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AN ESSAY ON THE COMPLETION OF QUANTUM THEORY.
I: GENERAL SETTING

WOLFGANG BERTRAM

Abstract. We propose a geometric setting of the axiomatic mathematical formalism of quantum theory. Guided by the idea that understanding the mathematical structures of these axioms is of similar importance as was historically the process of understanding the axioms of geometry, we complete the spaces of observables and of states in a similar way as in classical geometry linear or affine spaces are completed by projective spaces. In this sense, our theory can be considered as a “completion of usual linear quantum theory”, such that the usual theory appears as the special case where a reference frame is fixed once and for all. In the present first part, this general setting is explained. Dynamics (time evolution) will be discussed in subsequent work.

Dedicated to the memory of Tobias Brandes (1966 – 2017)

Preamble

Since our years of study in Göttingen, Tobias and I had a plan to write, some day, a book, “our” book on quantum mechanics. Our paths separated after the Diplom: Tobias became a physicist, and I, a mathematician. We believed that we had time to carry out our project. But we had not.

What I’m going to write up here, is my version of what might have been a draft for the first chapter of this book. There are a lot of excellent textbooks on quantum mechanics, and our aim cannot be, and never was, to copy them, or to cook a new one by mixing ingredients taken from them. Rather, by writing the book we would have wished to find answers to our own questions – the present version is certainly a biased choice of questions, the one of a mathematician, and Tobias is no longer here to correct and complete it by a physicist’s view. I’m well aware that the text is tentative, piecemeal, and possibly may appear altogether beside the point. My only excuse is that, from a purely mathematical point of view, the ideas exposed in the following seem natural, and are kind of unavoidable. I cannot claim that Tobias would have signed this text, but I’m sure that in the universe where he is now, he will forgive me for quoting his name in relation with ideas and speculations that, certainly, are not quite standard in our universe.

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1. Introduction

1.1. Quantum Mechanics: axioms versus interpretations. Whereas the interpretation of Quantum Mechanics is a hot topic – there are at least 15 different mainstream interpretations\(^1\), an unknown number of other interpretations, and thousands of pages of discussion –, it seems that the mathematical axioms of Quantum Mechanics are much less controversial: the Dirac-von Neumann axioms are generally accepted to be their definite version ([D, vN]). Although I find exiting and interesting the discussion on “interpretations”, I do not feel qualified to contribute to it. As a mathematician I feel more competent to comment on the axiomatic and formal structure of quantum mechanics: without being disrespectful towards Dirac and von Neumann, I find surprising that the “definite” form of the axioms has been fixed 85 years ago, shortly after the main discoveries of quantum theory had been made, and that since then essentially nothing has been changed. The whole discussion seems to turn around the “interpretation” of a theory whose formal mathematical structure is defined once and for all, without taking seriously into consideration that the axiomatic foundations may be questionable. This calls for comparison with the history of axioms of geometry: Euclide’s axiomatic construction of geometry is certainly among the greatest achievements of the human mind in ancient history; however, sticking to the axioms too closely prevented men for a long time from discovering non-Euclidean geometry. The rapid development of modern mathematics was possible only after mathematicians had questioned the structure of Euclide’s axioms. Could something similar occur with the axioms of quantum theory? I think this possibility cannot be completely excluded.

Of course, I neither claim that the Dirac-von Neumann axioms were “wrong” (they can be no more “wrong” than Euclide’s), nor to have a full-fledged counter-proposition, like Hilbert had when proposing his “Grundlagen der Geometrie”, putting Euclide’s axioms onto a rigorous and modern base. More modestly, I just want to point out that such possibilities may indeed exist, by presenting some tentative framework; it is then a matter of discussion between physicists and mathematicians to judge whether this deserves to be investigated further, and if so, to improve it and leading by iteration to a kind of optimal version, hopefully in less time then it took to progress from Euclide’s to Hilbert’s vision of geometry.

In a nutshell, my proposition is to “complete quantum theory”: since its present form is a linear theory, it calls for completion by some non-linear space, just like Euclidean geometry calls for completion by projective spaces. This proposition is presented in Section 4. Before presenting it, some more preliminary remarks.

1.2. The universe of mathematics, and the mathematical universe. Tobias was not the first and not the last to put forward the idea that “physics is mathematics” (I remember him exposing this idea to me on a paper napkin in Göttingen): Roland Omnès discussed such kind of idea in his book “Converging Realities” [O], saying: I suggest the name “physism” for the philosophical proposal that considers the foundations of mathematics as belonging to the laws of nature. More recently, this idea has been advanced by Max Tegmark (\([\text{\text{Teg}}]\)), who calls it

\(^{1}\) Hyperlinks are in grey in the electronic version of this text.
the Mathematical Universe Hypothesis (MUH), that is: Our external physical reality is a mathematical structure, and (loc. cit., p. 357): The MUH implies that mathematical existence equals physical existence. As a mathematician, I feel quite happy with this, and I like to take it as a heuristic principle, that is, as a welcome source of inspiration. The MUH suggests that physicists and mathematicians approach the same thing from different sides: physicists may call it the “Mathematical Universe”, and mathematicians may call it the “Universe of Mathematics”. Seen from the mathematician’s side, the axioms of quantum theory are part of the universe of mathematics, and finding their “optimal” form is not so much a matter of expedience, but rather an intrinsic mathematical question, whose importance is comparable to the one of the foundations of geometry. Indeed, my feeling is that these two questions are much more deeply related to each other than visible at present.

1.3. Form and content. Mathematicians tend to focus on the formal structure of the universe, on structures and relations, whatever their “meaning” or “content” may be. As Hilbert put it once, referring to his “Grundlagen der Geometrie” ([Reid], p. 57): “One must be able to say at all times – instead of points, straight lines, and planes – tables, chairs, and beer mugs.” Von Neumann (following the Göttingen spirit) defined the fundamental notions of quantum theory, state and observable, in a purely formal way as rays in a Hilbert space (“table”), respectively as self-adjoint operator in a Hilbert space (“beer mug”). Thus, passing from the “classical” to the “quantum world” is often presented by the following schema:

<table>
<thead>
<tr>
<th>Classical</th>
<th>Quantum</th>
</tr>
</thead>
<tbody>
<tr>
<td>state</td>
<td>point (element of a set)</td>
</tr>
<tr>
<td>observable</td>
<td>(real) function on the point set</td>
</tr>
</tbody>
</table>

Another version of this schema, in terms of $C^*$-algebras, reads as follows:

<table>
<thead>
<tr>
<th>Classical</th>
<th>Quantum</th>
</tr>
</thead>
<tbody>
<tr>
<td>state</td>
<td>point</td>
</tr>
<tr>
<td>observable</td>
<td>(real) function</td>
</tr>
</tbody>
</table>

This pattern is clean and neat, and there exist many excellent textbooks unfolding it in detail, both from the point of view of mathematics and of physics. As already said above, it is not our aim to reproduce them.

1.4. Plan. The pattern presented above looks clean and neat, but it is unsatisfying if you want to understand the “structure of the mathematical universe” – fundamental notions are defined via a construction (“take a Hilbert space or a $C^*$-algebra, and do this and that...”), and not via intrinsic properties and relations. In Section 2, we develop this criticism in more detail, and then present ingredients that might permit to formulate other axiomatizations (main sections: 3.2 and 4), essentially equivalent to the Dirac-von Neumann axioms, but opening a window towards possible new developments, by indicating what structure could be omitted or altered when wishing to start a trip into “non Dirac-von Neumannian quantum mechanics”. The
present text deals with the “general language” of quantum theory, whose main vocabulary is “state” and “observable”. In the subsequent second part, I will try to include dynamics into this theory (unitary time evolution).

1.5. **Geometry of quantum theory – history.** Before starting the mathematical discussion, let me very briefly sketch the history of our topic, from my own (admittedly subjective) viewpoint. Reading letters and texts by von Neumann ([Re, V]), I have the impression that nobody shared the dissatisfaction with his schema more than he himself. One the one hand, together with Jordan and Wigner, he investigated the possibility of constructing quantum mechanics by using only the “algebra of self-adjoint operators” – which is not an *associative* algebra, but (as we say nowadays) a *Jordan algebra*, with the symmetrized product $a \circ b := \frac{ab+ba}{2}$. I have been interested myself in the mathematical theory of Jordan algebras for a long time, and much of what follows is motivated by this research.

On the other hand, von Neumann writes in a letter to Garrett Birkhoff, in 1935 ([Re], p. 59): *I would like to make a confession which may seem immoral: I do not believe absolutely in Hilbert space any more.* He then attacks, together with Birkhoff, his deep and beautiful work on the lattice theoretic approach, completed later by contributions of other outstanding mathematicians, and presented in lectures by George Mackey giving rise to the monograph [V]. This monumental work is a major step in understanding mathematical structures underlying quantum mechanics, and it answers in many respects the criticism that I shall formulate below (cf. in particular the long notes to Chapter IV in [V], and [L17]).

However, reading [V], one ends up with the impression that the effect of this huge work is only to justify exactly the Dirac-von Neumann axioms as given before: we gain the satisfaction that they can be deduced from more general and more abstract principles. But nothing more – there seems to be no “window” that could be opened, comparably to opening Euclidean geometry towards non-Euclidean ones. Possibly, this feeling guided another generation of theoretical physicists, Aerts and his school on the one hand ([A99], where the term “completed quantum mechanics” is used in a sense different from ours, and [A09]), and on the other, Kibble, followed by Ashtekar and Shilling, and by Cirelli, Gatti, and Manià, and others (cf. references in [Be08a, Be08b]), who instead of lattice theory used (infinite dimensional) differential geometry to investigate the geometry of the “state manifold”, the projective Hilbert space $\mathbb{P}(H)$. This so-called “delinearization program” has also influenced my own approach [Be08a, Be08b], on which the present text is based. As far as I see, all of these authors pleading for a “geometric approach” to quantum theory have common aims and motivations, clearly formulated in [CGM]: “The delinearization program, by itself, is not related in our opinion to attempts to construct a non-linear extension of QM with operators that act non-linearly on the Hilbert space $H$. The true aim of the delinearization program is to free the mathematical foundations of QM from any reference to linear structure and to linear operators. It appears very gratifying to be aware of how naturally geometric concepts describe the more relevant aspects of ordinary QM, suggesting that the geometric approach could be very useful also in solving open problems in Quantum Theories.”
2. FROM CLASSICAL TO QUANTUM

Without going too much into details, here is what I would like to say as a mathematician, or as a “geometrician”, about the basic pattern presented above.

2.1. The “classical side”. Classical geometry deals with sets, say $M$, carrying additional structure having a “geometric flavor” (such as: manifold, symplectic or Lorentzian structure, and so on). The most elementary actors are the points of $M$, $p \in M$, which we call also pure states. Note that there is no “distinguished” point in $M$, no “origin”. However, one may object that “points” often appear to be a fiction, since they have no extension at all; it would be more realistic to replace points by probability measures $\mu$, also called mixed states, on $M$. Then it would be a matter of convenience to describe the correct topological, or measure-theoretic properties that one likes to impose. Anyhow, points $p$ may be identified with the corresponding point-mass, or Dirac measure, $\delta_p$, and finite convex combinations of Dirac measures represent mixed states coming from a finite number of pure states.

An observable is a real-valued function $f : M \to \mathbb{R}$ (in presence of additional structure, usually assumed to be continuous, or measurable, or smooth, and so on). Denote by $F(M)$ or $F(M, \mathbb{R})$ your space of observables (say, for the moment, the space of all real valued functions); then this space carries a rich structure: it is a vector space, by pointwise addition and multiplication by scalars, and a commutative algebra, by pointwise multiplication of functions, and there is a partial order: we may speak of positive functions. Note that all of these structures simply come from the corresponding ones of real numbers $\mathbb{R}$, since everything is defined pointwise. You just loose two things: $\mathbb{R}$ is a field, but $F(M, \mathbb{R})$ is not (it’s just a (commutative) ring), and the order on $\mathbb{R}$ is total, but the one on $F(M, \mathbb{R})$ is not (it’s only partial).

Next, states and observables naturally are in duality with each other: an observable $f$ can be evaluated at a point $p$, just by taking the value $f(p)$. If we work with mixed states (measures $\mu$), the same holds: if you take the view of defining a measure $\mu$ as a certain linear form on $F(M, \mathbb{R})$, then the value is denoted by $\mu(f)$; if you use classical measure theory, you will rather write $\int_M f d\mu$, but in the end this amounts to the same. In this context, $f$ may be called a random variable, and the value $\mu(f)$ is its expectation value. When $\mu$ is a Dirac measure $\delta_p$, then this value is always “sharp” (there is no variance), but in general we have to use the language of probability theory, as usual, e.g., in classical statistical mechanics. On a conceptual level, already at this point a serious problem becomes visible: the “problem of infinities” – certain measures attribute to certain functions the value $\infty$ (which is not a real number), or no value at all.

This is the basic set-up; much more could be said, and according to what you focus on, your theory will take different shapes. For instance, noticing that evaluation at a pure state is an algebra morphism $F(M, \mathbb{R}) \to \mathbb{R}$, you will be interested in kernels of the point evaluations, which are certain ideals of the algebra; pursuing this (and replacing $\mathbb{R}$ by $\mathbb{C}$ or other fields), you are lead towards formalisms used in algebraic geometry. On the other hand, keeping to real numbers, and noticing that measures are positive linear forms on $F(M)$, you are lead to look at the vector space $S(M)$ of signed measures, which is a subspace of the dual vector space $F(M)^*$, and to
realize that the Dirac measures are extremal points of the convex cone of positive functionals. This leads to duality of topological vector spaces, order structures, and to functional and convex analysis. Both viewpoints are extremely important in modern mathematics.

2.2. The “quantum side”, and the “superposition principle”. Concerning the “quantum side”, Varadarajan opens his book [V] by the phrase: _As laid down by Dirac in his great classic [D], the principle of superposition of states is the fundamental concept on which quantum theory is to be erected._ It is not easy to find a clear explanation of what this principle means – Dirac himself writes (in [D] p. 15): _The superposition process is a kind of additive process and implies that states can in some way be added._ Transposed to the classical picture drawn above, this would mean that we could “add” two pure states (points), and the result would be another pure state (point): that is, the manifold $M$ would be something like a linear space, with “addition” map assigning to a pair of points a third one. Thus, the passage from classical to quantum would resemble a procedure imposing some additional structure on the pure state space, turning it into something similar to a linear (=vector) space. This is not too far from a valid formal definition – today, we say much shorter: _a pure state is a ray in a (complex) Hilbert space $H$, so that the set of pure states is nothing but the projective space $\mathbb{P}(H)$ associated to $H$. Indeed, elements of a projective space cannot simply be “added”, but projective spaces do bear certain relations with linear spaces, and “superposition” refers to reminiscence of this kind of linearity in quantum theory._ For instance, two different points in a projective space define a unique projective line joining them, which is the set of superpositions of these two points (but the parametrisation of this line is not unique). Summing up, the state manifold $M$ becomes, on the quantum side, not quite a flat, linear space, but something related, a (complex) projective space $\mathbb{P}(H)$. _Projective geometry_ thus becomes part of quantum theory. This observation has triggered the geometric approaches to quantum mechanics mentioned above (subsection 1.5).

