $M_i$ is a $k^r$-module, so it can be decomposed as $M = \oplus M_i$, where $M_i := M_ι_i$. The structure map is therefore given as,
\[
\rho_0 : H(V) \to \text{End}_k(M) = (\text{Hom}_k(M_i, M_j)).
\]
Here, $\rho_0$ maps $H_{i,j}$ into $\text{Hom}_k(V_i, V_j)$, and therefore the formal family may be composed to give us the following $k$-algebra homomorphism,
\[
\rho : A \to (H_{i,j} \otimes \text{Hom}_k(V_i, V_j)) \to (\text{Hom}_k(M_i, M_j) \otimes_k \text{Hom}_k(V_i, V_j))
\]
\[
= (\text{Hom}_k(M_i \otimes_k V_i, M_j \otimes V_j)) = \text{End}_k(W).
\]
Here $W = \oplus_{i=1}^r (M_i \otimes V_i)$, and $k(M) := W$. Since $k[Γ] \subset \text{End}_k(P)$, where $P$ is $k[Γ]$, as $k[Γ]$-module, the first part of the theorem follows from Proposition 2.3. The rest is more or less evident.

□

To any, not necessarily finite, swarm $\mathcal{C} \subset \text{mod}(A)$ of right $A$-modules, we may use the above $O$-construction, to associated to $\mathcal{C}$ a $k$-algebra, see [La 3] and [La 4],
\[
O(\mathcal{C}, π) = \lim_{\mathcal{C} \subset \mathcal{C}_n} O(sVA),
\]
where $S = k[\mathcal{C}_0]$, is the $k$-algebra of the quiver associated to $\mathcal{C}_0$, where $\{V_i\}_{i=1}^r = |\mathcal{C}_0|$, and where $V := \oplus_i^r V_i$. There is a natural $k$-algebra homomorphism,
\[
η(\mathcal{C}) : A \to O(\mathcal{C})
\]
with the property that the $A$-module structure on $\mathcal{C}$ is extended to an $O$-module structure in an optimal way. We then defined an affine non-commutative scheme of right $A$-modules to be a swarm $\mathcal{C}$ of right $A$-modules, such that $η(\mathcal{C})$ is an isomorphism. In particular we considered, for finitely generated $k$-algebras, the swarm $\text{Simp}_{<\infty}(A)$ consisting of the finite dimensional simple $A$-modules, and the generic point $A$, together with all morphisms between them. The fact that this is a swarm, i.e. that for all objects $V_i, V_j \in \text{Simp}_{<\infty}$ we have $\dim_k \text{Ext}^1_{A}(V_i, V_j) < \infty$, is easily proved. We have in [La 4] proved the following result, (see (4.1), loc.cit.)

Proposition 2.4. Let $A$ be a geometric $k$-algebra, then the natural homomorphism,
\[
η(\text{Simp}^r(A)) : A \to O_π(\text{Simp}^r_{<\infty}(A))
\]
is an isomorphism, i.e. $\text{Simp}^r_{<\infty}(A)$ is a scheme for $A$.

In particular, $\text{Simp}^r_{<\infty}(k < x_1, x_2, ..., x_d >)$, is a scheme for $k < x_1, x_2, ..., x_d >$. To analyze the local structure of $\text{Simp}_n(A)$, we need the following, see [La 4],(3.23),
Lemma 2.5. Let \( \mathcal{V} = \{ V_i \}_{i=1,...,r} \) be a finite subset of \( \text{Simp}_{<\infty}(A) \), then the morphism of \( k \)-algebras,

\[
A \to O(\mathcal{V}) = (H_{i,j} \otimes_k \text{Hom}_k(V_i, V_j))
\]

is topologically surjective.

Proof. Since the simple modules \( V_i, i = 1, ..., r \) are distinct, there is an obvious surjection, \( \eta_0 : A \to \prod_{i=1,..,r} \text{End}_k(V_i) \). Put \( \tau = \ker \eta_0 \), and consider for \( m \geq 2 \) the finite-dimensional \( k \)-algebra, \( B := A/\tau^m \). Clearly \( \text{Simp}(B) = \mathcal{V} \), so that by the generalized Burnside theorem, see [La 2], (2.6), we find, \( B \cong O_B(\mathcal{V}) := (H_{i,j}^B \otimes_k \text{Hom}_k(V_i, V_j)) \). Consider the commutative diagram,

\[
\begin{array}{ccc}
A & \longrightarrow & (H_{i,j}^A \otimes_k \text{Hom}_k(V_i, V_j)) =: O^A(\mathcal{V}) \\
\downarrow & & \downarrow \\
B & \longrightarrow & (H_{i,j}^B \otimes_k \text{Hom}_k(V_i, V_j)) \longrightarrow O^A(\mathcal{V})/\mathfrak{m}^m
\end{array}
\]

where all morphisms are natural. In particular \( \alpha \) exists since \( B = A/\tau^m \) maps into \( O^A(\mathcal{V})/\text{rad}^m \), and therefore induces the morphism \( \alpha \) commuting with the rest of the morphisms. Consequently \( \alpha \) has to be surjective, and we have proved the contention.

\( \square \)

Localization, topology and the scheme structure on \( \text{Simp}(A) \). Let \( s \in A \), and consider the open subset \( D(s) = \{ V \in \text{Simp}(A) | \rho(s) \text{ invertible in } \text{End}_k(V) \} \). The Jacobson topology on \( \text{Simp}(A) \) is the topology with basis \( \{ D(s) | s \in A \} \). It is clear that the natural morphism,

\[
\eta : A \to O(D(s), \pi)
\]

maps \( s \) into an invertible element of \( O(D(s), \pi) \). Therefore we may define the localization \( A_{(s)} \) of \( A \), as the \( k \)-algebra generated in \( O(D(s), \pi) \) by \( \text{im } \eta \) and the inverse of \( \eta(s) \). This furnishes a general method of localization with all the properties one would wish. And in this way we also find a canonical (pre)sheaf, \( O \) defined on \( \text{Simp}(A) \).

Definition 2.6. When the \( k \)-algebra \( A \) is geometric, such that \( \text{Simp}^*(A) \) is a scheme for \( A \), we shall refer to the presheaf \( O \), defined above on the Jacobson topology, as the structure presheaf of the scheme \( \text{Simp}(A) \).

We shall now see that the Jacobson topology on \( \text{Simp}(A) \), restricted to each \( \text{Simp}^n(A) \) is the Zariski topology for a classical scheme-structure.

Recall first that a standard \( n \)-commutator relation in a \( k \)-algebra \( A \) is a relation of the type,

\[
[a_1, a_2, ..., a_{2n}] := \sum_{\sigma \in \Sigma_{2n}} \text{sign}(\sigma)a_{\sigma(1)}a_{\sigma(2)}...a_{\sigma(2n)} = 0
\]
where \( \{a_1, a_2, ..., a_{2n}\} \) is a subset of \( A \). Let \( I(n) \) be the two-sided ideal of \( A \) generated by the subset,

\[
\{[a_1, a_2, ..., a_{2n}] | \{a_1, a_2, ..., a_{2n}\} \subset A \}.
\]

Consider the canonical homomorphism,

\[
p_n : A \rightarrow A/I(n) =: A(n).
\]

It is well known that any homomorphism of \( k \)-algebras,

\[
\rho : A \rightarrow \text{End}_k(k^n) =: M_n(k)
\]
factors through \( p_n \), see e.g. [Formanek].

**Corollary 2.7.** (i). Let \( V_i, V_j \in \text{Simp}_{\leq n}(A) \) and put \( r = m_{V_i} \cap m_{V_j} \). Then we have, for \( m \geq 2 \),

\[
\text{Ext}^1_A(V_i, V_j) \simeq \text{Ext}^1_{A/rm}(V_i, V_j)
\]

(ii). Let \( V \in \text{Simp}_{n}(A) \). Then,

\[
\text{Ext}^1_A(V, V) \simeq \text{Ext}^1_{A(n)}(V, V)
\]

**Proof.** (i) follows directly from Lemma (2.5). To see (ii), notice that \( \text{Ext}^1_A(V, V) \simeq \text{Der}_k(A, \text{End}_k(V))/\text{Triv} \simeq \text{Der}_k(A(n), \text{End}_k(V))/\text{Triv} \simeq \text{Ext}^1_{A(n)}(V, V) \). The second isomorphism follows from the fact that any derivation maps a standard \( n \)-commutator relation into a sum of standard \( n \)-commutator relations.

\[\square\]

**Example 2.8.** Notice that, for distinct \( V_i, V_j \in \text{Simp}_{\leq n}(A) \), we may well have,

\[
\text{Ext}^1_A(V_i, V_j) \neq \text{Ext}^1_{A(n)}(V_i, V_j).
\]

In fact, consider the matrix \( k \)-algebra,

\[
A = \begin{pmatrix} k[x] & k[x] \\ 0 & k[x] \end{pmatrix},
\]

and let \( n = 1 \). Then \( A(1) = k[x] \oplus k[x] \). Put \( V_1 = k[x]/(x) \oplus (0), V_2 = (0) \oplus k[x]/(x) \), then it is easy to see that,

\[
\text{Ext}^1_A(V_i, V_j) = k, \ \text{Ext}^1_{A(1)}(V_i, V_j) = 0, i \neq j,
\]

but,

\[
\text{Ext}^1_A(V_i, V_i) = \text{Ext}^1_{A(1)}(V_i, V_i) = k, i = 1, 2.
\]
**Lemma 2.9.** Let $B$ be a $k$-algebra, and let $V$ be a vector space of dimension $n$, such that the $k$-algebra $B \otimes \text{End}_k(V)$ satisfies the standard $n$-commutator-relations, i.e. such that the ideal, $I(n) \subset B \otimes \text{End}_k(V)$ generated by the standard $n$-commutators $[x_1, x_2, ..., x_n]$, $x_i \in B \otimes \text{End}_k(V)$, is zero. Then $B$ is commutative.

*Proof.* In fact, if $b_1, b_2 \in B$ is such that $[b_1, b_2] \neq 0$, then the obvious $n$-commutator,

$$(b_1 e_{1,1})(b_2 e_{1,1}) e_{1,2} e_{2,2} ... e_{n-1,n} e_{n,n} = (b_2 e_{1,1})(b_1 e_{1,1}) e_{1,2} e_{2,2} ... e_{n-1,n} e_{n,n}$$

is different from 0. Here $e_{i,j}$ is the $n \times n$ matrix with all elements equal to 0, except the one in the $(i,j)$ position, where the element is equal to 1.

\(\Box\)

**Lemma 2.10.** If $A$ is a finite type $k$-algebra, then any $V \in \text{Simp}_n(A)$ is an $A(n)$-module. Let $V \subset \text{Simp}_n(A)$ be a finite family, then $H^{A(n)}(V)$ is commutative. In particular,

1. $\text{Ext}_{A(n)}(V, V_j) = 0$, for $V_i \neq V_j$
2. $H^{A(n)}(V) \simeq H^{A}(V)^{\text{com}} := H(V)/[H(V), H(V)]$.

*Proof.* Since

$$A(n) \rightarrow O(V) \simeq M_n(H^{A(n)}(V))$$

is topologically surjective, we find using (Lemma 2.9), that $H^{A(n)}(V)$ is commutative. This implies (1) and the commutativity of $H^{A(n)}(V)$. Consider for $V \in \text{Simp}_n(A)$, the natural commutative diagram of homomorphisms of $k$-algebras,

$$
\begin{array}{ccc}
Z(A(n)) & \longrightarrow & A(n) \\
\downarrow & \searrow & \downarrow \alpha \\
H(V)^{\text{com}} & \longrightarrow & H(V)^{\text{com}} \otimes_k \text{End}_k(V)
\end{array}
$$

where $Z(A(n))$ is the center of $A(n)$. The existence of $\alpha$ is a consequence of the ideal $I(n)$ of $A$ mapping to zero in $H(V)^{\text{com}} \otimes_k \text{End}_k(V) \simeq M_n(H(V)^{\text{com}})$. Therefore there are natural morphisms of formal moduli,

$$H^{A}(V) \rightarrow H^{A(n)}(V) \rightarrow H^{A}(V)^{\text{com}} \rightarrow H^{A(n)}(V)^{\text{com}}.$$

Since $H^{A(n)}(V) = H^{A(n)}(V)^{\text{com}}$ the composition,

$$H^{A(n)}(V) \rightarrow H^{A}(V)^{\text{com}} \rightarrow H^{A(n)}(V)^{\text{com}},$$

must be an isomorphism. Since, by Corollary (1.12), the tangent spaces of $H^{A(n)}(V)$ and $H^{A}(V)$ are isomorphic, the lemma is proved.

\(\Box\)
Corollary 2.11. Let $A = k < x_1, \ldots, x_d >$ be the free $k$-algebra on $d$ symbols, and let $V \in \text{Simp}_n(A)$. Then

$$H^A(V)_{\text{com}} \simeq H^{A(n)}(V) \simeq k[[t_1, \ldots, t_{(d-1)n^2+1}]]$$

This should be compared with the results of [Procesi 1], see also [Formanek]. In general, the natural morphism,

$$\eta(n) : A(n) \to \prod_{V \in \text{Simp}_n(A)} H^{A(n)}(V) \otimes_k \text{End}_k(V)$$

is not an injection, as it follows from the following,

Example 2.12. Let

$$A = \begin{pmatrix} k & k & k \\ k & k & k \\ 0 & 0 & k \end{pmatrix}.$$

The ideal $I(2)$ is generated by $[e_{1,1}, e_{1,2}, e_{2,2}, e_{2,3}] = e_{1,3}$. So

$$A(2) = \begin{pmatrix} k & k & k \\ k & k & k \\ 0 & 0 & k \end{pmatrix} / \begin{pmatrix} 0 & 0 & k \\ 0 & 0 & k \\ 0 & 0 & 0 \end{pmatrix} \simeq M_2(k) \oplus M_1(k).$$

However,

$$\prod_{V \in \text{Simp}_2(A)} H^{A(2)}(V) \otimes_k \text{End}_k(V) \simeq M_2(k),$$

therefore $\ker \eta(2) = M_1(k) = k$.

Let $O(n)$, be the image of $\eta(n)$, then,

$$O(n) \subseteq \prod_{V \in \text{Simp}_n(A)} H^{A(n)}(V) \otimes_k \text{End}_k(V)$$

and for every $V \in \text{Simp}_n(A)$,

$$H^{O(n)}(V) \simeq H^{A(n)}(V).$$

Put $B = \prod_{V \in \text{Simp}_n(A)} H^{A(n)}(V)$. Choosing bases in all $V \in \text{Simp}_n(A)$, then

$$\prod_{V \in \text{Simp}_n(A)} H^{A(n)}(V) \otimes_k \text{End}_k(V) \simeq M_n(B),$$

Let $x_i \in A, i = 1, \ldots, d$ be generators of $A$, and consider their images $(x_{i,p,q}) \in M_n(B)$. Now, $B$ is commutative, so the $k$-sub-algebra $C(n) \subset B$ generated by the elements $(x_{i,p,q})_{i=1,\ldots,d; \ p,q=1,\ldots,n}$ is commutative. We have an injection,

$$O(n) \to M_n(C(n)),$$
and for all $V \in \text{Simp}_n(A)$, with a chosen basis, there is a natural composition of homomorphisms of $k$-algebras,

$$\alpha : M_n(C(n)) \to M_n(H^{A(n)}(V)) \to \text{End}_k(V),$$

inducing a corresponding composition of homomorphisms of the centers,

$$Z(\alpha) : C(n) \to H^{A(n)}(V) \to k$$

This sets up a set theoretical injective map,

$$t : \text{Simp}_n(A) \to \text{Max}(C(n)),$$

defined by $t(V) := \ker Z(\alpha)$.

Since $A(n) \to H^{A(n)}(V) \otimes_k \text{End}_k(V)$ is topologically surjective, $H^{A(n)}(V) \otimes_k \text{End}_k(V)$ is topologically generated by the images of $x_i$, $i = 1, ..., d$. It follows that we have a surjective homomorphism,

$$\hat{C}(n)_{t(V)} \to H^{A(n)}(V).$$

Categorical properties implies, that there is another natural morphism,

$$H^{A(n)}(V) \to \hat{C}(n)_{t(V)},$$

which composed with the former is an automorphism of $H^{A(n)}(V)$. Since

$$M_n(C(n)) \subseteq \prod_{V \in \text{Simp}_n(A)} H^{A(n)}(V) \otimes_k \text{End}_k(V),$$

it follows that for $m_v \in \text{Max}(C(n))$, corresponding to $V \in \text{Simp}_n(A)$, the finite dimensional $k$-algebra $M_n(C(n)/m_v^2)$ sits in a finite dimensional quotient of some,

$$\prod_{V \in \mathcal{V}} H^{A(n)}(V) \otimes_k \text{End}_k(V).$$

where $\mathcal{V} \subset \text{Simp}_n(A)$ is finite. However, by Lemma (2.5), the morphism,

$$A(n) \to \prod_{V \in \mathcal{V}} H^{A(n)}(V) \otimes_k \text{End}_k(V)$$

is topologically surjective. Therefore the morphism,

$$A(n) \to M_n(C(n)/m_v^2)$$

is surjective, implying that the map

$$H^{A(n)}(V) \to \hat{C}(n)_{m_v},$$

is surjective, and consequently, $H^{A(n)}(V) \simeq \hat{C}(n)_{m_v}$.

We now have the following theorem, see Chapter VIII, §2, of the book [Procesi 2], where part of this theorem is proved.
Theorem 2.13. Let $V \in \text{Simp}_n(A)$, correspond to the point $m_n \in \text{Simp}_1(C(n))$.

(i) There exist a Zariski neighborhood $U_v$ of $v$ in $\text{Simp}_1(C(n))$ such that any closed point $m'_n \in U$ corresponds to a unique point $V' \in \text{Simp}_n(A)$.

Let $U(n)$ be the open subset of $\text{Simp}_1(C(n))$, the union of all $U_v$ for $V \in \text{Simp}_n(A)$.

(ii) $O(n)$ defines a non-commutative structure sheaf $\mathcal{O}(n) := \mathcal{O}_{U(n)}$ of Azumaya algebras on the topological space $U(n)$ (Jacobson topology).

(iii) The center $S(n)$ of $O(n)$, defines a scheme structure on $\text{Simp}_n(A)$.

(iv) The versal family of $n$-dimensional simple modules, $\tilde{V} := (C(n)) \otimes_k V$, over $\text{Simp}_n(A)$, is defined by the morphism,

$$\tilde{\rho} : A \to O(n) \subseteq \text{End}_{C(n)}(C(n)) \otimes_k V \simeq M_n(C(n)).$$

(v) The trace ring $\text{Tr} \tilde{\rho} \subseteq S(n) \subseteq C(n)$ defines a commutative affine scheme structure on $\text{Simp}_n(A)$. Moreover, there is a morphism of schemes,

$$\kappa : U(n) \longrightarrow \text{Simp}_n(A),$$

such that for any $v \in U(n)$,

$$H^A(n)(V) \simeq \tilde{S}(n)v \simeq (\text{Tr} \tilde{\rho})_{\kappa(v)} \simeq \tilde{C}(n)v$$

Proof. Let $\rho : A \longrightarrow \text{End}_k(V)$ be the surjective homomorphism of $k$-algebras, defining $V \in \text{Simp}_n(A)$. Let, as above $e_{i,j} \in \text{End}_k(V)$ be the elementary matrices, and pick $y_{i,j} \in A$ such that $\rho(y_{i,j}) = e_{i,j}$. Let us denote by $\sigma$ the cyclical permutation of the integers $\{1,2,\ldots,n\}$, and put,

$$s_k := [y_{\sigma^k(1),\sigma^k(2)}, y_{\sigma^k(2),\sigma^k(2)}, y_{\sigma^k(2),\sigma^k(3)} \cdots y_{\sigma^k(n),\sigma^k(n)}], \quad s := \sum_{k=0,1,\ldots,n-1} s_k \in A.$$

Clearly $s \in I(n - 1)$. Since $[e_{\sigma^k(1),\sigma^k(2)}, e_{\sigma^k(2),\sigma^k(2)}, e_{\sigma^k(2),\sigma^k(3)} \cdots e_{\sigma^k(n),\sigma^k(n)}] = e_{\sigma^k(1),\sigma^k(n)} \in \text{End}_k(V)$, $\rho(s) := \sum_{k=0,1,\ldots,n-1} \rho(s_k) \in \text{End}_k(V)$ is the matrix with non-zero elements, equal to 1, only in the $(\sigma^k(1),\sigma^k(n))$ position, so the determinant of $\rho(s)$ must be $+1$ or $-1$. The determinant $\det(s) \in C(n)$ is therefore nonzero at the point $v \in \text{Spec}(C(n))$ corresponding to $V$. Put $U = D(\det(s)) \subseteq \text{Spec}(C(n))$, and consider the localization $O(n)_{\{s\}} \subseteq M_n(C(n)_{\{\det(s)\}})$, the inclusion following from general properties of the localization. Now, any closed point $v' \in U$ corresponds to a $n$-dimensional representation of $A$, for which the element $s \in I(n - 1)$ is invertible. But then this representation cannot have a $m < n$ dimensional quotient, so it must be simple.

Since $s \in I(n - 1)$, the localized $k$-algebra $O(n)_{\{s\}}$ does not have any simple modules of dimension less than $n$, and no simple modules of dimension $> n$. In fact, for any finite dimensional $O(n)_{\{s\}}$-module $V$, of dimension $m$, the image $\tilde{s}$ of $s$ in $\text{End}_k(V)$ must be invertible. However, the inverse $\tilde{s}^{-1}$ must be the image of a polynomial (of degree $m - 1$) in $s$. Therefore, if $V$ is simple over $O(n)_{\{s\}}$, i.e. if the homomorphism $O(n)_{\{s\}} \rightarrow \text{End}_k(V)$ is surjective, $V$ must also be simple over $A$. Since now $s \in I(n - 1)$, it follows that $m \geq n$. If $m > n$, we may construct, in the same way as above an element in $I(n)$ mapping into a nonzero element of $\text{End}_k(V)$. Since, by construction, $I(n) = 0$ in $A(n)$, and therefore also in $O(n)_{\{s\}}$,
we have proved what we wanted. By a theorem of M. Artin, see [Artin], \( O(n) \) must be an Azumaya algebra with center, \( S(n) := Z(O(n)) \). Therefore \( O(n) \) defines a presheaf \( O(n) \) on \( U(n) \), of Azumaya algebras with center \( S(n) := Z(O(n)) \).

Clearly, any \( V \in Simp_n(A) \), corresponding to \( m \in Max(C(n)) \) maps to a point \( s(v) \in Simp(O(n)) \). Let \( m_{s(v)} \) be the corresponding maximal ideal of \( S(n) \). Since \( O(n) \) is locally Azumaya, it follows that,

\[
\hat{S}(n)_{m_{s(v)}} \cong H^{O(n)}(V) \cong H^{A(n)}(V).
\]

The rest is clear.

\[ \square \]

\( Spec(C(n)) \) is, in a sense, a compactification of \( Simp_n(A) \). It is, however, not the correct completion of \( Simp_n(A) \). In fact, the points of \( Spec(C(n)) \) \( - \) \( Simp_n(A) \) may correspond to semi-simple modules, which do not deform into simple \( n \)-dimensional modules. We shall return to the study of the (notion of) completion, together with the degeneration processes that occur, at infinity in \( Simp_n(A) \).

**Example 2.14.** Let us check the case of \( A = k < x_1, x_2 > \), the free non-commutative \( k \)-algebra on two symbols. First, we shall compute \( Ext^1_A(V, V) \) for a particular \( V \in Simp_2(A) \), and find a basis \( \{ t^i_+ \}_{i=1}^5 \), represented by derivations \( \partial_i := \partial_i(V) \in Der_k(A, \text{End}_k(V)) \), \( i = 1, 2, 3, 4, 5 \). This is easy, since for any two \( A \)-modules \( V_1, V_2 \), we have the exact sequence,

\[
0 \to \text{Hom}_A(V_1, V_2) \to \text{Hom}_k(V_1, V_2) \to \text{Der}_k(A, \text{Hom}_k(V_1, V_2)) \\
\to \text{Ext}^1_A(V_1, V_2) \to 0
\]

proving that, \( \text{Ext}^1_A(V_1, V_2) = \text{Der}_k(A, \text{Hom}_k(V_1, V_2))/\text{Triv} \), where \( \text{Triv} \) is the sub-vector space of trivial derivations. Pick \( V \in Simp_2(A) \) defined by the homomorphism \( A \to M_2(k) \) mapping the generators \( x_1, x_2 \) to the matrices

\[
X_1 := \begin{pmatrix} 0 & 1 \\ 0 & 0 \end{pmatrix} =: e_{1,2}, \quad X_2 := \begin{pmatrix} 0 & 0 \\ 1 & 0 \end{pmatrix} =: e_{2,1}.
\]

Notice that

\[
X_1X_2 = \begin{pmatrix} 1 & 0 \\ 0 & 0 \end{pmatrix} =: e_{1,1} = e_1, \quad X_2X_1 = \begin{pmatrix} 0 & 0 \\ 0 & 1 \end{pmatrix} =: e_{2,2} = e_2,
\]

and recall also that for any \( 2 \times 2 \)-matrix \( (a_{p,q}) \in M_2(k) \), \( e_i(a_{p,q})e_j = a_{i,j}e_{i,j} \). The trivial derivations are generated by the derivations \( \{ \delta_{p,q} \}_{p,q=1,2} \), defined by,

\[
\delta_{p,q}(x_i) = X_i e_{p,q} - e_{p,q} X_i.
\]

Clearly \( \delta_{1,1} + \delta_{2,2} = 0 \). Now, compute and show that the derivations \( \partial_i, i = 1, 2, 3, 4, 5 \), defined by,

\[
\partial_i(x_1) = 0, \text{ for } i = 1, 2, \quad \partial_i(x_2) = 0, \text{ for } i = 4, 5,
\]

by,

\[
\partial_1(x_2) = e_{1,1}, \quad \partial_2(x_2) = e_{1,2}, \quad \partial_3(x_1) = e_{1,2}, \quad \partial_4(x_1) = e_{2,2}, \quad \partial_5(x_1) = e_{2,1}
\]
and by,
\[ \partial_1(x_2) = e_{2,1}, \]
form a basis for \( \text{Ext}^1_k(V, V) = \text{Der}_k(A, \text{End}_k(V))/\text{Triv}. \) Since \( \text{Ext}^2_k(V, V) = 0 \) we find \( H(V) = k \ll t_1, t_2, t_3, t_4, t_5 >\) and so \( H(V)^\text{com} \simeq k[[t_1, t_2, t_3, t_4, t_5]]. \) The formal versal family \( \hat{V} \), is defined by the actions of \( x_1, x_2 \), given by,
\[ X_1 := \begin{pmatrix} 0 & 1 + t_3 \\ t_5 & t_4 \end{pmatrix}, \quad X_2 := \begin{pmatrix} t_1 & t_2 \\ 1 + t_3 & 0 \end{pmatrix}. \]
One checks that there are polynomials of \( X \) which are equal to \( t_1 e_{p,q} \), modulo the ideal \( (t_1, \ldots, t_5)^2 \subset H(V) \), for all \( i, p, q = 1, 2 \). This proves that \( \hat{C}(2) \) must be isomorphic to \( H(V) \), and that the composition,
\[ A \longrightarrow A(2) \longrightarrow M_2(C(2)) \subset M_2(H(V)) \]
is topologically surjective. By the construction of \( C(n) \) this also proves that
\[ C(2) \simeq k[t_1, t_2, t_3, t_4, t_5]. \]
locally in a Zariski neighborhood of the origin. Moreover, the Formanek center, in this case is cut out by the single equation:
\[ f := \det[X_1, X_2] = -(1 + t_3)^2 - t_2 t_5)^2 + (t_1(1 + t_3) + t_2 t_4)(t_4(1 + t_3) + t_1 t_5). \]
Computing, we also find the following formulas,
\[ \text{tr}X_1 = t_4, \quad \text{tr}X_2 = t_1, \]
\[ \det X_1 = -t_5 - t_3 t_5, \quad \det X_2 = -t_2 - t_2 t_3, \]
\[ \text{tr}(X_1 X_2) = (1 + t_3)^2 + t_2 t_5 \]
so the trace ring of this family is
\[ k[t_1, t_2 + t_2 t_3, 1 + 2 t_3 + t_3^2 + t_2 t_5, t_4, t_5 + t_3 t_5] =: k[u_1, u_2, \ldots, u_6], \]
with,
\[ u_1 = t_1, \quad u_2 = (1 + t_3) t_2, \quad u_3 = (1 + t_3)^2 + t_2 t_5, \quad u_4 = t_4, \quad u_5 = (1 + t_3) t_5, \]
and \( f = -u_3^2 + 4 u_2 u_5 + u_1 u_3 u_4 + u_1^2 u_5 + u_2 u_4^2. \) Moreover, the \( k[u] \) is algebraic over \( k[u_1] \), with discriminant, \( \Delta := 4 u_2 u_5 (u_3^2 - 4 u_2 u_5) = 4(1 + t_3)^2 t_2 t_5((1 + t_3)^2 - t_2 t_5)^2. \) From this follows that there is an étale covering
\[ \mathbb{A}^5 - V(f \Delta) \rightarrow \text{Simp}_2(A) - V(\Delta). \]
Notice that if we put \( t_1 = t_4 = 0 \), then \( f = \Delta. \) See the Example (3.7)
Completions of $\text{Simp}_n(A)$. In the example above it is easy to see that elements of the complement of $\text{Simp}_n(A)$ in the affine sub-scheme $\text{Spec}(C(n))$ will be represented by indecomposable, or decomposable representations. A decomposable representation $W$ will, however, not in general be deformable into a simple representation, since good deformations must conserve $\text{End}_A(W)$. Therefore, even though we have termed $\text{Spec}(C(n))$ a compactification of $\text{Simp}_n(A)$, it is a bad completion. The missing points at infinity of $\text{Simp}_n(A)$, should be represented as indecomposable representations, with $\text{End}_A(W) = k$. Any such is an iterated extension of simple representations $\{V_i\}_{i=1,2,..,s}$, with representation graph $\Gamma$ (corresponding to an extension type, see [La 4]), and $\sum_{i=1}^s \dim(V_i) = n$. To simplify the notations we shall write, $[\Gamma] := \{V_i\}_{i=1,2,..,s}$. In [La 2,4], see also [Jö-La-Sl], we treat the problem of classifying all such indecomposable representations, up to isomorphisms. Let us recall the main ideas.

Assume that the simple modules $\{V_i\}$ we shall talk about are such that all $\text{Ext}_A^r(V_i, V_j)$ are finite dimensional as $k$-vector spaces. Let $\Gamma$ be an ordered graph with set of nodes $|\Gamma| = |V_i|$. Starting with the the first node of $\Gamma$, we can construct, in many ways, an extension of the module corresponding to the end point of the first arrow of $\Gamma$, then continue, choosing an extension of the result with the module corresponding to the endpoint of the second arrow of $\Gamma$, etc. untill we have reached the endpoint of the last arrow. Any finite length module can be made in this way, the "oppositely ordered" $\Gamma$ corresponding to a decomposition of the module into simple constituencies, by peeling off one simple sub-module at a time, i.e. by picking one simple sub-module and forming the quotient, picking a second simple sub-module of the quotient and taking the quotient, and repeating the procedure untill it stops.

The "ordered" $k$-algebra $k[\Gamma]$ of the ordered graph $\Gamma$ is the quotient algebra of the usual algebra of the graph $\Gamma$ by the ideal generated by all admissible words which are not "intervals" of the ordered graph. Say $\ldots \gamma_{i,j}(n-1)\gamma_{j,k}(n)\gamma_{j,k}(n+1)\ldots$ is an interval of the ordered graph, then $\gamma_{i,j}(n-1)\gamma_{j,k}(n+1) = 0$ in $k[\Gamma]$.

Now, let $H([\Gamma])$ be the formal moduli of the family $[\Gamma]$. We show in [La 4], see Proposition 2. above, that any iterated extension of the $\{V_i\}_{i=1}$ with extension type, i.e. graph, $\Gamma$ corresponds to a morphism in $\underline{\alpha}$,

$$\alpha : H \rightarrow k[\Gamma].$$

Moreover the set of isomorphism classes of such modules is parametrized by a quotient space of the affine scheme,

$$\underline{A}(\Gamma) := \underline{\text{Mor}}_{\underline{\alpha}}(H([\Gamma]), k[\Gamma]).$$

Let $\alpha \in \underline{A}(\Gamma)$, and let $V(\alpha)$ denote the corresponding iterated extension module, then the tangent space of $\underline{A}(\Gamma)$ at $\alpha$ is,

$$T_{\underline{A}(\Gamma),\alpha} := \text{Der}_k(H([\Gamma]), k[\Gamma]_\alpha),$$

where $k[\Gamma]_\alpha$ is $k[\Gamma]$ considered as a $H([\Gamma])$-bimodule via $\alpha$. The obstruction space for the deformation functor of $\alpha$ is $HH^2(H([\Gamma]), k[\Gamma])$, and we may, as is explained in [La 0,1], compute the complete local ring of $\underline{A}(\Gamma)$ at $\alpha$. In particular we may decide whether the point is a smooth point of $\underline{A}(\Gamma)$, or not.
The automorphism group $G$ of $k[\Gamma]$, considered as an object of $\mathfrak{a}_r$, has a Lie algebra which we shall call $\mathfrak{g}$. Obviously we have,

$$\mathfrak{g} = \text{Der}_k(k[\Gamma], k[\Gamma]).$$

Clearly an iterated extension $\alpha$ with graph $\Gamma$ will be isomorphic as $A$-module to $\mathfrak{g}(\alpha)$, for any $\mathfrak{g} \in G$. In particular, if $\delta \in \mathfrak{g}$, then $\exp(\delta)(\alpha)$ is isomorphic to $\alpha$ as an iterated extension of $A$-modules, with the same graph as $\alpha$.

Consider the map,

$$\alpha^* : \text{Der}_k(k[\Gamma], k[\Gamma]) \to \text{Der}_k(H(\Gamma), k[\Gamma]).$$

The image of $\alpha^*$ is the subspace of the tangent space of $A(\Gamma)$ at $\alpha$ along which the corresponding module has constant isomorphism class.

Notice that if $\alpha$ is a smooth point, and $\alpha^*$ is not surjectiv then there is a positive-dimensional moduli space of iterated extension modules with graph $\Gamma$ through $\alpha$.

Clearly the kernel of $\alpha^*$ is contained in the Lie algebra of automorphisms of the module $\mathfrak{g}(\alpha)$, and should be contained in $\text{End}_A(\mathfrak{g}(\alpha))$. From this follows that if $\mathfrak{g}(\alpha)$ is indecomposable then $\ker \alpha^* = 0$. The Euler type derivations, defined by,

$$\delta_E(\gamma_{i,j}) = \rho_{i,j} \gamma_{i,j}, \rho_{i,j} \in k$$

are the easiest to check! Notice however, that there may be discrete automorphisms in $G$, not of exponential type, leaving $\alpha$ invariant. Notice also that an indecomposable module may have an endomorphism-ring which is a non-trivial local ring.

Assume now that we have identified the non-commutative scheme of indecomposable $\Gamma$-representation, call it $\text{Ind}_\Gamma(A)$. Put $\text{Simp}_{\Gamma}(A) := \text{Simp}_n(A) \cup \text{Ind}_{\Gamma}(A)$. Now, repeat the basics of the construction of $\text{Spec}(C(n))$ above. Consider for every open affine subscheme $D(s) \subset \text{Simp}_{\Gamma}(A)$, the natural morphism,

$$A \to \lim_{\varnothing \subset D(s)} O(\varnothing, \pi)$$

$\varnothing$ running through all finite subsets of $D(s)$. Put $B_s(\Gamma) := \prod_{V \in D(s)} H^A(n)(V)^{\text{com}},$ and consider the homomorphism,

$$A \to A(n) \to \prod_{V \in D(s)} H^A(n)(V)^{\text{com}} \otimes_k \text{End}_k(V) \simeq M_n(B_s(\Gamma)).$$

Let $x_i \in A, i = 1, ..., d$ be generators of $A$, and consider the images $(x^i_{p,q}) \in B_s(n) \otimes_k \text{End}_k(k^n)$ of $x_i$ via the homomorphism of $k$-algebras,

$$A \to B_s(\Gamma) \otimes M_n(k),$$

obtained by choosing bases in all $V \in \text{Simp}_{\Gamma}(A)$. Notice that since $V$ no longer is (necessarily) simple, we do not know that this map is topologically surjectiv. Now, $B_s(\Gamma)$ is commutative, so the $k$-sub-algebra $C_s(\Gamma) \subset B_s(\Gamma)$ generated by the elements $(x^i_{p,q})_{i=1, ..., d; p,q=1, ..., n}$ is commutative. We have a morphism,

$$I_s(\Gamma) : A \to C_s(\Gamma) \otimes_k M_n(k) = M_n(C_s(\Gamma)).$$
Moreover, these $C_s(\Gamma)$ define a presheaf, $C(\Gamma)$, on the Jacobson topology of $\text{Simp}_T(A)$. The rank $n$ free $C_s(\Gamma)$-modules with the $A$-actions given by $I_s(\Gamma)$, glue together to form a locally free $C(\Gamma)$-Module $E(\Gamma)$ on $\text{Simp}_T(A)$, and the morphisms $I_s(\Gamma)$ induce a morphism of algebras,

$$I(\Gamma) : A \to \text{End}_{C(\Gamma)}(E(\Gamma)).$$

As for every $V \in \text{Simp}_T(A)$, $\text{End}_A(V) = k$, the commutator of $A$ in $H^A(V)^{\text{com}} \otimes_k \text{End}_k(V)$ is $H^A(V)^{\text{com}}$. The morphism,

$$\zeta(V) : H^A(V)^{\text{com}} \to HH^0(A, H^A(V)^{\text{com}} \otimes_k \text{End}_k(V))$$

is therefore an isomorphism, and we may assume that the corresponding morphism,

$$\zeta : C(\Gamma) \to HH^0(A, \text{End}_{C(\Gamma)}(E(\Gamma)))$$

is an isomorphism of sheaves. For all $V \in D(s) \subset \text{Simp}_T(A)$ there is a natural projection,

$$\kappa := \kappa(\Gamma) : C_s(\Gamma) \otimes_k M_n(k) \to H^{A(n)}(V)^{\text{com}} \otimes_k \text{End}_k(V) \cong M_n(H^{A(n)}(V)^{\text{com}}),$$

which, composed with $I_s(\Gamma)$ is the natural homomorphism,

$$A \longrightarrow H^{A(n)}(V)^{\text{com}} \otimes_k \text{End}_k(V).$$

$\kappa$ defines a set theoretical map,

$$t : \text{Simp}_T(A) \longrightarrow \text{Spec}(C(\Gamma)),$$

and a natural surjective homomorphism,

$$\hat{C}(\Gamma)_{t(V)} \rightarrow H^{A(n)}(V)^{\text{com}}.$$

Categorical properties implies, as usual, that there is another natural morphism,

$$\iota : H^{A(n)}(V) \rightarrow \hat{C}(\Gamma)_{t(V)},$$

which composed with the former is the obvious surjection, and such that the induced composition,

$$A \longrightarrow H^{A(n)}(V)^{\text{com}} \otimes_k \text{End}_k(V) \rightarrow \hat{C}(\Gamma)_{t(V)} \otimes_k \text{End}_k(V),$$

is $I(\Gamma)$ formalized at $t(V)$. From this, and from the definition of $C(\Gamma)$, it follows that $\iota$ is surjective, such that for every $V \in \text{Simp}_T(A)$ there is an isomorphism $H^{A(n)}(V)^{\text{com}} \cong \hat{C}(\Gamma)_{t(V)}$. For $V \in \text{Simp}_T(A)$ there is also a natural commutative diagram,

$$\begin{array}{ccc}
Z A(n) & \longrightarrow & C(\Gamma) \\
| & & | \\
A(n) & \longrightarrow & \text{End}_{C(\Gamma)}(E(\Gamma)) \\
| & & | \\
H^{A(n)}(V) \otimes_k \text{End}_k(V) & \longrightarrow & \hat{C}(\Gamma)_{t(V)} \otimes_k \text{End}_k(V)
\end{array}$$

Formally at a point $V \in \text{Simp}_T(A)$, we have therefore proved that the local, commutative structure of $\text{Simp}_T(A)$ (as $A$ or $A(n)$-module), and the corresponding local structure of $\text{Spec}(C(\Gamma))$ at $V$, coincide. We have actually proved the following,
Theorem 2.15. The topological space $\text{Simp}_\Gamma(A)$, with the Jacobson topology, together with the sheaf of commutative $k$-algebras $\mathcal{C}(\Gamma)$ defines a scheme structure on $\text{Simp}_\Gamma(A)$, containing an open subscheme, étale over $\text{Simp}_n(A)$. Moreover, there is a morphism,

$$\pi(\Gamma): \text{Simp}_\Gamma(A) \to \text{Spec}(\mathcal{Z}A(n)),$$

extending the natural morphism,

$$\pi_0 : \text{Simp}_n(A) \to \text{Spec}(\mathcal{Z}A(n)).$$

Proof. As in Theorem (2.13) we prove that if $v = t(V)$, with $V \in \text{Simp}_n(A) \subseteq \text{Simp}_\Gamma(A)$, then there exists an open subscheme of $\text{Spec}(\mathcal{C}(\Gamma))$ containing only simple modules of dimension $n$. If $v$ is indecomposables with $\text{End}_A(V) = k$ we may look at the homomorphism of $\mathcal{C}(\Gamma)$-modules,

$$\text{End}_A(\mathcal{C}(\Gamma)) \otimes \text{End}_k(V) \longrightarrow \text{End}_A(V) = k.$$

Clearly there is an open neighborhood of $v$ in $\text{Spec}(\mathcal{C}(\Gamma))$ containing only indecomposables of dimension $n$.

These morphisms $\pi(\Gamma)$ are our candidates for the possibly different completions of $\text{Simp}_n(A)$. Notice that for $W \in \text{Spec}(\mathcal{C}(n)) - U_n$, the formal module $H^A(W)$ is not always pro-representing. If $W$ is semi-simple, but not simple then $\text{End}_A(W) \neq k$. The corresponding modular substratum will, locally, correspond to the semi-simple deformations of $W$, thus to a closed subscheme of $\text{Spec}(\mathcal{C}(n)) - U_n \subseteq \text{Spec}(\mathcal{C}(n))$. This follows from the fact that the substratum of modular deformations of any semisimple (but not simple) module $V$ will have a tangent space equal to the invariant space of the action of the $\text{End}_k(V)$ on $\text{Ext}^1_A(V, V)$, which must be the sum of the tangent spaces of the deformation spaces of the simple components of $V$.

$\text{Spec}(\mathcal{C}(n))$ is, in a sense, a compactification of $U_n$. It is, however not the correct completion of $U_n$. In fact, the points of $\text{Spec}(\mathcal{C}(n)) - U_n$ may correspond to semi-simple modules, which do not deform into simple $n$-dimensional modules. We shall in §4 return to the study of the (notion of) completion, in connection with the process of decay and creation of particles. Decay occur, at infinity, in $\text{Simp}_n(A)$, see the Introduction.

Morphisms, Hilbert schemes, Fields and Strings. Above we have studied moduli spaces of representations of finitely generated $k$-algebras. We might as well have studied the Hilbert functor, $\mathcal{H}_{A^r}$, of subschemes of length $r$ of the spectrum of the algebra $A$, or the moduli space $\mathcal{F}(A; R)$, of morphisms, $\kappa : A \to R$, for fixed algebras, $A$ and $R$. The difference is that whereas for finite $n$, the set $\text{Simp}_n(A)$ has a nice, finite dimensional scheme structure, this, in general, no longer true for the set, $\mathcal{H}_{A^r}$ nor for the set of fields, $\mathcal{F}(A; R)$, as the physicists call it, unless we put some extra conditions on the fields, so called decorations.

If $R$ is Artinian of length $n$, then the corresponding $\mathcal{F}(A; R)$ does exist and has a nice structure, both as commutative and as non-commutative scheme. The toy model of relativity theory, referred to in the introduction, is in fact modeled on $\mathcal{M}(k[x_1, x_2, x_3], k^2)$, i.e. on the set of surjective homomorphisms $k[x_1, x_2, x_3] \to R = k^2$. And, in all generality, the "space" $\mathcal{F}(A; R)$ has a tangent structure. I fact,
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depending on the point of view, the tangent space of a morphism \( \phi : A \to R \) is equal to,

\[
T_{\text{CalF}(A,R),\phi} = \text{Der}_k(A,R)/\text{Triv},
\]

where Triv either is 0 or the inner derivations induced by \( R \). Even though there is no obvious algebraic structure on \( \mathcal{F}(A;R) \) this general situation is important. It is the basis for Quantum Field Theory, as we shall see, in the next chapter. There it will be treated in combination with the notion of "clock" and/or string. Let us start here with the last notion.

**Definition 2.16.** A general string, a g-string, is an algebra \( R \) together with a pair of \( \text{Ph-points} \), i.e. a pair of homomorphisms \( \epsilon_i : \text{Ph}R \to k(p_i) \), corresponding to two points \( k(p_i) \in \text{Simp}_1(R) \) each outfitted with a tangent \( \xi_i \).

We might have considered any two points \( k(p_i) \in \text{Simp}_n(\text{Ph}R) \), but since the main properties of the g-strings will be equally well understood restricting to the case \( n = 1 \), we shall postpone this generalization.

For any g-string, consider the non-commutative tangent space of the pair of points,

\[
T(R,p_1,p_2) := \text{Ext}^1_{\text{Ph}R}(p_1,p_2).
\]

We shall call it the space of tensions, between the two points of the string.

Consider the space \( \text{String}_g(A) \) of \( g \)-strings in \( A \), i.e. the space of isomorphism classes of algebra homomorphisms \( \kappa : A \to R \) where \( R \) is a g-string, and where isomorphisms should correspond to isomorphisms of the g-string, thus conserving the two \( \text{Ph}R \)-points.

Any g-string in \( A, \kappa : A \to R \), induces a unique commutative diagram of algebras,

\[
\begin{array}{ccc}
A & \xrightarrow{\kappa} & R \\
| & | & | \\
\text{Ph}A & \xrightarrow{\text{Ph}\kappa} & \text{Ph}R
\end{array}
\]

The von Neumann condition imposed on a string \( \kappa \), is now the following,

\[
\epsilon_i \circ \text{Ph}\kappa \circ d = \kappa \ast \xi_i =: \xi_i = 0, \ i = 1 \lor i = 2,
\]

which, if \( x_j, j = 1, \ldots, n \) and \( \sigma_l, l = 1, \ldots, p \) are parameters of \( A \) respectively \( R \), is equivalent to the condition,

\[
\frac{\partial x_j}{\partial \sigma_l}(p_i) = 0, \ j = 1, \ldots, n, \ l = 1, \ldots, p, \ i = 1 \lor 2.
\]

Notice also that, since any derivation \( \xi \in \text{Der}_k(A,R) \) has a natural lifting to a derivation \( \bar{\xi} \in \text{Der}_k(\text{Ph}A,\text{Ph}R) \) defined by simply putting \( \bar{\xi}(a) = d(\xi(a)) \), we find, using the general machinery of deformations of diagrams, see [La 0], that any family of morphisms \( \kappa \) induces a family of the above diagram. If \( \tau_k, k = 1, \ldots, d \) are parameters of such a family, \( M = \text{Spec}(M) \), then \( d\tau_i \in \text{Ph}M \) corresponds to a derivation, \( \tau_i \in \text{Der}_k(A,R) \), and therefore to tangents \( \bar{\xi}_i, i = 1,2 \), of \( \text{Simp}_1(A) \) at the two points \( k(p_i) \). The Dirichlet condition on the string is now,

\[
\bar{\xi}_i = 0, \ i = 1 \lor 2,
\]
which is equivalent to the condition,

\[ \frac{\partial x_j}{\partial t_i}(p_i) = 0, j = 1, \ldots, n, \ t = 1, \ldots, p, \ i = 1 \lor 2. \]

These conditions will define new moduli spaces which we shall call \( \text{String}^N_R(A) \) and \( \text{String}^R(A) \), respectively. In the affine case the structure of these spaces is a problem, however we may of course do everything above for \( A \) and \( R \) replaced by projective schemes, and then all the moduli spaces exist as classical schemes. The volume form of the space the string is fanning out will give us an action functional, \( S \), defined on \( \text{String}_R(A) \), see next §.

Let us end this § by noticing that there is a non-commutative deformation theory for fields, just as there is one for representations of associative algebras. In fact, let \( \{ \kappa_i \}_{i=1, \ldots, r} \) be a finite family of fields, and consider for every pair \( (i, j)|1 \leq i, j \leq r \) the \( A \)-bimodule \( R_{i,j} \) where \( \kappa_i \) defines the left module structure, and \( \kappa_j \) the right hand structure. Then copying the definition for the non-commutative deformation functor for representations, replacing \( \text{Hom}_k(V_i, V_j) \) by \( R_{i,j} \), we may prove most of the results referred to at the beginning of this §. This may be of interest in relation with the problems of interactions treated in §4.

Example 2.17. (i) Let us go back to Example (1.1)(ii). It follows that the string of dimension 0, \( R = k^2 \), \( \text{Ph}(R) = k < x, dx > / ((x^2 - r^2), (xdx + dx^2)) \), has unique points, \( k(\pm r) \). The space of tensions is of dimension 1, the von Neumann condition is automatically satisfied, and the moduli space of \( k^2 \)-strings in \( A = k[x_1, x_2, x_3] \) is nothing but \( H := \text{Spec}k[t_1, \ldots, t_6] - \Delta \). If we consider the string with \( R = k[x]/(x^2), \text{Ph}R = k[x, dx]/(x^2, (xdx + dx^2)) \), then we see that there is just one point of \( R \), a line of point for \( \text{Ph}R \), all corresponding to \( x = 0 \) in \( R \). Therefore there is a 2-dimensional space of strings with the same \( R \). Compare this with the blow-up \( \bar{H} \), see [La 6]. (ii) In dimension 1 the simplest closed string is given by, \( R = k[x, y]/(f) \), with \( f = x^2 + y^2 - r^2 \), such that \( \text{Ph}R = k < x, y, dx, dy > / (f, [x, y], d[x, y], df) \), and with the two points, \( \xi_i : \text{Ph}R \to k(p_i) \), defined by the actions on \( k(p_i) := k \), given by, \( x_i, y_i, (dx)_i, (dy)_i, \ i = 1, 2 \). It is easy to see that the vectors, \( \xi_i := ((dx)_i, (dy)_i) \) are tangent vectors to the circle at the points \( p_i \), and if \( p_1 \neq p_2 \) we find that \( \text{Ext}^1_{\text{Ph}R}(k(p_1), k(p_2)) = k \). The von Neumann condition is, \( \xi_i = 0, \ i = 1 \lor i = 2 \), and this clearly means that \( \frac{\partial x}{\partial \sigma} = \frac{\partial y}{\partial \sigma} = 0 \) at one of the points \( p_i \). The 1-dimensional open string is now left as an exercise.

§3. GEOMETRY OF TIME-SPACES AND THE GENERAL DYNAMICAL LAW

Given a finitely generated \( k \)-algebra, and a natural number \( n \), we have in §2 constructed a scheme \( \text{Simp}_n(A) \), and a versal family,

\[ \tilde{\rho} : A \to \text{End}_{U(n)}(\bar{V}) \]

defined on an étale covering \( U(n) \) of \( \text{Simp}_n(A) \). \( U(n) \) is an open subscheme of an affine scheme \( \text{Spec}(C(n)) \), and the versal family is, in fact, defined on \( C(n) \).

Dynamical Systems.

We would now like to use this theory for the \( k \)-algebra \( \text{Ph}^\infty(A) \) of §1. However, \( \text{Ph}^\infty(A) \) is rarely of finite type. We shall therefore introduce the notion of dynamical structure, and the order of a dynamical structure, to reduce the problem to a
situation we can handle. This is also what physicists do, they invoke a parsimony principle, or an Action Principle, originally proposed by Fermat, and later by Mauerpertuis, with exactly this purpose, reducing the preparation needed to be able to see ahead, see (1.2).

**Definition 3.1.** A dynamical structure, $\sigma$, is a two-sided $\delta$-stable ideal $(\sigma) \subset Ph^\infty(A)$, such that

$$A(\sigma) := Ph^\infty(A)/(\sigma),$$

the corresponding, dynamical system, is of finite type. A dynamical structure, or system, is of order $n$ if the canonical morphism,

$$\sigma : Ph^{(n-1)}(A) \to A(\sigma)$$

is surjective. If $A$ is generated by the coordinate functions, $\{t_i\}_{i=1,2,\ldots,d}$, any dynamical system of order $n$ is defined by a Force Law, i.e. a system of equations,

$$\delta^n t_p = \Gamma^p(t_i, dt_j, d^2 t_k, \ldots, d^{n-1} t_l), \quad p = 1, 2, \ldots, d.$$

Put,

$$A(\sigma) := Ph^\infty(A)/(\delta^n t_p - \Gamma^p)$$

where $\sigma := (\delta^n t_p - \Gamma^p)$ is the two-sided $\delta$-ideal generated by the defining equations of $\sigma$. Obviously $\delta$ induces a derivation $\delta_\sigma \in Der_k(A(\sigma))$, also called the Dirac derivation, and usually just denoted $\delta$.

Notice that if $\sigma_i, \ i = 1, 2$, are two order $n$ dynamical systems, then we may well have,

$$A(\sigma_1) \simeq A(\sigma_2) \simeq Ph^{(n-1)}(A)/(\sigma_\sigma),$$

as $k$-algebras, see the Introduction.

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For any integer $n \geq 1$ consider the schemes $Simp_n(A(\sigma))$ and $Spec(C(n))$, and the corresponding (almost uni-) versal family,

$$\tilde{\rho} : A(\sigma) \rightarrow \text{End}_{Spec(C(n))}(\tilde{V}) \simeq M_n(C(n)).$$

The Dirac derivation $\delta \in Der_k(A(\sigma), A(\sigma))$, composed with $\rho$, corresponding to any element $v \in Simp_n(A(\sigma))$, defines, as explained above in §1, a tangent vector of $Simp_n(A(\sigma))$ at the point $v$, thus a distribution on $Simp_n(A(\sigma))$. The reason why the Dirac derivation, does not define a unique vector-field is, of course, that the structure morphisms of the simple modules can be scaled by any non-zero element of the field $k$. However, once we have chosen a versal family for the moduli space $Simp_n(A(\sigma))$, defined on $Spec(C(n))$, the Dirac derivation $\delta$ induces, by composition with $\tilde{\rho}$, an element,

$$\tilde{\delta} \in Der_k(A(\sigma), \text{End}_{C(n)}(\tilde{V})).$$

which obviously induces a well defined vector field $\xi \in \Theta_{U(n)}$, in the distribution defined by $\delta$. Now, to any vectorfield $\eta$ of $Spec(C(n))$, i.e. for any derivation $\eta \in Der_k(C(n))$, there is a unique element,

$$\eta' \in Der_k(A(\sigma), \text{End}_{C(n)}(\tilde{V})).$$
defined by,

\[ \eta'(a) = \eta(\tilde{\rho}(a)), \]

where we have identified \( \tilde{\rho}(a) \) with an element of \( M_n(C(n)) \). Suppose there exist a (rational) derivation \( [\delta] \in \text{Der}_k(C(n)) \), lifting the vector field \( \xi \) in \( \text{Simp}_n(\mathbf{A}((\sigma))) \), defined by the Dirac derivation, then by construction of \( C \), and of the versal family, \( \tilde{\rho} \), in general, we find that \( \delta - [\delta]^\pi \) is an inner derivation, defined by some \( Q \in \text{End}_{C(n)}(\hat{V}) \).

This is the situation which we shall find ourselves in, in the sequel, see the examples (3.5-3.7).

In general, we have the fundamental result:

\textbf{Theorem 3.2.} Formally, at any point \( v \in U(n) \), with local ring \( \text{C}(\hat{n})_v \), there is a derivation \( [\delta] \in \text{Der}_k(C(n)_v) \), and an Hamiltonian \( Q \in \text{End}_{C(\hat{n})}_v(\hat{V}_v) \), such that, as operators on \( \hat{V}_v \), we must have,

\[ \delta = [\delta] + [Q, -]. \]

This means that for every \( a \in \mathbf{A}((\sigma)) \), considered as an element \( \tilde{\rho}(a) \in M_n(C(\hat{n})_v) \), \( \delta(a) \) acts on \( \hat{V}_v \) as

\[ \tilde{\rho}(\delta(a)) = \xi(\tilde{\rho}(a)) + [Q, \tilde{\rho}(a)]. \]

\textbf{Proof.} Suppose the family,

\[ \tilde{\rho} : \mathbf{A}((\sigma)) \to \text{End}_{\text{Spec}(C(n))}(\hat{V}) \simeq M_n(C(n)). \]

had been the universal family of a fine moduli space. Then any (infinitesimal) automorphism of \( \mathbf{A}((\sigma)) \) would have been squared off by an (infinitesimal) automorphism of \( \text{End}_{\text{Spec}(C(n))}(\hat{V}) = M_n(C(n)) \). Now, any derivation of \( M_n(C(n)) \) is a sum of a derivation of \( C(n) \) plus an inner derivation \( \text{ad}(Q), Q \in M_n(C(n)) \), and we are through. Since our space \( U(n) \subset \text{Spec}(C(n)) \) is an étale covering of the fine moduli space \( \text{Simp}_n(\mathbf{A}((\sigma))) \), and since our versal family is only defined over \( U(n) \), we need to restrict to the formal case. Since \( C(\hat{n})_v \simeq \hat{O}_{\text{Simp}_n(\mathbf{A}((\sigma)))} \), this case is clear by general deformation theory.

\[ \square \]

As pointed out above, in the examples that we shall meet in the sequel, there are local (or even global) extensions of this result, where \( [\delta] \) and \( Q \) may be assumed to be defined (rationally) on \( C(n) \).

This \( Q \), the Hamiltonian of the system, is in the singular case, when \( [\delta] = 0 \), also called the Dirac operator, and sometimes denoted \( \delta \)-slashed, see e.g. [Schücker], or other texts on Connes’ spectral triples. In fact, a spectral triple is composed of a vector space like \( \hat{V} \), together with a Dirac operator, like \( Q \), and a complexification etc.

If \( [\delta] = 0 \), it is also easy to see that what we have observed implies that Heisenberg’s and Schrödinger’s way of doing quantum mechanics, are strictly equivalent.

In line with our general philosophy, we shall consider \( \xi \), or \( [\delta] \) as measuring time in \( \text{Simp}_n(\mathbf{A}((\sigma))) \), respectively in \( \text{Spec}(C(n)) \).

Assume for a while that \( k = \mathbb{R} \), the real numbers, and that our constructions go through, as if \( k \) were algebraically closed. Let \( \upsilon(\bar{\tau}) \in \text{Simp}_n(\mathbf{A}((\sigma))) \) be an element,
an event. Suppose there exist an integral curve $\gamma$ of $\xi$ through $v(\tau_0) \in \text{Simp}_1(C(n))$, ending at $v(\tau_1) \in \text{Simp}_1(C(n))$, given by the automorphisms $e(\tau) := \exp(\tau \xi)$, for $\tau \in [\tau_0, \tau_1] \subset \mathbb{R}$. The maximum $\tau$ for which the end point, $\xi$, of $\gamma$ is in $\text{Simp}_n(A(\sigma))$ should be called the lifetime of the particle. We shall see that it is relatively easy to compute these lifetimes, when the fundamental vector field $\xi$ has been computed.

This, however, is certainly not so easy, see the examples (3.4)-(3.8).

Let now $\psi(\tau_0) \in \hat{V}(v_0) \simeq V$ be a (classically considered) state of our quantum system, at the time $\tau_0$, and consider the (uni)versal family, $\tilde{\rho} : A(\sigma) \longrightarrow \text{End}_{C(n)}(\hat{V})$ restricted to $U(n) \subseteq \text{Simp}_1(C(n))$, the étale covering of $\text{Simp}_n(A(\sigma))$. We shall consider $A(\sigma)$ as our ring of observables.

What happens to $\psi(\tau_0) \in \hat{V}(0)$ when time passes from $\tau_0$ to $\tau$, along $\gamma$? This is obviously a question that has to do with whether we choose to consider the Heisenberg or the Schrödinger picture. In fact, if we consider the formal flow $\exp(t \delta)$ defined on the ring of observables, then putting,

$$u(\tau) := \exp(\tau \nabla_\xi),$$

where,

$$\nabla_\xi := \xi + Q \in \text{End}_k(\hat{V}),$$

we obtain for every $\psi \in \hat{V}$, and every $a \in A(\sigma)$, that the equation,

$$u(\tau)(\tilde{\rho}(\exp(-\tau \delta)(a))(\psi)) = \tilde{\rho}(a)(u(\tau)(\psi))$$

holds formally, at least up to first order. In fact, up to order one, in $\tau$, the left hand side is equal to

$$\tilde{\rho}(a)(\psi) - \tau \tilde{\rho}(\delta(a)(\psi)) + \tau \xi(\tilde{\rho}(a)(\psi)) + \tau Q \tilde{\rho}(a)(\psi),$$

and the right hand side is,

$$\tilde{\rho}(a)(\psi) + \tau \tilde{\rho}(a)(\xi(\psi)) + \tau \tilde{\rho}(a)(Q(\psi)).$$

Noticing that $\xi(\tilde{\rho}(a)(\psi)) = \xi(\tilde{\rho}(a))(\psi) + \tilde{\rho}(a)(\xi(\psi))$, and using (3.2) we find that the two sides are equal.

This means that $\exp(\tau \delta)$ keeps $\hat{V}$ fixed within its conjugate class, up to first order in $\tau$. Thus, an element $\psi \in \hat{V}$ which is flat with respect to the connection $\nabla_\xi$, above $\gamma$, has the property that,

$$\tilde{\rho}(\delta(a))\psi = \nabla_\xi(\tilde{\rho}(a)(\psi)),$$

for all $a \in A(\sigma)$.

It is therefore reasonable to consider any flat state, $\psi(l) \in \hat{V}$, as the time development of $\psi(0) \in V(0)$. Clearly, the flat states $\psi \in \hat{V}$, are solutions of the differential equation,

$$\xi(\psi) = -Q(\psi), \text{ i.e. } \frac{\partial \psi}{\partial \tau} = -Q(\psi).$$
which, if we accept that time is the parameter \( \tau \) of the integral curve \( \gamma \), is the Schrödinger equation.

Notice that, in the classical quantum-theoretical case, one works with one fixed representation, corresponding to what we have called a singular point of \( \xi \). This implies that we are looking at a representation \( V \) with \( \xi(v) = 0 \), and so we have no time. What we call time is then the parameter of the one-parameter automorphism group \( u(\tau) := \exp(\tau Q) \) acting on \( V \). This also leads to a Schrödinger equation, and to the next result, proving that \( \psi \) is completely determined, along any integral curve \( \gamma \) by the value of \( \psi(\tau_0) \), for any \( \tau_0 \in \gamma \).

**Theorem 3.3.** The evolution operator \( u(\tau_0, \tau_1) \) that changes the state \( \psi(\tau_0) \in \tilde{V}(v_0) \) into the state \( \psi(\tau_1) \in \tilde{V}(v_1) \), where \( \tau_1 - \tau_0 \) is the length of the integral curve \( \gamma \) connecting the two points \( v_0 \) and \( v_1 \), i.e. the time passed, is given by,

\[
\psi(\tau_1) = u(\tau_0, \tau_1)(\psi(\tau_0)) = \exp\left[ \int_\gamma Q(\tau)d\tau \right] (\psi(\tau_0)),
\]

where \( \exp \int_\gamma \) is the non-commutative version of the ordinary action integral, essentially defined by the equation,

\[
\exp\int_\gamma Q(\tau)d\tau = \exp\int_{\gamma_2} Q(\tau)d\tau \circ \exp\int_{\gamma_1} Q(\tau)d\tau
\]

where \( \gamma \) is \( \gamma_1 \) followed by \( \gamma_2 \).

**Proof.** This is a well known consequence of the Schrödinger equation above. In classical quantum theory one uses a chronological operator \( \tau \), to keep track of the intermediate time-steps that, in our case, are well defined by the integral curve \( \gamma \), the existence of which we assume. The formula above is related to what the physicists call the Dyson series, see [Weinberg], Vol I, Chap. 9, or [Elbaz], Chapitre 6. Since we have given the real curve \( \gamma \) parametrized by \( \tau \) we may look at \( \gamma \) as a closed interval of \( \mathbb{R} \), \( I := [0, \tau] \). Subdivide \( I \) into \( m \) equal intervals, \( [i\Delta \tau, (i+1)\Delta \tau] \), and see that the Schrödinger equation gives, formally,

\[
\psi((i + 1)\Delta \tau) = \exp(\Delta \tau Q)(\psi(i\Delta \tau)).
\]

Writing out the power series in \( \Delta \tau \), and summing up we find the formula above.

\( \square \)

The problem of integrating the differential equations above, i.e. finding algebraic geometric formulas for the integral curves of \( \xi = [\delta] \), is a classical problem, and we may use a technique already well known to Hamilton and Jacobi. In fact, assuming that \( A = k[t_1, ..., t_n] \), and that \( \sigma \) is determined by the following force-laws,

\[
d^2t_i = \Gamma^i(t_1, ..., t_n, dt_1, ..., dt_n)
\]

we have that,

\[
A(\sigma) = \text{Ph}^\infty(A)/(\sigma), \quad \delta = \sum_{i=1}^n (dt_i \frac{\partial}{\partial t_i} + \Gamma^i \frac{\partial}{\partial dt_i}).
\]

We may try to solve the equation,

\[
\delta \theta = 0
\]

in the ring \( A(\sigma) \). Obviously the set of solutions form a sub-ring of \( A(\sigma) \), the ring of invariants, and we have the following easy result,
Proposition 3.4. (i): Let $\Theta = \ker \delta$, be the subring of invariants in $A(\sigma)$, and let $\rho : A(\sigma) \to \text{End}_k(V)$ be an n-dimensional representation for which the tangent space of $\text{Simp}_n(A(\sigma))$, at $V$, $\text{Ext}_{A(\sigma)}^1(V, V) = 0$; or suppose $V$ corresponds to a point $t \in \text{Simp}_n(A(\sigma))$ for which $\xi(t) = 0$, then any $\theta \in \Theta$ is constant in $V$, i.e. $[Q, \rho(\theta)] = 0$, so that the eigenvectors of $Q$ are eigenvectors for $\theta$.

(ii): Consider for any $n$ the universal family, $\tilde{\rho} : A(\sigma) \to \text{End}_{C(n)}(\tilde{V})$, and let $\theta \in \Theta$, then $\text{Trace} \tilde{\rho}(\theta) \in C(n)$ is constant along any integral curve of $\xi$ in $\text{Simp}_n(A(\sigma))$, i.e.

$$[\delta](\text{Trace} \tilde{\rho}(\theta)) = 0$$

Proof. (i) Suppose $\delta(\theta) = 0$, and consider the dynamical equation,

$$\delta = [\delta] + [Q, -],$$

where we may assume $[\delta] = \xi$. If the tangent space of $V$ is trivial then, obviously $[\delta] = 0$ therefore $\delta(\theta) = 0$ implies $[Q, \rho(\theta)] = 0$.

(ii) If $\delta(\theta) = 0$, we must have, in $\text{End}_{C(n)}(\tilde{V})$,

$$0 = \text{Trace}_\xi(\tilde{\rho}(\theta)) + \text{Trace}[Q, \tilde{\rho}(\theta)] = \text{Trace}_\xi(\tilde{\rho}(\theta)) = \xi(\text{Trace} \tilde{\rho}(\theta)).$$

□

Notice that we find the same formulas if we extend the action of $A(\sigma)$ to $\tilde{V}_C := \tilde{V} \otimes_R C$. This is what turns out to be the interesting case in quantum physics. It is easy to see that if $A = k[x_1, ..., x_d] \subset A(\sigma)$ is a polynomial algebra, and $\sigma$ is a second order force-law, i.e. such that,

$$d^2x_i = \sum \Gamma^i_{j,k} dx_j dx_k, \ i, j, k = 1, 2, ..., d,$$

then, if we have chosen a versal family,

$$\tilde{\rho} : A(\sigma) \to \text{End}_{C(n)}(\tilde{V})$$

for the simple n-dimensional representations, we obtain another, complexified, versal family,

$$\tilde{\rho}_C : A(\sigma) \to \text{End}_{C(n)}(\tilde{V}_C)$$

with exactly the same formal properties by defining, $\tilde{\rho}_C(x_i) = \tilde{\rho}(x_i)$, $\tilde{\rho}_C(dx_i) = i\tilde{\rho}(dx_i)$, and putting $\xi_C = i\xi$, $Q_C = iQ$.