This apparently clear geometric picture suffers a setback when we wish to extend it to _mixed states:_ as on the classical side, one can speak of “mixed states”, again defined as formal convex combinations of pure states. However, one now must take care not to confuse such a formal convex combination with the superposition defined by the same coefficients! It is not clear what kind of “geometric object” the set of these general states then is: it is not a projective space, but still one would expect it to remember somehow the “superposition principle”, that is, to be some kind of geometry sharing properties with projective geometries – some kind of “generalized projective geometry”. Indeed, here we are lead to intrinsically mathematical questions concerning the structure of the Universe of Mathematics – and related to my own research.\(^2\)

Back to the quantum side, let’s now discuss the _observables:_ in the basic scheme, observables are represented by _self-adjoint operators on the Hilbert space $H$ (in general, unbounded operators_ – but let us, for the moment, prescind from this).

\(^2\) I have called, in [Be02], “generalized projective geometries” the precursors of the “Jordan geometries” from [Be14]. The approach is quite different from the lattice theoretic one developed in [V], cf. subsections 1.5 and 5.1.
States and observables are related with each other by a kind of duality, which in contrast to the classical case is now of quantum probabilistic nature: instead of a sharp “value of the observable \( A \) in the pure state \( \psi \), we just can speak of its expectation value, which is the number given by the formula (where \( \langle u,v \rangle \) is the scalar product in \( H \))

\[
\langle A \rangle_\psi = \frac{\langle \psi, A\psi \rangle}{\langle \psi, \psi \rangle},
\]

or, more generally, of the probability distribution of the values, including the second, \( \langle A^2 \rangle_\psi - \langle A \rangle_\psi^2 \), and higher moments. Although the expression \( \langle A \rangle_\psi \) looks more complicated than the classical \( f(p) \), it still is additive in \( A \), so that observables, just as in the classical case, form a vector space, with the usual operator sum being the same as “pointwise sum”. However, the formula is not “multiplicative” (i.e., not compatible with the composition of operators). The sum of operators thus seems to be the clear analog of the sum of functions from the classical case, whereas a clear interpretation of the product gets lost. The formula for \( \langle A \rangle_\psi \) is apparently not “linear” in the variable \( \psi \); all the more it is remarkable that the operator \( A \) itself acts linearly on \( \psi \) – the (complex) linearity of \( A \) is the surprising feature of quantum theory, and indeed it is the mathematical core of the “principle of superposition”. Whereas on the classical side there is just one source of linearity, on the quantum side there seem to be two such sources, one on the level of observables, the other on the level of states, which somehow appear to be compatible with each other. The precise formulation of this compatibility condition is rather subtle – there are at least two ways to formulate it, corresponding to the two ways of presenting the classical-quantum scheme given above, but in either way there is no such thing as “superposition of mixed states” (only of pure ones).

The first way is by identifying a mixed state \( W \), formal convex combination of orthonormal pure states \( \psi_i \) weighted by scalars \( w_i \in [0,1] \) such that \( \sum_i w_i = 1 \), with the corresponding “density matrix”, the operator represented by the diagonal matrix given by the \( w_i \) with respect to the \( \psi_i \). Then the expectation value of \( A \) in the mixed state \( W \) is given by

\[
\langle A \rangle_W = \text{trace}(WA).
\]

This formula is linear in \( W \), and even bilinear in \( (A,W) \). However, because of the normalizations, the density matrices do not form a linear space, but just a convex set, so the term “bilinear” has to be taken with some care.

The second way of interpreting these things, also going back to von Neumann, is to forget the Hilbert space \( H \) and to express everything in terms of the algebra \( A \) of (say, bounded) operators on \( H \), and in a next step taking for \( A \) more general types of associative algebras. Technically, one usually requires that \( A \) be a \( C^* \)-algebra. Most importantly, this means that \( A \) carries an involution, the adjoint map \( a \mapsto a^* \), and the observables then are the fixed points of this involution (\( a^* = a \), self-adjoint elements). The notion of state now becomes a derived notion: a state is a normalized positive linear functional \( \mu : A \to \mathbb{C} \), that is, a \( \mathbb{C} \)-linear map such that \( \mu(a^*a) \) is real and non-negative, for all \( a \in A \), and \( \mu(1) = 1 \), where \( 1 \) is the unit element of \( A \). The states form a convex cone, and the pure states are the extremal elements of this
cone. In this picture, the analogy with the classical picture appears very clearly: the commutative algebra $F(M)$ is replaced by the (complex) non-commutative algebra $A$, states are in both cases certain positive linear functionals, so there is a natural duality with the observables (interpreted in the classical case as exact value, and in the quantum case as expectation value, which is not sharp in general, even if the state is pure). Summing up, the following table refines the schema given in the preceding section:

<table>
<thead>
<tr>
<th></th>
<th>classical</th>
<th>quantum: Hilbert space $H$</th>
<th>quantum: $C^*$-algebra $A$</th>
</tr>
</thead>
<tbody>
<tr>
<td>pure state</td>
<td>point $p \in M$</td>
<td>ray $[\psi]$ in $H$</td>
<td>extremal $\mu : A \to \mathbb{C}$ normed, positive</td>
</tr>
<tr>
<td>(mixed) state</td>
<td>probability measure $\mu$</td>
<td>density matrix $W$</td>
<td></td>
</tr>
<tr>
<td>observable</td>
<td>function $f : M \to \mathbb{R}$</td>
<td>self-adjoint operator $A$</td>
<td>$a \in A$ with $a^* = a$</td>
</tr>
<tr>
<td>duality</td>
<td>$f(p), \mu(f)$</td>
<td>$\langle A \rangle_\psi, \langle A \rangle_W$</td>
<td>$\mu(a)$</td>
</tr>
</tbody>
</table>

This axiomatic framework proved to be powerful and robust since over 80 years. Why put it in question?

2.3. Questions. Stepping back and looking at the axioms from the point of view of the “universe of mathematics”, reasons of dissatisfaction may be:

1. Key notions (state, observable) are not introduced as primitive objects, but are defined by a construction using something else (Hilbert space $H$, $C^*$-algebra $A$). In other words, basic objects used in these constructions (vectors $\psi \in H$, elements $a \in A$) appear to be auxiliary: they do not have a physical interpretation (only the rays $[\psi]$, resp. the elements with $a^* = a$, do).

2. The axioms are “ungeometric”: this seems unavoidable when we define objects by their construction, and not by properties and relations. For instance, the linear structure of $H$, resp. the bilinear product of $A$ are imposed by decree, and the “superposition principle” comes out of this construction in a fairly indirect way.

3. Imposing by decree a linear structure also implies postulating the existence of an origin ($A = 0$, resp. $a = 0$), and of a unit ($A = I$, the identity operator, resp. $a = 1$, the unit of the algebra). These very distinguished “observables” do not look “physical”, but rather seem to reflect some kind of “convention”. What is their true status?

These three items are interwoven with each other. In the following, I shall use Item (3) as “line of attack”.

2.3.1. On the classical side. In the classical schema, clearly the two observables $f = 0$ and $f = 1$ (constant functions) play a very special rôle, and one may doubt if they deserve to be called “observable”: the function $f = 0$ defines the “origin” and the function $f = 1$ defines the “unit” with respect to which all other functions, and hence all “measurements”, are expressed. On the physics side, this raises deep questions about choices of unit systems, about existence of absolute zeros, and whether these values and choices are constant all the time, whether they are “canonical” or peculiar to our particular place in the multiverse, and so on. Since
these questions touch foundational issues of physics, I think that on the maths side, too, we should take them seriously. Thus, the least to say is that these two functions don’t look like observables of the same sort as the others: their status seems to be different.

Related to this item, a mathematician may regret that the duality between observables and states is somewhat unperfect: there are constant functions, but no “constant states”; there are pure states, but the “pure observables” (functions taking value 1 at one point and 0 else) play no rôle. Of course, there are analytic reasons for this: the “pure observables” are not continuous and would be of measure zero – at least, in all “continuous models”. Likewise, the “yes-no questions” (indicator functions $1_A$) are not continuous, hence out of scope of continuous models: are they “observables”, or “states”? Taking continuous models to be the only game in town excludes from the outset to encompass discrete models (see below, 2.4).

2.3.2. **On the quantum side.** The “scaling problem”, or “problem of the linear structure”, gets more involved on the quantum side than in the classical case, because of the double origin of its linear structure, see above. The “observables” 0 and 1 (the identity operator) play a very distinguished role also on the quantum side: again, they seem not to be “observables like the others”. In quantum logic, they represent truth values “always false” and “always true” – which clearly is a rather particular status, hardly comparable with observables like position or momentum.

2.4. **Continuous versus discrete.** Although this is not the main topic of the present work, let me say some words on this item, questioning the classical pattern. The question whether the universe should be seen as a “continuum”, or as a “granular (discrete) structure” is fundamental for choosing our mathematical model. The discussion already lasts over 2000 years, and sometimes one opinion prevailed, and sometimes the other: the hypothesis of a “granular” structure is attributed to Democritus; in “classical physics”, the universe appeared to be continuous; nowadays, in the “quantum era”, it appears to be discrete again, and maybe tomorrow, quantum-continuous... For a mathematician, the lesson to be drawn is that we should be ready to offer good models for both issues. But in practice we only have good models for the continuum model, and not for the discrete case (with its “worst case”: the one of a finite set). Indeed, in the continuous case, we have the whole of classical differential calculus at our disposition, which, combined with sophisticated functional analytic methods gives formidable strength to the approach of analyzing a “space” $M$ via function algebras and measure theory on $M$. Its power comes from duality, here: duality between function spaces and spaces of measures or distributions. In the finite case, the distinction between functions and measures becomes more or less a fiction: both $F(M)$ and $S(M)$ then are the same as $\mathbb{R}[M]$, the vector space with base indexed by elements of $M$, which carries a canonical scalar product

$$\langle f, g \rangle = \sum_{p \in M} f_p g_p,$$

so the same gadget may be considered as “function”, or as “signed measure”, whatever you prefer. This gives a purely algebraic model, which of course reflects nothing of the kind of properties of “infinitesimal calculus”. In the general discrete case, similar remarks hold. Thus, to get a sufficiently rich theory, we should ask: **is there some way to implement infinitesimal calculus in the discrete...**
case? I consider this to be a very interesting question – already in the realm of purely “classical” mathematics! Giving a serious answer would take too much place here. I have been doing research in this domain for a certain time – see [Be17a] for an overview; let me just say here that positive answers to this question do exist, and I believe they are relevant for the structure of the Universe of Mathematics.

3. The idea of “completion”

We explain the basic idea of how to complete a linear theory. In this chapter, we complete the space of classical observables, and in the next, the one of quantum observables and states.

3.1. Declaration of basic principles. In my own research, I regularly return to the “problem of choice of units and origins” (item (3) mentioned above). Any “geometric” theory starting by attributing a very special role to one, or two, or $n$, points, is flawed right from the beginning. I believe that somewhere in the constitution of the Universe of Mathematics it is written: All points are created equal. This axiom I hold for self-evident. It has corollaries, and to secure and formulate them it may be necessary to elect governments deriving their just powers from the consent of the governed. For instance, students learn in linear algebra courses that one may need to fix bases in order to define matrices; but we should be able to change bases (just like governments), or even to get entirely rid of them. Next, sometimes one wishes to get rid also of the origin of a vector space: this gives an affine space, which is nothing but a “vector space with forgotten origin”.$^3$ Likewise, in more abstract situations one may wish to get rid of “neutral elements” or “zero points” or “unit elements”: this is exactly the approach I advocate also for the “geometry of quantum mechanics”. For symmetry reasons, in quantum mechanics this strategy must be applied both to observables and states, featuring the duality between them. In this respect, duality in quantum mechanics appears to be more perfect than duality in the classical setting: the dual parts are of the same nature, they appear to be “self-dual”.

This proposition could be qualified “conservative”: no fancy new gadgets are introduced, but rather we renovate classical, if not old-fashioned, furnitures like affine spaces. More specifically, we will follow an old route pointed by the observation that the geometry of an affine space inevitably calls to be “completed” by adding “points at infinity”, the so-called “horizon”, to obtain some kind of more symmetric space, the projective space. My proposition is to place axiomatic quantum theory in the framework of a geometric space that “completes” an associative ($C^*$-) algebra, or a Jordan algebra, exactly like a projective space completes a usual vector of affine space. A good deal of my mathematical work has been devoted to such questions (see [Be00, Be02, BeKi, Be08a, Be08b, Be14, Be17b]), and to summarize, I can say that there is no mathematical obstruction to achieve this: such geometries do exist, and moreover, nothing is lost by the procedure of “completing”. What is gained?

- On the maths side, a more homogeneous and more symmetric picture, allowing

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$^3$ the idea is simple, but teaching it to students is not – the interested reader may look here for some remarks on this...
to look behind the horizon; on the physics side – I don’t know; but by comparison with classical geometry, and by the philosophy of the MUH, one may speculate that the gain could be non-zero. Future may tell.

3.2. The classical side revisited. Let’s start again by looking at the classical pattern, and by investigating more closely the special rôle of elements such as 1, 0 and ∞. First, the element 1 ∈ ℜ defines the “canonical” basis in ℜ with respect to its “canonical” origin 0 ∈ ℜ. Let’s forget the basis and look at ℜ just as an abstract one-dimensional vector space: there is a distinguished 0, but no distinguished 1. Then the dual space is also a one dimensional vector space, but both spaces should be distinguished from each other: to remind this, let us write f, v

we also know how to deal with brackets in such iterated products.

Summing up, forgetting the element 1 emphasizes the rôle of duality: it forces us to distinguish ℜ and ℜ′. On the level of functions, we now have so speak of “original functions” f : M → ℜ, and of “dual functions” φ : M → ℜ′, giving rise by paring to a function ⟨φ, f⟩ : M → ℜ. For the moment, let us think of both kinds of functions as “observables”. However, since obviously F(M, ℜ) and F(M, ℜ′) are sort of dual to each other (indeed, injecting a bit of language introduced in Appendix Б, the pair (F(M, ℜ), F(M, ℜ′)) is an archetypical example of an associative pair and of a Jordan pair), we will have to ask ourselves if it wouldn’t be more appropriate to identify one of the two spaces rather with some kind of “space of states”.

In a second step, let us forget both elements 0 and 1: this means to consider ℜ just as an affine line, that is, the one-dimensional affine space, ℜaf. By picking up any two distinct points a, b ∈ ℜaf, we can identify ℜaf with ℜ such that a corresponds to 0 and b to 1: namely, r ∈ ℜ corresponds to the “barycenter”

c = (1 − r)a + rb ∈ ℜaf. \tag{3.1}

Conversely, given a triple (a, b, c) of points in ℜaf, we recover r as (division) ratio

r = \frac{c − a}{b − a} = : R(c, b, a) ∈ ℜ, \tag{3.2}

and c = R(c, 1, 0). The ratio is an invariant of affine geometry. In the same way, our space of observables F(M, ℜ) is turned into an affine space F(M, ℜaf) by forgetting the functions 0 and 1. Relations (3.1) and (3.2) remain valid, pointwise: any two functions f0, f1 : M → ℜaf, taking different values at each point, can take the roles of “origin” and “unit”. Instead of f ∈ F(M, ℜ), we now consider the triple F = (f, f1, f0) as “observable”. The number describing the observable F in the pure state p is the ratio R(f(p), f1(p), f0(p)), that is, we define the “value of F at p” by

F(p) := R(f(p), f1(p), f0(p)) = \frac{f(p) − f0(p)}{f1(p) − f0(p)}. \tag{3.3}

We discover, then, that the inverse r−1 of a non-zero element r ∈ ℜ belongs to ℜ′: it is the unique element such that ⟨r, r−1⟩ = 1.
This formulation ensures that, even if values and choices of units and origins throughout the multiverse are uncommitted, the *mathematical form of laws* has a common description (as long as ratios are accepted as physical meaningful – which is possibly the oldest idea of exact science).\(^5\)

In the universe of mathematics, these first two steps force us to make a third one: the duality aspect from the first step has been lost in the second, and only by going to the *projective line* we can harvest the benefits of both steps together – indeed the *duality principle of projective geometry* is one of the highlights of classical geometry, and it makes projective geometry clearly superior to affine geometry. So, instead of \(F(M, \mathbb{R}^\text{af})\), let us look at the space \(F(M, \mathbb{P}^1)\) of functions \(f : M \to \mathbb{P}^1\) with values in the *real projective line* \(\mathbb{P}^1\). For the present purposes, it will be sufficient to define \(\mathbb{P}^1\) simply as the “one-point compactification” \(\mathbb{R} \cup \{\infty\}\) of \(\mathbb{R}\), by adding a single “point at infinity” (topologically, \(\mathbb{P}^1\) is a circle, but for the moment we are not interested in topology). We will use two basic facts about projective geometry (see Appendix C for some mathematical explanations):

1. removing an arbitrary hyperplane \(H\) from a projective space \(X\), an affine space \(X \setminus H\) remains (cf. Theorem C.1); in our case: removing an arbitrary point \(a\) from \(\mathbb{P}^1\), an affine line over \(\mathbb{R}\) remains, denote it by \(U_a = \mathbb{P}^1 \setminus \{a\}\);
2. picking up two different points \(a, b\) in the projective line, \(b\) may serve as origin in the affine space \(U_a\), and \(b\) may serve as origin in \(U_a\): thus we have two (one-dimensional) vector spaces \((U_a, U_b)\). Now, *these two vector spaces are dual to each other* (see subsection D.2). In other words, for any such choice of \((a, b)\), the pair \((U_a, U_b)\) is a “model” for \((\mathbb{R}, \mathbb{R}')\).

Both (1) and (2) can be used to associate to a quadruple \((a, b, c, d)\) of elements of \(\mathbb{P}^1\) a scalar in \(\mathbb{R}\): the ratio of \((a, b, c)\) in \(U_d\), and the duality \(\langle a, b \rangle\) in \((U_c, U_d)\). It is remarkable, then, that both procedures give the *same* number, namely the famous *cross-ratio* of the four values (see Appendix D):

\[
\text{CR}(a, b; c, d) = \frac{(c - a)(d - b)}{(c - b)(d - a)} = \frac{c - a}{c - b} : \frac{d - a}{d - b} = \frac{R(a, b, c)}{R(a, b, d)}.
\]

(To memorize notation: of the six possible differences, only \(a - b\) and \(c - d\) do not appear; the semicolon reminds this.) The cross-ratio is a rich and subtle projective invariant. It contains all information given by preceding constructions since

\[
R(a, b, c) = \text{CR}(a, b; c, \infty), \quad \frac{a}{b} = \text{CR}(a, b; 0, \infty), \quad a = \text{CR}(a, 1; 0, \infty).
\]

Again, the construction can be carried out pointwise: by (1), a completely arbitrary function \(h\), or \(f_\infty : M \to \mathbb{P}^1\) can serve as “horizon function”, or “infinity function”: the set of all functions never taking the same values as \(h\) forms an affine space, another copy of our \(F(M, \mathbb{P}^1)\). When \(h\) is the function \(f(x) = \infty\) for all \(x \in M\),

---

\(^5\) One may object that \(f = 0\) is distinguished by being a *constant* function, whereas \(f_0\) will in general not be constant. But my point is precisely that this distinction rather reflects a convention, and not a “fact of nature”: to take account of this, any choice of \(f_0\) also defines a modified *action of the group of bijections*, or of diffeomorphisms or whatever, on functions, such that \(f_0\) becomes invariant, i.e., “constant”, under this action. More generally, the notion of “constant section” of a vector bundle is not absolute, but depends on additional structure, such as, e.g., affine connections.
then we get back our old “standard realization”; but now we have also infinitely many other choices. Picking up three arbitrary functions, denoted by \( f_0, f_1, f_\infty : M \to \mathbb{RP}^1 \), subject to the condition that at any point they take pairwise different values, we use them as “reference triple”: \( f_1(p) \) as unit and \( f_0(p) \) as origin in the affine space \( U_{f_\infty(p)} \); thus given any function \( f : M \to \mathbb{RP}^1 \), we may define its “value at \( p \)” as the well-defined real number with respect to this reference triple, given by the pointwise cross-ratio of the quadruple \( F = (f, f_1, f_0, f_\infty) \):

\[
\langle F, p \rangle := CR\left((f(p), f_1(p); f_0(p), f_\infty(p))\right) = \frac{f(p) - f_0(p)}{f_1(p) - f_0(p)} : \frac{f_\infty(p) - f(p)}{f_\infty(p) - f_1(p)}.
\]

But, as said above, the same formula also realizes item (2)! It represents the bilinear pairing \( F(M, \mathbb{R}) \times F(M, \mathbb{R})' \to F(M, \mathbb{R}) \) when \( (f_0, f_\infty) \) is fixed, and in this context rather should be read

\[
(f, g) \mapsto CR\left((f(p), g(p); f_0(p), f_\infty(p))\right) = \frac{f(p) - f_0(p)}{g(p) - f_0(p)} : \frac{f_\infty(p) - f(p)}{f_\infty(p) - g(p)}
\]

Thus the pair \( (g, f_\infty) \) represents an object “dual” to \( (f, f_0) \), where the duality is given by the functional “pointwise cross-ratio”. This suggests a shift in the understanding of the notion of “state”: a state rather is a pair of dual functions, and an observable a pair of “original” functions, and this suggests to write the whole gadget \( F = (f, g; f_0, f_\infty) \) (call it “obstate”) rather as a matrix

\[
F = \begin{pmatrix} f & f_0 \\ g & f_\infty \end{pmatrix}.
\]

Its first row represents the “observable aspect”, and the second row the “state aspect”; the second column represents the “reference system aspect”, and the first column its “objective aspect”. The cross-ratio is invariant under exchange of rows, or of columns, and exchanging \( f \) and \( g \) (or \( f_0 \) and \( f_\infty \)) yields the inverse value. This shift of understanding furnishes a robust concept of “duality” and of “self-duality”, and it allows to separate this from the thorny problem of extracting a scalar valued pairing (using traces, integrals, measures – see Appendix E).

Summing up, classical mathematics, and classical mechanics and other classical theories, could equally well be described by replacing real valued functions by quadruples of \( \mathbb{RP}^1 \)-valued functions, and by working with cross-ratios instead of values of single functions. Of course, this looks heavy and unnecessarily complicated. And indeed, so it is, as long as origins, units and infinities are considered to be fixed once and for all. In classical mathematics, this assumption may seem reasonable; but even then it might be interesting to pursue this idea since it opens new views on certain fundamental issues. With this perspective in mind, we mention that a further property of \( \mathbb{R} \) generalizes rather nicely to \( \mathbb{RP}^1 \): the order relation of \( \mathbb{R} \) gives rise to a cyclic order on \( \mathbb{RP}^1 \). Namely,

- on the linear space \( \mathbb{R} \), the order is given by a unary relation: \( 0 < x \),
- on the affine space \( \mathbb{RAf} \), it is given by a binary relation, \( x < y \), as usual,
• on the projective space $\mathbb{RP}^1$, it is given by a ternary relation: a triple $(a, b, c)$ is cyclically ordered if $a < b$ in the affine space $U_c$.\footnote{Put differently and more formally: the group $\text{PGL}^+(\mathbb{R}^2)$ has two open orbits in $(\mathbb{RP}^1)^3$: one of them is the set of cyclically ordered triples.} We then write $b \in ]a, c[$, thus defining intervals on $\mathbb{RP}^1$.

Again, for functions, things carry over pointwise: what we get is a partial cyclic order on the space of functions from $M$ to $\mathbb{RP}^1$. The set of positive functions is generalised by intervals of this partial cyclic order (see \cite{Be17b} for more on this). Cyclic order and cross-ratio are related with each other: $\text{CR}(a, b; c, d)$ is negative iff $c$ lies in $]a, b[$ and $d$ in $]b, a[$, or vice versa, i.e., if the pair $(c, d)$ “separates” $(a, b)$. Some geometers (e.g., Coxeter) chose this separation relation as belonging to the structures appearing in axiomatic foundations of geometry.

3.3. From real to complex. Since quantum mechanics requires complex Hilbert spaces, and complex $*$-algebras, we may in a first step replace real functions from the classical picture by complex functions, $f : M \to \mathbb{C}$. One may agree that this is just a “trick”, since in the end the observables shall be real-valued. Everything said in the preceding section goes through (except the cyclic order, of course): it suffices to replace the real projective line by the complex projective line, $\mathbb{CP}^1 = \mathbb{C} \cup \{\infty\}$ (which now topologically is a 2-sphere, the Riemann sphere). The cross-ratio is defined in the same way, and it is invariant under complex conjugation. It follows that an “obstate” $F = (f, g; f_0, f_\infty)$ is real if, and only if, its cross-ratio is real. Now, it’s a classical fact that $\text{CR}(a, b; c, d)$ is real if, and only if, the four points $a, b, c, d$ lie on a generalised circle, that is, either lie on a circle in $\mathbb{C}$, or on a real affine line. Thus we have two possibilities to define “real obstates”:

1. a quadruple of $\mathbb{RP}^1$-valued functions, as in the preceding subsection,
2. a quadruple of $\mathbb{CP}^1$-valued functions such that, at every point $p \in M$, the four values lie on a generalised circle.

Let’s call an obstate “real” in the first sense, and “real-like” in the second one.

3.4. Antipode mapping. The cross-ratio is invariant under the full projective group: it is a projective invariant. On the other hand, the dynamics of quantum mechanics is governed by the unitary group, which is much smaller than the projective group. Thus at some point quantum mechanics requires to plug in some additional structure. For instance, we may fix a scalar product on $\mathbb{R}^2$, or on $\mathbb{C}^2$, say the standard scalar product $\langle x, y \rangle = x_1y_1 + x_2y_2$, and consider the induced polarity on $\mathbb{CP}^1$, that is, the orthocomplement map (where $J$ is the matrix given by (C.12))

$$\alpha : \mathbb{CP}^1 \to \mathbb{CP}^1, \quad [z] = \left[\begin{pmatrix} z_1 \\ z_2 \end{pmatrix}\right] \mapsto [z^\perp] = [Jz] = \left[\begin{pmatrix} -z_2 \\ z_1 \end{pmatrix}\right].$$

(3.9)

In the usual chart of $\mathbb{CP}^1$, this map is given by $z \mapsto -z^{-1}$; but if we identify $\mathbb{CP}^1$ with the Riemann sphere $S^2$, then $\alpha$ is rather represented by the antipode map sending a point of the sphere to its opposite, or antipode point. The projective maps commuting with $\alpha$ are exactly those coming from the projective unitary group.
PU(2). Thus we can reduce the projective invariant cross-ratio to a two-point invariant of \( \text{PU}(2) \): in formula (D.1), let \( a = \alpha(y), \ b = \alpha(x) \), then

\[
P(x, y) := \text{CR}(x, y; \alpha(y), \alpha(x)) = \frac{\langle x, y \rangle \cdot \langle y, x \rangle}{\langle x, x \rangle \cdot \langle y, y \rangle} = \cos^2(\phi(x, y)), \tag{3.10}
\]

where \( \phi \) measures the angle between the vectors \( x, y \in \mathbb{C}^2 \). Of course, the same holds for \( \mathbb{C} \) replaced by \( \mathbb{R} \).\(^7\) Again, applying everything pointwise, these definitions carry over to function spaces: we can define \( \alpha(f) \) and \( P(f, g) \) for functions.

The antipode mapping \( \alpha \) on \( \mathbb{C}P^1 \) is antiholomorphic, just like the complex conjugation \( \tau(z) = \overline{z} \) of \( \mathbb{C}P^1 \), whose fixed point set is \( \mathbb{R}P^1 \). Since \( \alpha \) and \( \tau \) commute, the composition \( \beta := \alpha \circ \tau \) is the map induced by the matrix \( J \), given in the usual chart by \( z \mapsto -z^{-1} \), which is a holomorphic map of order 2. It has precisely two fixed points: \( i \) and \( -i \). When picturing \( \mathbb{R}P^1 \) as equator of the sphere \( \mathbb{C}P^1 \), these two fixed points shall be pictured as north and south pole, and the points 0 and \( \infty \) on the equator could be called east and west pole, and 1 and \( -1 \) front and back pole. The “usual chart” is stereographic projection from the west pole onto the tangent plane of the sphere at the east pole (Figure 1). The four transformations \( \{ \tau, \alpha, \beta, \text{id} \} \) form an abelian group (a Klein four group) acting on \( \mathbb{C}P^1 \).