A section $\psi$ of the complex bundle $\tilde{V}$, a state is now a function on the moduli space $\text{Simp}(A)$, not a function on the configuration space, $\text{Simp}_1(A)$. The value $\psi(v) \in \tilde{V}(v)$ of $\psi$, at some point $v \in \text{Simp}_n(A)$, will be called a state, at the event $v$. 
EndC[n](\tilde{V}) induces also a complex bundle, of operators, on Simp_n(A(σ)). A section, φ of this bundle will be called a quantum field. In particular, any element \( a \in (A(σ)) \) will via versal family-map, \( \tilde{ρ} \), define a quantum field, and the set of quantum fields form a \( k \)-algebra.

Physicists will tend to be uncomfortable with this use of their language. A classical quantum field for any traditional physicist is, usually, a function \( ψ(p, σ, n) \), defined on configuration space, (which is not our Simp_n(A(σ))) with values, in the polynomial algebra generated by certain creation and annihilation operators in a Fock-space.

As we shall see, this interpretation may be viewed as a special case of our general notion.

**Classical Quantum Theory.**

Now, since the configuration-space coordinates \( x_i \) commute, we may find rational sections

\[
|x_ν(ς)⟩ ∈ \tilde{V}_C, \ ν = 1, \ldots, n,
\]

that are eigenvectors for all \( x_i \), such that a point in configuration space is given by the \( n \) possibilities, \( (x_1, ν(ς), \ldots, x_d, ν(ς)) \), where,

\[
\tilde{ρ}(x_ν(ς))(|x_ν(ς)⟩) = x_i, ν(ς)|x_ν(ς)⟩.
\]

In general, the observables \( dx_i, i=1, \ldots, d \), do not commute, but for every \( i \) we can still find eigenvectors,

\[
|dx_i, ν(ς disc)⟩ ∈ \tilde{V}_C, \ ν = 1, \ldots, n,
\]

such that,

\[
\tilde{ρ}(dx_i, ν(ς))(|dx_i, ν(ς)⟩) = dx_i, ν(ς)|dx_i, ν(ς)⟩.
\]

This will be explained in the section *Grand picture, Bosons, Fermions, and Supersymmetry*, where we also will focus on the notion of locality of interaction, see [Klein-Spiro], p. 104, where Cohen-Tannoudji gives a very readable explanation of this strange non-quantum phenomenon in the classical theory, see also in [Weinberg], the historical introduction.

Pick a point \( \tilde{L}_0 \in \tilde{U}(n) \subset Simp_n(A(σ)) \), and assume we have computed the integral curve \( γ \) through \( \tilde{L}_0 \), ending at \( \tilde{L}_1 \), parametrized by \( τ \). The evolution operator \( u(τ_0, τ_1) \) acts upon each \( |x_ν(ς)⟩, ν = 1, \ldots, n \). The result will have the form,

\[
u(τ_0, τ_1)(|x_ν(ς)⟩) = \sum_{μ=1}^{n} γ_{ν, μ}(τ)|x_μ(ς)⟩ \]

and,

\[
u(τ_0, τ_1)(|dx_i, ν(ς)⟩) = \sum_{μ=1}^{n} γ_{i, ν, μ}(τ)|dx_i, μ(ς)⟩, \]

where each \( γ_{ν, μ}(τ) \) and \( γ_{i, ν, μ}(τ) \) is a kind of action integral related to some "classical Lagrangian".

We might now consider the following laboratory situation, in which there are \( n^3 \) cells \( \{X_{q_1, q_2, q_3} \}_{q_1=1, \ldots, n, i=1, 2, 3} \), disposed in a structure like a grid of space, with coordinates \( (q_1, q_2, q_3) \), capable of clicking, if entered by a "particle". Each cell is outfitted with \( n^3 \) sub-cells, \( \{Y_{p_1, p_2, p_3} \}_{p_1, p_2, p_3=1, \ldots, n, i=1, 2, 3} \), capable of clicking if a particle there is outfitted with a certain momentum. Let us talk about these clicks as
a q-click and a p-click respectively. Interpreting \( \{ X_{q_1,q_2,q_3}, q_1 \leq n_1, q_2 \leq n_2, q_3 \leq n_3, n_1 + n_2 + n_3 = n \} \) as a basis of eigenvectors for the space-observables \( x_1, x_2, x_3 \), and considering \( \{ Y^{p_1}, Y^{p_2}, Y^{p_3}, p_i = 1, ..., n, i = 1, 2, 3 \} \), as a basis of eigenvectors for the momenta-observables \( dx_1, dx_2, dx_3 \), for some versal family of \( n \)-dimensional simple representations \( \tilde{V} \), defined on the \( k \)-algebra \( C(n) \). The possible outcomes of a measurement performed at time \( \tau \) are now limited to hearing a q-click in one of the \( n^3 \) points in space, corresponding to the eigenvalues of \( x_1, x_2, x_3 \), and for each such q-click, hearing a different p-click corresponding to one of the \( n^3 \) eigenvalues of \( dx_1, dx_2, dx_3 \). The experiment might consist of letting a beam of particles stream out of an outlet situated at one of the cells, say the one corresponding to the origin \( X_{0,0,0} \) of the q-grid. One checks the distribution of p-clicks from the sub-cells \( \{ Y_{0,0,0}^{p_1,p_2,p_3} \} \), say \( \{ \beta_1, \beta_2, \beta_3 \} \). Now, suppose we have chosen a simple representation \( V(\tau_0) \) such that,

\[
X_{0,0,0} = \sum_{E=1}^{n} \beta_{p_i} Y^{p_i}, \ i = 1, 2, 3,
\]

then we measure time \( \tau \) along the curve \( \gamma \) of \( U(n) \), starting at the point corresponding to \( V(\tau_0) \), and compute,

\[
U(\tau_0, \tau_1)(X_{0,0,0}) = \psi(\tau_1) = \sum_{2} \alpha_{q_1,q_2,q_3} X_{q_1,q_2,q_3}.
\]

We might then want to interpret the family \( |\alpha_{q_1,q_2,q_3}|^2 / |\psi(\tau_1)|^2 \) as the probability distribution, at time \( \tau_1 \), for finding the particle, at the corresponding point. And, correspondingly, one would be tempted to consider the normalized squares of the coefficients in the development,

\[
X_{q_1,q_2,q_3}(\tau_1) = \sum_{E=1}^{n} \beta_{p_i} Y^{p_i}(\tau_1), \ i = 1, 2, 3,
\]

as the probability distribution for momenta observed at the point \( q(\tau_1) \). However, we have to be careful, we have assumed that we might find an object \( \tilde{V} \) with the properties corresponding to our "preparation". This may be possible, as we shall see in an example, see \( (3.7) \), but the interpretation of the coefficients \( \alpha \) and \( \beta \) as probabilities, will probably depend upon the introduction of Hermitian norms on the representation \( \tilde{V} \). Anyway, this seems to lead to a kind of generalized Feynman’s "path"-integral. For a good exposition, for mathematicians, of path integrals, see [Faddeev].

**Planck’s constant(s) and Fock space.**

In [La 5] we treated the case of a conservative system, i.e. where the vector field \( \xi \) or \( [\delta] \)in \( \text{Simp}_n(\mathcal{A}(\sigma)) \) is singular, i.e. vanishes, at the point \( v \in \text{Simp}_n(\mathcal{A}(\sigma)) \) corresponding to the representation \( V \), and where therefore the Hamiltonian \( Q \) is both the time and energy operator, at the same "time". See examples \( (3.6) \) and \( (3.7) \) where we show how to compute these singularities in some classical cases.

We found, in this situation, see [La 6], or §1, that there is a notion of Planck’s constant \( \hbar \), with the ordinary properties.
Let \( \{v_i\}_{i \in I} \) be a basis of \( V \) (no longer assumed to be finite dimensional), formed by eigenvectors of \( Q \), and let the eigenvalues be given by,

\[
Q(v_i) = \kappa_i v_i.
\]

Consider the set \( \Lambda(\delta) \) of real numbers \( \lambda \) defined by,

\[
\Lambda(\delta) := \{ \lambda \in \mathbb{R} \mid \exists f_\lambda \in \text{Ph}^\infty(A), f_\lambda \neq 0, \rho_V(\delta(f_\lambda)) = \lambda \rho_V(f_\lambda) \in \text{End}_k(V) \}.
\]

Since \( \delta = [Q, -] \) is a derivation, if \( f_\lambda \) and \( f_\mu \) are eigenvectors for \( \delta \) in \( V \), then if \( f_\lambda f_\mu \) is non-trivial, it is also an eigenvector, with eigenvalue \( \lambda + \mu \), implying that if \( \lambda, \mu \in \Lambda(\delta) \), with \( f_\lambda f_\mu = 0 \), we must have \( \lambda + \mu \in \Lambda(\delta) \). Now,

\[
\lambda f_\lambda \cdot v_i = \delta(f_\lambda) \cdot v_i = (Q f_\lambda - f_\lambda Q)(v_i) = Q(f_\lambda \cdot v_i) - \kappa_i f_\lambda \cdot v_i,
\]

implying,

\[
Q(f_\lambda \cdot v_i) = (\kappa_i + \lambda) \cdot (f_\lambda \cdot v_i).
\]

If \( f_\lambda \cdot v_i \neq 0 \), it follows that: \( \kappa_i + \lambda = \kappa_j \) for some \( j \in I \). Therefore

\[
f_\lambda \cdot v_i = \alpha v_j, \quad \alpha \in \mathbb{R}, \text{ and } \lambda = \kappa_i - \kappa_j,
\]

and so,

\[
\Lambda(\delta) \subset \{ \kappa_i - \kappa_j \mid i, j \},
\]

Planck’s constant \( h \) should be a generator of the monoid \( \Lambda(\delta) \), when this is meaningful.

We can show, see example 3.5-3.7, that for the classical oscillator \( \Lambda(\delta) \) is infinite an additive monoid. See also that when \( \{f_\lambda\}_\lambda \) generate \( \text{End}_k(V) \) we must have \( \Lambda(\delta) = \{ \kappa_i - \kappa_j \mid i, j \} \), and that when \( h \) "tends to 0", any \( f \in \text{Ph}^\infty(A) \) maps every eigenspace \( V(\kappa_i) \) into itself, see §3. In the generic case when all \( \kappa_i \) are different, the image of \( \text{Ph}^\infty(A) \) into \( \text{End}_k(V) \) becomes commutative, a ring of functions on the spectrum of \( Q \).

The above introduction of Planck’s constant(s), also make sense in general, i.e. for the versal family of \( \text{Simp}(A(\sigma)) \). In fact, since

\[
[Q, \tilde{\rho}(\cdot)] = \tilde{\rho} \delta - [\delta] \tilde{\rho} : A(\sigma) \longrightarrow \text{End}_{\text{C}(n)}(\tilde{V})
\]

is a derivation, we show that the set,

\[
\Lambda(\sigma) := \{ \lambda \in C(n) \mid \exists f_\lambda \in A(\sigma), f_\lambda \neq 0, [Q, \tilde{\rho}(\delta(f_\lambda))] = \lambda \tilde{\rho}(f_\lambda) \},
\]

is a generalized additive monoid, i.e. if for \( \lambda, \lambda' \in \Lambda(\sigma) \) the product \( f_\lambda f_{\lambda'} \) is non-trivial, then \( \lambda + \lambda' \in \Lambda(\sigma) \).

Let \( h_l \in k \) be "generators" of \( \Lambda(\delta) \). These are our Planck’s constants, see examples (3.7) and (3.8). Now, assume there exists a \( C(n) \)-module basis \( \{\tilde{\psi}_i\}_{i \in I} \) of sections of \( \tilde{V} = C(n) \otimes V \), formed by eigenfunctions for the Hamiltonian, i.e. such that

\[
Q(\tilde{\psi}_i) = \kappa_i \tilde{\psi}_i, \quad i \in I,
\]
where \( \kappa_i \in C(n) \). An element such as \( \tilde{\psi}_i \in \tilde{V} \) is usually considered as a pure state
with energy \( \kappa_i \in C(n) \), depending on time, i.e., depending on \( \tau \), the length along the
integral curve \( \gamma \). It is also considered as an elementary particle (since \( \tilde{V} \) is,
by assumption, simple). As in \( \S I \) we find,
\[
\lambda \tilde{\rho}(f_{\lambda})(\tilde{\psi}_i) = Q(\tilde{\rho}(f_{\lambda})(\tilde{\psi}_i)) - \tilde{\rho}(f_{\lambda})(Q(\tilde{\psi}_i))
\]
\[
= Q(\tilde{\rho}(f_{\lambda})(\tilde{\psi}_i)) - \kappa_i \tilde{\rho}(f_{\lambda})(\tilde{\psi}_i)
\]
implying,
\[
Q(\tilde{\rho}(f_{\lambda})(\tilde{\psi}_i)) = (\kappa_i + \lambda) \tilde{\rho}(f_{\lambda})(\tilde{\psi}_i).
\]
By assumption, if \( \tilde{\rho}(f_{\lambda})(\tilde{\psi}_i) \neq 0 \) it must be an eigenvector of \( Q \), with eigenvalue, say
\( \kappa_j = \kappa_i + \lambda \). It follows that we have,
\[
\Lambda(\sigma) \subset \{ \kappa_j - \kappa_i \mid i, j \in I \}.
\]
To prove that the two sets are equal we need some extra conditions on the nature
of \( A(\sigma) \) and \( \tilde{\rho} \). If \( \{ \tilde{\rho}(f_{\lambda}) \}_{\lambda} \) generate \( \text{End}_{C(n)}(\tilde{V}) \), then the equality must hold,
since then \( \{ \tilde{\rho}(f_{\lambda})(\psi(0)) \}_{\lambda} \) must generate \( \tilde{V} \) as \( C(n) \)-module, and therefore contain
multiples of all \( \psi_i \), so that any \( \kappa_i \) must be equal to \( \kappa_0 + \lambda \) for some \( \lambda \).
Notice that if \( \hbar \rightarrow 0 \), meaning that \( [Q, \tilde{\rho}(a)] = 0 \), for all \( a \in A(\sigma) \), then all
\( a \in A(\sigma) \) must commute with \( Q \), and so act diagonally on the spectrum of \( Q \).
Notice also that if, at a point \( v \in \gamma \), \( h(v) \neq 0 \) as an element of \( k = \mathbb{R} \), it is
clearly reasonable to redefine \( \delta \) and \( Q(v) \) by dividing both with \( h(v) \). Then the
energy differences of \( 1/h(v)Q(v) \) will come up as integral values.
Assume that \( Q(\tau) = Q \) is constant along \( \gamma \), and use the theorem (3.3).
We find
\[
U(\tau_0, \tau_1)(\tilde{\psi}_i(\tau_0)) = \tilde{\psi}_i(\tau_1) = \exp[\int_\gamma Q d\tau] (\tilde{\psi}_i(\tau_0)) = \exp[\int_\gamma \kappa_i(\tau)d\tau] (\tilde{\psi}_i(\tau_0)),
\]
and in particular,
\[
\frac{\partial \tilde{\psi}_i(\tau)}{\partial \tau} = \kappa_i \exp(\int_\gamma \kappa_i(\tau)d\tau)(\tilde{\psi}_i(\tau_0)) = Q(\tilde{\psi}_i(\tau)),
\]
so, of course, again the Schrödinger’s equation, with \( \tau \) as time. For an example of
a non-constant Hamiltonian, see Example (3.8) and (3.9).
Above, \( \text{Simp}_n(A(\sigma)) \) is our "time-space", and \( \text{Simp}_1(A) \) is the analogue
of the classical configuration space. Given an element \( v \in \text{Simp}_n(A(\sigma)) \), corresponding
to a simple module \( V \) of dimension \( n \), there are for every \( a \in A(\sigma) \), a set of \( n \)
possible values, namely its eigenvalues, as operator on \( V \). Since \( V \) is simple, the
structure map,
\[
\rho_V : A(\sigma) \rightarrow \text{End}_k(V)
\]
is supposed surjective, and so in general (and, for order 2 dynamical systems, always) the operators \( \tilde{\rho}(a) \) and \( \tilde{\rho}(da) \), \( a \in A \), cannot all commute. In fact, if \( \dim V = \infty \), or \( \dim V \) "approaching" \( \infty \), see example (3.7), (3.8), any one \( a \in A(\sigma) \) would
tend to have a conjugate, i.e., an element \( b \in A(\sigma) \), such [\( \tilde{\rho}(a), \tilde{\rho}(b) \] = 1. Therefore,
if the values \( q_i \) of \( \tilde{\rho}(a) \) are determined, then the values \( p_i \) of \( \tilde{\rho}(b) \) will be totally
biased, and vice versa, giving us the Heisenberg indeterminacy problem. In general there is no way of fixing a point of \( \text{Simp}_1(A(\sigma)) \) as representing \( V \) or finding natural morphisms,
\[
\text{Simp}_n(A(\sigma)) \rightarrow \text{Simp}_m(A(\sigma)), m < n.
\]
However, as we know, see [La 4] and [Jø-La-Sl], there are partially defined decay maps,

\[ \operatorname{Simp}_n(A(\sigma))^{\infty} := \frac{\operatorname{Simp}_n(A(\sigma)) - \operatorname{Simp}_n(A(\sigma))}{\otimes_{n>m \geq 1} \operatorname{Simp}_n(A(\sigma))}. \]

In the very special case, where \( A = k[x_1, \ldots, x_p] \) is a commutative polynomial algebra, there exists moreover, for every linear form \( f : V \to k \), and every state \( \psi(\tau) \in \hat{V}|\gamma \) a curve \( \Psi(\gamma) \subset \text{Spec}(A) \cong A^p \) defined, by its coordinates, in the following way,

\[ x_i(\tau) = \int \tilde{\rho}(x_i) \psi(\tau) / \int \psi(\tau), \quad i = 1, \ldots, p. \]

Here \( \hat{V}|\gamma \) is identified with \( V \otimes_k O_\gamma \), \( \tau \) being a parameter of \( \gamma \). If we are able to pick common eigenfunctions \( \{ \phi_j \in \hat{V}_\gamma \}, \ j = 1, \ldots, n \) for \( \tilde{\rho}(x_i) \), \( i = 1, \ldots, p \), generating \( \hat{V}_\gamma \), with eigenvalues \( \kappa_j(\tau), \ j = 1, \ldots, n, \ i = 1, \ldots, p \), and if \( \psi(\tau) = \sum_j \lambda_j(\tau) \phi_j \), then picking the linear form defined by, \( \int \phi_j = 1, \ j = 1, \ldots, n \), we find,

\[ x_i(\tau) = \sum_j \lambda_j(\tau) \kappa_j(\tau) / \sum_j \lambda_j(\tau), \]

which is a general form of Ehrenfest’s theorem.

Suppose now that we have a situation where there is a unique non-trivial positive (as a real function) Planck’s constant, \( h \in C(n) \). Consider \( f_h \in A(\sigma) \), and assume that there are among the \( \{ f_\lambda \} \) a conjugate, i.e. a \( f_\mu \) such that \( [\tilde{\rho}(f_h), \tilde{\rho}(f_\mu)] = 1 \). This obviously cannot happen unless \( \dim V = \infty \), but see the examples (3.7) and (3.8) for what happens at the limit when \( \dim V \) goes to \( \infty \).

Then we easily find that \( \mu = -h \). Moreover, if \( \psi_0 \) is an eigenvector for \( Q \) with least energy (assumed always positive), \( \kappa_0 \), then \( \hat{N} := f_{-h} f_h \) is a quanta-counting operator, i.e. \( N(\psi_i) = i \), when \( \kappa_i = \kappa_0 + (i - 1)h \), is the \( i \)th energy level. It follows also that \( [Q, f_{-h} f_h] = 0 \). The algebra generated by \( \{ f_h, f_{-h} \} \) is a kind of a Fock representation, \( F \) on a Fock space. Its Lie algebra of derivations turns out to contain a Virasoro-like Lie-algebra. Since for \( Q_h := f_{-h} f_h \) we have that \( [Q_h, f_{-h}] = f_{-h}, [Q_h, f_h] = f_h \), it is often seen in physical texts that one identifies the Hamiltonian, \( Q \), with \( Q_h \), or with \( \sum_h Q_h \). We shall return to this in the examples (3.8), (3.9) and (3.16), at the end of this §.

We have seen that starting with a finitely generated \( k \)-algebra \( A \), and a dynamical system \( \sigma \), we have created an infinite series of spaces \( \operatorname{Simp}_n(A(\sigma)) \) and a quantum theory, on étale coverings \( U(n) \), of these spaces, with time being defined by the Dirac derivation \( \delta \).

Each \( C(n) \) is commutative and \( \hat{V} \) is a universal bundle on \( U(n) \subset \operatorname{Simp}_1(C(n)) \). Moreover, the elements of \( A(\sigma) \), the observables, become sections of the bundle of operators, \( \operatorname{End}_{C(\sigma)}(\hat{V}) \).

Clearly, if \( D \subset \operatorname{Simp}_1(C(n)) \) is a subvariety, say a curve parametrized by some parameter \( q \), then the universal family induces a homomorphism of algebras,

\[ \tilde{\rho}_D : A(\sigma) \longrightarrow \operatorname{End}_D(\hat{V}|D). \]

This is in many recent texts referred to as a quantification of the commutative algebra \( A(\sigma)/[A(\sigma), A(\sigma)] \), or to a quantum deformation, and the parameter \( q \) is
sometimes confounded with Planck’s constant. This is unfortunate, but probably unavoidable!

In quantum theory one attempts to treat the second quantification of an oscillator in dimension 1, as a certain representation on the Fock space, i.e. constructing observables acting on Fock space, with the properties one wants. This turns out to be related to the canonical representations of \( PhC := k < x, dx > \) on an \( n \)-bundle over the algebra, \( R := k[[n]] \). Here the \( p,q \)-deformed numbers \([n]_{p,q} \) are introduced as,

\[ [n]_{p,q} := q^{n-1} + pq^{n-2} + p^2 q^{n-3} + ... + p^{n-2} q + p^{n-1}, \]

and we may as well consider \( p,q \) as formal variables, so that \( R \subset k[q] \).

See (3.9) where we construct a homomorphism of \( A(\sigma) \) into an endomorphism ring of the form \( \text{End}_R(V \otimes_k R) \), see ([Elbaz], Appendice, on the “q-commutators”).

Picking representatives for \( x \) and \( dx \) in \( M_n(\mathbb{k} \{ x \}) \), it turns out that, instead of the classical defining relations for an oscillator, i.e. with \( a_+ := x + dx, a_- : x - dx \), and with a Hamiltonian \( Q \), such that in \( \text{End}_R(V \otimes_k R) \),

\[ [Q, x] = dx, \quad [Q, dx] = x, \quad [a_-, a_+] = 1 \]

one obtains,

\[ [Q, x]_q = dx, \quad [Q, dx]_q = x, \quad [a_-, a_+]_q = 1 \]

where \([a, b]_q := ab - qba\) is the “quantized” commutator. This holds in particular for \( p = 1 \), so for \( R = k[q] \), defining a curve \( D \) in \( \text{Simp}_n(Ph(k[x])) \).

However, this \( k[q]\)-parametrization is not parametrizing an integral curve of \( \xi \) in \( \text{Simp}_n(PhC) \). On the contrary, it is parametrizing a curve of anyons with \( q = -1, 1 \) corresponding to, respectively, Fermions and Bosons. A simple computation shows that it is transversal to \( \xi \), and therefore represent a phenomenon which takes place instantaneously, see the Examples (3.8), (3.9).

**General Quantum Fields, Lagrangians and Actions.**

Given algebras \( A \) and \( B \), supposed to be moduli spaces of interesting objects. Given dynamical structures \((\sigma), (\mu)\) of \( Ph^\infty(A) \) and \( Ph^\infty(B) \) respectively, corresponding to Dirac derivations, and corresponding vector fields, \( \delta_\sigma, \xi_\sigma \) and \( \delta_\mu, \xi_\mu \), respectively, we consider the set (or space, see §2.), \( \mathcal{F}(A,B) \), of iso-classes, of morphisms \( \phi : A \rightarrow B \). Any such induces a morphism,

\[ Ph(\phi) : PhA \rightarrow PhB. \]

and we shall assume also a morphism,

\[ \phi : A(\sigma) \rightarrow B(\mu). \]

This is actually what we shall call a field. The space of such is denoted \( \mathcal{F}(A,B) \).

The meaning of the term field, or its physical interpretation, is not obvious. In standard physics texts the notion is rarely well defined, see e.g. [Weinberg], I, (1.2), where one finds a nice historical introduction, and good reasons for lots of mathematical tears!

I would like to consider \( \phi : A \rightarrow B \) as a morphism of the moduli space \( \text{Spec}(B) \), of physical objects \( Y \), into the moduli space \( \text{Spec}(A) \), of physical objects \( X \), and in this way modeling composite physical objects \( (X,Y) \), as we shall see below.
Now, apply the deformation theory of categories of algebras, see e.g. [La 0]. From this theory follows readily that the tangent space of $\mathcal{F}(A, B)$ at a point, $\phi$, is,

$$T_{\mathcal{F}(A, B), \phi} = \text{Der}_k(A(\sigma), B(\mu))/\text{Triv},$$

where, as usual, $\text{Triv}$ depends upon the definition of $\mathcal{F}(A, B)$, i.e. upon the notion of isomorphisms among fields, see [La 0].

Unfortunately, this moduli space, $\mathcal{F}(A, B)$, is not, in general, a prescheme, neither commutative nor non-commutative. As we have, as a rule in this paper, identified any $k$-algebra with some space, we shall, never the less, at this point not hesitate to identify $\mathcal{F}(A, B)$ with the $k$-algebra of (reasonable) functions (or operators) defined on the space, and denote it by, $F(A, B)$. Then we are free to consider the versal (or maybe, universal) family of quantum fields,

$$\tilde{\phi} : A \to F(A, B) \otimes_k B.$$

Just in the same way as above, there is now a canonical vector field $[\delta]$ on the space $\mathcal{F}(A, B)$, defined by its value at $\phi$, given by,

$$[\delta](\phi) = \text{cl}(\delta_\sigma \phi - \phi \delta_\mu).$$

The set of stable, or singular fields, is now in complete analogy with the singular points in $\text{Simp}(A(\sigma))$ mentioned above, and treated in detail in the examples, (3.5)-(3.9) below,

$$M(A, B) := \{ \phi \in \mathcal{F}(A, B) | [\delta](\phi) = 0 \}.$$  

Here one sees where the Noether function $Q$ enters. In fact, if we are identifying fields up to automorphisms of $B$ defined by trivial derivations, $[\delta](\phi) = 0$ is equivalent to the existence of a "Hamiltonian" $Q \in B(\mu)$, such that for every $a \in A(\sigma)$

$$\delta_B(\phi(a)) - \phi(\delta_A(a)) = Q\phi(a) - \phi(a)Q = [Q, \phi(a)].$$

Consider this equation in rank 1, i.e. look at the commutativizations

$$\text{Ham} : A(\sigma) \to A(\sigma)/[A(\sigma), A(\sigma)] =: A_0(\sigma),$$

$$\text{Ham} : B(\mu) \to B(\mu)/[B(\mu), B(\mu)] =: B_0(\mu),$$

We find that in $\text{Simp}_1(B(\mu))$ the equation above reduces to,

$$\delta_B(\phi(a)) = \phi(\delta_A(a)),$$

which, geometrically, means the following: If $\gamma$ is an integral curve of $\delta_A$, in $U^A(n)$then the inverse image via $\phi$ is a union of integral curves for $\delta_B$ in $U^B(n)$.

The actual definition of a dynamical structure $(\sigma)$ has, up to now, been loosely treated. It may be defined in terms of "force laws", i.e. where $(\sigma)$ is the two-sided $\delta$-stable ideal generated by the equations $(\delta^n t_p - \Gamma^p)$, where,

$$\delta^n t_p := d^n t_p, \Gamma^p(t_i, d^2 t_k, ..., d^{n-1} t_l) \in \text{Ph}_\infty(A), p = 1, 2, ..., d.$$  

But, in general, force laws like these don’t pop up naturally. In fact, Nature seems to reveal its structure through some Action Principles. The physicists are looking
for a Lagrangian $L$, and an action functional $S(L)$ defined on $\mathcal{F}(A, B)$. In our setting, $L$ is simply an element,

$$L \in \text{Ph}(A).$$

For every field $\phi \in \mathcal{F}(A, B)$, the action, usually denoted,

$$S(\phi) := S(L)(\phi) \in k,$$

is constructed via some particularly chosen representation,

$$\rho : B(\mu) \rightarrow \text{End}_k(V).$$

Put, $\mathcal{L} := \rho(\phi(L))$ and let,

$$S(\phi) := \text{Tr}(\mathcal{L}).$$

In the classical case the trace, $\text{Tr}$, is an integral.

We may choose several canonical representations $\rho$, like the versal family of rank 1 representations, treated above, and called, $\text{Ham} : \text{Ph}(B) \rightarrow \text{Ph}_0(B) := \text{Ph}(B)/[\text{Ph}(B), \text{Ph}(B)]$, or the corresponding in rank $n$. We may also, for any derivation $\zeta$ of $B$, consider the canonical homomorphism $\rho_{\zeta} : \text{Ph}(B) \rightarrow B$, as a representation. In the first case it is clear that $\text{Tr}(\mathcal{L})$ is simply the image of $\mathcal{L}$ in $\text{Ph}_0(B)$. In the last case the Lagrangian density, i.e. $\mathcal{L}$, is now a function, $\mathcal{L}(\phi_i, \zeta(\phi_j))$, in $\phi_i := \phi(t_i)$, and in $\zeta(\phi_j)$ for some generators $t_i$ of $A$, and $\text{Tr}$ is an integral over some reasonably well defined subspace of $\text{Sim}_1(B)$. In this case one usually has to impose some boundary conditions on $\phi$.

Clearly, $S := S_\rho = \text{Tr}(\mathcal{L})$ is a function,

$$S : \mathcal{F}(A, B) \rightarrow k,$$

and $\nabla S \in \Theta_{\mathcal{F}(A, B)}$, is a vector field that corresponds to the fundamental vector field $\xi = [\delta]$, above. The equations defining the singularities of $\nabla S$, is usually written,

$$\delta S := \delta \int \mathcal{L} = 0,$$

since for most classical representations the dimension of $V$ is infinite, and the trace is an integral, see examples below.

Here is where the calculus of variation comes in. The corresponding Euler-Lagrange equations, the equations of motion, picks out the set of solutions, the singular fields, i.e. $\mathcal{M}(A, B) \subset \mathcal{F}(A, B)$.

The subspaces $\mathcal{M}(A, B)$ in $\mathcal{F}(A, B)$, defined by the Euler-Lagrange equations, are therefore uniquely defined by $L$, without reference to any dynamical structures of $A$ or $B$.

The problem with this is that, unless there actually exist a dynamical structure corresponding to $\nabla S$, we cannot know that our laws of nature are mathematically deterministic, see the Introduction, and compare [Weinberg], I, chapter 7.
Notice that the classical field theory corresponds to the situation where $A = k[t]$ and $B = k[x]$, and where $\phi$ is defined in terms of the fields, $\phi_i := \phi(t_i)$, and their time derivatives $\dot{\phi}_i := \phi(dt_i) := d\phi_i \in Ph(B)$. Choosing the representation $\rho_\zeta$ for some derivation $\zeta$ of $B$, we may assume the Lagrangian density has the form,

$$\mathcal{L} := L(\phi_i, \phi_{j,\alpha})$$

where $\phi_{j,\alpha} := \frac{\partial \phi_j}{\partial x_\alpha}$. The singular fields, are then picked out by the variation of the action integral, i.e.

$$\delta \int \mathcal{L} d\mu = 0,$$

or by the corresponding Euler-Lagrange equations,

$$\frac{\partial \mathcal{L}}{\partial \phi_i} - \sum_\alpha \frac{\partial}{\partial x_\alpha} (\frac{\partial \mathcal{L}}{\partial \phi_{i,\alpha}}) = 0.$$

In the light of the above, considering these equations as equations of motion, is now, maybe, a reasonable guess.

In fact, we shall show that this Lagrangian theory actually produce a dynamical structure, at least in special cases.

Pick a Lagrangian $L \in Ph(A)$, and assume $B = M_n(k)$, so that $\mathcal{F}(A, B) = \text{Rep}_n(A) \subset \text{Rep}_n(Ph(A))$. Restrict to $\text{Simp}_n(Ph(A)) \subset \text{Rep}_n(Ph(A))$, and consider the versal family, $\tilde{\rho} : Ph(A) \to M_n(C(n))$.

Put $\mathcal{L} := \tilde{\rho}(L) \in M_n(C(n))$ and $S := \text{Tr}\mathcal{L} \in C(n)$. If the choice of the Lagrangian $L$, is clever, the gradient, $\nabla S \in \Theta_{C(n)}$, restricted to $U(n)$ is a candidate for the vector field $\xi = [\delta]$, induced by the Dirac derivation $\delta$ defined by some dynamical structure, $A(\sigma)$. If the philosophy of contemporary physics is consistent, this is what we would expect.

Based on the parsimony principle involved in the theory of Lagrange, and given a dynamical system, with Dirac derivation $\delta$, we should expect that the Lagrangian $L$ is constant in time, i.e. that we have the Lagrangian equation,

$$\delta(L) = 0.$$ 

But then Theorem (3.4) tells us that, in the situation above, we have in $C(n)$,

$$[\delta](\text{Tr}\mathcal{L}) = 0,$$

i.e. for all $n \geq 1$, the equation $\nabla S = 0$, picks out the solutions $\gamma$ of the theory. Now we may try to turn the argument upside down, and ask whether, given $L$, we may construct a Dirac derivation, $\delta$, from the Lagrangian equation above. This is the purpose of the next examples. But be aware, this is not proving that the Lagrangian method for studying quantum fields, is equivalent to the one I propose above. Example (3.8) shows that there exists simple Lagrangians $L$ inducing unique force laws, but such that the set of solutions $\{\gamma\}$ is not determined by $\ker(\delta)$. Moreover, the classical relation between the Lagrangian and the Hamiltonian turns out to be more subtle in this non-commutative case. Notice that, above, we have
accepted fields $\phi : A \to B$ where $B = M_n(k)$ is a finite non-commutative space, for which the usual Euler-Lagrange equations do not apply.

Solving differential equations like the Lagrange equation, in non-commutative algebras, is not easy. However, if we reduce to the corresponding commutative quotient, things become much easier. In fact, as we mentioned in the Introduction, in the commutative situation we may write, in $Ph(A)$,

$$\delta = \sum_i (dt_i \frac{\partial}{\partial t_i} + d^2 t_i \frac{\partial}{\partial dt_i}),$$

and the Lagrange equation will produce order 2 dynamical structures, see Example (3.5). We may also consider the Euler-Lagrange equations, impose $\delta$, as time, and solve,

$$\delta \frac{\partial L}{\partial dt_i} - \frac{\partial L}{\partial t_i} = 0,$$

to find an order 2 force law, $d^2 t_i = \Gamma(t_i, dt_j)$.

The strategy will be to solve the equation in a representation like $Ham$, then try to lift it to $Ph(A)$ and then, eventually, map it back to say $Wey : Ph(A) \to Diff f_k(A, A)$. We shall now show that this strategy works in some interesting cases.

But first, let’s first have a look at the relationship, as we see it, between the picture we have drawn of QFT, and the one physicists presents in modern university textbooks.

**Grand Picture: Bosons, Fermions, and Supersymmetry.**

Now, go back to the QF-set-up. Consider a situation with a dynamical system, with Dirac derivation $\delta$, and fix the rank $n$ versal family, $A(\sigma) \to End_{C(n)}(\tilde{V})$.

Look at the singularities of the fundamental vectorfield $\xi \in Dcr_k(C(n))$. Let $V$ be a corresponding representation, the “particle”. Compute the set of eigenvalues $\Lambda$ of $adQ$ acting on $End_k(V)$, and the set of minimal elements, i.e. the set of Planck’s constants, $\{\hbar\}$, and the corresponding eigenvectors $f_\lambda \in End_k(V)$. We shall see in Examples (3.7) and (3.8), that if there exists a conjugate to $f_\hbar$, it must be $f_{-\hbar}$. If the Hamiltonian, $Q$, is diagonalized, with eigenvalues $\{q_0 \leq q_1 \leq ... \leq q_{n-1}\}$ it is of course easy to see that $\Lambda = \{(q_i - q_j)\}_{i,j=0,..,n-1}$, $f_\lambda = \sum_{q_i - q_j = \lambda} \epsilon_{i,j}$, and in particular, $f_{-\hbar} = f_{\hbar}^*$, the conjugate matrix.

Anyway, choosing a vacuum state $\phi_0 \in V$ for the Hamiltonian $Q$, i.e. an eigenvector with minimal positive-or zero-eigenvalue, we find that, for some $i > 0$, unless $f_\hbar^i = 0$, we have $Q(f_\hbar^i(\phi_0)) = i\hbar f_\hbar(\psi_0)$, i.e. the state $\phi_i := f_\hbar^i(\phi_0)$ may be occupied by $i$ quanta. If $\{\phi_i\}_{i=0,..,n-1}$ is a basis for $V$ then this is the purely Bosonic case, with $q_i - q_{i-1} = \hbar$, see (3.8), where we have treated the simple case of the harmonic oscillator.

What do I mean by state occupied by several quanta? The language is far from clear. Here we shall restrict ourself to the elementary language of quantum physics.. The phrase, ”the state $\phi$ is occupied by $n$ quanta” shall mean that $\phi$ is an eigenstate of the Hamiltonian $Q$, with eigen-value $n\hbar$. We shall also, as is explained above, assume $\hbar = 1$.
Physicists have come to the realization that there exist two types of particles, Bosons and Fermions, with different statistics, in the sense that states "containing" several identical Bosons are invariant upon permutations of these, but states "containing" several identical Fermions change sign with the permutation. This is another way of expressing that Fermions, like electrons, cannot all sink in to the lowest energy state in an atom, and stay there, killing chemistry.

We shall delay the discussion of collections of identical particles to §4.

Bosons can have states with an arbitrary occupation number, but Fermions have states only with occupation numbers 0 or 1.

If we know that no states are occupied by more than one quantum "at a time", then we must conclude that,
\[ f_h^2 = f_{-h}^2 = 0. \]

Moreover, we pose,
\[ \{ f_h, f_{-h} \} := f_{-h}f_h + f_hf_{-h} = 1, \]

implying, \((f_h + f_{-h})^2 = 1\). This is the purely Fermionic case.

These relations induce a split-up of the representation \(V\), i.e.
\[ V \simeq V_0 \oplus V_1, \]

In fact, put \( R_0 := f_h f_{-h}, \ R_1 := f_{-h}f_h, \) and see that
\[ R_0 + R_1 = 1, \ R_0 R_1 = R_1 R_0 = 0, \ R_i^2 = R_i, \ i = 0, 1, \]

and put \( V_0 = imR_0, V_1 = imR_1. \) Since \( R_i \) is the identity on \( V_i \), it is clear that the two linear maps,
\[ f_{-h} : V_0 \to V_1, \ f_h : V_1 \to V_0, \]

are isomorphisms, thus \( dim_k V_0 = dim_k V_1 = 1/2 \ n. \) Clearly, any endomorphism of \( V \) can be cut up into a sum of graded endomorphisms. Those of degree 0 we would like to call Bosonic. Those of degree 1, or -1, should then be called Fermionic. In dimension 4, this would look like:
\[
Q = \begin{pmatrix}
q_{1,1} & q_{1,2} & 0 & 0 \\
q_{2,1} & q_{2,2} & 0 & 0 \\
0 & 0 & q_{1,1} + 1 & q_{1,2} \\
0 & 0 & q_{2,1} & q_{2,2} + 1
\end{pmatrix},
\]

with,
\[
f_{-h} = \begin{pmatrix}
0 & 0 & 1 & 0 \\
0 & 0 & 0 & 1 \\
0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0
\end{pmatrix} \quad \text{and} \quad f_h = \begin{pmatrix}
0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 \\
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0
\end{pmatrix}.
\]

In the general case we may have a mix of Bosons and Fermions present, and this leads to the notion of super symmetry.

If we have a split up, as above, the modules \( V_i, i = 1, 2, \) having the Hamiltonians \( Q_0 \) and \( Q_1 := Q + 1, \) implying that, as Hamiltonians, \( Q_0 = Q_1, \) we see that \( End_k(V) \) is generated by the Bosonic operators, \( a := f_h, \ a^+ := f_{-h} \in End_k(V_i), \)
both defined for $Q_0 = Q_1$, together with the Fermionic operators $f^- := f_\hbar \in \text{Hom}_k(V_1, V_0)$, $f^+_\hbar \in \text{Hom}_k(V_0, V_1)$, for the Hamiltonian $Q$. In fact,

$$\text{End}_k(V) = \left( \begin{array}{cc} \text{End}_k(V_0) & \text{Hom}_k(V_0, V_1) \\ \text{Hom}_k(V_1, V_0) & \text{End}_k(V_1) \end{array} \right).$$

is generated by the $\text{End}_k(V_0)$ and $\text{End}_k(V_1)$, together with the isomorphisms $f^- : V_0 \rightarrow V_1 \in \text{Hom}_k(V_0, V_1)$, $f_\hbar : V_1 \rightarrow V_0 \in \text{Hom}_k(V_1, V_0)$.

Put, $f = i(f^- + f^+)$, and see that there are two eigenstates of $f$, the Fermion with eigenvalue $1$, and the anti-Fermion with eigenvalue $-1$.

The general situation is much like the one above. We may assume that we have a Hamiltonian $Q$, split up as above, with corresponding Bosonic operators, $a_1$, $a_1^+$, and Fermionic operators, $f_p$, $f^+_p$, generating $\text{End}_k(V)$. We may also assume that, given the vacuum state $\phi_0 \in V$, there is a basis of $V$, given by $\{ \phi_i := (a_j^+)^i(\phi_0) \}_{i=0}^{n-1}$. Moreover $a_l$ kills $\phi_0$, and $a_l^+$ kills $\phi_{n-1}$. Considering the versal family $\tilde{p}$, and the global Hamiltonian $Q \in \text{End}_{C(n)}(V)$, the situation becomes more subtle. Here we have the étale morphism $\pi : U(n) \rightarrow \text{Simp}_n(A(\sigma))$. Fix the field $k = \mathbb{R}$, and assume no harm is made by this choice. Clearly we have a monodromy homomorphism,

$$\mu(v) : \pi_1(v; \text{Simp}_n(A(\sigma))) \rightarrow \text{Aut}(V) \simeq \text{Gl}_n(\mathbb{R}).$$

One would be tempted to define Bosons, Fermions, and Anyons, with respect to this monodromy map. The component of $\text{Simp}_n(A(\sigma))$ where $\mu(v)$ is trivial are Bosonic, the one with $i\mu(v) = \{ +1, -1 \} \simeq \mathbb{Z}_2$ is Fermionic, and the rest are Anyonic. Notice that the fiber of $\pi$ is composed of identical particles. The treatment of such are, as mentioned above, postponed until the next paragraph, §4.

If $Q$ is constant, we may of course assume that the Bosonic operators, $a_1$, $a_1^+$, and Fermionic operators, $f_p$, $f^+_p$, are elements of $\text{End}_{C(n)}(V)$, generating $\text{End}_{C(n)}(V)$, as $C(n)$-module. Then any quantum field would look like,

$$\psi(v) = \psi(v; a, a^+, f, f^+),$$

the functions being polynomials in the operator variables. In particular $\tilde{p}(t_j)$ and $\tilde{p}(dt_j)$ would have this form, where, in most cases relevant for physics, the polynomial function would be linear, see the case of the harmonic oscillator, Example 3.8.

This is very close to the form one finds in physics books, the only problem is that the function $\psi$ is a function on $\text{Simp}(A(\sigma))$, not on the configuration space, with fixed momentum, as is usually the case in physics.

Suppose we have a classical case, where the algebra $A = k[t_1, ..., t_r]$ is the commutative affine algebra of the "configuration" variety, $X := \text{Spec}(A)$. Then an element $v \in \text{Simp}_n(A(\sigma))$ will correspond to up to $n$ different points in $q \in X$. If one imposes commutation rules on the $dt_j$, as physicists do, then to $v$, there corresponds also up to $n$ values of the momenta, $p_l \in \text{Spec}(k[dt_1, ..., dt_r])$. However, there is no way to pinpoint the representation $v$, by fixing $q$ and $p$. Because $t_i$ and $dt_i$ do not commute, which imposes the Heisenberg uncertainty relation with respect to determining $q's$ and $p's$, at "the same time", the physicists will have to introduce some mean values, using different versions of the spectral theorem for Hermitian operators, to obtain reasonable definitions of the notion of quantum
field. Usually the generators, of the algebra of quantum fields, are expressed in the form of an integral, like,
\[
\psi_l(x) = \sum_\sigma (2\pi)^{-3/2} \int d^3p[p(u(p, \sigma, n)a(p, \sigma, n)\exp(ipx) \\
+ v(p, \sigma, n^c)a^+(p, \sigma, n)\exp(-ipx)],
\]
see [Weinberg], I, p.260. Here \( \sigma \) means spin, \( n \) the particle "species", and \( n^c \) the antiparticle of the species \( n \). Integration is on different domains, depending on the situation, and the whole thing is deduced from ordinary quantum theory, imposing relativistic invariance.

The corresponding interpretation of interaction, implies that interaction takes place at points in configuration space. This is the so called locality of action. See the very readable article of Gilles Cohen-Tannoudji in [Klein-Spiro], p.104.

Of course, the interesting Hamiltonians \( Q \in \text{End}_{C(n)}(\hat{V}) \) will not be constant, therefore physicists introduce what is called Perturbation Theory, which amounts to assuming that there exists a background situation, defined by an essentially constant Hamiltonian \( Q_0 \), such that for the real situation, given by the versal family \( \hat{\rho} : \mathcal{A}(\sigma) \to \text{End}_{C(n)}(\hat{V}) \), and the Dirac derivation \( \delta \), the Hamiltonian \( Q \) may be considered as a perturbation of \( Q_0 \), with an Interaction \( I \in \text{End}_{C(n)}(\hat{V}) \), such that,
\[
Q = Q_0 + I.
\]

Then, using the basis \( \{\phi_i\} \), given for \( \hat{V} \), defined by the Creation operators \( \{a^+_i\} \), see above, one may apply Theorem 3.3, and obtain formulas for the evaluation operator, along the curve \( \gamma \), applied to any \( \phi_i \). If we have a Hermitian metric on the bundle \( \hat{V} \), then we obtain formulas for the so called \( S = (S_{ij}) \)-matrix, calculating the probability for a \( \phi_i \) observed at the start-point \( v_0 \) of \( \gamma \), to be observed as changed into \( \phi_j \), at the end-point \( v_1 \). The same types of formulas as one finds in elementary physics books, like [Weinberg], I, p.260, pop up. And the computation is again made easier by chopping up the formulas, by introducing Feynman diagrams. In our case, the integrals along the compact \( \gamma \), are, of course, easily seen to be well defined, but then we have not explained how we may know that our preparation gave us the start-point \( v_0 \).

This is where the problem of locality of action enters. Suppose we have fixed a basis, \( \{\phi_i\}_{i=0}^{n-1} \), of the \( C(n) \)-module of sections of \( \hat{V} \), composed of common eigenvectors for the commuting operators, \( \hat{\rho}(t_j) \in \text{End}_{C(n)}(\hat{V}) \). Suppose also that the operator \( \hat{\rho}(t_j), \hat{\rho}(dt_j) \) is sufficiently close to the identity, see examples (3.8) and (3.9). Notice again that it must have a vanishing trace, since we are in finite dimension. Fix an index \( l \) and let \( \{\kappa_i(l)\}_{i=0}^{n-1} \) be a basis of the \( C(n) \)-module of sections of \( \hat{V} \), composed of eigenvectors for the operator \( \hat{\rho}(dt_l) \). Then, given a \( v \in \text{Simp}_n(\mathcal{A}(\sigma)) \), represented by the \( n \)-dimensional \( \mathcal{A}(\sigma) \)-module \( V \), we have,
\[
t_j(\phi_i) = q^i_j \phi_i, \quad dt_l(\kappa_i) = \mu^i_l,
\]
where, for \( i = 0, ..., n-1 \), \( q_j(v) := q^i_j = (q^i_0, ..., q^i_{n-1}) \in X \) are the possible configuration positions of \( v \), and, for \( j = 0, ..., n-1 \), the possible values of the \( l \)-component of the momenta, are given by, \( \mu^i_l(v) := p^i_l \). Consider now the base-change matrices, \( (\lambda^i_j) \), and \( (\mu^i_l) \), such that, \( \phi_j = \sum \lambda^i_j \kappa_i \), and \( \kappa_j = \sum \mu^i_l \phi_i \), and compute,
\[
q^i_l dt_l(\phi_i) = dt_l(q^i_l \phi_i) = dt_l t_l(\phi_i) = (t_l dt_l + [dt_l, t_l])(\phi_i).
\]
We obtain,
\[
[dt_1, t_2](\phi) = \sum_{j,k} \lambda^j_i \mu^l_k (q^l_k - q^l_i) \phi_k .
\]

By assumption, the operator \([dt_1, t_2]\) is close to the identity, which implies that when the \(l\)-coordinate of the configuration positions are clustered tightly about a certain point, then the \(l\)-coordinate of the corresponding momenta cannot be kept bounded. This is the analogy of the Heisenberg uncertainty relation of the classical quantum theory.

With this in mind, one would be tempted to formulate the task of the experimenter in physics, as follows.

She should test out the possibilities of the laboratory technology, to prepare the situation by bounding the configuration positions of the phenomenon she is interested in, to a subset \(D(q) \subset X := \text{Spec}(k[t_1, \ldots, t_r])\), and at the same time bound the corresponding \(l\)-component of the momenta, to a subset \(D(p, l) \subset Y := \text{Spec}(k[dt_1, \ldots, dt_r])\) by performing repetitive experiments. Each experiment, setting up the preparation would have to be performed within a short time interval \(\Delta \tau\).

Then she should compute the subset,
\[
D(q, p, l) = \{ v \in \text{Simp}(A(\sigma)) | q^l_i(v) \in D(q), p^l_j(v) \in D(p, l), i, j = 1, \ldots, r \},
\]
and finally, she should compute for each \(v \in D(q, p, l)\) the solution curve \(\gamma_v(\tau)\) through \(v\) with length \(\tau\), that is, with time-development \(\tau\), ending at \(v(\tau)\). The end-points of all of these curves, would form a subset \(D(q, p, l : \tau)\), and one would expect that the result of letting the phenomenon develop through the time-interval \(\tau\), would give position and momenta results within the subsets,
\[
D(q : \tau) = \{ q^l_i(v) | v \in D(q, p : \tau) \}, \quad D(p, l : \tau) = \{ p^l_j(v) | v \in D(q, p, l : \tau) \}.
\]

The philosophically interesting result would be that no interaction is really local.

One interesting consequence of the above assumption, our Heisenberg uncertainty relation, is that if we are considering a natural phenomenon related to a macroscopic object, i.e. such that all \(|q^l_k - q^l_i|\) are allowed to be, relatively, very big, then we may prepare the object in such a way that all \(|p^l_k - p^l_i|\) are very small. We then have classical situation, where the result would be, relatively, unique! The Big Bang, see the last subsection of this \(\S\), would in this respect be the extreme opposite situation, where we are totally incapable to trace unique curves, \(\gamma\), from the assumed unique point in "configuration space" where BB happens. And, of course, the End of it all, would correspond to a totally homogenous universe, with a uniquely given future!

**Example 3.5.** Let \(C\) be a finite type commutative \(k\)-algebra, say parametrizing an interesting moduli space, and assume it is non-singular, and pick a system of regular coordinates \(\{t_1, t_2, \ldots, t_r\}\) in \(C\). The problem of constructing a dynamical system of interest to physics, has been discussed in the Introduction, and above. We may consider an element \(L \in \text{Ph}(C)\), a Lagrangian, and try to find a force law, with Dirac derivation \(\delta\), such that,
\[
\delta(L) = 0.
\]
We could start with the trivial Lagrangian, \( L := g = \sum_{i=1}^{r} dt_i^2 \in \text{PhC} \). The Lagrange equation becomes, \( 0 = \delta(g) = \sum_{i=1}^{r} (d^2 t_i dt_i + dt_i d^2 t_i) \in \text{PhC} \). with the obvious solution, 
\[
d^2 t_i = 0, \quad i = 1, \ldots, r.
\]
inducing a dynamical structure \((\sigma)\) in \( \text{Ph}(C) \), generated by the relations,
\[
[dt_i, t_j] + [t_i, dt_j], [dt_i, dt_j], i \neq j.
\]
The corresponding dynamic system, \( C(\sigma) \), is the dynamical system for a free particle. Notice however, that, classically, one imposes also the relations, \( [dt_i, t_j] = 0 \) for \( i \neq j \), and \( [dt_i, t_i] = 1 \).
Consider the representations of dimension 1, corresponding to \( \rho = \text{Ham} \), and use Theorem (3.2), with \( n = 1 \). Then, obviously, the Hamiltonian \( Q \) must be a function, and we find,
\[
\tilde{\rho}(dt_i) = [\delta](t_i), 0 = \tilde{\rho}(\delta^2(t_i)) = [\delta](\delta(t_i)).
\]
This fits well with,
\[
[\delta] = \sum_{i=1}^{r} dt_i \frac{\partial}{\partial t_i},
\]
which gives us the canonical symplectic structure on the commutativization of \( \text{Ph}(C) \), the \( C(1) \) for this situation. Notice that the corresponding Poisson bracket now give us,
\[
\{dt_i, t_j\} = \delta_{i,j},
\]
defining a "deformation" of the commutative phase space which is the quotient of \( C(\sigma) \) defined above.

**Connections and the Generic Dynamical Structure.**

Now, let, \( L := g = 1/2 \sum_{i,j=1}^{r} g_{i,j} dt_i dt_j \in \text{PhC}, \) be a Riemannian metric. Recall the formula for the Levi-Civita connection, 
\[
\sum_l g_{l,k} \Gamma^l_{j,i} = 1/2(\frac{\partial g_{k,i}}{\partial t_j} + \frac{\partial g_{j,k}}{\partial t_i} - \frac{\partial g_{i,j}}{\partial t_k}).
\]
Since,
\[
\delta(g) = \sum_{i,j,k=1}^{r} \frac{\partial g_{i,j}}{\partial t_k} dt_k dt_i dt_j + \sum_{i,j=1}^{r} g_{i,j} (d^2 t_i dt_j + dt_i d^2 t_j),
\]
we may plug in the formula,
\[
\delta^2 t_i = -\Gamma^i := -\sum_i \Gamma^i_{i,j} dt_i dt_j,
\]
on the right hand side, and see that we have got a solution of the Lagrange equation,
\[
\delta(L) = 0,
\]
in the commutative situation. This solution has the form of a *force law*,
\[
d^2 t_i = -\Gamma^i := -\sum_i \Gamma^i_{i,j} dt_i dt_j,
\]
generating a dynamical structure \( \langle \sigma \rangle := \langle \sigma(g) \rangle \) of order 2. The dynamic system is, of course, as an algebra,
\[
C(\sigma) = k[t, \xi]
\]
where \( \xi \) is the class of \( dt_j \). The Dirac derivation now has the form,
\[
\delta = \sum_l (\xi_l \frac{\partial}{\partial t_l} - \Gamma^l_i \frac{\partial}{\partial \xi_l}),
\]
and the fundamental vector field \([\delta]\) in \( \text{Simp}_1(C(\sigma)) = \text{Spec}(k[t, \xi]) \), is, of course, the same. Use Theorem 3.4,(ii), and see that \([\delta](g) = 0\), which means that \( g \) is constant along the integral curves of \([\delta]\) in \( \text{Simp}_1(Ph(C)) \), which projects onto \( \text{Simp}_1(C) \) to give the geodesics of the metric \( g \), with equations,
\[
\ddot{t}_l = - \sum_{i,j} \Gamma^l_{i,j} \dot{t}_i \dot{t}_j.
\]

Put \( \delta_i := \frac{\partial}{\partial t_i} \), and consider the Levi-Civita-connection,
\[
\nabla : \Theta_C \rightarrow \text{End}_k(\theta_C)
\]
expressed in coordinates as,
\[
\nabla_{\delta_i}(\delta_j) = \sum_l \Gamma^l_{i,j} \delta_l
\]

Classsically we define the curvature tensor \( R_{i,j}(\delta_k) = \sum_l R^l_{i,j,k} \delta_l \), of a connection \( \nabla \), as the obstruction for \( \nabla \) to be a Lie-algebra homomorphism. We find,
\[
([\nabla_{\delta_i}, \nabla_{\delta_j}] - \nabla_{[\delta_i, \delta_j]})(\delta_k) = \sum_l R^l_{i,j,k} \delta_l.
\]
This, we shall see, is a commutative version of the more precise notion of curvature, related to a more general dynamic system, to be studied below. Recall that the Ricci tensor is given as,
\[
Ric_{i,k}(g) = \sum_j R^j_{i,j,k}
\]
and that, assuming the metric is non-degenerate with inverse \( g^{k,i} \), one defines the scalar curvature of \( g \), as,
\[
S(g) := \sum_{k,i} g^{k,i} Ric_{i,k}.
\]
These are fundamental metric invariants. Recall also Einstein’s equation,
\[
Ric - 1/2S(g)g = U,
\]
where \( U \) is the stress-mass tensor.
A non-degenerate metric, \( g \in Ph(C) \) induces an isomorphism of \( C \)-modules
\[
\Theta_C = Hom_C(\Omega_C, C) \simeq \Omega_C.
\]
Assume first that \( g = 1/2 \sum_{i=1}^d dt_i^2 \), i.e. assume that the space is Euclidean, and pick a basis \( \{ \delta_i := \frac{\partial}{\partial t_i} \} \) of \( \Theta_C \), and a basis \( \{ dt_j \} \) of \( \Omega_C \) such that,

\[
\delta_j(dt_i) = \delta_{i,j}.
\]

Consider a \( C \)-module, \( V \). Any connection \( \nabla \) on \( V \) induces a homomorphism,

\[
\rho := \rho_\nabla : \text{Ph}(C) \to \text{End}_k(V),
\]

with, \( \rho(dt_i) := \nabla_{\delta_i} = \frac{\partial}{\partial t_i} + \nabla_i \). To see this we just have to check that the relations, \([dt_i, t_j] + [t_i, dt_j] = 0\), in \( \text{Ph}(C) \) are not violated. Since we obviously have,

\[
\rho([dt_i, t_j]) = [\nabla_{\delta_i}, \rho(t_j)] = \delta_{i,j},
\]

the homomorphism \( \rho \) is well defined. We are therefore led to consider the dynamical structure on \( C \), generated by the ideal,

\[
(\sigma) := ([dt_i, t_j] - \delta_{i,j}) \subset \text{Ph}^\infty(C).
\]

Since \( \delta(t_i) = [g, t_i] = dt_i \), the Dirac derivation is given by,

\[
\delta = \text{ad}(g).
\]

(\( \sigma \)) is clearly invariant under isometries. Moreover, in \( C(\sigma) \) we have,

\[
\delta^2(t_i) = -1/2 \sum_k (dt_k[dt_i, dt_k] + [dt_i, dt_k]dt_k).
\]

Given any connection, \( \nabla \), on an \( C \)-module, \( V \), and consider the corresponding representation, \( \rho : \text{Ph}(C) \to \text{End}_k(V) \). If \( V \) is of infinite dimension as \( k \)-vector space, we cannot prove that there is a useful moduli space in which \( V \) is a point. However we now know that \( \rho \) is singular. This follows since there exist a Hamiltonian, \( Q := \rho(g) \in \text{End}_k(V) \), such that for all \( a \in \text{Ph}(C) \),

\[
\rho(da) = [Q, \rho(a)].
\]

In particular we have, \( \rho(dt_i) = \nabla_{\delta_i} = [Q, t_i] \). This imply,

\[
Q = 1/2 \sum_i \nabla_{\delta_i}^2.
\]

Thus for any connection \( \nabla \) we find a force law, in \( \text{End}_k(V) \), given by,

\[
\rho_\nabla(d^2t_i) = -1/2 \sum_{j=1}^d \nabla_{\delta_j} [\nabla_{\delta_i}, \nabla_{\delta_j}] - 1/2 \sum_{j=1}^d [\nabla_{\delta_i}, \nabla_{\delta_j}] \nabla_{\delta_j}.
\]

We shall in this situation use the notations,

\[
R_{i,j} := [dt_i, dt_j] \in \text{Ph}(C), \ F_{i,j} := [\nabla_{\delta_i}, \nabla_{\delta_j}] \in \text{End}_C(V),
\]
GEOMETRY OF TIME-SPACES.

\( F_{i,j} \) being the curvature tensor, of the connection. Below we shall come back to these notions in the general situation.

Since we now have,

\[
\nabla \delta_j F_{i,j} = F_{i,j} \nabla \delta_j + (\frac{\partial F_{i,j}}{\partial t_j} + [\nabla_j, F_{i,j}])
\]

we find the following equation,

\[
\rho \nabla (d^2 t) = -F \rho (dt) - q,
\]

where \( q \) (by definition) is the charge of the field. \( q \) is a vector, the coordinates of which,

\[
q_i = \frac{1}{2} \sum_{j=1}^{r} (\frac{\partial F_{i,j}}{\partial t_j} + [\nabla_j, F_{i,j}]),
\]

are endomorphisms of the bundle.

As we shall see in several examples, the dynamic structure defined above is sufficiently general to serve as basis for what is usually called quantization, of the electromagnetic field. For the gravitational field, we have to do some more work.

Let us look at the last first. We then have to consider a general, non-degenerate, metric, \( g = \frac{1}{2} \sum_{i=1}^{d} g_{i,p} dt_i dt_j \), and the corresponding dynamical system, \( (\sigma) = ([dt_i, t_j] - g^{i,j}) \). Again it is easy to see that this is not violating the relations, \([dt_i, t_j] + [t_i, dt_j] = 0\) of \( Ph(C) \). Notice also that in \( C(\sigma) \) we have,

\[
[[dt_i, dt_j], t_k] = g^{il} \frac{\partial g_{j,p}}{\partial t_l} - g^{jk} \frac{\partial g_{i,p}}{\partial t_k},
\]

meaning that the curvature does not commute with the action of \( C \). Introducing \( dt_i := \sum g_{i,p} dt_p \), we find that \([dt_i, t_j] = \delta_{i,j}\). Moreover, if we let,

\[
\bar{g} := \frac{1}{2} \sum_{i=1}^{d} dt_i^2,
\]

we find, \( ad(\bar{g})(t_i) = dt_i \). Using the above, we find that there is a one-to-one correspondence between connections \( \nabla \) on a \( C \)-module \( V \) and morphisms,

\[
\rho : C(\sigma) \to \text{End}_k(V),
\]

defined by,

\[
\rho(\sigma_i) = \sum_j g^{j,i} \nabla \delta_i = \nabla \xi_i,
\]

where \( \xi_i = \sum_j g^{i,j} \delta_j \) is the dual to \( dt_i \).

Consider now the Levi-Civita connection \( \nabla \delta_i = \frac{\partial}{\partial t_i} + \nabla_i \), where,

\[
\nabla_i \in \text{End}_C(\Theta_C),
\]
is given by the matrix formula, $\nabla_i = (\Gamma^q_{p,i})$. Put,

$$T := \frac{1}{2} \sum_{j,k} \frac{\partial g^{j,k}}{\partial t_j} dt_k = \frac{1}{2} \sum_{j,k,l} \frac{\partial g^{j,k}}{\partial t_j} g_{k,l} dt_l$$

and consider the inner derivation,

$$\delta := ad(g - T),$$

then after a dull computation, using the well known formula for Levi-Civita connection,

$$\frac{\partial g_{i,j}}{\partial t_k} = \sum_l (\Gamma^l_{k,i} g_{l,j} + \Gamma^l_{k,j} g_{i,l}),$$

we obtain, in $C(\sigma)$,

$$T := -\frac{1}{2} \sum_{k,l} \Gamma^k_{k,l} dt_l + \sum_{k,p,q} g^{k,q} \Gamma^p_{k,q} g_{p,l} dt_l,$$

$$\delta(t_i) := ad(g - T)(t_i) = dt_i, \ i = 1, ..., d.$$

Therefore we have a well-defined dynamical structure $(\sigma)$, with Dirac derivation $\delta := ad(g - T)$. It is easy to see that $(\sigma)$ is invariant w.r.t. isometries.

Moreover, the representation, $\rho$ of $C(\sigma)$, defined on $\Theta_C$, by the Levi-Civita connection, has a Hamiltonian,

$$Q := \rho(g - T) = \frac{1}{2} \sum_{i,j} g^{ij} \nabla_i \nabla_j,$$

i.e. the generalized Laplace-Beltrami-operator, which is also invariant w.r.t. isometries, although the proof demands some algebra. Put

$$\Gamma^i_{p,q} := \sum_{l,r} g^{r,i} \Gamma^l_{r,p} g_{l,q},$$

then,

$$T = \sum_i T_i dt_i$$

$$T_i = -\frac{1}{2} \sum_j (\Gamma^j_{i,j} + \Gamma^j_{j,i}) = -\frac{1}{2} (\text{trace} \nabla_i + \text{trace} \nabla_i).$$

Since $\delta(t_i) := ad(g - T)(t_i) = dt_i$, the general force law, in $C(\sigma)$, looks like,

$$d^2 t_i = [g - T, dt_i] = -\frac{1}{2} \sum_{p,q} (\Gamma^i_{i,p} + \Gamma^i_{q,p}) dt_p dt_q$$

$$+ \frac{1}{2} \sum_{p,q} g_{p,q} (R_{p,i} dt_q + dt_p R_{q,i})$$

$$+ [dt_i, T],$$

where, as above, $R_{i,j} = [dt_i, dt_j]$. We shall want to write this in the following form,
**General Force Law.** In $C(\sigma)$ we have the following force law,

\[
d^2 t_i = - \sum_{p,q} \Gamma^i_{p,q} dt_p dt_q - 1/2 \sum_{p,q} g_{p,q}(F_{i,p} dt_q + dt_p F_{i,q}) + 1/2 \sum_{l,p,q} g_{p,q}[dt_p, (\Gamma^q_{l,-} - \Gamma^q_{l,i})] dt_l + [dt_l, T].
\]

**Proof.** As we have seen, the dual of $dt_i$ is $\xi_i = \sum_l g^{i,l} \frac{\partial}{\partial t_l}$, therefore

\[
[\xi_i, \xi_j] = \sum_{l,k} g^{i,l} \frac{\partial g^{j,k}}{\partial t_l} \frac{\partial}{\partial t_k} - g^{j,k} \frac{\partial g^{i,l}}{\partial t_k} \frac{\partial}{\partial t_l}
\]

is dual to

\[
\sum_{l,k,p} g^{i,l} \frac{\partial g^{j,k}}{\partial t_l} g_{k,p} dt_p - g^{j,k} \frac{\partial g^{i,l}}{\partial t_k} g_{k,p} dt_p.
\]

Using the above equations relating the derivatives of $g^{i,j}$ to the Levi-Civita connection, we find,

\[
\sum_{l,k,p} g^{i,l} \frac{\partial g^{j,k}}{\partial t_l} g_{k,p} dt_p - g^{j,k} \frac{\partial g^{i,l}}{\partial t_k} g_{k,p} dt_p = \sum_p (\Gamma^{j,i}_p - \Gamma^{i,j}_p) dt_p
\]

where $\Gamma^{j,i}_p = \sum_s g^{j,k} \Gamma^{i,k}_p$. Let now,

\[
F_{i,j} := R_{i,j} - \sum_p (\Gamma^{j,i}_p - \Gamma^{i,j}_p) dt_p.
\]

For every connection $\nabla$ on a $C$-module $E$, given by a representation, $\rho_E$, we obtain,

\[
\rho_E(F_{i,j}) = [\nabla_{\xi_i}, \nabla_{\xi_j}] - \nabla_{[\xi_i, \xi_j]},
\]

i.e. the curvature of the connection, $F(\xi_i, \xi_j)$.