**Figure 1.** The Riemann sphere \( \mathbb{C}P^1 \) with six poles.

4. **Completion of Quantum Theory**

This is the main section: we are going to explain the general setting “completing” usual, linear quantum theory. By “usual” formulation we mean the one in terms of a \( C^*\)-algebra \( \mathbb{A} \) (but we will not use all properties of a \( C^*\)-algebra, only those which define a \( P^*\)-algebra, see Appendix A). For some mathematical constructions and definitions we shall refer to the appendices. The algebra \( \mathbb{A} \), respectively, its

\(^7\) The formula for \( P(x, y) \) defines a transition probability, in the sense of [L98], p.80, or [L17], p. 31. Using (D.1), the same procedure can be applied to any projective space.
real subspace $\text{Herm}(A)$, are completed by the following “geometric spaces”

\[
\mathcal{G} := \text{Gras}(A^2) = \{x \subset A^2 \mid x \text{ (right) submodule}, x \neq 0, x \neq A^2\},
\]

(4.1)

\[
\mathcal{S} := \text{Gras}_A(A^2) = \{x \in \text{Gras}(A^2) \mid x \cong A\},
\]

(4.2)

\[
\mathcal{S}' := \text{Gras}^A(A^2) = \{x \in \text{Gras}(A^2) \mid A^2/x \cong A\},
\]

(4.3)

\[
\mathcal{S} := \mathbb{A}^1 := \text{Gras}_A^A(A^2) := \mathcal{S} \cap \mathcal{S}',
\]

(4.4)

\[
\mathcal{R} := \{x \in \mathcal{S} \mid x_\perp = Jx\} (= \text{Lagrangian variety of } \omega(u,v) = \langle Ju,v \rangle),
\]

(4.5)

\[
\mathcal{R}' := \{x \in \mathcal{R} \mid A^2 = x \oplus x_\perp\} = \{x \in \mathcal{R} \mid A^2 = x \oplus Jx\}
\]

(4.6)

where $J$ is given by (C.12) and $\perp$ the orthogonal complement with respect to the usual “scalar product” $\langle u,v \rangle = \sum_i u_i^* v_i$ on $A^2$. In the classical case $A = \mathbb{C}$, the spaces $\mathcal{G}, \mathcal{S}, \mathcal{S}'$ and $\mathcal{S}$ all coincide with the Riemann sphere $S^2$, and the spaces $\mathcal{R}$ and $\mathcal{R}'$ both coincide with the unit circle (equator) $S^1$. When $A$ is infinite dimensional, the inclusions

\[
\mathcal{R}' \subset \mathcal{R} \subset \mathcal{S} \subset \mathcal{S}' \subset \mathcal{G}
\]

are in general all strict, and one may consider them as inclusions of “universes”, sitting inside each other like Matryoshka dolls. The “base points” $0 = [(1,0)]$ and $\infty = [(0,1)]$ belong to all of them, and so does the “affine part” $\{[(1,a)] \mid a \in \text{Herm}(A)\} \cong \text{Herm}(A)$, which represents the “algebra of (bounded) observables” from “usual” quantum mechanics, so that the nested sequence arises by adding more and more “points at infinity” to the “usual” space. Having fixed the pair $(0, \infty)$, the “natural chart” $A \subset \mathcal{S}$ generalizes stereographic projection; but the “set at infinity” (the part of $\mathcal{S}$ not covered by $A$) now is in general quite a big set: it contains a distinguished point $\infty$, but also many others. If $A$ is finite-dimensional, then $A$ will always be dense in $\mathcal{S}$, but if $A$ is infinite-dimensional, then this need not be the case. We will describe two versions of the setting:

(1) a weak, or projective setting (subsection 4.1): the structure is just given by $(\mathcal{S}, \tau)$; its symmetry group is big (the whole projective group of $\mathcal{R}$). This setting suffices to define expectation values (first moments),

(2) a strong, or unitary setting (subsection 4.2): the structure is given by $(\mathcal{S}, \tau, \alpha)$; its symmetry group is smaller (essentially, unitary), and it permits to recast the whole of quantum theory, including higher moments.

As said in the introduction, this text is still preliminary and experimental: at present, it is not entirely clear to me which parts of quantum theory really belong to the “weak setting”, and which to the “strong setting”, or maybe to some intermediate setting yet to be defined.

4.1. The weak (projective) setting. This setting is given by the data $(\mathcal{S}, \mathcal{R}, \tau)$. Its symmetry group, generalizing the usual projective group $\mathbb{P}(\text{SL}_2(\mathbb{R}))$, is described in the appendix, equation (C.14). In this setting, it is appropriate to distinguish the projective line from its dual:

4.1.1. Duality and self-duality. The projective line over $A$, as well as the Hermitian projective line, are self-dual: they agree with their dual projective line, $\mathcal{S} = \mathcal{S}'$, $\mathcal{R} = \mathcal{R}'$, see Appendix C. However, both for mathematical and for philosophical
reasons, we shall separate, whenever possible, two copies \((\mathcal{R}, \mathcal{R}')\), resp. \((\mathcal{S}, \mathcal{S}')\), considered to be “dual to each other”. In more technical terms, this means that we treat, whenever possible, the associative algebra \(A\) as an associative pair \((A^+, A^-) = (A, A)\) (and the Jordan algebra \(\text{Herm}(A)\) as a Jordan pair); see Appendix B. Still put differently, we try, as long as possible, to avoid using the unit element \(1\) of \(A\).

4.1.2. Basic terminology: complete obstates. We use the term “obstate” for “observable-state” to denote an entity incorporating “observables” and “states”.

Definition 4.1. A complete obstate, or obstate for short, is a quadruple

\[
\mathbf{O} := (A, W; A_0, W_{\infty})
\]

such that: \(A, A_0 \in \mathcal{R}\), and \(W, W_{\infty} \in \mathcal{R}'\), and \(A_0 \top W_{\infty}\), \(A_0 \top W\), \(ATW_{\infty}\) (where \(\top\) means “transversal”, see Appendix C.3). The pair \((A, W)\) is called the object part of the obstate, or objective obstate, and the pair \((A_0, W_{\infty})\) is called the reference part, or reference obstate. The pair \((A; A_0)\) is called complete observable, and the pair \((W; W_{\infty})\) complete state. This terminology is summarized by the “obstate matrix”:

<table>
<thead>
<tr>
<th>complete observable</th>
<th>object part</th>
<th>reference part</th>
</tr>
</thead>
<tbody>
<tr>
<td>(A)</td>
<td>(A_0)</td>
<td></td>
</tr>
<tr>
<td>(W)</td>
<td>(W_{\infty})</td>
<td></td>
</tr>
</tbody>
</table>

In “usual” quantum physics, the reference part is fixed once and for all, and then is denoted just by \((0, \infty)\). The idea of “complete quantum physics” should be that complete obstates are the intrinsic objects to be studied. As long as the reference part is fixed, no deviation from usual quantum physics should appear, that is, we postulate the mathematical “conservation rule”: for a fixed reference part \((A_0, W_{\infty})\), the rules of complete quantum theory shall reduce to the rules of usual quantum theory. (To ensure this, we have included the transversality assumption \(A_0 \top W_{\infty}\) in the definition.) Since the unit element \(1\) of the algebra does not appear in the reference part, this “weak setting” comprises all aspects of the usual theory that do not depend on, or do not require the choice of, a unit element.

We do not make any claims about “interpretations” or “reality” corresponding to a possible change of reference parts. Indeed, the reader may safely assume that \((A_0, W_{\infty})\) is fixed once and for all. Maybe s-he will find, later on, that it is much more convenient to assume that this is not the case, and that this is in much better keeping with some of the current “interpretations” of usual quantum theory. This, possibly, could open the hypothetical window towards “non Dirac-von Neumannian quantum theory” – which should never be in contradiction with usual quantum physics. (There might be apparent contradictions due to unclear terminology: before projective geometry was invented, a phrase like “two parallel lines intersect at infinity” sounded contradictory.)

4.1.3. Pure states; rank. Under our “conservation rule”, for a fixed reference system \((0, \infty)\), states \(W\) shall correspond to density matrices, and thus pure states shall correspond to density matrices of rank one. The aim of the paper \([\text{BeL}]\) is to give a “geometric interpretation” of such concepts, in the general context of Jordan geometries. Let me try to summarize the main ideas in the present, more special
context: given two elements $x, y \in \mathcal{R}$ that belong to some affine part $U_a$ of $\mathcal{R}$, the locus of the real affine line $[x, y]_a = \{ tx + (1 - t)y \mid t \in \mathbb{R} \}$ does in general heavily depend on the choice of $a$. But for certain choices of the pair $(x, y)$, this locus does not depend on the choice of the affinization $a$:

**Definition 4.2.** A pair $(x, y) \in \mathcal{R}^2$ is said of rank 1, or of arithmetic distance 1, if $x \neq y$ and for all $a, b \in \mathcal{R}$ with $x, y \in U_a \cap U_b$: 

$$ [x, y]_a \cap U_b = [x, y]_b \cap U_a. $$

Then $\overline{[x, y]} := [x, y]_a \cup \{ \infty \}$ is a copy of the real projective line $\mathbb{RP}^1$ in $\mathcal{R}$ that depends only on $x$ and $y$, called an intrinsic line in $\mathcal{R}$. A complete state $(W; W_\infty)$ is said pure (and then we shall often write $(\psi; \psi_\infty)$, following a venerable tradition) if the pair $(W, W_\infty)$ is of rank 1.

For instance, in a projective space $\mathbb{RP}^n$, every pair $(x, y)$ with $x \neq y$ is of rank 1: every pair of distinct points defines a unique intrinsic projective line joining them. In sharp contrast, for higher rank geometries, such as $\mathcal{R}$, such lines can only follow very special directions (these directions lie on the extreme boundary of the “light cones” that define the “generalized conformal structure” of $\mathcal{R}$, see below, 4.1.7). Algebraically, saying that $(0, x)$ is of rank 1 corresponds to saying that $x$ is von Neumann regular, or that we can find $\infty$ such that $x$ becomes an idempotent (see Def. B.3), or yet that $x$ generates a minimal inner ideal. Likewise, in [BeL] it is explained that higher rank is related to inner ideals that need not be minimal. Their geometric counterpart has been christiandel intrinsic subspace.

**Remark 4.1.** The term arithmetic distance is due to L.-K. Hua, who studied it for all series of finite-dimensional matrix geometries. There are interesting links between the arithmetic distance between $(x, y)$ and algebraic invariants of the torsors $U_{xy}$. For instance, when $(x, y)$ is of rank 1, then $U_{xy}$ is a solvable group with derived series having one non-trivial term.

4.1.4. **Expectation value of an obstate.** Assume $(A, W; A_0, W_\infty)$ is a complete obstate. We have to extract a real number from these data, which for the fixed reference system $(A_0, W_\infty)$ shall coincide with the one given by (2.1) or (2.2). Imperatively, this scalar has to be given by a scalar valued cross-ratio:

**Definition 4.3.** The expectation value of the complete obstate $(A, W; A_0, W_\infty)$ is

$$ \langle A, W; A_0, W_\infty \rangle := \text{trace}(K_{A_0, W_\infty}(A, W)) $$

$$ = \text{trace}(\text{CR}(A, W; A_0, W_\infty)) $$

where $K$ and the generalized cross-ratio $\text{CR}$ are defined by eqn. (D.5) – (D.7), and trace denotes a trace on $\mathcal{A}$ in the sense of Def. E.3.

This definition is natural, in the sense that it is invariant under the automorphism group Aut($\mathcal{S}, \tau$). However, mind:

1. **Traces** of linear operators always exist in finite dimension over a field, but not always in very general situations. Indeed, this is not a “quantum” problem, but already appears in the “classical case” (section 3.2): associating a scalar to a pair (function, dual function) is some kind of integration, and already classical integrals may lead to infinite values (see Appendix E).
(2) For the formula from Definition 4.3 to reduce to (2.2), in case \((A_0, W_\infty) = (0, \infty)\), we have to carefully distinguish a space from its dual space. If one misses that point, one would read the expression as \(\text{trace}(AW^{-1})\).

(3) We cannot define “second moments” in the same way, since the definition of an operator \(A^2\) depends on the choice of a unit element \(1\), which is not given in the present setting. As far as I see, it is not possible to define such higher moments in the present “weak setting”: one needs more, and more rigid, structure to define them, see below.

The second item is related to the normalization which, in the usual theory, is necessary to write formula (2.2); in our “intrinsic” formula in def. 4.3 no normalization is necessary (think of \(W_\infty\) as the “zero matrix”, which of course is not a density matrix itself, in the usual theory).

Let us re-interprete this construction in a more geometric way for pure states \((\psi, \psi_\infty)\): in this case, the intrinsic projective line \(L \cong \mathbb{RP}^1\) determined by the pure state contains already two distinguished elements, \(\psi\), and \(\psi_\infty\). The observable \((A_0, A)\) defines two other distinguished points \((a_0, a)\) on \(L\): namely, \(a\) is the unique point of \(L\) completing the affine line \(L \cap U_A\), and likewise for \(a_0\). Now, the expectation value is the (classical) cross-ratio of these four points on the line \(L \cong \mathbb{RP}^1\):

\[
\langle A, \psi; A_0, \psi_\infty \rangle = \text{CR}(a, \psi; a_0, \psi_\infty).
\]

This is the analog of (2.1). If \((A, A_0)\) happens to be already on \(L\) (so \(a = A, a_0 = A_0\)), then the measurement is “sharp”, but in general, this will not be the case, and there will be higher moments (cf. below). Note that for pure states we do not have to bother about problem (1) mentioned above, since traces exist for rank-one operators.

4.1.5. Axiomatic setting; superposition principle. As said above, in the usual setting, Dirac’s “superposition principle of quantum theory” corresponds to assuming that observables are operators acting linearly on a linear space, or that the “(Jordan) algebra of observables” carries a bilinear product. In our setting, this property can be translated into the form of geometric axioms (cf. [Be02, Be14]): it means that the geometry \((\mathcal{R}, \mathcal{R}')\) is an affine pair geometry – every element \(a \in \mathcal{R}\) defines an affine part \(U_a\) of \(\mathcal{R}'\), and vice versa, every \(w \in \mathcal{R}'\) defines an affine part \(U_w\) of \(\mathcal{R}\). In other words, the geometry is covered by “affine charts”, which are part of its structure. In an axiomatic “geometrically complete quantum theory”, this property should be part of the axioms. It then becomes a theorem (cf. [BeL]) that the intrinsic lines form, in turn, another geometry, that is, they also have a local linear structure, corresponding to the superposition principle. Thus a truly axiomatic presentation of “completed quantum theory” should be possible; but for pedagogical reasons it must be postponed.