Now, plug this in the force law above, i.e. write,

\[
1/2 \sum_{p,q} g_{p,q}(R_{p,i} dt_q + dt_p R_{q,i}) =
\]

\[
1/2 \sum_{p,q} g_{p,q}((R_{p,i} - \sum_l (\Gamma^{i,p}_l - \Gamma^{p,i}_l) dt_l) dt_q + dt_p(R_{q,i} - \sum_l (\Gamma^{i,q}_l - \Gamma^{q,i}_l) dt_l))
\]

\[
+ 1/2 \sum_{l,p,q} g_{p,q}(\sum_l (\Gamma^{i,p}_l - \Gamma^{p,i}_l) dt_l) dt_q
\]

\[
+ 1/2 \sum_{p,q} g_{p,q}(\sum_l (\Gamma^{i,q}_l - \Gamma^{q,i}_l) dt_l),
\]

and use

\[
+ 1/2 \sum p g_{p,q}(\sum_l (\Gamma^{i,p}_l - \Gamma^{p,i}_l) dt_l) dt_q = 1/2 \Gamma^{i,q}_{l,q} dt_l dt_q - 1/2 \Gamma^{i,q}_{q,l} dt_l dt_q
\]

\[
+ 1/2 \sum q g_{p,q}(\sum_l (\Gamma^{i,q}_l - \Gamma^{q,i}_l) dt_p dt_l) = 1/2 \Gamma^{i,q}_{l,p} dt_p dt_l - 1/2 \Gamma^{i,q}_{p,l} dt_p dt_l.
\]
Finally use,
\[ dt_p(\Gamma^{i,q}_l - \Gamma^{q,i}_l) = (\Gamma^{i,q}_l - \Gamma^{q,i}_l)dt_p + [dt_p, (\Gamma^{i,q}_l - \Gamma^{q,i}_l)]. \]

□

We shall consider the above formula as a general Force Law, in \( Ph(C) \), induced by the metric \( g \). As explained before, this means the following: Let \( \mathfrak{c} \) be the \( \delta \)- stable ideal generated by this equation in \( Ph^\infty(C) \). Since the force law above holds in the dynamical system defined by \( (\sigma) \), we obviously have \( \mathfrak{c} \subset (\sigma) \), and we may hope this new dynamical system leads to a Quantum Field Theory, as defined above, with new and interesting properties. We know that this dynamical structure reduces to the generic structure for connections, i.e. for the singular cases.

Notice that this force law reduces to an equation of motion in General Relativity, in the representation-dimension 1 case, i.e. in the commutative case. More interesting is that it leads to both Lorentz force law, and to Maxwell’s field-equations for Electro-Magnetism in the classical flat-space-situation, see Examples (3.16) and (3.17).

An easy calculation in \( C(\sigma) \), shows that,
\[ [T, dt_i] = 1/2 \sum_j T_j R_{j,i} - 1/2 \sum_{j,l} \frac{\partial T_j}{\partial t_l} g^{j,i} dt_j =: q_i. \]

But, be careful, these \( q_i \)'s no longer vanish in the classical phase-space, i.e. in the commutativization of \( Ph(C) \).

Now, choose a representation \( \rho_E : C(\sigma) \rightarrow \text{End}_k(E) \), i.e. a connection \( \nabla \), on a \( C \)-module \( E \). The generalized curvature \( F_{i,j} =: F(\xi_i, \xi_j) \in \text{End}_C(E) \) maps to the classical one, and we observe that there is an “interaction” between the geometry, defined by the metric \( g \) and the ”geometry” defined by the connection \( \nabla \). Our Force Law above will now take the form,
\[ \rho_E(d^2t_i) + \sum_{p,q} \Gamma^{i}_{p,q} \nabla_{\xi_p} \nabla_{\xi_q} \]
\[ = 1/2 \sum_p F_{p,i} \nabla_{\delta_p} + 1/2 \sum_p \nabla_{\delta_p} F_{p,i} + 1/2 \sum_{i,q} \delta_q(\Gamma^{i,q}_l - \Gamma^{q,i}_l) \nabla_{\xi_l} + [\nabla_{\xi_i}, \rho_E(T)]. \]

See also Example (3.16).

Notice also, that for the Levi-Civita connection, there is a possible relationship between this formula and the Einstein field equation. See [Sachs-Wu], Proposition 4.2.2., p.114. If, above we assume that we are in a geodesic reference frame, i.e. along a geodesic \( \gamma \) in our space \( \text{Sim}p_1(C) \), then an average of the excess-relative-acceleration, i.e. of \( d^2t_i + 1/2 \sum_{p,q}(\Gamma^{i}_{p,q} + \Gamma^{q}_{p,i})dt_p dt_q \), evaluated in \( \Theta_C|\gamma \), is proved to be given by the \( \text{Ric} \) tensor. But, above this relative-acceleration is, for any representation corresponding to a connection \( \nabla \), equal to
\[ 1/2 \sum_p F_{p,i} \nabla_{\delta_p} + 1/2 \sum_p \nabla_{\delta_p} F_{p,i} + 1/2 \sum_{i,q} \delta_q(\Gamma^{i,q}_l - \Gamma^{q,i}_l) dt_l + [\nabla_{\xi_i}, \rho_E(T)]. \]
Since this excess-relative-acceleration, representing a *tidal force*, should be a measure of the inertial mass present, it is tempting to consider this force law as a generalized, quantized, Maxwell-Einstein’s equation. The reference to Maxwell here is natural, since if the bundle $E = \Theta_C$ above is the tangent bundle, and we consider the connection, given by the potential $A = (A_1, \ldots, A_n), A_i \in C$, then the resulting curvature is the electro-magnetic force field. See Example (3.17) for the notion of Charge, and see Example (3.16), where the problem of Mass will be addressed.

In this generality, it is not really meaningful to ask for invariance of this general Force Law, w.r.t. isometries. This is linked to the fact that, in general, this force law, considered as a dynamical structure on $C$, may have non-singular finite-dimensional representations, and then invariance under isometries of $Simp_1(C)$ is not the proper question to pose. We shall come back to this later, but see Example (3.16) for relations to Newton and Kepler’s laws.

Notice that applying $\rho$, corresponding to the Levi-Civita connection, the above translate into,

$$\rho(\delta t_i) = \sum_{j=1}^d [Q, g^{ij} \nabla \delta_j]$$

where $Q$ is the Laplace-Beltrami operator.

Before we turn to situations requiring a general quantum theoretical treatment, let us go back to the discussion above, about how to look at parsimony, via Lagrange functions or via dynamical systems. We claimed that the integral curves of the vector field

$$\delta = \sum_l (Q \frac{\partial}{\partial t_l} - \Gamma_l^{ij} \frac{\partial}{\partial \xi_l})$$

in $Simp_1 Ph(C)$, projects onto the geodesics of the metric $g$ in $Simp_1(C)$. These geodesics are assumed to be trajectories of free test particles in the geometric space $Simp_1(C)$ outfitted with the the metric $g$. As such they must be curves parametrized by some clock parameter $\tau$. Since quantum field theory is assumed to model such movements, we now have two different methods to pick out such trajectories, i.e. to find the solutions $M(C, k[\tau]) \subset F(C, k[\tau])$. One, using dynamical systems, the force law $\sum \Gamma_{ij} dx_i dx_j = - \sum \Gamma_{ij} dt_i dt_j$, deduced above for $C$, and the obvious $d^2 \tau = 0$, for the free particle modeled by $B := k[\tau]$, the other using the Euler-Lagrange equations as described above.

In the first case we have the equation,

$$\delta_B(\phi) - \phi(\delta_C) = 0,$$

which evaluated at $d^2 t_i$ give us,

$$\phi(- \sum \Gamma_{p,q} dt_p dt_q) = (\frac{\partial}{\partial \tau})^2(\phi) d\tau^2,$$

with the resulting equation,

$$\ddot{\phi} = - \sum \Gamma_{p,q} \dot{\phi}_p \dot{\phi}_q,$$

i.e. the equations for a geodesic.
In the second case, we should use the obvious derivation $\zeta = \frac{\partial}{\partial \tau}$, the corresponding representation $\rho_\zeta : \text{Ph}(k[\tau]) \to k[\tau]$, pick the Lagrangian $L := g$, and look at the resulting action and corresponding Euler-Lagrange equations. We obtain,

$$
L = \sum g_{p,q}\dot{\phi}_p \dot{\phi}_q
$$

$$
S = \int \sum g_{p,q}\dot{\phi}_p \dot{\phi}_q d\tau,
$$

together with the Euler-Lagrange equations,

$$
\frac{\partial g}{\partial \phi_i} - \frac{\partial}{\partial \tau} (\frac{\partial g}{\partial \dot{\phi}_i}) = 0
$$

which reduces to the same equations for geodesics.

**Example 3.6.** With this done, let us consider some easy examples of quantum theory, first in dimension 1, and still in rank 1. That is, we start with the $k$-algebra $C = k \langle x \rangle = k[x]$, and consider the classical Lagrangians,

$$
L = \frac{1}{2} dx^2 - V(x) \in \text{Ph}C.
$$

The corresponding dynamical system $\sigma$, deduced from the Lagrange equations, as above, is given by the force law,

$$
d^2 x = \frac{\partial V}{\partial x},
$$

and is of order 2, so the algebra of interest is,

$$
C(\sigma) = \text{Ph}C = k < x, dx > \simeq k < x_1, x_2 >.
$$

Notice that the classical Hamiltonian $H := dx^2 - L$, is not an invariant, i.e. $\delta(H) \neq 0$.

Let us first compute the particles in rank 1 for some cases, and let us start with $V(x) = 1/2 x^2$, i.e. the classical oscillator. The fundamental equation of the dynamical system is,

$$
\delta = [\delta] + [Q, -],
$$

where, in dimension 1, the endomorphism $Q$ obviously commutes with the actions of $x_i$, $i = 1, 2$. To solve the equation above, we may therefore forget about $Q$, so we are left with the vector fields,

$$
[\delta] = \xi.
$$

The space, $\text{Simp}_1(C(\sigma))$, is just the ordinary phase space, $\text{Simp}_1(k[x, dx])$. Put as above, $x_1 := x, x_2 := dx$. We must solve the equations,

$$
\delta(x) = [\delta](x) = [\delta](x_1)
$$

$$
\delta^2(x) = [\delta](dx) = [\delta](x_2)
$$

We can obviously pick,

$$
\delta_i = \chi_i = \frac{\partial}{\partial x_i}.
$$
so we must have
\[ [\delta] = \xi_1 \frac{\partial}{\partial x_1} + \xi_2 \frac{\partial}{\partial x_2}. \]
In the case of the potential, \( V = \frac{1}{2}x^2 \), we get the equations,
\[
\begin{align*}
    x_2 &= [\delta](x) = [\delta](x_1) = \xi_1 \\
    x_1 &= [\delta](dx) = [\delta](x_2) = \xi_2
\end{align*}
\]
Therefore the fundamental vector field is,
\[ \xi = x_2 \frac{\partial}{\partial x_1} + x_1 \frac{\partial}{\partial x_2}, \]
i.e. we find hyperbolic motions in the phase space, with general solutions,
\[
    x = x_1 = r \cosh(t + c), \quad dx = x_2 = r \sinh(t + c)
\]
which is what we expected.

In the case of the oscillator, \( V = -\frac{1}{2}x^2 \), we get the equations,
\[
\begin{align*}
    x_2 &= [\delta](x) = [\delta](x_1) = \xi_1 \\
    -x_1 &= [\delta](dx) = [\delta](x_2) = \xi_2
\end{align*}
\]
Therefore the fundamental vector field is,
\[ \xi = x_2 \frac{\partial}{\partial x_1} - x_1 \frac{\partial}{\partial x_2}, \]
i.e. we find circular motions in the phase space, with general solutions,
\[
    \gamma : x = x_1 = r \cos(t + c), \quad dx = x_2 = -r \sin(t + c),
\]
which is also what we expected.

Consider now the versal family restricted to \( \gamma \),
\[ \tilde{\rho}_\gamma : k < x, dx > \rightarrow \text{End}_\gamma(\tilde{\mathcal{V}}|\gamma), \]
and a state \( \psi(t) \in \tilde{\mathcal{V}}|\gamma \). If \( Q \), restricted to \( \gamma \), is multiplication by \( \kappa(t) \), (in physics, one usually puts \( \kappa(t) = i\kappa \)), then the Schrödinger equation becomes,
\[ \frac{\partial}{\partial t} \psi = \kappa(t)\psi \]
so that we should have,
\[ \psi(t) = \exp(i \int_\gamma \kappa). \]
This will turn out much nicer if we extend the action of \( k < x_1, x_2 > \) to \( \tilde{\mathcal{V}}_C \), and put \( Q \), restricted to \( \gamma \), equal to multiplication by \( i\kappa \). Then we find the reasonable result,
\[ \psi(t) = \exp(i \int_\gamma \kappa). \]
See again [La 5].

In the repulsive, resp. attractive, Newtonian case, with $V = \pm 1/x$, we find,

$$x_2 = [\delta](x) = [\delta](x_1) = \xi_1$$
$$\epsilon(1/x_1^2) = [\delta](dx) = [\delta](x_2) = \xi_2, \quad \epsilon = +, - .$$

Therefore the fundamental vector field is,

$$\xi = x_2 \frac{\partial}{\partial x_1} + \epsilon(1/x_1^2) \frac{\partial}{\partial x_2}$$

with the classical solution,

$$x = \epsilon(9/2)t^{2/3} .$$

In higher dimensions, say in the case of our "toy" model $H$, of the Introduction, this rank 1 theory reduces to the wave-mechanics of de Broglie. Recall that $Ph(H)$ has a natural action of $U(1)$, and there is therefore a natural complex structure on the tangent space $T_H$. Consider the trivial versal family,$$Ph(H) \rightarrow \text{End}_{C(1)}(C(1), C(1)),$$where we may assume $C(1)$ is a complex vector space. Any order 2 dynamical structure defined on $H$, will induce a vector field in $C(1) = Ph(H)$, which in the case of the Levi-Civita connection, considered as force law, makes the integral curves geodesics. From this Klein-Gordon follows in a natural way, and here we may also discuss interference and diffraction of light, see [La 6].

**Example 3.7.** Now let us go back to the case of $A = k < x_1, x_2 >$, the free non-commutative $k$-algebra on two symbols, and the rank $n = 2$, see (2.14). We found,

$$C(2) \simeq k[t_1, t_2, t_3, t_4, t_5].$$

locally, in a Zariski neighborhood of the origin. The versal family $V$, is defined by the actions of $x_1, x_2$, given by,

$$X_1 := \begin{pmatrix} 0 & 1 + t_3 \\ t_5 & t_4 \end{pmatrix}, \quad X_2 := \begin{pmatrix} t_1 & t_2 \\ 1 + t_3 & 0 \end{pmatrix}.$$  

The Formanek center, in this case, is cut out by the single equation:

$$f := \det[X_1, X_2] = -((1 + t_3)^2 - t_2 t_5)^2 + (t_1(1 + t_3) + t_2 t_4)(t_4(1 + t_3) + t_1 t_5).$$

and

$$\text{tr}X_1 = t_4, \quad \text{tr}X_2 = t_1,$$
$$\det X_1 = -t_5 - t_3 t_5, \quad \det X_2 = -t_2 - t_2 t_3,$$
$$\text{tr}(X_1X_2) = (1 + t_3)^2 + t_2 t_5,$$

so the trace ring of this family is

$$k[t_1, t_2 + t_2 t_3, 1 + 2t_3 + t_3^2 + t_2 t_5, t_4, t_5 + t_3 t_5] =: k[u_1, u_2, u_3, u_4, u_5],$$
with,
\[ u_1 = t_1, \ u_2 = (1 + t_3)t_2, \ u_3 = (1 + t_3)^2 + t_2t_5, \ u_4 = t_4, \ u_5 = (1 + t_3)t_5, \]
and \( f = -u_2^4 + 4u_2u_5 + u_1u_4 + u_1^2u_5 + u_2u_4^2 \). Moreover, \( k[u] \) is algebraic over \( k[u] \), with discriminant, \( \Delta := 4u_2u_5(u_2^3 - 4u_2u_5) = 4(1 + t_3)^2t_2t_5((1 + t_3)^2 - t_2t_5)^2, \) and there is an \( \acute{e}tale \) covering,
\[ A^5 - V(\Delta) \rightarrow \text{Simp}_2(A) - V(\Delta). \]

Notice that if we put \( t_1 = t_4 = 0 \), then \( f \) divides \( \Delta \).

**Example 3.8.** For the oscillator, given by the Lagrangian, \( L = \frac{1}{2}dx^2 - \frac{1}{2}x^2 \), and with the force law, \( d^2x = x \), in rank 2, things are more difficult. As we have computed above, we have found a (partial) versal family of \( \text{Simp}_2(Ph k[x]) \), given by,
\[ x = \begin{pmatrix} 0 & 1 + t_3 \\ t_5 & t_4 \end{pmatrix}, dx = \begin{pmatrix} t_1 & t_2 \\ 1 + t_3 & 0 \end{pmatrix} \]
and we may chose,
\[ \chi_i = \frac{\partial}{\partial t_i}, \ i = 1, 2, ..., 5. \]
and, obviously,
\[ \delta_i = \frac{\partial}{\partial t_i}, \ i = 1, 2, ..., 5. \]
The fundamental vector fields will have the form,
\[ [\delta] = \sum \xi_i \delta_i, \ \xi = \sum \xi_i \frac{\partial}{\partial t_i}, \]
with 5 unknowns, \( \xi_i, i = 1, 2, ..., 5 \) Moreover,
\[ Q = \begin{pmatrix} q_{1,1} & q_{1,2} \\ q_{2,1} & q_{2,2} \end{pmatrix}, \]
with 4 unknowns \( q_{i,j}, i = 1, 2, j = 1, 2 \). Now, recall that \( Q \) can only be determined up to a central element from \( M_2(C) \), i.e. we have 8 essential unknowns, \( \xi_i, i = 1, 2, 3, 4, 5 \) and \( (q_{1,1} - q_{2,2}), q_{1,2}, q_{2,1} \) in the two matrix equations,
\[ \delta(x) = dx = [\delta](x) + [Q, x] \]
\[ \delta^2(x) = x = [\delta](dx) + [Q, dx] \]
On the right hand side of the equations we have the terms,
\[ [\delta](x) = \sum \xi_i \delta_i(\begin{pmatrix} 0 & 1 + t_3 \\ t_5 & t_4 \end{pmatrix}) = \begin{pmatrix} 0 & \xi_3 \\ \xi_5 & \xi_4 \end{pmatrix} \]
\[ [\delta](dx) = \sum \xi_i \delta_i(\begin{pmatrix} t_1 & t_2 \\ 1 + t_3 & 0 \end{pmatrix}) = \begin{pmatrix} \xi_1 & \xi_2 \\ \xi_3 & 0 \end{pmatrix} \]
and the terms,

\[ [Q, x] = \begin{pmatrix} t_5q_{1,2} - (1 + t_3)q_{2,1} & (1 + t_3)q_{1,1} + t_4q_{1,2} - (1 + t_3)q_{2,2} \\ t_5q_{2,2} - t_5q_{1,1} - t_4q_{2,1} & (1 + t_3)q_{2,1} - t_5q_{1,2} \end{pmatrix} \]

\[ [Q, dx] = \begin{pmatrix} (1 + t_3)q_{1,2} - t_2q_{2,1} & t_2q_{1,1} - t_1q_{1,2} - t_2q_{2,2} \\ t_1q_{2,1} + (1 + t_3)q_{2,2} - (1 + t_3)q_{1,1} & t_2q_{2,1} - (1 + t_3)q_{1,1} \end{pmatrix}, \]

and on the left side, we have,

\[ \delta(x) = dx = \begin{pmatrix} t_1 & t_2 \\ 1 + t_3 & 0 \end{pmatrix} \]

\[ \delta^2(x) = x = \pm \begin{pmatrix} 0 & 1 + t_3 \\ t_5 & t_4 \end{pmatrix}. \]

Writing up the matrix for the corresponding linear equation, we find that the determinant of the $8 \times 8$ matrix turns out to be easily computed, it is,

\[ D = 2(1 + t_3)(t_2t_5 - (1 + t_3)^2). \]

Notice that $D$ is a divisor in the discriminant, $\Delta = 4(1 + t_3)^2t_2t_5((1 + t_3)^2 - t_2t_5)^2$, see (3.5). Moreover we find,

\[ q_{1,1} - q_{2,2} = D^{-1}(-(1 + t_3)(t_1^2 + t_4^2) + (t_2 - t_5)(t_2t_5 - (1 + t_3)^2 - t_1t_4)) \]

\[ q_{1,2} = D^{-1}(2(1 + t_3)(t_1t_2 + (1 + t_3)t_4)) \]

\[ q_{2,1} = D^{-1}(2(1 + t_3)(t_4t_5 + t_1(1 + t_3))) \]

\[ \xi_1 = t_2q_{2,1} - (1 + t_3)q_{1,2} \]

\[ \xi_2 = -t_2(q_{1,1} - q_{2,2}) + t_1q_{1,2} + (1 + t_3) \]

\[ \xi_3 = (1 + t_3)(q_{1,1} - q_{2,2}) + t_1q_{2,1} + t_5 \]

\[ \xi_4 = t_5q_{1,2} - (1 + t_3)q_{2,1} \]

\[ \xi_5 = t_5(q_{1,1} - q_{2,2}) + t_4q_{2,1} + (1 + t_3) \]

See that $\xi_1 = \xi_4 = 0$ imply,

\[ ((1 + t_3)^2 - t_2t_5)q_{1,2} = ((1 + t_3)^2 - t_2t_5)q_{2,1} = 0, \]

and, since we assume that $\Delta \neq 0$, therefore, $((1 + t_3)^2 - t_2t_5) \neq 0$, and so $q_{1,2} = q_{2,1} = 0$, this also implies that $t_1 = t_4 = 0$. Therefore the singularities of $\xi$ are given, by,

\[ t_2 = -(1 + t_3), \ t_5 = +(1 + t_3), \]

or, up to isomorphisms, uniquely, by the representation,

\[ x = \begin{pmatrix} 0 & 1 \\ 1 & 0 \end{pmatrix} \]

\[ dx = \begin{pmatrix} 0 & -1 \\ 1 & 0 \end{pmatrix} \]

\[ Q = \begin{pmatrix} q_{1,1} & 0 \\ 0 & q_{1,1} + 1 \end{pmatrix}. \]
corresponding to \( t_1 = 0, t_2 = -1, t_3 = 0, t_4 = 0, t_5 = 1 \). Notice that in this case we find, in all ranks, that \( f_\hbar := \rho(x + dx) \), is an eigenvector for \([Q, -] \) with \( f_\hbar \) is the quantum counting operator.

Let us pause a little, to compute the gradient of the action, \( S = \text{Trace}(\mathcal{L}) \). Since \( L = 1/2dx^2 + 1/2x^2 \), this is easy, and we find,

\[
S = 1/2(t_1^2 + 2t_2(1 + t_3) + 2t_3(1 + t_3) + t_4^2), \nabla S = (t_1, (1 + t_3), (t_2 + t_5), t_4, (1 + t_3)),
\]

which, clearly is different from the vector field \( \xi \) above, see the Introduction. But the singularities, obtained by solving \( \nabla S = 0 \), for,

\[
\begin{pmatrix}
  0 & 1 + t_3 \\
  t_5 & t_4
\end{pmatrix}
\]

\[
\begin{pmatrix}
  t_1 & t_2 \\
  1 + t_3 & 0
\end{pmatrix}
\]

gives us,

\[
\begin{pmatrix}
  0 & 0 \\
  -t_2 & 0
\end{pmatrix}
\]

\[
\begin{pmatrix}
  0 & t_2 \\
  0 & 0
\end{pmatrix}
\]

which is isomorphic to the singularity for \( \xi \), after the coordinate change, \( a^+ := 1/2(x + dx), a := 1/2(x - dx) \), with the same Hamiltonian \( Q \). Notice, however, that even though the Formaneck center, \( f \), is non-vanishing, our family is not good at this point. Since \( 1 + t_3 = 0 \), the discriminant \( \Delta = 0 \), and so our family is not étale at this point.

Now, to find the integral curves of the vector field \( \xi \), we must solve the obvious system of differential equations,

\[
\frac{\partial t_i}{\partial \tau} = \xi_i, i = 1, \ldots, 5.
\]

It turns out that we are mostly interested in the solutions for which there exists singular point, corresponding to \( t_1 = t_4 = 0 \). If they exist they look like,

\[
\frac{\partial t_1}{\partial \tau} = \xi_1 = 0
\]

\[
\frac{\partial t_2}{\partial \tau} = \xi_2 = -t_2(t_2 - t_5)(2 + 2t_3)^{-1} + (1 + t_3)
\]

\[
\frac{\partial t_3}{\partial \tau} = \xi_3 = 1/2(t_2 - t_5) + t_5
\]

\[
\frac{\partial t_4}{\partial \tau} = \xi_4 = 0
\]

\[
\frac{\partial t_5}{\partial \tau} = \xi_5 = t_5(t_2 - t_5)(2 + 2t_3)^{-1} + (1 + t_3).
\]

And these equations are obviously consistent with the conditions \( t_1 = t_4 = 0 \).

Introducing new variables,

\[
y_1 = (t_2 - t_5)
\]

\[
y_2 = (t_2 + t_5)
\]

\[
y_3 = (2 + 2t_3)
\]

so that,

\[
t_2 = 1/2 (y_2 + y_1)
\]

\[
t_5 = 1/2 (y_2 - y_1)
\]

\[
t_3 = 1/2 y_3 - 1.
\]
things look nicer. We find,
\[ \xi_2 = -\frac{1}{2} (y_1 + y_2) + \frac{1}{2} y_3 \]
\[ \xi_3 = \frac{1}{2} y_2 \]
\[ \xi_5 = \frac{1}{2} (y_2 - y_1) y_1 y_3^{-1} \]

In the new coordinates the system of equations above reduces to,
\[ y_1 \frac{\partial y_1}{\partial \tau} - y_2 \frac{\partial y_2}{\partial \tau} + y_3 \frac{\partial y_3}{\partial \tau} = 0 \]
\[ y_1^{-1} \frac{\partial y_1}{\partial \tau} + y_3^{-1} \frac{\partial y_3}{\partial \tau} = 0. \]

The integral curves are therefore intersections of the form,
\[ C(c_1, c_2) := V(y_1^2 - y_2^2 + y_3^2 = c_1) \cap V(y_1 y_3 = c_2). \]

Moreover, the stratum at infinity, given by \( f = 0 \), where \( f \) is the Formanek center, is now easily computed, in terms of the new coordinates it is given as,
\[ f = -1/16(y_1^2 - y_2^2 + y_3^2)^2 \]

This shows that a particle corresponding to an integral curve \( \gamma := C(c_1, c_2) \), with \( c_1 \neq 0 \) lives eternally, as it should. Its completion does not intersect the Formanek center, the stratum at infinity.

An easy calculation gives us, see (3.6),
\[ 16(y_1^2 - y_2^2 + y_3^2)^2 = -u_3^2 + 4u_2u_5 \]
\[ y_1 y_3 = 2(u_2 - u_5), \]

where the \( u \)-coordinates are those of the trace ring, see (3.6). The integral curves of the harmonic oscillator will be therefore be plane conic curves in the part of \( \text{Simp}_2(Phk[x]) \), where \( \Delta \neq 0 \), \( u_1 = u_4 = 0 \), given by,
\[ u_3^2 - 4u_2u_5 = c_3, \quad (u_2 - u_5) = c_4. \]

Here \( c_3 \neq 0, c_4 \) are constants. Notice also that our special point, the singularity for \( \xi \), given by \( y_1 = -2, y_2 = 0, y_3 = 2 \), sits on the curve defined by \( c_1 = 8, c_2 = -4 \), corresponding to \( c_3 = 32, c_4 = -2 \).

In the new, \( y \)-coordinates, the versal family of \( \text{Simp}_2(Phk[x]) \), lifted to \( U(2) \), and restricted to \( t_1 = t_2 = 0 \), is given by,
\[ x = \begin{pmatrix} 0 & \frac{1}{2}y_3 \\ \frac{1}{2}(y_2 - y_1) & 0 \end{pmatrix}, \quad dx = \begin{pmatrix} 0 & \frac{1}{2}(y_1 + y_2) \\ \frac{1}{2}y_3 & 0 \end{pmatrix}. \]

Moreover along the curve \( \gamma \), defined by \( c_1 = 8, c_2 = -4 \), which is given by the equations,
\[ y_3 = -4y_1^{-1}, \quad y_2^2 = y_1^2 + 16y_1^{-2} - 8 = (y_1^2 - 4)^2y_1^{-2}, \]
the vectorfield $\xi$ is given by,

$$
\xi = -\frac{1}{4}(y_1 + 2)(y_1 - 2)y_1 \frac{\partial}{\partial y_1}
$$

or,

$$
\xi = \frac{3}{4}(y_1 + 2)(y_1 - 2)y_1 \frac{\partial}{\partial y_1},
$$

depending on which root we choose for $y_2$ above. The corresponding time along $\gamma$, is then given as, $\tau = -\log(y_1) + 1/2\log(\frac{y_1}{y_1 - 2})$, respectively $\tau = 1/3\log(y_1) - 1/6\log(y_1 + 2) - 1/6\log(\frac{y_1}{y_1 - 2})$, both with a singularity at $y_1 = -2$, $y_1 = 2$, corresponding to the same unique singularity of $\xi$, in $\text{Simp}_2(\text{Ph}(k[x]))$.

This shows that to reach the singularity, from outside, would take infinite time.

The versal family is not defined at $y_1 = 0$, see above.

Example 3.9. (i) We shall not treat oscillators in rank $\geq 3$, in general, but only look at the singularities, in all ranks. This is all well known in physics, see [Elbaz], section 16, although in most books in physics, it is treated rather formally, in relation with the second quantification and the introduction of Fock-spaces, and their associated representations of the algebra of observables. We shall see that this second quantification is a natural quotient of the algebra of observables $\text{Ph}_C$, in line with the general philosophy of this paper. Although we may work in a very general setting, we shall, as above, restrict our attention to the classical oscillator $L = 1/2dx^2 - 1/2x^2$, in dimension 1.

As above we find,

$$
d^2x = x
$$

and the Dirac derivation has therefore,

$$
a_+ := 1/2(x + dx), \quad a_- := 1/2(x - dx)
$$

as eigenvectors, with eigenvalues 1 and -1 respectively. Since $\text{Ph}_C = k < x, dx >$ is generated by the elements $a_+ := 1/2(x + dx), \quad a_- := 1/2(x - dx)$, it is clear that Planck’s constant $\hbar = 1$. Notice also that the classical Hamiltonian is given by,

$$
Q := dx^2 - L = 2a_+a_-
$$

Using the method above it is easy to see that for any rank $n = \text{dim}V$, a singular point $v \in \text{Simp}_n(\text{Ph}_C)$ corresponds to a $k < x, dx >$-module $V$, with $x$ and $dx$ acting as endomorphisms $X, dX \in \text{End}_k(V)$ for which there exists an endomorphism, the Hamiltonian, $Q \in \text{End}_k(V)$ with,

$$
dX := \rho(dx) = [Q, \rho(x)] =: [Q, X]
$$

$$
X = \rho(d^2x) = [Q, \rho(dx)] =: [Q, dX]
$$

Let $\psi_0$ be any eigenvector for $Q$ with eigenvalue $\kappa_0$. Since $V$ is simple, the family $\{a_+^m a_-^n (\psi_0)\}$ must generate $V$. Moreover, if $a_+^m a_-^n (\psi_0) \neq 0$, we know it must be an eigenvector for $Q$, with eigenvalue $\kappa_0 + (m - n)$. We can, by adding $\lambda 1$ to $Q$, 

assume that there is a basis for $V$ of eigenvectors for $Q$, with eigenvalues of this form. This means that $Q$ can be assumed to have the form,

$$Q = \begin{pmatrix}
\kappa_0 & 0 & 0 & 0 & \ldots & 0 \\
0 & \kappa_0 + \lambda_1 & 0 & 0 & \ldots & 0 \\
0 & 0 & \kappa_0 + \lambda_2 & 0 & \ldots & 0 \\
\ddots & \ddots & \ddots & \ddots & \ddots & 0 \\
0 & 0 & 0 & 0 & \ldots & \kappa_0 + \lambda_{n-1}
\end{pmatrix},$$

where $0 \leq \lambda_1 \leq \lambda_2 \leq \ldots \leq \lambda_{n-1}$ are all integers. Moreover, since $V$ is simple, and $[Q, a_+] = a_+$, $[Q, a_-] = -a_-$, an easy computation shows that,

$$a_+ = \begin{pmatrix}
0 & 0 & 0 & 0 & \ldots & 0 \\
0 & a_{2,1} & 0 & 0 & \ldots & 0 \\
\ddots & \ddots & \ddots & \ddots & \ddots & 0 \\
0 & 0 & 0 & 0 & \ldots & a_{n,n-1}
\end{pmatrix},$$

$$a_- = \begin{pmatrix}
0 & a_{1,2} & 0 & 0 & \ldots & 0 \\
0 & 0 & a_{2,3} & 0 & \ldots & 0 \\
\ddots & \ddots & \ddots & \ddots & \ddots & 0 \\
0 & 0 & 0 & 0 & \ldots & 0
\end{pmatrix},$$

where all $a_{i,i-1}, a_{i,i+1} \neq 0$. We also find,

$$[a_+, a_-] = \begin{pmatrix}
-a_{1,2}a_{2,1} & 0 & 0 & 0 & \ldots & 0 \\
0 & a_{2,1}a_{1,2} - a_{2,3}a_{3,2} & 0 & 0 & \ldots & 0 \\
0 & 0 & a_{3,2}a_{2,3} - a_{3,4}a_{4,3} & 0 & \ldots & 0 \\
\ddots & \ddots & \ddots & \ddots & \ddots & 0 \\
0 & 0 & 0 & 0 & \ldots & a_{n,n-1}a_{n-1,n}
\end{pmatrix},$$

obviously with vanishing trace.