4.1.6. Real versus complex. Expectation values shall be real, and not complex. This can be achieved by a purely real theory, and the setting presented so far does not (yet) really explain why complex numbers play such an important rôle in quantum theory, compared to the classical theory. Indeed, everything said so far makes sense more generally when \(\mathcal{R}\) is the Jordan geometry corresponding to an abstract ordered...
Jordan algebra, cf. [Be17b] (except that the definition of the generalized cross-ratio becomes more involved if no associative structure is around). As far as I see, the true role of the complex numbers appears more clearly in the “strong setting”. For the moment, we have the same two options for formulating the “complete” theory as mentioned in subsection 3.3, and so far both of them appear to be reasonable:

1. “real”: we work in the universe of the Hermitian projective line \( \mathcal{R} = \mathcal{R}' \); that is, all four components of \((A,A_0,W,W_\infty)\) shall belong to \( \mathcal{R} \);
2. “real-like”: we work in \( \mathcal{S} = \mathbb{A} \mathbb{P}^1 \), but we require that all four components of a complete obstate belong to a “generalised circle” (conjugate of the Hermitian projective line under the projective group \( \mathbb{P} \text{Gl}(2,\mathbb{A}) \)). Since expectation values of quadruples are invariant under the projective group, this still ensures that all expectation values are real.

4.1.7. Positivity: cyclic order. In the \( C^* \)-algebra setting, it is part of the definition of states that they are positive linear functionals. We have not included positivity in our definition of a complete state, since the precise formulation of such an assumption is related to the question of “interpretations” of the formalism. First of all, in the projective setting, the binary order relation generalizes to a ternary relation. As starting point, we use the binary partial order \( \leq \) on \( \text{Herm}(\mathbb{A}) \), which exists by definition of a \( P^* \)-algebra (def. A.6), and then define a partial order \( \leq_c \) on each affine part \( U_c \) (cf. [Be17b], Theorem 4.1), defining the ternary relation:

Definition 4.4. A complete state \((W;W_\infty)\) is called positive with respect to a reference part \((A_0,W_\infty)\), if the triple \((A_0,W,W_\infty)\) is cyclically ordered, that is, if \( A_0 \leq W \) in the ordered vector space \( U_{W_\infty} \). We say that \((A,A_0;W,W_\infty)\) is a cyclically ordered obstate if \((A,A_0,W_\infty)\) and \((A_0,W,W_\infty)\) are cyclically ordered triples. This implies that the expectation value \( \langle A,W;A_0,W_\infty \rangle \) is positive.

As explained in [Be17b], the intervals on \( \mathcal{R} \) define a kind of generalized conformal, or causal, structure, modelled on the positive cone of \( \text{Herm}(\mathbb{A}) \).

4.2. The strong (unitary) setting. Now we add the following datum to the “weak setting”: the standard scalar product on \( \mathbb{A}^2 \) defines an orthocomplementation map \( \alpha : \mathcal{S} \to \mathcal{S} \), \( x \mapsto x^\perp \) which is antiholomorphic and commutes with \( \tau \), so that the holomorphic map \( \beta := \alpha \circ \tau : \mathcal{S} \to \mathcal{S} \) is again of order 2. The data \((\mathcal{S},\tau,\alpha)\) define the “strong setting”. There are no “closed” formulae for \( \tau \) and \( \alpha \), but as in the classical case, the map \( \beta \) is induced by the matrix \( J = \begin{pmatrix} 0 & -1 \\ 1 & 0 \end{pmatrix} \), that is, \( \beta = [J] \), so that in the usual chart, for \( z \in \mathbb{A} \),

\[
\begin{align*}
\tau(z) &= z^*, \\
\beta(z) &= -z^{-1}, \\
\alpha(z) &= -(z^*)^{-1}. 
\end{align*}
\] (4.8)

4.2.1. North and south pole. The map \( J : \mathbb{A}^2 \to \mathbb{A}^2 \) is diagonalizable over \( \mathbb{A} \): it has two eigenvectors \((i,1)\) and \((-i,1)\) with eigenvalues \(i, -i\), so

\[
\begin{pmatrix}
C^{-1}JC
\end{pmatrix} = \begin{pmatrix}
i & 0 \\
0 & -i
\end{pmatrix} = \begin{pmatrix}
1 & 0 \\
0 & -1
\end{pmatrix}, \quad \text{where } C = \begin{pmatrix}i & -i \\
1 & 1 \end{pmatrix}
\] (4.9)

(the matrix \( C \) describes the Cayley transform, see below). Thus the map \( \beta = [J] : \mathcal{S} \to \mathcal{S} \) has precisely two fixed points, called north pole and south pole,

\[
N := [(i,1)], \quad S := [(-i,1)].
\] (4.10)
Since one eigenvalue is the negative of the other, the map $\beta$ acts by multiplication with $-1$ on the linear spaces $(S_N, S)$ and $(S_S, N)$, i.e., $\beta = (-1)_{N,S}$ (using notation (C.8)). Thus $\beta$, and hence also $\alpha$, can be recovered from $(N, S)$, and we see that the data $(S, \tau, \alpha)$ and $(S, \tau, N, S)$ are essentially equivalent.

4.2.2. The canonical $S^1$-action. Since the data $(N, S)$ are canonical, not only the reflection map $(-1)_{N,S}$ is canonical, but every map of the form $\lambda_{N,S}$ with $\lambda \in S^1$. Indeed, these maps commute with $\tau$ since $\tau \circ \lambda_{N,S} \circ \tau = \lambda_{\tau(N), \tau(S)} = (-\lambda)_{S,N} = \lambda_{N,S}$, hence preserve $\mathcal{R}$. Thus we get a canonical action

$$S^1 \times \mathcal{R} \to \mathcal{R}, \quad (\lambda, x) \mapsto \lambda_{N,S}(x). \quad (4.11)$$

In particular, $\beta$ has $i_{N,S}$ as a canonical square root: $i^2_{N,S} = (-1)_{N,S} = \beta$.

4.2.3. Automorphism groups: unitary structure. The strong setting is more rigid than the weak one, hence its automorphism group is smaller. Nevertheless, this group still has “big orbits”. Let’s explain this: The automorphism group of $S$ is $\mathbb{P}GL(2, A)$; the one of the weak setting is $\text{Aut}(S, \tau) = \{g \in \text{Aut}(S) \mid g \circ \tau = \tau \circ g\}$ (cf. Appendix C.4), and the one of the strong setting is

$$U := \text{Aut}(S, \tau, \alpha) = \text{Aut}(S, \tau) \cap \text{Aut}(S, \alpha)$$

$$= \text{Aut}(S, \alpha) \cap \text{Aut}(S, \beta) = \mathbb{P}U(2; A) \cap \text{Aut}(S, \beta)$$

$$= \mathbb{P}\{f = \begin{pmatrix} a & b \\ -b & -a \end{pmatrix} \mid a^*a + b^*b = 1, a^*b - b^*a = 0\}$$

$$= \mathbb{P}\{f = \begin{pmatrix} a & b \\ -b & -a \end{pmatrix} \mid a, b \in A, (a + ib) \in U(A)\},$$

given by unitary operators $f : A^2 \to A^2$ such that $fJ = Jf$. Via the Cayley transform, this group is isomorphic to $\mathbb{P}(U(A) \times U(A))$.

4.2.4. The real unitary universe $\mathcal{R}_{N,S}$, and strong obstates.

**Definition 4.5.** We call real unitary universe the subset of $\mathcal{R}$ given by all elements that are both transversal to $N$ and to $S$,

$$\mathcal{R}_{N,S} := \mathcal{R} \cap U_{N,S} = \{x \in \mathcal{R} \mid x \perp N, x \perp S\}. \quad (4.12)$$

With $\mathcal{R}'$ given by (4.6), we always have $\mathcal{R}_{N,S} \subset \mathcal{R}' \subset \mathcal{R}$. The real unitary universe is the space where strong obstates live:

**Definition 4.6.** A strong obstate is an obstate $(A, W; A_0, W_\infty)$ such that

1. $W_\infty$ and $A_0$ are antipodes of each other: $W_\infty = \alpha(A_0)$,
2. $A_0 \in \mathcal{R}_{N,S}$ [by Theorem 4.9 below, all 4 components then belong to $\mathcal{R}_{N,S}$].

4.2.5. On the structure of the real unitary inverse. Here are the most important results on the structure of $\mathcal{R}_{N,S}$. They are special cases of more general and abstract results from [BeKi2] (except for Theorem 4.9); we will give more computational and down-to-earth proofs (using the Cayley transform) in [Bexy]. The first result, contained in [BeKi2], says that $\mathcal{R}_{N,S}$ “is” the unitary group. This will be basic for our interpretation of unitary time evolution [Bexy].
Theorem 4.7. The real unitary universe \( \mathcal{R}_{N,S} \) carries a canonical torsor structure, that is, for any choice of origin \( a \in \mathcal{R}_{N,S} \) this set carries a group structure with unit element \( a \) and product \( xz = x \cdot_a z \), such that with respect to any other origin \( y \), the product is given by \( x \cdot_y z = xy^{-1}z \). Moreover, any of the groups thus obtained is isomorphic to the unitary group \( U(\mathcal{A}) = U(\mathcal{A}, \ast) \) (cf. Def. A.4).

Theorem 4.8. The automorphism group \( U \) acts transitively on the real unitary universe \( \mathcal{R}_{N,S} \). The stabilizer group of a point \( o \) is isomorphic to \( U(\mathcal{A}) \), so that as homogeneous space, \( \mathcal{R}_{N,S} = U \sim = (U(\mathcal{A}) \times U(\mathcal{A})) / U(\mathcal{A}) \).

Indeed, any torsor acts transitively on itself by left or right translations, and these always belong to the automorphism group ([BeKi2]).

Theorem 4.9. The real unitary universe contains affine parts defined by all of its points: for all \( a \in \mathcal{R}_{N,S} \), the set \( \mathcal{A}_a = \{ x \in S \mid x \oplus \alpha(a) = \mathcal{A}^2 \} \) carries the structure of an associative algebra, with zero vector \( a \) and unit element \( b := i_{N,S}(a) \). This algebra is isomorphic to the associative algebra \( \mathcal{A} \) with unit 1. Likewise, its real form \( \mathcal{A}_a^\tau \) is a Jordan algebra with unit element \( b \) and zero vector \( a \), isomorphic to \( \text{Herm}(\mathcal{A}) \).

Proof. Since the stabilizer group \( U(\mathcal{A}) \) acts (via conjugation) by automorphisms on the algebra \( \mathcal{A} \), we may transport the algebra structure from the algebra at the base point 0 to any other point of \( \mathcal{R}_{N,S} \), by transitivity. The pair \((0, \infty)\) is mapped to \((a, \alpha(a))\), the unit 1 is then mapped to a point \( b \). The only thing which is not
quite obvious is that then, necessarily, \( b = i_{N,S}(a) \). This, again, is proved using the Cayley transform, see [Bexy].

Remark 4.2. On may think of \( A_a \), or rather of \( \text{Herm}(A_a) \) as a “tangent algebra of the geometry at the point \( a \)” (see [Be14]). The product in the algebra with neutral element \( b \), and the unitary group law \( \cdot \) with unit \( b \), are dual to each other, in the sense of Cartan duality of symmetric spaces: the Jordan cone at \( b \) is kind of “non-compact dual” of the “compact-like” unitary group at \( b \). This is related to the topic of Jordan-Lie algebras, see 4.2.7.

The preceding results permit to reduce “strong completed quantum theory” to “business as usual”: since all algebras \( A_a \) are equivalent, we may (as long as \( a \) is considered to be fixed), “without loss of generality”, assume that \( a = 0 = [(0,1)] \) is the “usual base point. For instance:

4.2.6. Second and higher moments. Given a strong obstacle \((A,W; A_0, W_\infty)\), we may compute in the algebra \( A_a \) with \( a = A_0 \). Let \( AW \) be the product of \( A \) and \( W \) in this algebra, and \( L_X : A_a \to A_a \), \( Y \mapsto XY \) the operator of left multiplication by \( X \). Then the operator valued cross ratio and \( L_{AW} \) coincide:

\[
\text{CR}(A,W; A_0, W_\infty) = L_{AW},
\]

and hence also their traces: \( \langle A,W; A_0, W_\infty \rangle = \text{trace}(AW) \), so that expectation values are calculated in the algebra \( A_a \) in the usual way. Since \( A_a \) is an algebra with binary product (and not just an associative pair), we now can form also all expressions of the form \( A^k, k \in \mathbb{N} \), and in particular we can define as usual the second moment (variance) of the strong obstacle, by

\[
V(A,W; A_0, W_\infty) := \langle A^2 \rangle_W - \langle A \rangle_W^2 = \text{trace}(AWA) - (\text{trace}(AW))^2.
\]

If we assume that \( A \) is a \( C^* \)-algebra, we can also define the probability distribution on \( \mathbb{R} \) induced by the complete obstacle, via the spectral theorem, in the usual way (cf. e.g., [vN, L17]). In a similar way, all other properties and constructions can be carried over from \( A \) to \( A_a \).

4.2.7. Conceptual approach: geometry of Jordan-Lie algebras. Presenting things by simply transferring everything to “business as usual”, as phrased above, is not very conceptual, nor satisfying, but at least we see that a geometric, base-point-free setting for the geometry of quantum theory exists (and this is all I wanted to show at present). Possibly, a better understanding of what is going on here can only be achieved in connection with studying dynamics: the unitary (Schrödinger) evolution on the one hand (Part II [Bexy]), and, much more difficult, a mathematical analysis of the measurement process from a geometric viewpoint (Part III ?). Mathematically, as far as I see, the “strong setting” is the geometric counterpart of the algebraic structure of a Jordan-Lie algebra (cf. Appendix A, and Part II [Bexy] for a more detailed introduction). In [E] and in [L98], Jordan-Lie algebras are taken as mathematical starting point to develop quantum theory; thus on purely mathematical grounds, I think it should be important to fully understand what the “geometry of a Jordan-Lie algebra” really is. In particular, the interplay between the weak, projective, setting, and the strong, unitary, setting is rather subtle, and the explanations given above are certainly insufficient.
4.2.8. **Completed qubits.** The smallest non-commutative real universe is the **qubit-space**, the completion of the 4-dimensional Jordan algebra \( \text{Herm}(2, \mathbb{C}) \) (cf. also the table in Appendix C.4). This Jordan algebra is isomorphic to Minkowski space \( \mathbb{R}^{3,1} \), and its positive cone is the Lorentz cone. Its completion is precisely the **conformal compactification of Minkowski space**, often used in relativity theory. (Of course, this is just a pure coincidence, isn’t it?)

4.2.9. **Towards the second chapter: dynamics.** At this point, the first chapter of our book would end. Almost everything the reader is waiting for is still missing, so for sure, s-he would be impatient to start reading the second chapter: so far, there is no Schrödinger equation (no dynamics, no time at all), no Heisenberg relation, not even \( \hbar \) did show up. So, I hope to meet you soon again.

## 5. SOME CONCLUDING REMARKS

I don’t know how many chapters the book may have, and if it will ever be finished. From a mathematical point of view, I think the ideas presented here are kind of inevitable, and should be pursued until they are fully understood. Meanwhile, here are a few more mathematical remarks.

5.1. **Duality, self-duality, and von Neumann.** For my taste, one of the most interesting aspects of the theory explained so far is the interplay between duality, and self-duality: in order to understand and to organize projective geometry, or geometry of quantum theory, duality is a necessary principle; but then it turns out that certain structures are self-dual. The self-dual structures are an important part of the landscape. If I’m not mistaken, it is precisely the feature of self-duality that distinguishes our approach fundamentally from the lattice theoretic Birkhoff–von Neumann approach \([V]\): both are rooted in projective geometry, but self-duality is uninteresting in classical approaches, which deal with projective geometries over **fields**. The lattice structure of a projective line over a field is trivial, hence uninteresting (subspaces are just individual points); over **rings**, this changes drastically.

Seen from a different angle, working with geometries over rings, as opposed to those over fields, also allows to integrate aspects of “fuzzy”, or “intuitionistic logic”, into our theory, without having to use abstract tools like topos theory (cf. \([L17]\), Chapter 12): namely, in projective geometries over fields, there is just one “incidence relation” – a point belongs to a line, or not: **tertium non datur**. In geometries over rings, there are a lot of shades of gray, between white (the point has nothing in common with a line), and black (the point is totally included in the line).\(^8\) Maybe this viewpoint could add a new facet to the topic of “quantum logic”.

5.2. **Infinities; completeness.** The completion of a linear space, such as the linear space of quantum theory, by “points at infinity” provides a convenient and geometric language to speak about “infinities”. Remarkably, it allows to give some sense to “infinities” that seem untractable without the geometric framework (cf. Section 2 of \([BeKi]\)). The bigger the “set at infinity” is, the more it carries structure reflecting

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\(^8\) A one-dimensional submodule over a ring may intersect another subspace non-trivially, without being totally included in it.
complicated analytic or arithmetic structures – for instance, the “most difficult base ring” for our theory is $\mathbb{K} = \mathbb{Z}$ (it has few invertible elements, so we have to add a lot of points at infinity). Could this be a good piece of language for speaking about “problems of infinities” arising in physics?

Paradoxically, while insisting on “geometric completeness”, we relax demands on analytic completeness: we prefer to use the more general $P^*$-algebras rather than $C^*$-algebras; they need not be complete in the analytic or metric sense. For instance I try to avoid using Banach space norms altogether (their geometric meaning in the non-linear context is unclear to me). And although it is analytically very convenient, on physical grounds it seems hard to justify that all Cauchy sequences must converge. Statements of the kind “if something converges, then...” should suffice to cast the logical structure. The uncountable set of all possible limits of all possible Cauchy sequences forms another “infinity” in the geometric sense: both notions of completeness and of infinity have non-trivial relations with each other.

5.3. Completion of commutative and of non-commutative geometry. According to the basic pattern, commutative algebras $A$ correspond to classical systems. This remains true on the level of the “completed” theory: commutative $C^*$-algebras are function algebras, and when $A$ is a function algebra, our formalism of “complete quantum theory” corresponds exactly to what has been proposed in Section 3.2: the projective line over $F(M, \mathbb{R})$ really is the space $F(M, \mathbb{RP}^1)$, and hence “in the classical limit”, we shall get back a classical system. In other words, Section 3.2 describes the “completion of commutative geometry”.

In the same way, the general “complete quantum theory” can be seen as “completion of Non-Commutative Geometry”, where Non-Commutative Geometry (NCG) here is understood in its technical sense defined by A. Connes. Although methods and aims of NCG appear to be quite different from what is proposed here, I see no principal obstruction for asking about transferring certain of its methods and results to the “completed” setting. After all, motivation of NCG by physics is often emphasized, so it might turn out that NCG and “complete quantum theory” are complementary, approaching the same reality from different sides.

5.4. Composed systems. If $A\mathbb{P}^1$ corresponds to one system and $B\mathbb{P}^1$ to another, then the composed system can be described by $(A \otimes B)\mathbb{P}^1$ – composition of systems corresponds to tensor product of algebras. This idea works well for Jordan-Lie algebras, and it even distinguishes them among general Jordan algebras, for which a tensor product of algebras is missing. In fact, this observation was the historical origin for developing the concept of Jordan-Lie algebra in [GP], going back to ideas on “composition classes” by Niels Bohr – see [Be08b] for references and some more remarks. This, again, motivates to develop a theory describing the geometry corresponding to Jordan-Lie algebras ([Bexy]).

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9 If we take algebras of continuous, or smooth, functions, then some non-trivial analysis is needed to describe the precise relationship between the projective line over this algebra, and the space of all (continuous or smooth) functions with values in $\mathbb{RP}^1$. 
Appendix A. $P^*$-algebras

One cannot do mathematics without using formulas. We have avoided them as much as possible in the main text, but in the appendices we give precise definitions and formulas for some of the objects mentioned in the main text. First of all, some definitions related to algebra. Algebraists have the habit to work with algebras defined over a general commutative field or ring $\mathbb{K}$. We will do the same here; for physicists this may be motivated by the fact that in view of understanding discrete models (cf. subsection 2.4) it may interesting to have formalisms that are valid beyond the “usual choice” $\mathbb{K} = \mathbb{R}$.

Definition A.1. A (binary) algebra (over a commutative field or ring $\mathbb{K}$) is a linear space over $\mathbb{K}$, together with a bilinear product map $\mathbb{A} \times \mathbb{A} \rightarrow \mathbb{A}$. If the product is associative, we call $\mathbb{A}$ an associative algebra, and often write the product as $a \cdot b$, or by simple juxtaposition $ab$. If the product is skew-symmetric and satisfies the Jacobi-identity, $\mathbb{A}$ is called a Lie algebra, and the product is often denoted by $[a,b]$. If the product, written $\bullet$, is commutative and satisfies the Jordan identity, $\forall a,b \in \mathbb{A} : a \bullet (b \bullet a^2) = (a \bullet b) \bullet a^2$, where $a^2 = a \bullet a$, then $\mathbb{A}$ is called a Jordan algebra.

Every associative algebra gives rise to a family of associative, Lie, and of Jordan algebras, sometimes called homotopes of each other:

Lemma A.2. Let $\mathbb{A}$ be an associative algebra, and fix $u \in \mathbb{A}$. Then

$$ a \cdot u \ b = aub, \quad [a,b]_u := aub - bua, \quad a \bullet_u b := \frac{1}{2}(aub + bua). $$

are associative, Lie, respectively Jordan algebra products on $\mathbb{A}$.

When $u = 1$ is a neutral element, we get the “usual” products $ab$, $[a,b]$ and $a \bullet b$.

Definition A.3. An associative algebra is called unital if it has a unit element $1$. Then an element $a \in \mathbb{A}$ is called invertible if there is $b \in \mathbb{A}$ with $ab = 1 = ba$. The set $\mathbb{A}^\times$ of invertible elements then is a group.

It is easy to show that $a$ is invertible, if and only if, both the left and right multiplication operators $L_a(x) = ax$ and $R_a(x) = xa$ are invertible, iff the operator $Q_a(x) = L_a \circ R_a(x) = axa$ is invertible. This may serve to define invertible elements even in non-unital associative algebras (see below, def. B.2).

Definition A.4 ($\ast$-algebra). A $\ast$-algebra is an associative complex algebra $\mathbb{A}$ together with an involution $\mathbb{A} \rightarrow \mathbb{A}$, $a \mapsto a^\ast$ (that is, a complex anti-linear map such that $(a^\ast)^\ast = a$, $(ab)^\ast = b^\ast a^\ast$, $1^\ast = 1$). An element $a \in \mathbb{A}$ is called Hermitian if $a^\ast = a$, and skew-Hermitian if $a^\ast = -a$. The sets of (skew) Hermitian elements are denoted by

$$ \text{Herm}(\mathbb{A}) = \{ a \in \mathbb{A} \mid a^\ast = a \}, \quad \text{SHer}(\mathbb{A}) = \{ a \in \mathbb{A} \mid a^\ast = a \}. $$

The unitary group of a (unital) $\ast$-algebra is the subgroup of $\mathbb{A}^\times$ given by

$$ U(\mathbb{A}, \ast) := \{ a \in \mathbb{A} \mid aa^\ast = 1 = a^\ast a \} = \{ a \in \mathbb{A}^\times \mid a^{-1} = a^\ast \}. $$
By decomposing \( a = \frac{a^+ + a^-}{2} + \frac{a^+ - a^-}{2} \), we see that \( A = \text{Herm}(A) \oplus \text{S Herm}(A) \), and since \( * \) is antilinear, we get \( \text{S Herm}(A) = i \text{Herm}(A) \), whence

\[
A = \text{Herm}(A) \oplus i \text{Herm}(A).
\]  

(A.2)

**Lemma A.5.** Assume \( A \) is a *-algebra, and consider the products given by (A.1). If \( u^* = u \), then \( \text{Herm}(A) \) is a Jordan algebra and \( \text{S Herm}(A) \) a Lie algebra, and if \( u^* = -u \), then \( \text{Herm}(A) \) is a Lie algebra and \( \text{S Herm}(A) \) is a Jordan algebra. In particular, \( \text{Herm}(A) \) is a Jordan algebra for the product \( \bullet \) and a Lie algebra for the product \([-\cdot]\), for any constant \( h \in \mathbb{R} \).

**Proof.** It suffices to check that the spaces are stable under the product in question, and this follows directly from \((a b)^* = b^* u^* a^* \).

The Jordan and Lie products on \( \text{Herm}(A) \) satisfy certain natural compatibility conditions: they define a *Jordan-Lie algebra*. This structure is important for the theory of *time evolution*, and will be investigated in more detail in Part II [Bexy].

**Definition A.6 (\( P^* \)-algebra).** A \( P^* \)-algebra is a positive *-algebra, that is, a *-algebra such that \( \text{Herm}(A) \) carries a structure of ordered vector space (over the real field, or over some other partially ordered ring), such that,

1. Whenever \( b \in \text{Herm}(A) \) is positive (i.e., \( b \geq 0 \)), then \( \forall a \in A : a b^* \geq 0 \),
2. For all \( a \in A \), and all invertible \( b \in A \), the element \( a^* a + b^* b \) is invertible.

If this holds, then \( (\text{Herm}(A), \leq) \) is an ordered Jordan algebra (see [Be17b]). The second condition is a weakening of the well-known condition of being formally real.

**Lemma A.7.** Any \( C^* \)-algebra \((where \( x \geq 0 \) iff \( \text{Spec}(x) \geq 0 \)) is a \( P^* \)-algebra.

Indeed, in a \( C^* \)-algebra, (2) is vacuous if the algebra has no invertible elements; else, \( \text{Spec}(b^* b) \geq \lambda \) for some constant \( \lambda > 0 \) if \( b \) is invertible, hence the same holds for \( a^* a + b^* b \), hence \( a^* a + b^* b \) is invertible. We prefer to work with \( P^* \)-algebras, since they are more general than \( C^* \)-algebras, and their defining properties have a clear geometric meaning.

**Appendix B. Associative pairs**

Square matrices are generalized by rectangular matrices (including the important special cases of row and column vectors). In the same way, usual (binary) algebras (associative, or Jordan) are generalized by associative pairs, resp. Jordan pairs. Both concepts are not very well known among mathematicians. The Jordan pair concept, as introduced by Loos in [Lo], is quite technical, and we will not use it in this text. The concept of associative pair, on the other hand, is very simple (see [Lo], or Appendix B in [BeKi]):

**Definition B.1.** An associative pair (over a commutative ring \( \mathbb{K} \)) is a pair \((\mathbb{A}^+, \mathbb{A}^-)\) of \( \mathbb{K} \)-modules together with two trilinear maps

\[
\langle \cdot, \cdot, \cdot \rangle^\pm : \mathbb{A}^+ \times \mathbb{A}^+ \times \mathbb{A}^+ \to \mathbb{A}^+, \quad (x, y, z) \mapsto \langle xyz \rangle^\pm
\]
satisfying the following para-associative law:

\[
\langle xy(zu)v \rangle^\pm = \langle \langle xyz \rangle^\pm uv \rangle^\pm = \langle x(uy)z^\pm v \rangle^\pm.
\]
It is called commutative if always \( \langle xyz \rangle^\pm = \langle zyx \rangle^\pm \). Fixing the middle element \( a \in A^\pm \), we get a binary associative product on \( A^\pm \), denoted by \( A_a \) and called the \( a \)-homotope:

\[
xz := x \cdot_a z := \langle xaz \rangle^\pm.
\]

(B.1)

Examples of associative pairs.

1. Every associative algebra \( A \) gives rise to an associative pair \( A^+ = A^- = A \) via

\[
\langle xyz \rangle^+ = xyz, \quad \langle xyz \rangle^- = zyx.
\]

2. For \( K \)-modules \( E \) and \( F \), let \( A^+ = \text{Hom}(E,F) \), \( A^- = \text{Hom}(F,E) \), and

\[
\langle XYZ \rangle^+ = X \circ Y \circ Z, \quad \langle XYZ \rangle^- = Z \circ Y \circ X.
\]

Taking \( F = K \), a linear space \( E \) and its dual \( E' \) form an associative pair.

3. \((\mathbb{R}, \mathbb{R}')\) is an associative pair, and so is \((F(M,'\mathbb{R}), F(M, \mathbb{R}'))\) (cf. section 3.2).

4. Let \( \hat{A} \) be an associative algebra with unit 1 and idempotent \( e \) (that is, \( e^2 = e \)) and \( f := 1 - e \) its “opposite idempotent”. Let

\[
\hat{A} = f\hat{A}f \oplus f\hat{A}e \oplus e\hat{A}e \oplus e\hat{A}f = A_{00} \oplus A_{01} \oplus A_{11} \oplus A_{10}
\]

with \( A_{ij} = \{ x \in \hat{A} \mid ex = ix, xe = jx \} \) the associated eigenspace (Peirce) decomposition. Then

\[
(A^+, A^-) := (A_{01}, A_{10}), \quad \langle xyz \rangle^+ := xyz, \quad \langle xyz \rangle^- := zyx
\]

is an associative pair.

It is easy to show that every associative pair arises from an associative algebra \( \hat{A} \) with idempotent \( e \) in the way just described ([Lo], Notes to Chapter II).

**Definition B.2** (invertible elements). We call an element \( x \in A^\pm \) invertible if

\[
Q_x : A^\pm \to A^\pm, \quad y \mapsto \langle xyx \rangle
\]

is an invertible operator.