Now to have the classical formulas, see ([Elbaz], p.377-380), we just have to impose the condition that $a_+$ and $a_-$ be conjugate operators, i.e. that

$$[a_+, a_-] = \begin{pmatrix}
-1 & 0 & 0 & \ldots & 0 \\
0 & -1 & 0 & \ldots & 0 \\
0 & 0 & -1 & \ldots & 0 \\
\ddots & \ddots & \ddots & \ddots & 0 \\
0 & 0 & 0 & \ldots & (n-1)
\end{pmatrix}.$$

Then, introducing a base change, corresponding to an inner automorphism defined by a diagonal matrix, we find that we may assume $a_{i,i+1} = a_{i+1,i}$. It follows that,

$$X = \begin{pmatrix}
0 & \sqrt{1} & 0 & 0 & \ldots & 0 \\
\sqrt{1} & 0 & \sqrt{2} & 0 & \ldots & 0 \\
0 & \sqrt{2} & 0 & \sqrt{3} & \ldots & 0 \\
\ddots & \ddots & \ddots & \ddots & \ddots & \sqrt{(n-1)} \\
0 & 0 & 0 & \ldots & \sqrt{(n-1)} & 0
\end{pmatrix}.$$
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\[ dX = \begin{pmatrix}
0 & -\sqrt{1} & 0 & 0 & \cdots & 0 \\
\sqrt{1} & 0 & -\sqrt{2} & 0 & \cdots & 0 \\
0 & \sqrt{2} & 0 & -\sqrt{3} & \cdots & 0 \\
0 & 0 & \cdots & \cdots & \cdots & 0 \\
\cdots & \cdots & \cdots & \cdots & \cdots & \cdots \\
0 & 0 & 0 & \cdots & \sqrt{(n-1)} & 0
\end{pmatrix} \]

with associated Hamiltonian,

\[ Q = \begin{pmatrix}
1/2 & 0 & 0 & 0 & \cdots & 0 \\
0 & 3/2 & 0 & 0 & \cdots & 0 \\
0 & 0 & 5/2 & 0 & \cdots & 0 \\
0 & 0 & 0 & \cdots & \cdots & 0 \\
\cdots & \cdots & \cdots & \cdots & \cdots & \cdots \\
0 & 0 & 0 & \cdots & 0 & (2n-1)/2
\end{pmatrix} \]

Clearly, we cannot impose, \([a-, a+] = 1\), in finite rank. If, however, we let \(n = \dim V\) tend to \(\infty\), then we find exactly the classical formulas for the oscillator in the "second quantification", see the reference above. In particular it follows that \([a-, a+] = 1\) is the only relation between the operators \(a_-\) and \(a_+\) in this classical limit representation.

On the basis of the examples above, in particular (3.8), it is tempting to conjecture that all integral curves of \(\xi\) are intersections of hypersurfaces of \(\text{Spec}(C(n))\), of the form \(\text{Trace} \xi(\tilde{\rho}(\theta)) = \text{const.}\). However, this is not true, as we can see by going back to (3.7). Here we have

\[ A = k[x],\ A(\sigma) = PhA = k < x, dx > = k < x, y >, \ y = dx, \ \delta = y \frac{\partial}{\partial x} + x \frac{\partial}{\partial y}. \]

There are only two obvious invariants, \(\theta_1 = x^2 - y^2\), i.e. the Hamiltonian, and \(\theta_2 = xy - yx\). Moreover the universal family on \(C(2) = k[t_1, ..., t_5]\), is given by,

\[ \tilde{\rho}(x) = \begin{pmatrix} 0 & 1 + t_3 \\ t_5 & t_4 \end{pmatrix}, \tilde{\rho}(y) = \begin{pmatrix} t_1 \\ 1 + t_3 \end{pmatrix}. \]

We find, see (3.5), that the invariants expressed in the coordinates \((u_1, ..., u_5)\), looks like,

\[ \text{trace}(\tilde{\rho}(\theta_1)) = -u_1 - 2u_2 + u_4 + 2u_5 \]
\[ \text{det}(\tilde{\rho}(\theta_1)) = (u_5 - u_2 - u_1^2)(u_5 - u_2 + u_4^2) - u_1 u_3 u_4 - u_2^2 u_2. \]
\[ \text{det}(\tilde{\rho}(\theta_2)) = -u_3^2 + 4u_3 u_5 + u_1 u_3 u_4 + u_1^2 u_5 + u_2 u_4^2 \]
\[ \text{det}(\tilde{\rho}(\theta_1)\tilde{\rho}(\theta_2)) = 0. \]

Recall from above that,

\[ u_1 = t_1, \ u_2 = (1 + t_3)t_2, \ u_3 = (1 + t_3)^2 + t_2 t_5, \ u_4 = t_4, \ u_5 = (1 + t_3)t_5, \]

and,

\[ \xi_2 = -1/2 (y_1 + y_2) + 1/2 y_3 \]
\[ \xi_4 = 1/2 y_2 \]
\[ \xi_5 = 1/2 (y_2 - y_1)y_1 y_3^{-1}. \]
If we put \( t_1 = t_4 = 0 \), we find the result of (3.7), namely
\[
\text{Trace}(\tilde{\rho}(\theta_1)) = y_1y_3 = 2(u_3 - u_2), \quad \text{det}(\tilde{\rho}(\theta_1)) = 1/4(y_1^2u_2^2) = (u_5 - u_2)^2, \quad \text{det}(\tilde{\rho}(\theta_2)) = -1/16(y_1^2 - y_2^2 + y_3^2) = -u_3^2 + 4u_2u_5.
\]
However, the fact that \( \text{det}(\tilde{\rho}(\theta_1)\tilde{\rho}(\theta_2)) = 0 \) indicates that there are non-algebraic integral curves sitting on an algebraic surface of \( A^5 \). This is related to the problem of hyperbolicity of complex algebraic surfaces. In fact, we see that any integral curve of \( \xi = [\delta] \) is sitting on an algebraic surface, and we may find one for which \( \xi \) have no singularities. Is the integral curve algebraic, or may it be dense on the surface, in the Zariski topology? Exact conditions on algebraic surfaces for being hyperbolic seems not to be known. Notice moreover that the non-commutative invariant \( \theta_2 \) is essential in the integration of \( \xi \) in this case. Notice also that when \( A = k[x_1, x_2, x_3] \), and if the Lagrangian \( L = 1/2(dx_1^2) + 1/2(dx_2^2) + 1/2(dx_3^2) + U \), has a potential \( U \), such that \( \partial U/\partial x_i = \partial U/\partial x_j \), i.e. concerns a central force, then the angular momenta \( L_{i,j} := x_idx_j - x_jdx_i, \) are constants, i.e. \( \delta(L_{i,j}) = 0 \), in rank 1, which of course have the classical consequences one knows. Combining this with the representations discussed in the Example(1.1), (iii), we find interesting results, see next subsection.

Emmy Noethers theorem is, in this context, reduced to the following observation. Suppose a non-trivial inner derivation, \( ad(T) \) of \( A(\sigma) \simeq Ph(A) \) leaves the dynamical structure of the versal family \( \tilde{\rho} \)-invariant, i.e. suppose,
\[
\tilde{\rho}(\text{ad}(T), [\delta]) = 0,
\]
then for all \( a \in A(\sigma) \) we must have, \( \tilde{\rho}(\delta(T), a) = 0 \), so \( \tilde{\rho}(\delta(T)) \), is a central element of \( M_n(C(n)) \), therefore simply an element \( c(T) \in C(n) \). Let \( C(T) \in C(n) \) be such that \( \xi(C(T)) = c(T) \). Since \( \tilde{\rho} \) is surjective, there exists an element \( C(T) \in A(\sigma) \) such that \( \tilde{\rho}(C(T)) = C(T) \). Applying (3.2) we find \( [Q, \tilde{\rho}(T)] = c(T) - \xi(\tilde{\rho}(T)) \), and \( \tilde{\rho}(\delta(T - C(T))) = 0 \), therefore,
\[
\xi(\text{Trace}(\tilde{\rho}(T - C(T)))) = 0,
\]
so \( \text{Trace}(\tilde{\rho}(T)) - nC(T) \) is a constant for the theory.

**Remark 3.10.** We might try to find functions, or formal power series, \([n] \in k[[\tau]]\) such that the representation,
\[
x(n) = \begin{pmatrix}
0 & \sqrt{1} & 0 & 0 & \ldots & 0 \\
\sqrt{1} & 0 & \sqrt{2} & 0 & \ldots & 0 \\
0 & \sqrt{2} & 0 & \sqrt{3} & \ldots & 0 \\
\vdots & \vdots & \vdots & \vdots & \ddots & \vdots \\
0 & 0 & 0 & \ldots & \sqrt{(n-1)} & 0
\end{pmatrix},
\]
\[
dx(n) = \begin{pmatrix}
0 & -\sqrt{1} & 0 & 0 & \ldots & 0 \\
\sqrt{1} & 0 & -\sqrt{2} & 0 & \ldots & 0 \\
0 & \sqrt{2} & 0 & -\sqrt{3} & \ldots & 0 \\
\vdots & \vdots & \vdots & \vdots & \ddots & \vdots \\
0 & 0 & 0 & \ldots & -\sqrt{(n-1)} & 0
\end{pmatrix}
\]
with associated Hamiltonian,

\[
Q = \begin{pmatrix}
\frac{1}{2} + [0] & 0 & 0 & 0 & \ldots & 0 \\
0 & \frac{1}{2} + [1] & 0 & 0 & \ldots & 0 \\
0 & 0 & \frac{1}{2} + [2] & 0 & \ldots & 0 \\
\vdots & \vdots & \vdots & \ddots & \ddots & \vdots \\
0 & 0 & 0 & \ldots & 0 & \frac{1}{2} + [n-1]
\end{pmatrix}
\]

satisfyi the fundamental dynamical equation,

\[
\delta = [\delta] + [Q, -].
\]

We may, of course choose \([\delta] := \frac{\partial}{\partial \tau}\) as the generator of the vector fields on the \(\tau\)-line. We find the following system of differential equations,

\[
\begin{align*}
\frac{\partial}{\partial \tau} f_n + (f_n^2 - f_{n-1}^2) f_n &= f_n \\
\frac{\partial}{\partial \tau} f_n + (-f_n^2 - f_{n-1}^2) f_n &= -f_n
\end{align*}
\]

where \(f_n := \sqrt{[n]}\), and with boundary conditions,

\(f_n(0)^2 = n\).

These equations immediately lead to \(\delta(f_n) = 0\), so to constant \(f_n\)'s, and therefore proves that the curve in \(\text{Simp}_n(PhC)\) defined by the family \(\{x(n), dx(n)\}\) is transversal to the fundamental vector field \(\xi\). The introduction of the \((p,q)\) commutators, and their treatment in physics, makes it possible to treat the fermions and the bosons in a common structure. Letting the parameter \(q\) in the above family slide from 1 to -1, the \(q\)-commutator \([-,-]_q\) changes from the ordinary Lie product to the Jordan product. The computation above shows that this change takes place transversal to time, i.e. instantaneously!

**Remark 3.11.** For the harmonic oscillator in dimension \(n = 2\) we have \(A = k[x_1, x_2]\), and, \(Ph(A) = k < x_1, x_2, dx_1, dx_2 > /[x_1, x_2], [x_1, dx_2] - [x_2, dx_1]\), and,

\[
A(\sigma) = k < x_1, x_2, dx_1, dx_2 > /[x_1, x_2], [x_1, dx_2] - [x_2, dx_1], [dx_1, dx_2].
\]

Moreover, in rank 2 we find a simple representation of \(A(\sigma)\), given by,

\[
X_1 = \begin{pmatrix} 1 & 0 \\ 0 & 0 \end{pmatrix}, X_2 = \begin{pmatrix} 0 & 0 \\ 0 & 1 \end{pmatrix}
\]

\[
dX_1 = \begin{pmatrix} 0 & -1 \\ 1 & 0 \end{pmatrix}, dX_2 = \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix}
\]

with,

\[
[X_1, dX_1] = [X_2, dX_2] = \begin{pmatrix} 0 & -1 \\ -1 & 0 \end{pmatrix}.
\]
Example 3.12. For the quartic anharmonic oscillator, given by \( L = 1/2 \, dx^2 - 1/4 \, \alpha x^4 \) we may easily compute the rank 2 and 3 versal families. In rank 2 we find that there is a one-dimensional singular family of dimension 2 simple modules, with,

\[
X = \begin{pmatrix} \alpha t^5 \\ 0 \end{pmatrix}, \quad dX = \begin{pmatrix} 0 \\ -\alpha^2 t^5 \end{pmatrix}, \quad Q = X = \begin{pmatrix} 0 & 0 \\ 0 & \alpha \end{pmatrix}.
\]

In rank 3 we find that there are no simple singular module with corresponding diagonal Hamiltonian. This may be one reason why the energy levels of the quartic anharmonic oscillator is not known to the physicists.

Example 3.13. Now, let us consider the infinite rank case. In particular we may consider the representation given in the above example, when \( n = \dim V \) “becomes” \( \infty \). Notice that this is given as the limit case of the singular simple representation of the classical oscillator in dimension \( n \), with an obvious conjugation condition imposed. For \( k = \mathbb{R} \), we have a real Planck’s constant which we obviously may assume equal to \( \hbar = 1 \).

Moreover, we now have, \([a_+, a_-] = 1\), and we have a representation of \( PhC \) onto the algebra \( F \), generated by \( \{a_+, a_-\} \). Notice that in each finite rank, this algebra generate the whole \( \text{End}_k(V) \). The commutation relations is given by a classical formula,

\[
a^m a^n_+ = a^m a^n_+ - mn a^{m-1}_+ a^{n-1}_+ + 1/2! m(m-1)n(n-1) a^{m-2}_+ a^{n-2}_- + 1/3! m(m-1)(m-2)n(n-1)(n-2) a^{m-3}_+ a^{n-3}_- + ... 
\]

and the Lie algebra \( f \) of derivations of \( F \) are easily seen to be generated by the derivations \( \{\delta_{p,q}\}_{p,q} \), defined as,

\[
\delta_{p,q}(a_+) = a^p a^q_+ \quad \delta_{p,q}(a_-) = -p/(q+1) a^{p-1}_+ a^{q+1}_-.
\]

If we put, for \( m, n \geq 0 \)

\[
\chi_{m,n} := \delta_{m+1,n}, \quad \chi_{m} := \chi_{m,0}
\]

then we find the Witt-algebra, with the classical relations,

\[
[\chi_{m}, \chi_{n}] = (n - m) \chi_{m+n}.
\]

Moreover we find,

\[
[\chi_{0}, \chi_{m,n}] = (m - n) \chi_{m,n} =: \text{deg}(\chi_{m,n}) \chi_{m,n}.
\]

Clearly the Lie algebra \( \text{Der}_k(\mathcal{F}) \) has an ascending filtration with respect to the degree, \( \text{deg} \), defined above, and it is easy to see that the corresponding graded Lie algebra \( \mathfrak{g} := \text{gr}(\text{Der}_k(\mathcal{F})) \) has the following products,

\[
[\chi_{p,q}, \chi_{r,s}] = (r - p + (s + 1)^{-1}(r + 1)q - (q + 1)^{-1}(p + 1)s) \chi_{p+r,q+s}.
\]

In particular the degree zero component of \( \mathfrak{g} \) is Abelian.
Example 3.14. Finally let $C := \mathbb{R}[x]$, and let $\mathcal{C} := C \otimes \mathbb{R} C$, and consider some representation on $V = \mathcal{C}$ of $\text{PhC} = \mathbb{R} < x, dx >$. Clearly,

$$\text{Ext}^1_C(V, V) = 0,$$

but, in general,

$$\text{Ext}^1_{\text{PhC}}(V, V)$$

is infinite dimensional.

(i) Consider the free particle, i.e. the dynamical system, $\sigma$ given by,

$$L = \frac{1}{2} dx^2, \quad \sigma : \delta^2 x = 0,$$

and let $V$ be defined by letting $dx$ act as the identity. Then we find that,

$$[\delta] = 0, Q = \frac{\partial}{\partial x}.$$

This means that $[\delta]$ does not move $V$ in the moduli space of $V$. The Hamiltonian $Q$ defines time, and

$$\exp(tQ)(f(x)) = f(x + t).$$

(ii) Consider the same dynamical system, and let $V$ be defined by letting $dx$ act as $\frac{\partial}{\partial x}$. Then we find that,

$$[\delta] = 0, Q = \left(\frac{\partial}{\partial x}\right)^2.$$

As above, $[\delta]$ does not move $V$ in the moduli space. The Hamiltonian $Q$ defines time, and the time evolution looks like,

$$U(t, \psi) = \exp(tQ)(\psi).$$

Introducing the Fourier transformed $\hat{\psi}$, we obtain a time evolution given by,

$$U(t, \hat{\psi}) = \exp(tp^2)(\hat{\psi}).$$

(iii) Consider again the harmonic oscillator, and let the representation $V := k[x^{-1}]$ be defined by letting $x$ act as multiplication by $x^{-1}$, and $dx$ act as $\frac{\partial}{\partial x}$. Then we find that,

$$[\delta] = 0, Q = (x \frac{\partial}{\partial x}).$$

As above, $[\delta]$ does not move $V$ in the moduli space. The eigenvectors of the Hamiltonian $Q$ are the monomials $x^{-n}$, $n \geq 0$, with eigenvalues $-n$, and the time evolution looks like,

$$U(t, x^{-n}) = \exp(-nt)x^{-n}.$$

Notice that,

$$[x, dx] = x^2,$$
as operators on $V$. Notice also that $V$ in this case is not simple. It is, however, a limit of the finite representations, $V_n := k[x^{-1}]/(x^{-1})^n$. The representation $V_2$ is given by the actions,

$$x = \begin{pmatrix} 0 & 1 & 0 \\ 1 & 0 & 0 \end{pmatrix},$$
$$dx = \begin{pmatrix} 0 & 0 \\ 0 & 0 \end{pmatrix}$$

where we have chosen the basis $\{1, x^{-1}\}$ in $V_2$. It is clearly not simple, but it sits as a point at infinity, $t_1 = t_2 = 1 + t_3 = t_4 = 0, t_5 = 1$, for the (almost) versal family,

$$x = \begin{pmatrix} 0 & 1 + t_3 \\ t_5 & t_4 \end{pmatrix},$$
$$dx = \begin{pmatrix} t_1 & t_2 \\ 1 + t_3 & 0 \end{pmatrix}.$$

**Example 3.15.** Given a dynamic system, $A(\sigma)$ and a versal family for simple representations of dimension $n$. Let $\xi$ be the fundamental vectorfield defined on $U(n)$. Recall Theorem (2.13) and Theorem (3.2). There is a commutative diagram of generalized schemes,

$$U(n) \rightarrow \text{Simp}_n(A) \rightarrow \text{Simp}_n(A)/<\xi> =: \mathcal{X}_n.$$

The quotient “spaces” $\mathcal{Y}_n$ and $\mathcal{X}_n$ are orbit spaces, where each orbit is a curve. Completing, when necessary, $U(n)$ and/or $\text{Simp}_n(A)$, we may assume these curves complete. Restricting to an integrable part of $U(n)$, resp. of $\text{Simp}_n(A)$, we may then hope to find natural morphisms,

$$\Gamma_n : \mathcal{Y}_n \rightarrow \mathcal{M},$$

where $\mathcal{M}$ is the moduli space of the complete algebraic curves. Moreover, Theorem (3.3) should produce a rank $n$ bundle $\mathcal{V}_n$ on $\mathcal{Y}_n$, $n \geq 1$, and one might ask for conditions for the existence of universal bundles $U_n$ on $\mathcal{M}$, such that $\mathcal{V}_n = \Gamma_n^* U_n$.

These are questions related to vertex algebras (bundles), see e.g. [F]. There is a large literature on the subject. Seen from our point of view, the hidden agenda of the vertex algebra framework, seems to be to construct the relevant algebra $A(\sigma)$ of observables for a given quantum (field) theoretic situation.

In our language, let $t_{i_0} \in \text{Simp}_n(A(\sigma))$ be a singularity for $\xi$. Consider the Planc’s constant, $\hbar(t_{i_0})$, and the corresponding operators, $a_i^+, a_i^- \in A(\sigma)$, together with the vacuum state $\omega(t_{i_0}) \in \tilde{V}(t_{i_0}) =: V$ (any flat section $\omega$ of $\tilde{V}$ along $\gamma$ will produce a vacuum state), such that the action of $A(\sigma)$ induces an isomorphism,

$$k[a_i^+, \ldots, a_i^+] \simeq V,$$
a situation that we have seen realized in the case of the harmonic oscillator in dimension 1, but which is easily seen to generalize to any dimension, then there pops up a family of generalized vertex algebras. In fact, consider the restriction of the versal family

$$\tilde{\rho} : A(\sigma) \to \text{End}_{C(\sigma)}(\tilde{V}),$$

to the integral curve $\gamma$ through the point $\xi_0 \in \text{Simp}_n(A(\sigma))$. It is singular at $\xi_0$, so parametrized with time, $\tau$, the completion will produce a map,

$$Y : V = \tilde{V}(\xi_0) \simeq k[a_1^+, \ldots, a_n^+] \subset A(\sigma) \to \text{End}_k(V) \otimes_k k[[\tau]][\tau^{-1}],$$

see (3.6), which will be a kind of generalized vertex algebra. In particular, the localization axiom of vertex algebras imply that $\tilde{\rho}(a_i^+)$ and $\tilde{\rho}(a_i^-)$ commute, which here is obvious. Moreover we observe that the exponentiating formula of Y.-Z. Huang, see [F], p.18, (16),

$$Y(a, t) = R(\rho)Y(R(\rho^{-1}a, \rho(t))R(\rho)^{-1},$$

for $a \in A$, and for any $\rho \in \text{Aut}(\hat{O}_{\gamma, 0}) \simeq \text{Aut}(C[[t]])$, follows from Theorem (3.3) above. We shall, hopefully, return to this in a later paper.

Clocks and Classical Dynamics.

Going back to General Quantum Fields, Lagrangians and Actions, we shall study the 1-dimensional case from a different perspective.

When we talk about a clock, we obviously do not talk about the clock. We just think of a device that can measure the changes that we choose to study, in a most objective way.

The western way of thinking about time has always been cyclical, life, death, reincarnation, new death, etc. This way of viewing the world would be more comfortable with a Big Crunch turning into a new Big Bang, and so on.

The measuring device for this kind of time-concept might therefore be a one-dimensional harmonic oscillator i.e. $k[\tau]$ with dynamical structure $d^2\tau = \tau$.

Our representations of a western clock in rank 1 is easy. We have a $k$-algebra $k < \tau, d\tau >$ with a one-dimensional automorphism, given by, $exp(t\delta)(f)(\tau, d\tau) = f(t_0 + td\tau, d\tau)$. In rank 2 we may look at the situation in (3.7), and we find that the western clock has no singularities in rank 2. It never stops, in contrast to what we found for the eastern clock, see (3.7).

Let now $A := k[\tau], \ B := k[t_1, \ldots, t_m]$, and let the dynamical system $\sigma$ defined on $A$ be the Eastern Clock, and the dynamical system $\mu$ on $B$ be the free particle, so that $A(\sigma) = Ph(A), \ d^2\tau = \tau, \ B(\mu) = Ph(B), d^2t_i = 0, \ i = 1, \ldots, n$. A field, $\phi : A \to B$, should be considered as a model for the physical phenomenon that can be observed everywhere on $\text{Simp}_1(B)$, having a value that oscillate, (real or complex situation). The equation of motion for this field is,

$$d\phi(a) - \phi(da) = [Q, \phi(a)], \ a \in A(\sigma), \ Q \in B(\mu).$$
Plugging in $a = d\tau$, we obtain,

$$d^2\phi(\tau) - \phi(\tau) = [Q, \phi(d\tau)],$$

in $B(\mu)$. In the representation $Ham: B(\mu) \to B_0(\mu) = k[t_i, dt_j]$, this gives us,

$$\sum_{i,j} \frac{\partial^2 \phi}{\partial t_i \partial t_j} k_i k_j = \phi,$$

where $k_i = Ham(dt_i)$, necessarily are constants, since $d(dt_i) = 0$. This choice of vector $k = (k_1, ..., k_m)$, is now arbitrary, and corresponds to the vector that pop up in physicists text, in Fourier expansions. We shall show that this immediately leads to the quantized Klein-Gordon equation, but first we have to make clear where we are.

**Time-space and Space-times.**

Go back to our basic model, $Hilb^2(\mathbb{E}^3) := H = \mathbb{H}/\mathbb{Z}_2$, classifying the family of pairs of points $(o, x)$ of the Euclidean 3-space, $\mathbb{E}^3$. As above we shall first consider the structure of $\mathbb{H}$ before we extend the results to $Hilb^2$. Recall from the Introduction, and above, that the tangent bundle $T(H)$, outside of $\underline{A}$ is decomposed into the sum of the basic tangent bundles $B_o, B_x, A_o, x$, each of rank 2. Recall also that $A_{o,x}$ is decomposed into a unique 0-velocity, dual to $<dt_o>$ and a light velocity dual to $<dt_3>$. The sub-bundle $S_{o,x}$, given by the triples $(\psi, -\psi, \phi) \in B_o \oplus B_x \oplus A_{o,x}$, in which the pair $(\psi, -\psi)$ corresponds to a light-velocity, is at each point $H$ a 4-dimensional tangent sub-space, with a unique 0-velocity. Choosing a line $l \subset \mathbb{E}^3$, the subscheme $H(l) \subset H$, has a much simpler structure than $H$. The sub-bundle $S_{o,x}$, restricted to $H(l)$ can be integrated in $H$, and we obtain a 4-dimensional subspace $S(l) \subset H$, in which we choose coordinates $t_0, dt_1, dt_2, t_3$, where $dt_0$ and $dt_3$ are as above, and $dt_1, dt_2$ are dual coordinates for the transverse bundle, $B_0$, (isomorphic to the inversely oriented bundle $B_x$), normal to $l$ in $\mathbb{E}^3$.

This subspace $S(l)$ of $H$ may be identified with a natural moduli-subspace $M(l) \subset H$. In fact, let $M(l) = \{(o, x) \in H | 1/2(o + x) \in l\}$. Since $o = 1/2(o + x) + 1/2(o - x), x = 1/2(o + x) - 1/2(o - x)$, it is of dimension 4, and contains $H(l)$. At every point $(o, x) \in M(l)$ the tangent space is easily identified with $S_{o,x}$, therefore identifying $S(l)$ and $M(l)$, as spaces. See also the subhead Cosmology, Big Bang and all that, for another characterization of $S(l)$.

Moreover, the usual Minkowsky space, is recovered as the restriction, $U(l)$, of the universal family $U$, to $M(l)$. The metric is deduced from the energy-function of $M(l)$.

Now substitute $S(l)$ for $B$ above, and see that we have two solutions for $\phi$, respectively, $\phi = exp(kt), \phi^1 = exp(-kt)$, in the real case, and with $\delta = ad$ substituting for $d$, $\phi = exp(ikt), \phi^1 = exp(-ikt)$, in the complex case. Now, in $S(l)(\mu)$ the elements $\phi$ and $\phi^1$ commute, but they do not commute with $d\phi$ and $d\phi^1$. However, if we now look at the representation,

$$Wey : S(l)(\mu) \to Diff_k(S(l)),$$

where we recall that,

$$Diff_k(S(l)) = k \mathbb{R}^{2n} > \{[t_i, dt_j] | [t_i, dt_j], [t_i, dt_j]_{i \neq j}, [t_i, dt_i] - 1\},$$
which, restricted to
is outside the diagonal,
in a
at \( t \in \tilde{\eta} \subset \eta \), \( \eta \) spin the corresponding oriented line. The length of the spin vector is called the
space \( H \). This above becomes particularly nice if integrated in our "toy" model, where
laws.
\( c \) above
\( \xi \) defines a unique point
unique point \( t \).
Example (3.16), Relativistic movement in \( \tilde{H} \), Newton’s and Kepler’s
laws. This above becomes particularly nice if integrated in our "toy" model, where
QFT, see e.g. [Mandl]. Notice also that, from (3.7) we know that the energy of this
"particle" is given in terms of the number operator, here \( Q = -1/2(\phi d\phi^\dagger + d\phi^\dagger \phi) = 1/2(\phi^\dagger d\phi + d\phi^\dagger) \).

Example (3.16), Relativistic movement in \( \tilde{H} \), Newton’s and Kepler’s
laws. This above becomes particularly nice if integrated in our "toy" model, where
the space \( H := \tilde{H}/\mathbb{Z}_2 \) is the Hilbert scheme of length 2 subschemes in the Eucledean
space \( E^3 \). But let us first study the geometry of \( H \).
Recall that \( \tilde{H} \to H \), is the (real) blow up of the diagonal \( \Delta \subset H \), where \( H \) is the space of pairs of points in \( E^3 \). Clearly any point \( t \in H \) outside the diagonal, determines a vector \( \xi(o,x) \) and an oriented line \( l(o,x) \subset E^3 \), on which both the observer \( o \) and the observed \( x \) sits. This line also determines a subscheme \( H(l) \subset H \), see (3.7) or [La 6], and in \( H(l) \) there is unique light velocity curve \( l(t) \), through \( t \), an integral curve of the distribution \( \tilde{c} \), and this curve cuts the diagonal \( \Delta \) in a unique point \( c(o,x) \), the center of gravity of the observer and the observed, and thus defines a unique point \( \xi(t) \), of the blow-up of the diagonal, in the fiber of \( \tilde{H} \to H \), above \( c(o,x) \). Any tangent \( \eta := (\eta_1, \eta_2), \eta_2 = -\eta_1 \), of \( H \) at \( t \) in \( \tilde{c} \) normal to \( l(t) \) corresponds to a light velocity, to a spin vector, \( \eta \times \xi(o,x) \), in \( E^3 \), with spin axis, the corresponding oriented line. The length of the spin vector is called the spin of \( \eta \). We have shown in [La 6] that there exists a metric on \( H \) which, restricted to every 3-space \( c(t) \sim \{c(o,x)\} \), has the form,

\[
\begin{align*}
    ds^2 = dt_3^2 + (1 + r^{-2})dt_1^2 + (1 + r^{-2})dt_2^2 - r^{-4}(t_1 dt_1 + t_2 dt_2)^2,
\end{align*}
\]
where we have chosen the coordinates such that $dt_3$ corresponds to the oriented line $l(o, x)$, and $dt_i$, $i = 1, 2$, correspond to the spin-momenta, assuming $t_3 \neq 0$. The nice property of this metric is the following. Consider a spin momentum $\eta := (\eta_1, -\eta_1)$ and its corresponding spin vector $\xi := (\eta_1 \times \xi(o, x), -\eta_1 \times \xi(o, x))$ along the line $l(t)$. Clearly the length of this vector, when $r = t_3$ is large, is just the classical spin. When $r$ tends to zero, $\eta$ defines a tangent vector of the exceptional fiber of the blow up at $c(o, x)$, i.e. of the projective 2-space, and of the covering 2-sphere. And we see that the length of $\eta \times \xi(o, x)$ tends to the length of this tangent vector in the Fubini-Studi metric of $\mathbb{P}^2$.

To see what this may lead us to, we need a convenient parametrization of $\tilde{H}$. Consider, as above, for each $t \in \tilde{H}$ the length $\rho$, in $E^3$, the Euclidean space, of the vector $(o, x)$. Given a point $\lambda \in \lambda_\xi$ and a point $\xi \in E(\lambda) = \pi^{-1}(\lambda)$, the fiber of,

$$\pi : \tilde{H} \to H,$$

at the point $\lambda$, for $o = x$. Since $E(\lambda)$ is isomorphic to $S^2$, parametrized by $\phi$, any element of $\tilde{H}$ is now uniquely determined in terms of the triple $t = (\lambda, \phi, \rho)$, such that $c(t) = c(o, x) = \lambda$, and such that $\xi$ is defined by the line $\alpha_\xi$. Here $\rho > 0$, see also the section Cosmology, Big Bang and all that.

Consider any metric on $\tilde{H}$, of the form,

$$g = h_1(\lambda, \phi, \rho)d\rho^2 + h_2(\lambda, \phi, \rho)d\phi^2 + h_3(\lambda, \phi, \rho)d\lambda^2,$$

where $d\phi^2$ is the natural metric in $S^2 = E(\lambda)$. It is reasonable to believe that the geometry of $(\tilde{H}, g)$, might explain the notions like energy, mass, charge, etc., In fact, we tentatively propose that the source of mass and charge etc. is located in the black holes $E(\lambda)$. This would imply that mass, charge, etc. are properties of the 5-dimensional superstructure of our usual 3-dimensional Euclidean space, essentially given by a density, $h(\lambda, \phi, \theta)$. This might bring to mind Kaluza-Klein-theory. However, it seems to me that there are important differences, making comparison very difficult.

Here we shall just consider the following simple case, where

$$h_1 = \left(\frac{\rho - h}{\rho}\right)^2, \quad h_2 = (\rho - h)^2, \quad h_3 = 1,$$

where $h$ is a positive real number. It clearly reduces to the Euclidean metric far away from $\lambda_\xi$, and it is singular on the horizon of the black hole, given by $\rho = h$.

Moreover, we shall reduce to a plane in the light directions, i.e. we shall just assume that $S^2 = E(\lambda)$, is reduced to a circle, with coordinate $\phi$. Consider the Lagrangian $L = g$, see (Example 3.5), we find the Euler-Lagrange equations,

$$0 = \delta\left(\frac{\partial L}{\partial d\phi}\right) - \frac{\partial L}{\partial \phi} = 2\left(\frac{\rho - h}{\rho}\right)^2d^2\rho + 2\left(\frac{\rho - h}{\rho}\right)\left(\frac{h}{\rho}\right)d\rho^2 - 2(\rho - h)d\phi^2$$

$$0 = \delta\left(\frac{\partial L}{\partial d\phi}\right) - \frac{\partial L}{\partial \phi} = 2(\rho - h)^2d^2\phi + 4(\rho - h)d\rho d\phi$$

$$0 = \delta\left(\frac{\partial L}{\partial d\lambda}\right) - \frac{\partial L}{\partial \lambda} = 2d^2\lambda$$

where $\lambda = \lambda_\xi$, as above.
We solve these equations, and find,

\[ \frac{d^2 \rho}{dt^2} = -\left(\frac{h}{\rho(\rho - h)}\right) \frac{d\rho}{dt} \rho^2 + \left(\frac{\rho^2}{(\rho - h)}\right) \frac{d\phi}{dt} \rho^2 \]

\[ \frac{d^2 \phi}{dt^2} = -2/(\rho - h) \frac{d\rho}{dt} \frac{d\phi}{dt} \rho \]

\[ \frac{d^2 \lambda}{dt^2} = 0. \]

This, of course, is the same solutions as what we would have found by solving the Lagrange equation.