As shown in [Lo], associative pairs with invertible elements correspond to unital associative algebras: namely, \( x \) is invertible if and only if the homotope algebra \( A_x \) has a unit (which is then \( x^{-1} := Q_x^{-1}x \)).

**Definition B.3** (idempotent). An idempotent in an associative pair is a pair \( (e^+, e^-) \in A^+ \times A^- \) such that

\[
\langle e^+, e^-, e^+ \rangle = e^+ \quad \text{and} \quad \langle e^-, e^+, e^- \rangle = e^-.
\]

Idempotents are a tool to start to “glue together” the two spaces \( A^+ \) and \( A^- \).

**Appendix C. Projective spaces, projective lines**

In this appendix we describe the construction of basic “geometric spaces” by using rings and algebras. For sake of generality, in this appendix \( R \) is a \( (possibly \text{non-commutative}) \) ring with unit 1 (we reserve the letter \( K \) to commutative rings), and \( W \) a right module over \( R \). For a first reading, think of \( R = \mathbb{R} \) or \( \mathbb{C} \), and \( W \) a vector space, say \( W = \mathbb{C}^n \), or a Hilbert space; but for a second reading, it will be important to allow for \( R \) a non-commutative ring: namely, the role of \( R \) may be taken by some \( \ast \)-algebra \( A \). For a general ring \( R \), recall that a module over \( K \) is
defined like a vector space, except that there is no commutativity of scalars, and therefore we agree to write scalars always on the right of vectors.

C.1. **Grassmannians.** Let $W$ be a right $R$-module, together with a direct sum decomposition $W = A \oplus Z$. We define the Grassmannian of type $A$ and co-type $Z$ to be the set of all submodules $E$ that are isomorphic to $A$ and admit a complement $E'$ isomorphic to $Z$:

$$\text{Gras}_A^Z(W) := \{ E \subset W \text{ submodule} \mid E \cong A, \exists E' \cong Z : W = E \oplus E' \}.$$  \hspace{1cm} (C.1)

The pair of Grassmannians

$$(\mathcal{X}, \mathcal{X}') = (\text{Gras}_A^Z(W), \text{Gras}_Z^A(W))$$  \hspace{1cm} (C.2)

is said to be in duality. For instance, the pair $(\text{Gras}_p(\mathbb{K}^{p+q}), \text{Gras}_q(\mathbb{K}^{p+q}))$, where for a field $\mathbb{K}$, $\text{Gras}_k(\mathbb{K}^n)$ is the Grassmannian of $k$-dimensional subspaces of $\mathbb{K}^n$, is such a pair. The general linear group $GL(W)$ acts on $\mathcal{X}$ and on $\mathcal{X}'$. This action is transitive: by definition, there are linear isomorphisms $g_1 : A \to E$, $g_2 : Z \to E'$, whence $g := g_1 \oplus g_2 \in GL(W)$ sends $A$ to $E$ (and $Z$ to $E'$). If $R$ is commutative, then the scalars act trivially on $\mathcal{X}$ (but else not).

C.2. **Projective spaces, projective lines, self-duality.** We say that $A$ is a line if $A$ is isomorphic to the base ring ($A \cong R$), and then call $Z$ a hyperplane if $W = A \oplus Z$. The projective space of $W$ is the space of all lines admitting a hyperplane complement, and its dual projective space is the space of all hyperplanes: we write

$$(\mathcal{X}, \mathcal{X}') = (\text{Gras}_A^Z(W), \text{Gras}_Z^A(W)) = (\mathbb{P}(W), \mathbb{P}(W)').$$  \hspace{1cm} (C.3)

In the special case $W = R \oplus R = R^2$, with $A$ the first and $Z$ the second factor, this defines the projective line $R\mathbb{P}^1$ over $R$, together with its dual projective line $(R\mathbb{P}^1)'$:

$$(R\mathbb{P}^1, (R\mathbb{P}^1)') = (\text{Gras}_R^R(R \oplus R), \text{Gras}_R^R(R \oplus R)).$$  \hspace{1cm} (C.4)

As sets, $R\mathbb{P}^1$ and $(R\mathbb{P}^1)'$ agree: the projective line is self-dual. Both copies may be distinguished by taking different base points: in the first copy, the base point 0 is the first factor $R \times 0 = [(1,0)]$, and in the second copy, the base point $\infty$ is the second factor $0 \times R = [(0,1)]$, where we write $[(x,y)] = (x,y)R$ for the right module generated by $(x,y)$. Given these base points, there is a natural imbedding

$$R \to R\mathbb{P}^1, \quad z \mapsto [(z,1)],$$  \hspace{1cm} (C.5)

and $GL(2,R) = GL(R \oplus R)$ acts, just as in the classical case $R = \mathbb{C}$, on the affine part $R$ by “fractional linear transformations”:

$$\begin{pmatrix} a & b \\ c & d \end{pmatrix} \begin{pmatrix} z \\ 1 \end{pmatrix} = \begin{pmatrix} (az + b) \\ (cz + d) \end{pmatrix} = \begin{pmatrix} (az + b)(cz + d)^{-1} \\ 1 \end{pmatrix}.$$  \hspace{1cm} (C.6)

C.3. **Transversality, projection operators.** Let $(\mathcal{X}, \mathcal{X}')$ be as in (C.2). A pair $(x,a) \in \mathcal{X} \times \mathcal{X}'$ is called transversal, and we write $a \perp x$, if $W$ is the direct sum of $a$ and $x$: $W = a \oplus x$. We denote the set of all complementary subspaces of $a$ by

$$U_a = \{ x \in \mathcal{X} \mid x \perp a \}.$$  \hspace{1cm} (C.7)

**Theorem C.1.** The set $U_a$ carries a natural structure of an affine space over $R$. 

This result is classical, and easily proved. For instance, for any scalar \( r \in R^\times \), multiplication by \( r \) in the linear space \((U_a, x)\),
\[
    r_{a,x} : \mathcal{X} \to \mathcal{X}, \quad y \mapsto r_{a,x}(y)
\]
is given by the matrix \( \left( \begin{smallmatrix} r & 0 \\ 0 & 1 \end{smallmatrix} \right) \) with respect to the direct sum decomposition \( W = x \oplus a \). The following is less classical (to my knowledge, [BeKi] is the first time it appeared):

**Theorem C.2.** Fix \((a, b) \in \mathcal{X}' \times \mathcal{X}'\) and let \( U_{ab} := U_a \cap U_b \) (set of common complements of \( a \) and \( b \)). If \( U_{ab} \) is not empty, fix an arbitrary element \( y \in U_{ab} \). Then \( U_{ab} \) carries a natural group structure with neutral element \( y \). In case \( a = b \), this group is isomorphic to the general linear group \( GL(a) \).

Let us, following [BeKi], describe the group law on \( U_{ab} \). Its product shall be denoted by \( x \cdot y, z, \) or \((xyz)_{ab}\). The main tool for defining it are the projection operators: when \( a \cap x \), denote by \( P^a_x : W \to W \) the linear projector having kernel \( a \) and image \( x \). Then, clearly, \( P^a_x \circ P^b_x = P^b_x \circ P^a_x = P^b_x \), and \( P^a_x \circ P^b_x = 0 \). From these rules it follows by direct computation that, whenever \( x, y, z \in U_{ab} \), the linear operator
\[
    M_{xyz} := P^a_x - P^b_z : W \to W
\]
is invertible, with inverse \( M_{zyx} \). Applying it to \( y \) gives the group law:
\[
    x \cdot y \cdot z = (xyz)_{ab} = M_{xyz}(y) = (P^a_x - P^b_z)(y).
\]
And \( r_{a,x} = rP^a_x + P^x \). Much more can be said about this (cf. loc. cit.).

**Example C.1.** When \( X = \mathbb{P}(W) \) is a projective space and \( a \) a hyperplane, then \( U_a \) is identified with the set of all points of \( X \) that do not belong to the hyperplane defined by \( a \), \( U_a = \mathcal{X} \setminus H_a \), and \( H_a \) is the “horizon of \( U_a \)” (set of points at infinity of the affine space \( U_a \)). In this case, there exists a simple lattice-theoretic formula describing the group law of \( U_{ab} \).

**Example C.2.** If \( X = R\mathbb{P}^1 = \text{Gras}^R(R \oplus R) \) is a generalized projective line, then the image of the imbedding (C.5) is \( U_\infty \), isomorphic to \((R, +)\) as a group. Likewise, \( U_0 \cong (R, +) \) as a group, and \( U_{0\infty} = U_\infty \cap U_0 = R^\times \) is the multiplicative group of the ring \( R \) (which is possibly non-commutative, according to our assumption).

**C.4. The Hermitian projective line.** Now consider the case where the ring \( R \) is a *-algebra \( A \), say over \( \mathbb{K} = \mathbb{C} \). Since \( A\mathbb{P}^1 \) generalizes the Riemann sphere, we also write \( S := A\mathbb{P}^1 = \text{Gras}^A(A^2) \). The involution \( * : A \to A \) induces an involution \( \tau : S \to S \). This involution is given by taking the orthocomplement \( \tau(x) := x^{\perp \omega} \) of the submodule \( x \) with respect to the “Poisson form”, that is, the skew-Hermitian sesquilinear form \( \omega : A^2 \times A^2 \to A \) given by
\[
    \omega((u_1, u_2), (v_1, v_2)) = (u_1^* u_2^*, v_1^* v_2) J \begin{pmatrix} v_2 \\ v_1 \end{pmatrix} = u_1^* v_2 - u_2^* v_1.
\]
where \( J \) is the matrix
\[
    J := \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix}.
\]
Indeed, the orthocomplement of \([1, a]\) is \([1, a^*]\), because for all \(x, y \in \mathbb{A}\),

\[
\omega((x, ax), (y, a^* y)) = x^*(a^* y) - (ax)^* y = x^* a^* y - x^* a^* y = 0,
\]

i.e., \(\tau(a) = a^*\) on \(\mathbb{A} \subset \mathcal{S}\). (For general elements of \(\mathcal{S}\), there is no “closed formula” describing \(\tau\), unless \(\mathbb{A}\) is commutative.) The fixed point set of this involution is called the *Hermitian projective line over \((\mathbb{A}, *)\); it generalizes the “equator of the Riemann sphere”, and it is the set of \(\omega\)-Lagrangian subspaces which we denote by

\[
\mathcal{R} := \text{Herm}(\mathbb{A}) \mathbb{P}^1 = \text{Lag}_\omega(\mathbb{A}^2) = \{ x \in \mathcal{S} \mid x^{1/\omega} = x \}.
\] (C.13)

The self-dual geometry \((\mathcal{R}, \mathcal{R})\) is the total space for “completed quantum theory”. Its automorphism group is \(\text{Aut}(\mathcal{R}) = \text{Aut}(\mathcal{S}, \tau) = \mathbb{P}\text{Aut}(\omega) = \text{Aut}(\omega)/\mathbb{Z}\), the “symplectic group” of \(\omega\), modulo its center. It acts transitively on \(\mathcal{R}\). In matrix form, the group \(\text{Aut}(\omega)\), and its Lie algebra \(\text{Der}(\omega)\) can be described by \(2 \times 2\)-matrices

\[
\text{Aut}(\omega) = \{ g \in \text{Gl}(2, \mathbb{A}) \mid ^t g^* J g = J \}, \quad \text{Der}(\omega) = \{ M \in \text{M}(2, 2; \mathbb{A}) \mid ^t M^* J = -JM \}
\] (C.14) (C.15)

\[
= \{ \left( \begin{array}{cc} a & b \\ c & -a^* \end{array} \right) \mid a, b, c \in \mathbb{A}, b^* = b, c^* = c \}. \quad \text{(C.16)}
\]

As a real vector space, \(\text{Der}(\omega) = \text{Herm}(\mathbb{A}) \oplus \mathbb{A} \oplus \text{Herm}(\mathbb{A})\), and when \(\mathbb{A} = \mathbb{C}\), we see that \(\text{Der}(\omega) = \mathfrak{sl}_2(\mathbb{R}) \oplus i\mathbb{R}\), where \(i\mathbb{R}\) is the center. Finally, in Subsection 4.2 we consider also the orthocomplementation map \(\alpha\) with respect to the positive “scalar product” on \(\mathbb{A}^2\), whose automorphism group is the unitary group \(\text{U}(2, \mathbb{A})\).

C.4.1. *Toy model.* In the following table (second column) we give formulae describing the special case of the algebra \(\mathbb{A} = M(n, n; \mathbb{C})\) with involution \(A^* = \overline{A}\) (conjugate transpose); this case is a good finite dimensional “toy model” for quantum mechanics (\(n = 1\) is “classical”; \(n = 2\) is the “qubit”).

<table>
<thead>
<tr>
<th>complex associative algebra</th>
<th>Hilbert space setting: (\mathcal{H} = \mathbb{C}^n), (\langle u, v \rangle = \sum_i u_i v_i)</th>
<th>(*)-algebra setting: ((\mathbb{A}, \ast))</th>
</tr>
</thead>
<tbody>
<tr>
<td>involution (*)</td>
<td>(a^* = \overline{a})</td>
<td>(*)</td>
</tr>
<tr>
<td>observables (Jordan algebra)</td>
<td>(\text{Herm}(\mathcal{H}) = \text{Herm}(n, \mathbb{C}))</td>
<td>(\text{Herm}(\mathbb{A}))</td>
</tr>
<tr>
<td>projective line (\mathbb{A}^1)</td>
<td>(\text{Gras}_n(\mathcal{H} \oplus \mathcal{H}) = \text{Gras}_n(\mathbb{C}^{2n}))</td>
<td>(\text{Gras}_n^{\ast}(\mathbb{A} \oplus \mathbb{A}))</td>
</tr>
<tr>
<td>automorphism group (\text{Aut}(\mathbb{A}^1))</td>
<td>(\mathbb{P}\text{Gl}(2n, \mathbb{C}))</td>
<td>(\text{Gl}(2, \mathbb{A})/\text{Center})</td>
</tr>
<tr>
<td>finite part (\mathbb{A} \subset \mathbb{A}^1)</td>
<td>({ \text{Graph}_a \mid a \in \text{End}(\mathcal{H}) })</td>
<td>({([1, a]) \mid a \in \mathbb{A} }), ({([u, v]) = ([u, v])_\mathcal{A} })</td>
</tr>
<tr>
<td>base point 0</td>
<td>(\text{Graph}_0 = \mathcal{H} \times 0)</td>
<td>((1, 0))</td>
</tr>
<tr>
<td>base point (\infty)</td>
<td>(0 \times \mathcal{H})</td>
<td>((0, 1))</td>
</tr>
<tr>
<td>unit 1</td>
<td>(\text{Graph}_1 = \text{diag}(\mathcal{H} \times \mathcal{H}))</td>
<td>((1, 1))</td>
</tr>
<tr>
<td>Poisson form (\omega((u_1, u_2), (v_1, v_2)))</td>
<td>(\ldots = \langle u_1, v_2 \rangle - \langle u_2, v_1 \rangle)</td>
<td>(\ldots = u_1^* v_2 - u_2^* v_1)</td>
</tr>
<tr>
<td>involution of (\mathbb{A}^1)</td>
<td>(\omega)-orthocomplement</td>
<td>(x \mapsto x^{1/\omega})</td>
</tr>
<tr>
<td>Hermitian projective line (\mathcal{R})</td>
<td>(\text{Lag}_\omega(\mathcal{H} \oplus \mathcal{H}))</td>
<td>(\text{Lag}_\omega(\mathbb{A}^2))</td>
</tr>
<tr>
<td>(G = \text{Aut}(\mathcal{R}))</td>
<td>(\text{Aut}(\omega) \cong \text{U}(n, n)/\text{Center})</td>
<td>(\text{Aut}(\omega)/\text{Center})</td>
</tr>
</tbody>
</table>
The formulae in the first column are obtained from the general formulae of the
second column by considering a \((2n) \times (2n)\)-matrix as a \(2 \times 2\)-matrix with entries
in the algebra \(A = M(n, n; \mathbb{C})\). Note, however, that in the general case (second
column) \(A\) may be an infinite dimensional algebra over \(\mathbb{C}\), and all the preceding
groups are infinite dimensional Lie groups (and they are much “bigger” than those
usually considered in quantum mechanics: they contain many “hidden variables”,
that is, degrees of freedom that are not activated in usual quantum mechanics).
With some care, formulae from the “toy model” generalize to the infinite dimensional
Hilbert space setting of quantum mechanics (to get conceptual formulae, it may be
useful to replace \(\mathbb{C}^{2n}\) by \(\mathcal{H} \oplus \mathcal{H}'\), the direct sum of a Hilbert space and its dual space).

Appendix D. Cross-ratio

The cross-ratio is the most important invariant of a projective line. We approach
it in two steps: first of all, we recall the classical definition for the projective line
over a commutative field or ring, and second, we discuss generalizations to the
non-commutative case.

D.1. The classical cross-ratio. Retain notation from the preceding appendix,
and assume \(R\) is a commutative field henceforth denoted by \(\mathbb{K}\). We denote by
\(W' := \text{Hom}(W, \mathbb{K})\) the algebraic dual space of \(W\). Then \((X, X') = (\mathbb{P}(W), \mathbb{P}(W'))\)
are projective spaces in duality (a hyperplane \(H\) in \(W\) corresponds to \(\ker(A)\) where
\(A : W \to \mathbb{K}\) is determined up to a scalar). We define the cross-ratio of a quadruple
\((x, y, a, b) = ([\xi], [\eta], [A], [B]) \in X^2 \times (X')^2\) (so vectors \(\xi, \eta\) and linear forms \(A, B\)
defined up to a scalar) by

\[
\text{CR}(x, y; a, b) := \frac{A(\xi)}{A(\eta)} : \frac{B(\xi)}{B(\eta)} = \frac{A(\xi) \cdot B(\eta)}{A(\eta) \cdot B(\xi)} . \tag{D.1}
\]

Note that this is well-defined (independent of scaling of \(\xi, \eta, A, B\)), and it clearly is
an invariant under the natural action of the general linear group \(\text{GL}(2, \mathbb{K})\), acting as
usual on \(W\) and on \(W'\). In other words, this definition defines a natural invariant
of projective spaces and their duals.

Now assume that \(W\) is two-dimensional, say \(W = \mathbb{K}^2\). The special feature of this
case is the existence of a canonical symplectic form \(\omega\): it is given by the same
formula as (C.11), with involution the identity map:

\[
\omega : \mathbb{K}^2 \times \mathbb{K}^2 \to \mathbb{K}, \quad ((x_1, x_2), (a_1, a_2)) \mapsto x_1a_2 - x_2a_1 . \tag{D.2}
\]

This form is \(\mathbb{K}\)-bilinear (commutativity of \(\mathbb{K}\) is crucial here!), and up to a factor
it is invariant under the whole linear group \(\text{GL}(2, \mathbb{K})\). We may use it in order to
identify \(X\) and \(X'\), so that our invariant (D.1) is turned into a function defined on
\(X^4\) and given by

\[
\text{CR}(x, y; a, b) = \frac{\omega(A, \xi)}{\omega(A, \eta)} : \frac{\omega(B, \xi)}{\omega(B, \eta)} = \frac{(a_1x_2 - a_2x_1)(b_1y_2 - b_2y_1)}{(a_1y_2 - a_2y_1)(b_1x_2 - b_2x_1)} . \tag{D.3}
\]

Letting in this formula \(x_2 = 1 = a_2 = y_2 = b_2\), we get the value \(\frac{(a_1 - x_1)(b_1 - y_1)}{(a_1 - y_1)(b_1 - x_1)}\),
which corresponds to the “usual” definition of the cross-ratio \(\text{CR}(x_1, y_1; a_1, b_1)\), as given by
formula (3.4). Hence both definitions are in keeping. The one given here has several advantages: it features duality, and it shows where self-duality and commutativity of \( K \) enter into the definition.

D.2. Operator valued cross-ratio. We wish to define analogs of the cross-ratio in the general setting of C.1. In fact, it is always possible to define an “operator valued cross-ratio”, but it may be problematic to extract from it a scalar valued invariant – to do this, one needs things like determinants, or traces, and these may not exist in infinite dimension. Let us explain this: assume \( a \perp x \). Then the linear spaces \((U_x, U_a)\), with origin \((a, x)\), are in duality with each other, in the following sense: with respect to the decomposition \( W = a \oplus x = a \times x \), every element \( b \in U_x \) can be written as the graph of a unique linear map \( \beta : a \to x \), and every element \( y \in U_a \) as graph of a linear map \( \eta : x \to a \). Thus, as pair of linear spaces,

\[
(U_x, U_a) = (\text{Hom}_R(a, x), \text{Hom}_R(x, a)).
\]  

(D.4)

Note that this pair is again an associative pair. This observation leads to define two “canonical kernel functions”

\[
K_{x,a} : U_x \times U_a \to \text{End}(x), \quad (\beta, \eta) \mapsto \beta \circ \eta,
\]

(D.5)

\[
K_{a,x} : U_a \times U_x \to \text{End}(a), \quad (\eta, \beta) \mapsto \eta \circ \beta.
\]

(D.6)

Now we define the generalized (operator valued) cross-ratio by

\[
\text{CR}(y, b; x, a) := K_{x,a}(b, y) \in \text{End}(x)
\]

(D.7)

(so \( \text{CR}(y, b; 0, \infty) = by \)). The construction is natural, that is, invariant under the symmetry group \( \text{Gl}(W) \):

\[
\forall g \in \text{Gl}(W) : \quad K_{gx,ga}(gb, gy) = gK_{x,a}(b, y)g^{-1}.
\]

(D.8)

The operators now live in a space depending on \( x \) (or on \( a \)). In technical terms, they define an invariant section of a vector bundle. Moreover, at least when \( a \) and \( b \) are transversal, the operator valued cross-ratio is closely related to the left, right and middle multiplication operators defined in [BeKi].

D.3. Scalar valued cross-ratio (expectation value). To extract a well defined scalar from the operator valued cross-ratio, we have to compose the \( \text{End}(x) \)-valued cross-ratio with a function \( \text{End}(x) \to \mathbb{K} \) that is conjugation invariant, such as determinant or trace functions (see Appendix E for more on traces),

\[
\det(gAg^{-1}) = \det(g), \quad \text{trace}(gAg^{-1}) = \text{trace}(A),
\]

(D.9)

giving two candidates to define a scalar valued cross-ratio:

\[
(x, y; a, b) \mapsto \det(K_{x,a}(b, y)), \quad (x, y; a, b) \mapsto \text{trace}(K_{x,a}(b, y)).
\]

(D.10)

More generally, \( \det \) and \( \text{trace} \) could be replaced here by any map \( \chi : \text{End}(x) \to \mathbb{K} \) that satisfies (D.9). In the same way, when \( A \) is a non-commutative algebra over \( \mathbb{R} \), there is an operator valued cross-ratio on \( \mathbb{A}^\mathbb{P}^1 \); but to extract from it an \( \mathbb{R} \)-scalar valued one, we again need some conjugation invariant map \( \chi : A \to \mathbb{R} \). Defining such maps is, in infinite dimension, closely related to integration theory, see the following appendix.
Appendix E. Pairings, densities, and traces

Recall from subsection 3.2 that \( \mathbb{R} \) and \( \mathbb{R}' \) are just two copies of \( \mathbb{R} \), without a fixed base, and \( \mathbb{R}' \) is the dual space of \( \mathbb{R} \). If \( M \) is, say, a topological space, or a general measurable space, we denote in this appendix by \( F(M, \mathbb{R}) \) and \( F(M, \mathbb{R}') \) the set of all measurable functions on \( M \).

**Definition E.1.** A pairing on \( (F(M, \mathbb{R}), F(M, \mathbb{R}')) \) is a map

\[
\Pi : F(M, \mathbb{R}) \times F(M, \mathbb{R}') \to \mathbb{R} \cup \{\infty\} = \mathbb{RP}^1, \quad (f, g) \mapsto \Pi(f, g)
\]

having the following properties:

1. it is \( \mathbb{R} \)-bilinear (in the sense of extended arithmetic operations including rules like \( x + \infty = \infty, \lambda \infty = \infty \), for \( x, \lambda \in \mathbb{R} \)),
2. for all \( f, g, h \in F(M, \mathbb{R}) \), we have \( \Pi(fh, g) = \Pi(f, hg) \),
3. it is \( \sigma \)-continuous: if \( f_n \downarrow 0 \) (pointwise monotone convergence), then \( \Pi(f_n, g) \downarrow 0 \) (whenever \( \Pi(f_N, g) \neq \infty \) for at least some \( N \in \mathbb{N} \)), and likewise in the second argument.

If \( \mu \) is a measure on \( M \), we write \( \mu(f) = \int_M f \, d\mu \), and then the formula

\[
\Pi_{\mu}(f, g) := \mu(fg) \quad \text{(E.1)}
\]

defines a pairing. This pairing is “invariant” in the following sense: let \( G \) be the group of measurable bijections of \( M \) preserving the collection of sets of measure zero. An element \( \phi \in G \) need not preserve \( \mu \), but \( \phi_* \mu(h) := \mu(h \circ \phi) \) defines another measure \( \phi_* \mu \), which is absolutely continuous with respect to \( \mu \). Thus, by the Radon-Nikodym theorem, there is a function \( \phi' \) with \( \phi_* \mu = \phi' \mu \), i.e., for all \( h \),

\[
\mu(h \circ \phi) = \mu(\phi' \cdot h). \quad \text{(E.2)}
\]

(to be precise, \( \phi' \) is defined up to sets of measure zero). Applying (E.2) twice, we get the “chain rule”: for all \( \phi, \psi \in G \), we have \( (\phi \circ \psi)' = \phi' \cdot (\psi' \circ \phi^{-1}) \). Now we let \( \phi \) act \( G \) in the “usual” way by \( \phi.f = f \circ \phi^{-1} \) on “usual” functions \( F(M, \mathbb{R}) \), and via

\[
\phi.h := \phi' \cdot (h \circ \phi^{-1}) \quad \text{(E.3)}
\]
on \( F(M, \mathbb{R}') \). By the “chain rule”, this is indeed an action. When equipped with this action, we call the space \( F(M, \mathbb{R}') \) the space of \( \mu \)-densities. Now, using (E.2),

\[
\Pi_{\mu}(\phi.f, \phi.h) = \mu((f \circ \phi^{-1}) \cdot (h \circ \phi^{-1})) \circ \phi = \mu(fh) = \Pi_{\mu}(f, h).
\]

This proves:

**Proposition E.2.** With notation as above, the pairing (E.1) is invariant under the group \( G \): for all \( \phi \in G \), we have \( \Pi_{\mu}(\phi.f, \phi.h) = \Pi_{\mu}(f, h) \).

The following examples illustrate that, under natural assumptions, the pairing can be considered as “canonical” – it does not really depend on \( \mu \), but only on the class of measures having the same sets of measure zero as \( \mu \):
(1) If \( M \) is a differentiable manifold, then we define integration with respect to volume forms, as usual. Then densities in our sense coincide with those in the sense of differential and our trace is the pairing between densities and (say) continuous functions. This pairing is invariant under the group of all diffeomorphisms, which is a subgroup of \( G \).

(2) Generalizing the preceding item, whenever we have a partition of unity subordinate to an atlas of \( M \), then pairings defined with respect to chart domains can be glued together to give a pairing on \( M \).

(3) If \( M \) is a finite set, then trace\((f,g) := \sum_{p \in M} f(p)g(p)\) defines a pairing. It is invariant under the group of all bijections of \( M \).

Once we have a theory of pairings on the “classical pair” \((F(M, \mathbb{R}), F(M, \mathbb{R}'))\), one would like to develop such a theory on more general associative pairs \((A^+, A^-)\). In case \( A^+ = A^- = A \) is a \( C^*\)-algebra, this should more or less correspond to spectral theory, and hence developing such a theory is a big task clearly exceeding the scope of the present work. We shall just state the following definitions:

**Definition E.3.** Let \( (A, \ast) \) be a \( \mathbb{P}^*\)-algebra (def. A.6). A trace on \( A \) is a linear map in the generalized sense, as above)\[
\text{trace} : A \rightarrow \mathbb{CP}^1
\]
that is symmetric: \( \text{trace}(ab) = \text{trace}(ba) \) and positive: whenever \( b \in \text{Herm}(A) \) is positive, then \( \text{trace}(b) \geq 0 \). It may be normalized by the following condition: if \( a \in \text{Herm}(A) \) is an idempotent of rank one, then \( \text{trace}(a) = 1 \).

More geometrically, the trace map could be defined as an ordered morphism of projective lines \( \mathbb{AP}^1 \rightarrow \mathbb{CP}^1 \), and possibly normalized by the condition that on each “intrinsic projective line” \( \mathbb{CP}^1 \) contained in \( \mathbb{AP}^1 \), it should induce the identity mapping. Given such a trace function, the binary map \( (a, b) \mapsto \Pi(a, b) := \text{trace}(ab) \) then should define what one might call a “pairing on \((A, A)\)”, and \( (x, y; a; b) \mapsto \text{trace}(K_{x,a}(b, y)) \) an “expectation value” (scalar valued cross-ratio).

**References**


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