According to (3.5) the corresponding equations for the geodesics in \( \tilde{H} \) are,

\[ \frac{d^2 \rho}{dt^2} = -\left(\frac{h}{\rho(\rho - h)}\right) \left(\frac{d\rho}{dt}\right)^2 + \left(\frac{\rho^2}{(\rho - h)}\right) \left(\frac{d\phi}{dt}\right)^2, \]

\[ \frac{d^2 \phi}{dt^2} = -2/(\rho - h) \frac{d\rho}{dt} \frac{d\phi}{dt} \rho \]

\[ \frac{d^2 \lambda}{dt^2} = 0. \]

where \( t \) is time. But time is, by definition, the distance function in \( \tilde{H} \), so we must have,

\[ \left(\frac{\rho - h}{\rho}\right)^2 \left(\frac{d\rho}{dt}\right)^2 + (\rho - h)^2 \left(\frac{d\phi}{dt}\right)^2 + \left(\frac{d\lambda}{dt}\right)^2 = 1, \]

from which we find,

\[ \left(\frac{d\rho}{dt}\right)^2 = \rho^2((\rho - h)^{-2}(1 - (\frac{d\lambda}{dt})^2) - (\frac{d\phi}{dt})^2). \]

Put this into the first equation above, and obtain,

\[ \frac{d^2 \rho}{dt^2} = -h(1 - (\frac{d\lambda}{dt})^2)(\frac{\rho}{\rho - h}) \frac{1}{(\rho - h)^2} + \left(\frac{\rho + h}{\rho - h}\right) \rho \left(\frac{d\phi}{dt}\right)^2. \]

From the third equation above, we find that \( \frac{d\lambda}{dt} \), \( j = 1, 2, 3 \), are constants, and \( |\frac{d\lambda}{dt}| \) is the rest-mass of the system. Put \( K^2 = (1 - |\frac{d\lambda}{dt}|^2) \), then \( K \) is the kinetic energy of the system.

Assume now \( r := \rho - h \cong \rho \), we find ,

\[ \frac{d^2 r}{dt^2} = \frac{hK^2}{r^2} + r \left(\frac{d\phi}{dt}\right)^2, \]

i.e. Keplers 1. equation (law). The constant \( h \) is thus also related to mass. In the same way, the second equation above gives us Keplers second law,

\[ r \left(\frac{d^2 \phi}{dt^2}\right) + 2r \left(\frac{dr}{dt}\right) \left(\frac{d\phi}{dt}\right) = 0. \]

Notice that with the chosen metric, time, in light velocity direction, is standing still on the horizon \( \rho = h \), of the black hole at \( \lambda \in \mathbb{A} \). Therefore no light can escape from the black hole. In fact, no geodesics can pass through \( \rho = h \). Notice also that,
for a photon with light velocity, $K = 1$, so we may measure $h$, by measuring the trajectories of photons in the neighborhood of the black hole.

Let us now go back, and consider, in this case, the generic dynamical structure ($\sigma$), of the subsection Connections, and the Generic Dynamical Structure, related to the above metric. Put $\rho = t_1, \phi = t_2, \lambda = t_3$, then the Euler-Lagrange equations above give us immediately the following formulas,

$$
\Gamma_{1,1}^1 = h/\rho(\rho - h), \quad \Gamma_{1,2}^1 = -\rho^2/(\rho - h)
$$

$$
\Gamma_{1,2}^2 = 1/(\rho - h), \quad \Gamma_{2,1}^2 = 1/(\rho - h)
$$

$$
\Gamma_{3, i,j}^3 = 0
$$

All other components vanish. From this we find the following formulas,

$$
\nabla_1 = \begin{pmatrix} h/\rho(\rho - h) & 0 & 0 \\ 0 & 1/(\rho - h) & 0 \\ 0 & 0 & 0 \end{pmatrix}
$$

$$
\nabla_2 = \begin{pmatrix} 0 & 1/(\rho - h) & 0 \\ -\rho^2/(\rho - h) & 0 & 0 \\ 0 & 0 & 0 \end{pmatrix}
$$

$$
[\nabla_1, \nabla_2] = \begin{pmatrix} 0 & -1/(\rho - h) & 0 \\ -\rho/(\rho - h) & 0 & 0 \\ 0 & 0 & 0 \end{pmatrix}
$$

$$
\frac{\partial}{\partial \rho} \nabla_2 = (\rho - h)^{-2} \begin{pmatrix} 0 & -1 & 0 \\ -\rho(\rho - 2h) & 0 & 0 \\ 0 & 0 & 0 \end{pmatrix}
$$

$$
D_\rho := \nabla_{\delta_1} = \frac{\partial}{\partial \rho} + \nabla_1, \quad D_\phi := \nabla_{\delta_2} = \frac{\partial}{\partial \phi} + \nabla_2, \quad D_\lambda := \nabla_{\delta_3} = \frac{\partial}{\partial \lambda} + \nabla_3
$$

$$
Q = \sum_{i=1}^3 1/h_i \nabla_{\delta_i}^2
$$

$$
\rho(\delta^2(t_i)) = [Q, \rho(dt_i)] = 1/h_i [Q, \nabla_{\delta_i}].
$$

Here the $h_i$ is the function defined above, i.e. $g_{i,i}$ in our metric.

Recall our general Force Law for the proposed metric on $\tilde{H}$, and the Levi-Civita connection,

$$
\rho E(d^2t_i) + \sum_{p,q} \Gamma_{p,q}^i \nabla_{\xi_p} \nabla_{\xi_q} = 1/2 \sum_p F_{p,i} \nabla_{\delta_p} + 1/2 \sum_p \nabla_{\delta_p} F_{p,i} + 1/2 \sum_{l,q} \delta_q (\Gamma_{l,q}^i - \Gamma_{q,l}^i) \nabla_{\xi_l} + [\nabla_{\xi_i}, \rho E(T)].
$$

A dull, and lengthy computation gives the formulas,

$$
\xi_\rho = 2\rho^2(\rho - h)^{-2} \frac{\partial}{\partial \rho}, \quad \xi_\phi = 2(\rho - h)^{-2} \frac{\partial}{\partial \phi}, \quad \xi_\lambda = 2 \frac{\partial}{\partial \lambda},
$$

$$
[\xi_\rho, \xi_\phi] = -8\rho^2(\rho - h)^{-5} \frac{\partial}{\partial \phi}
$$

$$
F(\xi_\rho, \xi_\phi) = 4\rho^2(\rho - h)^{-4} \left( [\nabla_{\rho}, \nabla_{\phi}] + \frac{\partial}{\partial \rho} \nabla_{\phi} \right).
$$
Putting these formulas together gives the expressions,

\[
\rho_\Theta(d^2 \rho) + h \rho^{-1}(\rho - h)^{-1}\nabla^2 \xi_\rho - \rho^2(\rho - h)^{-1}\nabla^2 \xi_\phi
\]

\[
= F(\xi_\rho, \xi_\phi)(\frac{\partial}{\partial \phi} + \nabla \phi) + 2 \rho^2(\rho - h)^{-5}(\frac{\partial}{\partial \phi} + \nabla \phi) + [\nabla \xi_\rho, \rho_\Theta(T)],
\]

and,

\[
\rho_\Theta(d^2 \phi) + (\rho - h)^{-1}\nabla \xi_\phi \nabla \xi_\rho - \nabla \xi_\rho \nabla \xi_\phi
\]

\[
= \frac{1}{2} F(\xi_\rho, \xi_\phi)((\frac{\partial}{\partial \phi} + \nabla \phi) + \frac{1}{2}(\frac{\partial}{\partial \phi} + \nabla \phi) F(\xi_\rho, \xi_\phi)
\]

\[
- (4 \rho(\rho - h)^{-5} + 6 \rho^2(\rho - h)^{-6}(\frac{\partial}{\partial \phi} + \nabla \phi) + [\nabla \xi_\phi, \rho_\Theta(T)].
\]

**Example 3.17: Classical Electro-Magnetism.** There are at least two possible models of an electromagnetic field.

First, given a potential,

\[
\phi = \sum_{j=0}^{3} \phi_j d_1 t_j \in PhS(l),
\]

considered as a field \( \phi : k[\tau] \to Ph(S(l)) \), and, say, the trivial metric \( g \) on \( S(l) \). The interpretation is a plane wave, corresponding to a linear form on the tangent space of each point of \( S(l) \). Notice that we have a canonical derivation \( d_1 : S(l) \to Ph(S(l)) \), and that we have the following relations,

\[
[d_1 t_i, t_j] + [t_i, d_1 t_j] = 0.
\]

Let \( A := k[\tau] \) be the Western clock, with Dirac derivation \( d \) with \( d^2 \tau = 0 \). It is easy to check that the following relations actually define a dynamic system on \( B \), with Dirac derivation, \( d \):

\[
[t_i, dt_j] = 0, i, j \geq 0,
\]

\[
dt_i d_1 t_j = -dt_j d_1 t_i, i \neq j,
\]

\[
dt_i d_1 t_i = dt_j d_1 t_j, i, j \geq 0,
\]

\[
d^2 t_i = dd_1 t_i = 0, i = 0, 1, 2, 3.
\]

If we put, for \( i, j, k = 1, 2, 3 \), with \( sgn(i, j, i \times j) = 1 \),

\[
E_i := \frac{\partial \phi_i}{\partial t_0} - \frac{\partial \phi_0}{\partial t_i}, B_k := \frac{\partial \phi_j}{\partial t_i} - \frac{\partial \phi_i}{\partial t_j}, k = i \times j,
\]

we find, modulo these relations,

\[
d\phi = E_1 dt_0 d_1 t_1 + E_2 dt_0 d_1 t_2 + E_3 dt_0 d_1 t_3 + B_3 dt_1 d_1 t_2 + B_2 dt_3 d_1 t_1 + B_1 dt_2 d_1 t_3
\]

\[
+ \nabla \phi dt_0 d_1 t_0 + \sum_{i=0}^{3} \phi_i dd_1 t_i.
\]
Computing we find,
\[
\delta^2(\phi) = \left( \frac{\partial}{\partial t_0} + \frac{\partial}{\partial t_1} + \frac{\partial}{\partial t_2} + \frac{\partial}{\partial t_3} \right) \cdot \frac{\partial^2 \phi}{\partial t_0 \partial t_1 \partial t_2 \partial t_3} dt_0 dt_1 dt_2 dt_3 + D.
\]

Here,
\[
D = \left( \frac{\partial \phi_0}{\partial t_1} + \frac{\partial \phi_1}{\partial t_0} \right) dt_1 dt_2 + \left( \frac{\partial \phi_2}{\partial t_1} - \frac{\partial \phi_0}{\partial t_2} \right) dt_2 dt_3 + \left( \frac{\partial \phi_3}{\partial t_1} - \frac{\partial \phi_2}{\partial t_2} \right) dt_3 dt_0 + \sum_{i=0}^{3} \frac{\partial \phi_i}{\partial t_i} dt_1 dt_2 dt_3.
\]

vanish. The equation of motion is, \( d^2 \phi = 0 \), which implies the following equations,
\[
\begin{align*}
\frac{\partial}{\partial t_0} \left( \nabla \cdot \phi \right) + \nabla_s E &= 0 \\
\nabla_s (\nabla \cdot \phi) + \frac{\partial E}{\partial t_0} + \nabla \times B &= 0 \\
\nabla \times E + \frac{\partial B}{\partial t_0} &= 0 \\
\n\nabla_s B &= 0,
\end{align*}
\]

where \( \nabla_s \) is the space-gradient.

These are Maxwell’s equations, with electrical charge equal to \( \rho := \frac{\partial (\nabla \cdot \phi)}{\partial t_0} \), and electrical current equal to \( j := \nabla_s \cdot (\nabla \phi) \). The last two equations are Bianchi’s equations and are trivial, given that we start with a potential.
In our language, we see that charge becomes a rest-mass, and electric current a kinetic energy.

Notice also that in Space-Time coordinates our potential satisfies, \( \nabla^2 \phi_i = 0 \), which explains the extra conditions necessary in the classical case. A classical "free" particle, clocked by a Western clock, is now, according to (3.17) a field

\[
\kappa : S(l) \rightarrow k[\tau]
\]
such that,

\[
\delta S = 0, \quad S := \int L, d\tau
\]

where, putting \( \dot{\kappa}_j := \frac{\partial \kappa_j}{\partial \tau}, \)

\[
L := Ph(\kappa)(\phi) = \sum_{j=0}^{3} \phi_j(\kappa_0, \kappa_1, \kappa_2, \kappa_3) \dot{\kappa}_j.
\]

Classically, where the representations one considers are \( L^2 \)-spaces of functions, this is interpreted as,

\[
\delta \int Ph(\kappa)(\phi) d\tau = 0.
\]

The Euler-Lagrange equations applies and one gets the system of equations,

\[
\frac{\partial \phi_j}{\partial \tau} - \sum_i \frac{\partial \phi_i}{\partial t_j} \dot{\kappa}_i = 0, \forall j,
\]

which reduces to,

\[
\sum_{i=0}^{3} \frac{\partial \phi_i}{\partial t_i} \dot{\kappa}_i - \sum_{i=0}^{3} \frac{\partial \phi_i}{\partial t_j} \dot{\kappa}_i = 0, \forall j.
\]

Put, as above,

\[
E_i := \frac{\partial \phi_i}{\partial t_0} - \frac{\partial \phi_0}{\partial t_i}, \quad B_k := \frac{\partial \phi_i}{\partial t_k} - \frac{\partial \phi_k}{\partial t_i}, k = i \times j,
\]

and define the 3-vectors,

\[
\psi := (\kappa_1, \kappa_2, \kappa_3), \quad \dot{\psi} := (\dot{\kappa}_1, \dot{\kappa}_2, \dot{\kappa}_3) \quad E := (E_1, E_2, E_3), \quad B := (B_1, B_2, B_3).
\]

Then the equations above simply says the following,

\[
B \times \dot{\psi} + E.\dot{\kappa}_0 = 0, \quad E.\dot{\psi} = 0.
\]

Now let us denote by \( \gamma \) the curve defined by \( \kappa = (\kappa_0, \kappa_1, \kappa_2, \kappa_3) \). Then in our space \( S(l) \subset H \) the squared energy, \( m^2 \), of the particle, measured by the clock we are using, is \( \sum_{i=0}^{3} \dot{\kappa}_i^2 \). Recalling that,

\[
\dot{\gamma} = m.\mathbf{v},
\]
where the unit vector $v$ is the velocity, and $\dot{\psi} = m v_{\text{space}}$, $\dot{k}_0 = m v_0 := \text{restmass}$. We find the following equations,

$$B \times v_{\text{space}} + E v_0 = 0.$$ 

If we, together with the electromagnetic potential, also take into consideration gravitation, i.e. time, say just the trivial metric, $g := \sum_{i=0}^{3} dt_i^2$, i.e. if we consider the Lagrangian,

$$L := \phi + g,$$

then, the Euler-Lagrange equations look like,

$$\frac{\partial k_j}{\partial \tau} + \frac{\partial \phi_j}{\partial \tau} - \sum_i \frac{\partial \phi_i}{\partial t_j} k_i = 0, \forall j.$$ 

Referring to Example (3.16), we find for the real time, $t$, with respect to the clock, $\tau$,

$$\left(\frac{\partial t}{\partial \tau}\right)^2 = \sum_{i=0}^{3} k_i^2, \frac{\partial v}{\partial \tau} = v = \frac{\partial v}{\partial t} \frac{\partial t}{\partial \tau} = a m,$$

where $a$ is the relativistic acceleration, and the equations above become,

$$ma = m \begin{pmatrix} 0 & E_1 & E_2 & E_3 \\ -E_1 & 0 & B_3 & -B_2 \\ -E_2 & -B_3 & 0 & B_1 \\ -E_3 & B_2 & -B_1 & 0 \end{pmatrix} \begin{pmatrix} v_0 \\ v_1 \\ v_2 \\ v_3 \end{pmatrix}.$$ 

This is, basically, what one finds in any textbook in physics, see again [Mandl]. The mass, or relativistic energy of the "test particle" here is $m$, and the charge density becomes a kind of relativistic energy density in the (unique) 0-velocity direction $dt_0$, both measured with our clock. Notice that this kind of rest-energy is the only energy or mass that we seem to be able to observe via electromagnetic interaction. The "missing" two other 0-velocity directions might hide black energy/mass? Notice also that $v$ and therefore $a$ are just dependent upon the structure of $H$ and the curve $\gamma$, not on the clock.

With this done, one may consider the Weyl representation,

$$Ph(S(l)) \rightarrow Diff(S(l)),$$

which simply means to introduce the new relations, $[t_i, d_1 t_j] = 0, i \neq j, [t_i, d_1 t_i] = 1$. Then, going about as above, in the case of the Klein-Gordon equation, one obtains the classical quantization of Electro Magnetism, see also the next Example.

For later use, see that the last formula is a kind of Schrödinger equation,

$$\frac{\partial}{\partial t} (\psi) = Q(\psi),$$

where $\Psi \in \Theta_{S(l)}$ and $Q \in \text{End}_{S(l)}(\Theta_{S(l)})$.

Let us now go back to the section, Connections and the Generic Dynamical Structure (3.5), and do all this via the interpretation of an electromagnetic field, as given, for a trivial metric, by the connection $\nabla$, of the tangent bundle $\Theta_{S(l)}$,

$$\nabla_{\dot{k}_i} = \frac{\partial}{\partial t_i} + A_i.$$
Here $A_i \in \operatorname{End}_{S(l)}(\Theta_{S(l)})$, in the above, classical situation, are just a functions, acting as a diagonal matrix. Notice, however, that this connection now has torsion.

As we have seen, there is a canonical associated representation, $\rho_{\nabla}$, for which the curvature is given by,

$$F_{i,j} := \rho_{\nabla}([dt_i, dt_j]) = \frac{\partial A_i}{\partial t_j} - \frac{\partial A_j}{\partial t_i}.$$ 

Since the metric is flat, the generic dynamical system, $(\sigma)$, will give us,

$$d^2 t_i = -\frac{1}{2} \sum_{j=1}^{r} (dt_j [dt_i, dt_j] + [dt_i, dt_j] dt_j),$$

so we have, the force law,

$$d^2 t_i = -\sum_{j=1}^{r} F_{i,j} dt_j - q_i,$$

where the vector $q$ is the charge-current density. In this case we have,

$$q_i = 1/2 \sum_{j=1}^{r} \frac{\partial F_{i,j}}{\partial t_j}.$$ 

The Maxwell equations, the 2 first non-trivial ones, are then equivalent to,

$$q = \nabla(\nabla A) - \nabla^2 (A).$$

We see here that the charge occur also in the equation of motion for fields. It disappeared in the essentially commutative QF-version, presented above. As in the case of General Relativity, where the problem was to explain the notion of Mass, as a property of the geometry of the time-space, the problem here is to explain how the notion of Charge is related to the the geometry of the same time-space. This will be treated in a forthcoming paper by Olav Gravir Imenes, see also [GI].

**Example 3.18.** Gauge groups, Invariant theory, Spin, Isospin, Hypercharge and Quarks. Above we have seen that $Ph(H)$ and therefore also $Ph(S(l))$ are moduli spaces of interest. We know that $U(1)$ acts on $Ph(H)$, conferring a complex structure on the tangent bundle of $H$. Moreover, the fundamental gauge group,

$$G := SU(2) \times SU(3),$$

acts on the complexified tangent bundle of $H$, i.e. there is a principal $G$-bundle $\tilde{G}$ defined over $H$, acting on $\Theta_{H}$. Notice that if we choose a velocity $\nu$, i.e. a directed line $l$ in a tangent space of $H$, and the corresponding Minkowski space-time defined by this directed tangent-line, then the action of $G$ on the tangent space of this space-time, is trivially invariant under Lorentz boosts, since it, of course, leaves $H(l)$ fixed. This suggests that if we had had a natural extension of the action of $\tilde{G}$ to the non-commutative algebraic scheme, $\operatorname{Simp}(Ph(H))$, then the non-commutative orbit spaces,

$$\mathcal{L} := \operatorname{Simp}(Ph(H) : \tilde{G}),$$
or
\[ L := \text{Simp}(PH : \tilde{g}), \]
would have been a prime target for mathematical physics, see [La 4]. However, this is not possible unless one gives up the commutativity of the base space \( H \). We shall therefore, at this moment, restrict ourselves to a rather trivial special case, which is sufficiently general to explain the Dirac equation, see below.

First, any Lie group \( G \), acting on a \( k \)-algebra \( A \), induces a homomorphism of Lie algebras, \[ \eta : g \rightarrow \text{Der}_k(A). \]

For a fixed integer \( n \), there is a versal family, \[ \tilde{\rho} : A \rightarrow \text{End}_{C(n)}(\tilde{V}). \]

Any element \( \chi \in g \), considered as a derivation of \( A \), acts, according to Theorem (3.2), like, \[ \tilde{\rho}(\chi(a)) = [\chi](\tilde{\rho}(a)) + [\nabla_\chi, \tilde{\rho}(a)]. \]

for \( a \in A \), and for some "Hamiltonian" \( \nabla_\chi \). This looks like a connection, \[ \nabla : g \rightarrow \text{End}_k(\tilde{V}). \]

Clearly, the condition \( [\chi] = 0 \) for all \( \chi \in g \), implies that the action \( \eta \) induces an \( A - g \)-module structure on \( \tilde{V} \),

\[ \nabla : g \rightarrow \text{End}_{C(n)}(\tilde{V}), \]

as in [La 4], p.563. This means that \( \nabla \) is a Lie-algebra homomorphism, and that for all \( a \in A \), \( \chi \in g \), \( \psi \in \tilde{V} \), we have

\[ \nabla_\chi(\tilde{\rho}(a)(\psi)) = \tilde{\rho}(a)\nabla_\chi(\psi) + \tilde{\rho}(\chi(a))(\psi). \]

In particular the curvature vanishes, i.e.

\[ R(\chi, \eta) := [\nabla_\chi, \nabla_\eta] - \nabla_{[\chi, \eta]} = 0. \]

Moreover, we find that the sub-scheme,

\[ C(n; g) := \{ c \in \text{Spec}(C(n)) | \forall \chi \in g, [\chi](c) = 0 \}, \]

is a sub-scheme of \( \text{Simp}(A : g) \), the non-commutative invariant space defined by the \( g \)-action.

In fact, what we do by imposing the conditions \( [\chi] = 0 \), for all \( \chi \in g \), is to construct a slice in \( C(n) \) cutting all orbits of \( [g] \) at a single point. The result is obviously an orbit space.

Picking a line \( l \subset E^3 \) and, as above, considering the subspace \( S(l) \), we see immediately that the trivial principal bundle \( SU(2) \), as well as \( su(2) \), and therefore also the complexified \( su(2) \), i.e. \( sl_2(C) \), acts on the restriction of the tangent bundle of \( H \) to \( S(l) \).
The tangent bundle of $H$, restricted to $H - \Delta$ decomposes as,
\[ CB_o \oplus CB_x \oplus CA_{o,x}, \]
and on on the exceptional fibers of $H$, it decomposes as,
\[ CC_o \times CA_o \times C\Delta. \]

Restricted to $S(l)$ these bundles are natural representation of $su(2)$ on $CB_o \oplus CB_x$, and of $su(3)$ on $C\Delta$. But, of course, $su(2)$ acts trivially on $Ph(H(l))$. The above discussion concerning invariant spaces imply that we should be interested in $I = Simp(Ph(S(l)) : su(2))$ or the representation theory of $Ph(S(l)) \times U(su(2))$, where $U(g)$ is the universal enveloping algebra of the Lie algebra $g$.

As an example, pick $\alpha \in su(2) \subset \text{End}(C^2)$. Using the parity-operator $P$, see the Introduction, it acts naturally on $CB_o$ and on $CB_x$, as $\alpha$ and $-\alpha$, respectively. So it acts as,
\[ \gamma_0(\alpha) = \begin{pmatrix} \alpha & 0 \\ 0 & -\alpha \end{pmatrix}, \]
on the complex rank 4 vector bundle, $CB_o \oplus CB_x$, the sections of which are the spinors of the physicists. Composed with the parity operator $P$, $\gamma(\alpha)$ acts like,
\[ \gamma(\alpha) = \begin{pmatrix} 0 & \alpha \\ -\alpha & 0 \end{pmatrix}, \]
i.e. like Diracs representation.

Fix now a representation of $su(2)$, given in terms of the generators,
\[ \alpha_1 = \begin{pmatrix} i & 0 \\ 0 & -i \end{pmatrix}, \alpha_2 = \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix}, \alpha_3 = \begin{pmatrix} 0 & i \\ i & 0 \end{pmatrix}, \]
the analogues of the Pauli matrices, and put, $L = (\alpha_1, \alpha_2, \alpha_3)$.

Notice that, picking the basis of $CB_o \oplus CB_x$ given by $e_1 = (1, 0, -1, 0), e_2 = (0, 1, 0, -1), e_4 = (1, 0, 1, 0), e_5 = (0, 1, 0, 1)$, then $e_1, e_2$ correspond to light velocities, and $e_4, e_5$ correspond to zero velocities. Let $S_{o,z}$ be the subspace generated by the light-velocities $e_1, e_2$, and let $I_{oz}$ be the subspace generated by the 0-velocities $e_4, e_5$, then the parity operator $P$ will permute $S_{o,z}$ and $I_{oz}$.

Any particle, i.e. any simple representation of $Ph(S(l)) \ltimes U(su(2))$, is now canonically a $su(2)$-representation, as in classical quantum theory, where one imposes such an action, a spin structure, and in particular a Casimir element $S = L^2 \in su(2), L^2 = \sum \alpha_i^2$. Notice also that we have here a double situation, a spin structure for the action on the light-velocities, $S_{o,z}$, and an iso-spin structure for the action on the 0-velocities, $I_{oz}$. Let us explain this a little better: The Lie algebra $su(3)$ has a 2-dimensional Cartan subalgebra $\mathfrak{h}$, and two copies of $su(2)$, $\mathfrak{g}_1$, and $\mathfrak{g}_2$.

We may pick, in an essentially unique way, $\mathfrak{g}_1$, such that it leaves the $dt_0$-direction of the tangent space invariant, together with a non-zero element $s \in \mathfrak{h} \cap \mathfrak{g}_1 \subset \mathfrak{h}_1$, where the last Lie algebra is the 1-dimensional Cartan Lie-algebra of $\mathfrak{g}_1$. We may, with an obvious matrix notation, pick
\[ s = \begin{pmatrix} 1/2 & 0 \\ 0 & -1/2 \end{pmatrix}, \]
and assume that \( h \) has a basis given by \( s \) and the element,

\[
y = \begin{pmatrix}
1/3 & 0 & 0 \\
0 & 1/3 & 0 \\
0 & 0 & -2/3
\end{pmatrix}.
\]

Physicists call \( s \) the \textit{isospin}, and \( y \) the \textit{hypercharge}. They are commonly denoted \( T_z \) and \( Y \), and of course, are dependent upon the choice of the direction, \( l \), and therefore of \( H(l) \). The common eigenvectors for \( s \) and \( y \) are called \textit{quarks}. They come as up-quark, as down-quark, or as strange-quark. Obviously, there is room for far more strange and colorful occupants of the tangent bundle of \( S(l) \), by combining the many observables that we have at hand.

Now, check that the following equations define a dynamical structure on the algebra \( Ph(S(l)) \otimes U(su(2)) \),

\[
\delta(t_i) = dt_i + \alpha_i, \quad \delta(dt_i) = 0, \quad i = 0, 1, 2, 3, \quad \delta \alpha_i = \sum_{j=0}^3 [\alpha_j, \alpha_i] dt_j.
\]

Notice that in the universal algebra of \( su(2) \), we have \( 1 = -\alpha_3^2 \), and put,

\[
\gamma(\alpha_0) = \begin{pmatrix} 1 & 0 \\ 0 & -1 \end{pmatrix}.
\]

Let \( \alpha_i \) act on \( CB_o \times CB_x \) as \( \gamma_i := \gamma(\alpha_i) \), and \( dt_i \) act as \( \nabla_{\delta_i} := \frac{\partial}{\partial t_i} \). Put \( \nabla := \sum_{i=0}^3 \nabla_{\delta_i} \) and \( L := 1/2 \sum_{i=0}^3 \gamma_i \). Assuming that this gives us a singular representation \( \rho \) of \( Ph(S(l)) \otimes U(su(2)) \) on \( CB_o \times CB_x \), we compute the Hamiltonian, and find,

\[
Q = \Delta + 2D
\]

where \( D = 1/2 \sum_{i=0}^3 \nabla_{\delta_i} \gamma_i = \nabla L \) is the Dirac operator. Since for any ”simple” representation \( \rho \) of \( Ph(S(l)) \otimes U(su(2)) \) on \( CB_o \times CB_x \), we compute the Hamiltonian, and find,

\[
Q = \Delta + 2D + S,
\]

where all 3 components commute. The equation for a particle in state \( \psi \) with mass \( m \) and spin \( s \) is,

\[
Q(\psi) = (\Delta + 2D + S)(\psi) = (m + s)^2 \psi,
\]

and so, we find,

\[
(D - sm)\psi = 0,
\]

i.e. the Dirac equation.

However, there are more information here. Computing we find the following force law:

\[
\delta^2(t_i) = \delta \alpha_i = \sum_{j=0}^3 [\alpha_j, \alpha_i] dt_j
\]

\[
\delta^2(\alpha_i) = \sum_{j,k=0}^3 ([[\alpha_k, \alpha_j], \alpha_i] + [\alpha_j, [\alpha_k, \alpha_i]]) dt_k dt_j = \sum_{j,k=0}^3 [\alpha_k, [\alpha_j, \alpha_i]] dt_k dt_j.
\]
If we go back to the situation of (3.18), and consider the moduli space $H = \tilde{H}/\mathbb{Z}_2$, see the Introduction, then it is clear that the gauge group, $\mathbb{Z}_2$, maps $B_o$ isomorphically onto $B_x$ and vice versa. Let $f_+: B_o \oplus B_x \rightarrow B_o \oplus B_x$ and $f_- : B_o \oplus B_x \rightarrow B_o \oplus B_x$ be the corresponding projections. Obviously we have,

$$\{f_-, f_+\} = 1.$$ 

We therefore have the 4-dimensional situation described in Grand Picture, and we see that the generator $\Sigma \in \mathbb{Z}_2$, not to be confused with the Parity operator $P$ discussed above, reverses spin and isospin, but conserves the hypercharge. Notice that for a point in $\tilde{\Delta}$ the decomposition of the tangent bundle of $\mathbb{S}(l)$ is different.

There we have,

$$T_{\tilde{H},o} = C_o \oplus A_o \oplus \tilde{\Delta},$$

where $C_o \oplus A_o$ are light-velocities, and here su(2) act only upon the first factor. Thus here we have no isospin, and no quarks! To end this sketch, notice also that the canonical metric $g$ defined on $\tilde{H}$, is very similar to the Schwartzchild metric of a black hole, with horizon equal to the exceptional fibre. The area of this "black hole horizon", bounding nothing, as a function on $\tilde{\Delta}$, is a candidate for mass-density of the Universe. See also that the symmetry, i.e. the gauge group, corresponding to points in $\tilde{\Delta}$ is different from the gauge group outside $\tilde{\Delta}$. This is analogous to deformation theory, where the automorphism group of an object sees a spontaneous reduction along a deformation outside the modular stratum. It therefore seems to me that, in this purely mathematical model, one might find a correlation between a notion of mass-distribution, and a kind of Higgs mechanism.

Cosmology, Big Bang and all that.

In the paper [La 6], we discussed the possibility of including a cosmological model in our "toy-model" of Time-Space. The 1-dimensional model we presented there was created by the deformations of the trivial singularity, $O := k[x]/(x)^2$. Using elementary deformation theory for algebras, we obtained amusing results, depending upon some rather bold mathematical interpretations of the, more or less accepted, cosmological vernacular. Here we shall go one step further on, and show that our "toy-model", i.e. the moduli space, $H$, of two points in the Euclidean 3-space, or its ´ etale covering, $\tilde{H}$, is created by the (non-commutative) deformations of the obvious singularity in 3-dimensions, $U := k < x_1, x_2, x_3 > / (x_1, x_2, x_3)^2$.

In fact, it is easy to see that the versal space of the deformation functor of the $k$-algebra $U$, contains a flat component (a room in the (commutativized) modular suite, see [La-Pf]) isomorphic to $\tilde{H}$, and that the modular stratum (the inner room) is reduced to the base point. Notice that we are working with the non-commutative model of the 3-dimensional space. This is, of course, not visible in dimension 1, and therefore not highlighted in the above mentioned paper.

The tangent space of the versal base $W_o$ of the deformation functor of $U$, is given, see e.g. [La 0], by the cohomology,

$$T_{W, o} = H^1(k, U; U) = \text{Hom}_F(J, U)/\text{Der},$$

where $\pi : F \rightarrow U$, is a surjection of a (non-comutative) free $k$-algebra $F$ onto $U$, $J = \ker(\pi)$, and $\text{Der}$, the vector space of restrictions of derivations $\delta \in \text{Der}_k(F; U)$
to \( J \). One easily obtains a basis, \( \{ \alpha_k^{i,j} \}_{i,j,k} \) for \( \text{Hom}_F(J,U) \), given in terms of the expression of the values of any \( F \)-linear maps \( \alpha \in \text{Hom}_F(J,U) \),

\[
\alpha(x_i x_j) = \sum_k \alpha_k^{i,j} \epsilon_k,
\]

where \( \epsilon_k \) is the class of \( x_k \) in \( U \).

The \( k \)-vector space of derivations, \( \text{Der}_k(F,U) \), is of dimension 9, but the restrictions to \( J = (x_i x_j) \) are given in terms of,

\[
\delta(x_i x_j) = \delta(x_i) \epsilon_j + \epsilon_i \delta(x_j),
\]

so that, with the notations above, \( \delta_k^{i,j} = \delta_i \) if \( k = j \), \( \delta_k^{i,j} = \delta_j \) if \( k = i \), and \( \delta_k^{i,j} = 2 \delta_i \) if \( k = j \), and 0 otherwise. In particular \( \text{Der} \) is of dimension 3, determined by the values \( \delta^i := \delta(x_i) \).

A point of \( \tilde{\mathbf{H}} \) is an ordered pair \((o,x)\) of two points \( o = (\alpha_1^1,\alpha_2^1,\alpha_3^1), \ x = (\alpha_1^2,\alpha_2^2,\alpha_3^2) \). Consider now the sub vector space \( T(2) \) of \( \text{Hom}_F(J,U) \), generated by the linear maps defined by,

\[
\alpha(x_i x_j) = \alpha_i^1 \epsilon_j + \alpha_j^2 \epsilon_i.
\]

The expressions above show that \( \text{Der} \subset T(2) \). Moreover, the quotient space, i.e. the subspace of the tangent space \( T_{\tilde{\mathbf{W}}^*} \), defined by \( T(2) \), can be represented by the maps,

\[
\alpha(x_i x_j) = (\alpha_i^1 - \alpha_j^1) \epsilon_i,
\]

i.e. by the vectors \( o \bar{x} \) in Euclidean 3-space.

This shows that \( \tilde{\mathbf{H}} - \tilde{\Delta} \) is a natural subspace of \( \tilde{\mathbf{W}} \), defined by the equations,

\[
(x_i - \alpha_i^1)(x_j - \alpha_j^1) = 0, \ i,j = 1,2,3.
\]

Notice that this system of equations is not equivalent to the commutative algebra equations,

\[
(x_1 - \alpha_1^1, x_2 - \alpha_2^1, x_3 - \alpha_3^1)(x_1 - \alpha_1^2, x_2 - \alpha_2^2, x_3 - \alpha_3^2) = 0
\]

which gives as the set of pairs of unordered points in 3-space, i.e. part of \( \mathbf{H} \).

The Lie algebra of infinitesimal automorphisms of \( U \), i.e. \( \text{Der}_k(U,U) \) acts on the tangent space of \( \tilde{\mathbf{W}} \) at the base point, \( * \), leaving only 0 invariant, proving that the modular substratum of \( \tilde{\mathbf{W}} \) is reduced to \( * \). We may identify \( * \), i.e. the singularity \( U \), with the point \( \alpha_i^p = 0, i = 1,2,3, p = 1,2 \), contained in \( \tilde{\Delta} \), but recall that the points of \( \Delta \) do not sit in \( \tilde{\mathbf{H}} \). Moreover, the fact that the tangent space of \( * \) does not contain any 0-velocity vectors, proves that \( * \) is a fixed point of the versal space \( \tilde{\mathbf{W}} \) and therefore of \( \Delta \).

Given any point \( \xi = (o,x) \in \tilde{\mathbf{H}} \) there is a translation \( \omega(o,x) \in \tilde{\Delta}(\xi) \), the meaning of which should be clear, translating the vector \( o \bar{x} \) so that its middle point coincides with \( * \). The length, \( \omega \) of \( \omega(o,x) \) is our Cosmological Time.

For every point \( (o,x) \in \tilde{\mathbf{H}} \) there is a vector \( \xi \in T_{\tilde{\mathbf{H}},0} = \bar{\epsilon}(*) \) such that \( o \bar{x} \) is a translation by \( \omega(o,x) \) of \( \omega \xi \). We may express this by saying that \( (o,x) \) is created by the tangent vector \( \xi \) of \( U \) at the cosmological time \( \omega \).
Since $*$ is a fixed point in $\Delta \simeq \mathbb{E}^3$, the geometry of our space $\tilde{H}$ changes. The natural symmetry (gauge) group operating on $\mathbb{E}^3$ is now $G := \text{GL}(3)$. The action of $g := \text{lie}G$ defines a distribution $\tilde{\mathfrak{g}}$, in $\tilde{H}$, which is different from $\tilde{\Delta}$. In particular, the diagonal group, $K^* \subset G$, generates a 1-dimensional distribution, $\tilde{\omega} \subset \tilde{\mathfrak{g}}$, in $\tilde{H}$.

Notice that the subspace $U(\omega) \subset \tilde{H}$ with a fixed cosmological time $\omega$, is of dimension 5. The tangent space of any point $t \in U(\omega)$ has a 2-dimensional subspace, $\tilde{\beta}(t)$, normal to $\tilde{\omega}(t)$, which is not a subspace of $\tilde{\Delta}(t)$.

The metric $g$ defined on $\tilde{H}$ may now be written,

$$g = \lambda(t)d\omega^2 + g(\omega)$$

Here $g(\omega)$ is the metric induced on $U(\omega)$.

Notice that the maximal distribution $\tilde{\mathfrak{c}}$, the light-velocities is no longer defined. From the above equations, we see that cosmological time is therefore of the same form as Einstein’s proper time. This should be compared with the Einstein-de Sitter metric of the elementary cosmological model, see [Sachs-Wu],

$$-g = du^4 \otimes du^4 - (R^2(u^4)) \sum_{i,j=1}^{3} du^i \otimes du^j,$$

and the Friedmann-Robertson-Walker metric, see [Elbaz],

$$-ds^2 = dr^2 = dt^2 - R^2(t)\left(\frac{dr^2}{1 - kr^2} + r^2 d\Omega^2\right).$$

It seems to us, that our model has an advantage over these, ad-hoc, models. Accepting Big Bang as the creator of space, via the (uni)versal family of the primordial $U$, changes the geometry of $\tilde{H}$. In this picture, the point $(o, x) =: t \in \tilde{H}$ is created from $* = U$, via the tangent vector $v := 1/\omega \vec{o}x$ in the tangent space of $* = U$, in the direction, defined by the middle point of $(o, x = t)$, of $\Delta$, and in the time-span defined by the cosmological time $\omega$. This defines a half-line $\bar{\pi}_{\omega}$ extending from $*$ through $t \in \tilde{H}$.

The visible universe $\mathcal{V}(o)$, for $\text{ME} = (o, o)$ should be the union of all world curves of photons, leaving $*$, reaching $\text{ME}$. This should coincide with the union of all solutions of the Lagrangian $\lambda(t)d\omega^2 = dt^2 - g(\omega)$, or with the corresponding geodesics of this non-positive definite metric. This implies that the universe must be curved, just like the picture drawn in [La 6], page 262, suggests. In that paper we just worked with a 1-dimensional real space. Never the less, the situation here, in our 3-dimensional picture, is basically the same, complete with a Hubble-constant, that is not really constant, etc.

If we want a catchy way to express these basic properties of our model, we might say that, $U$, the Big Bang, that created our (visible) world, is equivalent to any
center of the exceptional fiber in \( \widehat{H} \), and so the infinite small has the same structure as the origin of the world, \( U \).

In particular, see Example 3.18, we find that the symmetry brake of our model for quarks, in \( S(l) \), where the \( su(3) \) symmetry is broken by the unique 0-velocity \( dt_0 \), now could have been made global, using the cosmological distribution instead of \( dt_0 \).

Making all this fit with contemporary quantum theory and cosmology is, however, not an easy task. There are serious interpretational difficulties here, as well as in most papers we have seen, on cosmology. We shall therefore leave it for now, and hopefully be able to return to this later.

\section*{§4 Interaction and Non-commutative Algebraic Geometry}

Given a dynamical system \( \sigma \) of, order 2. A particle, \( \tilde{V} \), that we know "happened" at some point \( t \in \text{Simp}_n(A(\sigma)) \), producing a simple representation \( V := \tilde{V}(t) \) will after some time \( \tau \) have developed into the particle sitting at a point on the integral curve \( \gamma \) defined by the vectorfield \( \xi \) of \( \sigma \), at a "distance" \( \tau \) in \( \text{Simp}_n(A(\sigma)) \) (we are of course assuming the field \( k \) is contained in the real numbers). Now, this may well be a point on the border of \( \text{Simp}_n(A(\sigma)) \), i.e. in \( \Gamma_n = \text{Simp}(C(n)) - U(n) \), where it decays into an indecomposable, or into a semi-simple, representation, i.e. into two or more new particles \( \{ V_i \in \text{Simp}_n(A(\sigma), n = \sum n_i \} \).

What happens now is taken care of by the following scenario: If the different particles we have produced are not interacting, each of the new particles should be considered as an independent object, evolving according to the Dirac derivation \( \delta \). However, if the particles we have produced are interacting, we have a different situation.

Notice first that for \( n = 1 \), we have a canonical morphism of schemes,

\[ \text{Simp}_1(A(\sigma)) \longrightarrow \text{Simp}_1(\mathcal{A}) \]

and a canonical vector-field \( \xi \) in \( \text{Simp}_1(A(\sigma)) \), the phase space. Given any point of \( \text{Simp}_1(\mathcal{A}) \), the configuration space, and any tangent-vector at this point, there is an integral curve of \( \xi \) in \( \text{Simp}_1(A(\sigma)) \), through the corresponding point, projecting down to a fundamental curve in the configuration space.

For \( n \geq 2 \) the spaces \( \text{Simp}_n(A(\sigma)) \) and \( \text{Simp}_n(\mathcal{A}) \) are, however, totally different and without any easy relations to each other.

Let now \( v_i \in \text{Simp}_n(A(\sigma)) \), \( i = 1, 2 \) be two points of \( \text{Simp}(A(\sigma)) \) corresponding to representations \( V_1, V_2 \), maybe in different components, and/or ranks. Consider their components, i.e. the universal families in which they are contained,

\[ \tilde{\rho}_i : A(\sigma) \longrightarrow \text{End}_{C(n_i)}(\tilde{V}_i) \]

The Dirac derivation, \( \delta \), defines derivations,

\[ [\delta_i] : A(\sigma) \longrightarrow \text{End}_{C(n_i)}(\tilde{V}_i) \]

and therefore also the fundamental vector-fields, \( \partial_i \in Ext^1_{A(\sigma) \otimes C(n_i)}(\tilde{V}_i, \tilde{V}_i) \), and \( \xi_i \in \text{Der}(C(n_i)) \).
4.1 Definition. Let $B$ be any finitely generated $k$-algebra. We shall say that the components, $C_1 \subseteq \text{Simp}_{n_1}(B)$, $C_2 \subseteq \text{Simp}_{n_2}(B)$, or the corresponding particles $\tilde{V}_i$, $i=1,2$, are non-interacting if

$$\text{Ext}_B^1(V_1, V_2) = 0, \forall v_1 \in C_1, \forall v_2 \in C_2.$$ 

Otherwise they interact.

Suppose now that the points $v_1$ and $v_2$, sit in $\text{Simp}_{n_1}(A(\sigma))$ and $\text{Simp}_{n_2}(A(\sigma))$, respectively. "Physically", we shall consider this as an "observation" of two particles, $\tilde{V}_1$ and $\tilde{V}_2$ in the "state" $V_1$ and $V_2$, at some instant. If the two particles are non-interacting, the resulting entity, considered as the sum $V := V_1 \oplus V_2$, of dimension $n := n_1 + n_2$, as module over $A(\sigma)$, will stay, "as time passes", a sum of two simples.

If $V_1$ and $V_2$ interact, this may change. To explain what may happen, we have to take into consideration the non-commutativity of the geometry of $\text{PhA}$. In particular, we have to consider the non-commutative deformation theory, see §2, and [La 2,3,4]. Consider the deformation functor,

$$\text{Def}_{\{V_1, V_2\}} : a_2 \rightarrow \text{Sets},$$

or, if we want to deal with more points, say a finite family $V_i, i = 1,2,...,r$, the deformation functor,

$$\text{Def}_{\{V_i\}} : a_r \rightarrow \text{Sets},$$

and its formal moduli,

$$H := \begin{pmatrix} H_{1,1} & \cdots & H_{1,r} \\ \vdots & \ddots & \vdots \\ H_{r,1} & \cdots & H_{r,r} \end{pmatrix},$$

together with the versal family, i.e. the essentially unique homomorphism of $k$-algebras,

$$\tilde{\rho} : A(\sigma) \rightarrow \begin{pmatrix} H_{1,1} \otimes \text{End}_k(V_1) & \cdots & H_{1,r} \otimes \text{Hom}_k(V_1, V_r) \\ \vdots & \ddots & \vdots \\ H_{r,1} \otimes \text{Hom}_k(V_r, V_1) & \cdots & H_{r,r} \otimes \text{End}_k(V_r) \end{pmatrix}.$$

This is, in an obvious sense, the universal interaction. However, we need a way of specifying which interactions we want to consider. This is the purpose of the following, tentative, definition.

4.2 Definition. Given $v_i \in \text{Simp}(A(\sigma)), i = 1,...,r$. An interaction mode for the corresponding family of modules $\{V_i\}, i=1,...,r$, is a right $H(\{V_i\})$-module $M$.

An interaction mode is a kind of higher order preparation, see (1.2). It consists of a rule, telling us, for the given family of $r$ points, $v_i \in \text{Simp}_{n_i}(A(\sigma))$ how to prepare their interactions. The structure morphism $\phi : H(\{V_i\}) \rightarrow \text{End}_k(M)$, fixes all relevant higher order momenta, i.e. $\phi$ evaluates all the tangents between these modules, and by the Beilinson-type Theorem, see §2, creates a new $A(\sigma)$-module.

In fact, an interaction mode induces a homomorphism,

$$\kappa(M) : A(\sigma) \rightarrow \text{End}_k(\tilde{V}),$$
where \( \tilde{V} := \oplus_{i=1,...,r} M_i \otimes V_i \).

Thus, we have constructed a new \( n \)-dimensional \( A(\sigma) \)-module, which may be decomposable, indecomposable or simple, depending on the interaction mode we choose, and, of course, depending upon the tangent structure of the moduli space \( Simp(A(\sigma)) \).

Assuming the impossible, that our \( k \)-algebra of observables, \( A(\sigma) \), consisted of all observables that our curiosity fancied. Assume moreover that the notion of interaction mode, i.e. a right \( H(Simp(A(\sigma))) \)-module, made sense, then we might talk about a QFT-theory of the UNIVERSE. The piece-wise defined curves, \( \gamma_0, \gamma_1, ..., \gamma_l, ... \), corresponding to successive choices at each "moment" of decay, of a new interaction mode, and therefore of a new integral curve in the relevant time-space \( Simp_N(A(\sigma)) \), might be called a History.

Determining the conditions for de-coherence, and the assignment of probabilities for these choices, will be left for now. See [Gell-Mann], and [Gell-Mann and Hartle].

The problem of time, in this context, observed from inside or outside of the universe, will also be postponed, maybe à la calendes grecques.

**Example 4.3.** Let us consider the case of two point-objects interacting in 3-space. Go back to the Example (1.1), (iii), and let us look at two fields, 

\[
\phi(i) : A_0 := k[x_1, x_2, x_3] \to k[\tau], \quad i = 1, 2,
\]

inducing,

\[
Ph(\phi(i)) : A := Ph(k[x_1, x_2, x_3]) \to Ph(k[\tau]), \quad i = 1, 2.
\]

This corresponds to two curves \( \gamma(i), \quad i = 1, 2 \), in the phase space, \( A \), parametrized by the same clock-time \( \tau \). Start with \( \tau = 0 \), and let this time correspond to the two points, 

\[
\phi(i, 0) = v_i := (q_i, p_i) \in A, \quad i = 1, 2.
\]

Let \( V_i := k(v_i), \quad i = 1, 2 \), assume \( q_1 \neq q_2 \), and use (1.1)(iii), to find that the formal moduli of the family \( \{ V_1, V_2 \} \) of \( A \)-modules, has the form,

\[
\left( \begin{array}{c}
\hat{A}_{p_1} \\
< \tau_{2,1} > \
\hat{A}_{p_2}
\end{array} \right)
\]

where \( \tau_{i,j} \) is a generator for \( Ext^1_A(V_i, V_j)^* \simeq k \).

Since \( Hom_k(V_i, V_j) = k, \quad i, j = 1, 2 \), we have a natural, versal family, i.e. a morphism,

\[
A \to \left( \begin{array}{c}
\hat{A}_{p_1} \\
< \tau_{2,1} > \
\hat{A}_{p_2}
\end{array} \right),
\]

An interaction mode between the two point-particles, defined by \( \phi(i), \quad i = 1, 2 \), is now given by evaluating the \( \tau_{i,j} \), together with expressing the two fields \( \phi_i \) by the corresponding morphism,

\[
\phi : A \to \left( \begin{array}{c}
Ph(k[\tau]) \\
0 \\
0
\end{array} \right).
\]

A force law should now be given by some elements,

\[
\psi_{i,j}(\tau) \in Ext^1_A(\phi_i(\tau), \phi_j(\tau)).
\]
by, say, putting,
\[ \psi_{i,j}(\tau) = \phi(i, j : \tau) \cdot \xi_{i,j}, \]
where \( \xi_{i,j} \) is the generator of \( \text{Ext}^1_A(\phi_i, \phi_j) \) found in (1.1), (iii).

In the general case, we need a notion of hypermetric to formulate a reasonable theory of interaction. This will, I hope, be the topic of a forthcoming paper.

**Example 4.4.** Let us consider the notion of interaction between two particles, \( V_i := k(v_i), \ i = 1, 2, \) as above, in the spirit of §1, i.e. by looking at the \( A_0 := k[x] \)-module \( V := V_1 \oplus V_2, \) i.e. the homomorphism of \( k \)-algebras, \( \rho_0 : A_0 \to \text{End}_k(V), \) and let us try to extend this module-structure to a
\[ \rho : \mathcal{P} h^\infty A_0 \to \text{End}_k(V). \]

We have the following relations in \( \mathcal{P} h^\infty A_0: \)
\[ [x_i, x_j] = 0 \]
\[ dx_i, x_j + [x_i, dx_j] = 0 \]
\[ \ldots \]
\[ \sum_{l=0}^{p} \binom{p}{l} [d^l t_i, d^{p-l} t_j] = 0. \]

Put,
\[ \rho_0(x_i) = \rho_0(dx^0 x_i) = \begin{pmatrix} x_i(1) & 0 \\ 0 & x_i(2) \end{pmatrix} =: \begin{pmatrix} \alpha^0_i(1) & 0 \\ 0 & \alpha^0_i(2) \end{pmatrix}, \]
and, \( \alpha^0_i(r, s) := x_i(r) - x_i(s), \ r, s = 1, 2. \) Let, for \( q \geq 0, \)
\[ \rho(dx^q x_i) = \begin{pmatrix} \alpha^q_i(1) & r^q_i(1, 2) \\ r^q_i(2, 1) & \alpha^q_i(2) \end{pmatrix}. \]

Now, compute, for any \( p \geq k, \)
\[ [\rho(d^k x_i), \rho(d^{p-k} x_j)] = \]
\[ \begin{pmatrix} r^k_i(1, 2) r^{p-k}_j(2, 1) - r^{p-k}_j(1, 2) r^k_i(2, 1) & r^{p-k}_j(1, 2) \alpha^k_i(1, 2) + r^k_i(1, 2) \alpha^{p-k}_j(2, 1) \\ r^k_i(2, 1) \alpha^{p-k}_j(1, 2) + r^{p-k}_j(2, 1) \alpha^k_i(2, 1) & r^k_i(2, 1) r^{p-k}_j(1, 2) - r^{p-k}_j(2, 1) r^k_i(1, 2) \end{pmatrix} \]
and observe that,
\[ [\rho(dx^q x_i), \rho(dx^r x_j)] + [\rho(x_i), \rho(dx^p x_j)] = \]
\[ \begin{pmatrix} r^q_i(1, 2) \alpha^0_j(2, 1) + r^0_j(1, 2) \alpha^q_i(1, 2) & 0 \\ 0 & r^q_j(1, 2) \alpha^0_i(2, 1) + r^0_i(1, 2) \alpha^q_j(1, 2) \end{pmatrix}. \]

After some computation we find the following condition for these matrices to define a homomorphism \( \rho, \) independent of the choice of diagonal forms,
\[ r^k_i(r, s) = \sum_{l=0}^{k} \binom{k}{l} \sigma_{k-l}^l \alpha^0_i(r, s), \ r, s = 1, 2, \]
where the sequence \( \{ \sigma_l \} \), \( l = 0, 1, \ldots \) is an arbitrary sequence of coupling constants, with \( \sigma_0 = 0 \) and \( \sigma_l \) of order \( l \). By recursion, we prove that this is true, for \( k \leq p - 1 \), therefore \( r_k^p(1, 2) = r_p(2, 1) \), and so the diagonal elements above vanish, i.e.

\[
r_k^p(1, 2)r_j^{p-k}(2, 1) - r_j^{p-k}(1, 2)r_k^p(2, 1) = 0.
\]

The general relation is therefore proved if we can show that with the above choice of \( r_k^p(r, s) \) we obtain, for every \( p \geq 0 \),

\[
\sum_{k=0}^{p} \left( \begin{array}{c} p \\ k \end{array} \right) \left( \sum_{l=0}^{p-k} \left( \begin{array}{c} p-k \\ l \end{array} \right) \sigma_{p-k-l}^l(1, 2) \alpha_k^p(1, 2) + \sum_{l=0}^{p} \left( \begin{array}{c} p \\ k \end{array} \right) \left( \sum_{l=0}^{k} \left( \begin{array}{c} k \\ l \end{array} \right) \sigma_{k-l}^l(1, 2) \alpha_j^{p-k}(2, 1) = 0,
\right.
\]

and this is the formula,

\[
\sum_{k=0}^{p} \left( \begin{array}{c} p \\ k \end{array} \right) \left( \sum_{l=0}^{k} \left( \begin{array}{c} k \\ l \end{array} \right) \sigma_{k-l}^l(1, 2) \alpha_j^{p-k}(2, 1) - \alpha_j^l(1, 2) \alpha_j^{p-k}(2, 1) \right) = 0,
\]

Notice that the relations above are of the same form for any commutative coefficient ring \( C \), i.e. they will define a homomorphism,

\[
\rho : Ph^\infty A_0 \to M_2(C),
\]

for any commutative \( k \)-algebra \( C \).

Now, consider the Dirac time development \( D(\tau) = exp(\tau \delta) \) in \( D \), the completion of \( Ph^\infty A_0 \). Composing with the morphism \( \rho \) defined above, we find a homomorphism,

\[
\rho(\tau) : Ph^\infty (A_0) \to M_2(k[[\tau]]),
\]

where,

\[
X_i := \rho(\tau)(t_i) = \begin{pmatrix} \Phi_i(1) & \Phi_i(1, 2) \\ \Phi_i(2, 1) & \Phi_i(2) \end{pmatrix}
\]

and

\[
\Phi_i(r) = \sum_{n=0}^{\infty} \frac{1}{(n!)^2} \tau^n \cdot \alpha_i^n(r), \ r = 1, 2,
\]

\[
\Phi_i^0(r, s) = \sum_{n=0}^{\infty} \frac{1}{(n!)^2} \tau^n \cdot \alpha_i^n(r, s), \ r, s = 1, 2,
\]

\[
\sigma = \sum_{n=0}^{\infty} \frac{1}{(n!)^2} \tau^n \cdot \sigma_n,
\]

\[
\Phi_i(r, s) = \sigma \cdot \Phi_i^0(r, s), \ r, s = 1, 2.
\]
This describes the most general, Heisenberg model, of motion of our particles, clocked by $\tau$. Observe that the interaction acceleration $\Phi_i(r,s)$ is pointed from $r$ to $s$, just like in physics!

The formula above, is now seen to be a consequence of the obvious equality of the two products of the formal power series, $\sigma \cdot (\Phi_0^0(1,2)\Phi_0^0(1,2))$ and $\sigma \cdot (\Phi_0^0(1,2)\Phi_0^0(1,2))$, just compare the coefficients of the resulting power series.

What we have got is nothing but a formula for commuting matrices $\{X_i\}_{i=1}^d$ in $M_2(k[[\tau]])$, since for such matrices we must have,

$$\frac{d^n}{d\tau^n}[X_i, X_j] = 0, \quad n \geq 1.$$

The eigenvalues $\lambda_i(1,\tau)$, and $\lambda_i(2,\tau)$ of $X_i$ describes points in the space that should be considered the trajectories of the two points under interaction. This is OK, at least as long as we are able to label them by $1$ and $2$ in a continuous way with respect to the clock time $\tau$. If all coupling constants, $\sigma_n$, $n \geq 0$, vanish, then the system is simply given by the two curves $\Phi(r) := (\Phi_1(r), \Phi_2(r), \ldots, \Phi_d(r))$, $r = 1, 2$, where $d$ is the dimension of $A_0$. In general, the eigenvalues of $X_i$ are given by,

$$\lambda_i(r) = 1/2 \cdot (\Phi_i(1) + \Phi_i(2))$$

$$(-1)^{1/2} \sqrt{(\Phi_i(1) + \Phi_i(2))^2 - 4(\Phi_i(1) \cdot \Phi_i(2) + \sigma^2(\Phi_i(1) - \Phi_i(2))^2)},$$

for $r = 1, 2$. Clearly,

$$1/2(\lambda_i(1) + \lambda_i(2)) = 1/2(\Phi_i(1) + \Phi_i(2))$$

$$\lambda_i(1) - \lambda_i(2) = (\Phi_i(1) + \Phi_i(2))^2 - 4(\Phi_i(1) \cdot \Phi_i(2) + \sigma^2(\Phi_i(1) - \Phi_i(2))^2)$$

$$= (1 - 4\sigma^2)(\Phi_i(1) - \Phi_i(2))^2.$$

Denote by, $\lambda(r) = (\lambda_1(r), \lambda_2(r), \ldots, \lambda_d(r))$, $r = 1, 2$, the vectors corresponding to the eigenvalues, and by,

$$o := 1/2(\lambda(1) + \lambda(2)) = 1/2(\Phi(1) + \Phi(2))$$

the common median, and put,

$$R_0 := |(\Phi(1) - \Phi(2)|$$

$$R := |\lambda(1) - \lambda(2)|,$$

then

$$R = \sqrt{(1 - 4\sigma^2)}R_0.$$

Choose coupling constants such that,

$$\frac{d^2}{d\tau^2} R = r R^{-2},$$

where $r$ is a constant we should expect Newton like interaction. If $\sigma \geq 1/2$, the point particles are confounded, the eigenvalues of $X_i$ become imaginary, and the result is no longer obvious. If we pick $\sigma = \sqrt{1 - r^2 R_0^{-2}}$, the relative motion will be circular, with constant radius $r$ about $o$. 
Example 4.5. Let $B$ be the free $k$-algebra on two non-commuting symbols, $B = k < x_1, x_2 >$, and see Example (2.14). Let $P_1$ and $P_2$ be two different points in the $(x_1, x_2)$-plane, and let the corresponding simple $B$-modules be $V_1, V_2$. Then, $\text{Ext}^1_B(V_1, V_2) = k$. Let $\Gamma$ be the quiver,

$$V_1 \bullet \longrightarrow \bullet V_2,$$

then an interaction mode is given by the following elements: First the formal moduli of $\{V_1, V_2\}$,

$$H := \begin{pmatrix} k < u_1, u_2 > \quad < t_{1,2} > \\ < t_{2,1} > \quad k < v_1, v_2 > \end{pmatrix},$$

then the $k$-algebra,

$$k\Gamma := \begin{pmatrix} k & k \\ 0 & k \end{pmatrix},$$

and finally a homomorphism,

$$\phi : H \longrightarrow k\Gamma$$

specifying the value of $\phi(t_{1,2}) \in \text{Ext}^1_B(V_1, V_2)$. Since $\text{Hom}_k(V_i, V_j) = k$, we obtain $V = k^2$, and we may choose a representation of $\phi(t_{1,2})$ as a derivation, $\psi_{1,2} \in \text{Der}_k(B, \text{Hom}_k(V_1, V_2))$, such that the $B$-module $V = k^2$ is defined by the actions of $x_1, x_2$, given by,

$$X_1 := \begin{pmatrix} \alpha_1 & 1 \\ 0 & \alpha_2 \end{pmatrix}, \quad X_2 := \begin{pmatrix} \beta_1 & 0 \\ 0 & \beta_2 \end{pmatrix},$$

where $P_1 = (\alpha_1, \beta_1)$ and $P_2 = (\alpha_2, \beta_2)$. $V$ is therefore an indecomposable $B$-module, but not simple.

If we had chosen the following quiver,

$$V_1 \bullet \xleftarrow{\epsilon_{1,2}} \epsilon_{2,1} \bullet V_2, \quad \epsilon_{i,j} \epsilon_{j,i} = 0, \quad i, j = 1, 2,$$

then the resulting $B$-module $V = k^2$ would have been simple, represented by,

$$X_1 := \begin{pmatrix} \alpha_1 & 1 \\ 0 & \alpha_2 \end{pmatrix}, \quad X_2 := \begin{pmatrix} \beta_1 & 0 \\ 1 & \beta_2 \end{pmatrix}.$$
If they do, this situation is analogous to the case which in physics is referred to as the "super-selection rule".

Or, if \([\delta] \in \text{Ext}^1_{\text{A}(\sigma)}(V, V)\) does not sit (or stay) in the modular stratum, the particle \(V\) looses automorphisms, and may become indecomposable, or simple, instantaneously.

We may thus create new particles, and we have in Example (3.7) discussed the notion of lifetime for a given particle. In particular we found that the harmonic oscillator had ever-lasting particles of \(k\)-rank 2. If, however, we forget about the dynamical system, and adopt the more physical point of view, picking a Lagrangian, and its corresponding action, we may easily produce particles of finite lifetime.

**Example 4.6.** Let, as in (3.6) \(A := \text{PhA}_0 = k < x, dx >\), with \(A_0 = k[x]\) and put \(x =: x_1, \, dx =: x_2\). Consider the curve of two-dimensional simple \(A\)-modules,

\[
X_1 = \begin{pmatrix} 0 & 1 + t \\ 0 & t \end{pmatrix}, \quad X_2 = \begin{pmatrix} t & 0 \\ 1 + t & 0 \end{pmatrix},
\]

either as a free particle, with Lagrangian \(1/2dx^2\), or as a harmonic oscillator with Lagrangian \(1/2dx^2 + 1/2x^2\). The action is, in the first case, \(S = 1/2TrX_1^2 = t^2\), and in the second case, \(S = 1/2Tr(X_2^2 + 1/2X_1^2) = 2t^2\). Thus the "Dirac"-derivation becomes \(\nabla S = t \frac{\partial}{\partial t}\), or \(\nabla S = 2t \frac{\partial}{\partial t}\). Computing the Formanek center \(f\), see (3.6), we find,

\[
f(t) = t^2(1 + t)^2 - (1 + t)^4.
\]

The corresponding particle, born at \(t < 0\), decays at \(t = -1/2\), and thus has a finite lifetime. Of course, the parameter \(t\) in this example, is not our time, and the curve it traces is not an integral curve of the dynamic system of the harmonic oscillator, see (3.7). This shows that one has to be careful about mixing the notions of dynamic system, and the dynamics stemming from a Lagrangian-, or from a related action-principle.

**Example 4.7.** Suppose we are given an element \(v \in \text{Simp}(A(\sigma))\), and consider the monodromy homomorphism,

\[
\mu(v) : \pi_1(v; \text{Simp}_n(A(\sigma))) \to \text{Gl}_n(k).
\]

If \(v\) is Fermionic, then there exist a loop in \(\text{Simp}_n(A(\sigma))\) for which the monodromy is non-trivial. Assume the tangent of this loop at \(v\) is given by \(\xi \in \text{Ext}^1_{A(\sigma)}(V, V)\). Since this tangent has no obstructions it is reasonable to assume that there is a quotient,

\[
H([V, V]) \to \begin{pmatrix} k & f^- \\ f^+ & k \end{pmatrix},
\]

with \(f^- f^+ = f^+ f^- = 1\). This would give us an \(A(\sigma)\)-module, with structure map,

\[
A(\sigma) \to \begin{pmatrix} \text{End}_k(V) & f^- \otimes_k \text{End}_k(V) \\ f^+ \otimes_k \text{End}_k(V) & \text{End}_k(V) \end{pmatrix},
\]

i.e. a simple \(A(\sigma)\)-module of Fermionic type, see §3, Grand Picture.
Notice that, for a given connection on the vector-bundle $\tilde{V}$, the correct monodromy group to consider for the sake of defining Bosons and Fermions, etc., should probably be the infinitesimal monodromy group generated by the derivations of the curvature tensor, $R$.

In physics, interactions are often represented by tensor products of the representations involved. For this to fit into the philosophy we have followed here, we must give reasons for why these tensor products pop up, seen from our moduli point of view. It seems to me that the most natural point of view might be the following: Suppose $A$ is the moduli algebra parametrizing some objects $\{X\}$, and $B$ is the moduli for some objects $\{Y\}$, then considering the product, or rather, the pair, $(X,Y)$, one would like to find the moduli space of these pairs. A good guess would be that $A \otimes_k B$ would be such a space, since it algebraically defines the product of the two moduli spaces. However, this is, as we know, too simplistic. There are no reasons why the pair of two objects, should deform independently, unless we assume that they do not fit into any ambient space, i.e. unless the two objects are considered to sit in totally separate universes, and then we have done nothing, but doubling our model in a trivial way.

In fact, we should, for the purpose of explaining the role of the product, assume that our entire universe is parametrized by the moduli algebra $A$, and accepting, for two objects $X$ and $Y$ in this universe, that the superposition, or the pair, $(X,Y)$, correspond to a collection of, maybe new, objects parametrized by $A$.

This is basically what we do, when we assume that $X$ and $Y$ are of the same sort, represented by modules $V$ and $W$ of some moduli $k$-algebra $A$, and then consider some tensor product of the representations, say $V \otimes_k W$, as a new representation, modeling a collection of new particles.

As above we observe that the obvious moduli space of tensor-products of representations, or of the pairs $(V,W)$ is $A \otimes A$. But since these representations should be of the same nature as any representation of $A$, this would, by universality, lead to a homomorphism of moduli algebras,

$$\Delta : A \to A \otimes A,$$

i.e. to a bialgebra structure on the moduli algebra $A$.

This is just one of the reasons why mathematical physicists are interested in Tannaka Categories, and in the vast theory of quantum groups. For an elementary introduction, and a good bibliography, see [Kassel].

See now that, if $A_0$ is commutative, and if we put $A = Ph(A_0)$, then there exist a canonical homomorphism,

$$\Delta : A \to A \otimes_{A_0} A.$$

In fact, the canonical homomorphism $i : A_0 \to Ph(A_0)$ identifies $A_0$ with the a sub-algebra of $A \otimes_{A_0} A$. Moreover, $d \otimes 1 + 1 \otimes d$ is a natural derivation, $A_0 \to A \otimes_{A_0} A$, so by universality, $\Delta$ is defined. Thus, for representations of $A := Ph(A_0)$ there is a natural tensor product, $-$ $\otimes_{A_0}$ $-$ . Thus, in (3.18), the tensor product of the fiber bundles defined on $S(l)$, $PN := \tilde{\Delta} \otimes_H \tilde{\Delta} \otimes_H \tilde{\Delta}$, is, in a natural way, a new representation of $Ph(S(l))$, the fibers of which is the triple tensor product,

$$(3) \otimes_k (3) \otimes_k (3).$$
of the Lie-algebra su(3). The representation $PN$ therefore splits up in the well-known “swarm” of elementary particles, among which, the proton and the neutron, see (3.18).

Notice, finally, that the purpose of my notion of swarm, see [La 4], is to be able to handle a more complicated situation than the one above. One should be prepared to sort out the ”swarm” of those representations of the known observables, that one would like to accept as models for physical objects, and then compute the parameter algebra best fitting this ”swarm”. This is, in my opinion, one of the main objects of a fully developed future non-commutative algebraic geometry.

References


Usefull readings.


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