The correct formulation of Gleason’s theorem in quaternionic Hilbert spaces

Valter Moretti\textsuperscript{a}, Marco Oppio\textsuperscript{b}

\textsuperscript{a} Department of Mathematics University of Trento, and INFN-TIFPA
via Sommarive 15, I-38123 Povo (Trento), Italy.
valter.moretti@unitn.it

\textsuperscript{b} Faculty of Mathematics, University of Regensburg
Universitätsstrasse 31, 93053 Regensburg, Germany
marco.oppio@ur.de

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Abstract. Quantum Theories can be formulated in real, complex or quaternionic Hilbert spaces as established in Solé\’r’s theorem. Quantum states are here pictured in terms of $\sigma$-additive probability measures over the non-Boolean lattice of orthogonal projectors of the considered Hilbert space. Gleason’s theorem proves that, if the Hilbert space is either real or complex and some technical hypotheses are true, then these measures are one-to-one with standard density matrices used by physicists recovering and motivating the familiar notion of state. The extension of this result to quaternionic Hilbert spaces was obtained by Varadarajan in 1968. Unfortunately, the formulation of this extension \cite{Va68} is partially mathematically incorrect due to some peculiarities of the notion of trace in quaternionic Hilbert spaces. A minor issue also affects Varadarajan’s statement for real Hilbert space formulation. This paper is devoted to present Gleason-Varadarajan’s theorem into a technically correct and physically meaningful form valid for the three types of Hilbert spaces. In particular, we prove that only the \textit{real part} of the trace enters the formalism of quantum theories (also dealing with unbounded observables and symmetries) and it can be safely used to formulate and prove a common statement of Gleason’s theorem.
1 Introduction

1.1 Gleason’s theorem and troubles with the quaternionic formulations

The idea to formulate Quantum Theories starting from the partially ordered set of elementary propositions can be traced back to Birkhoff and von Neumann [BiNe36] with crucial contributions by Mackay [Ma63] (see [EGL09] for a review on these issues). From a physical point of view, these elementary propositions are the statements which can be physically tested on a quantum system, receiving either the outcome 0 (not valid) or the outcome 1 (valid). In particular, every physical quantity \( A \) which can be measured on the system can be pictured as a collection of such elementary propositions \( P_E(A) \) labelled by sets of reals \( E \) like an interval \( E = (a, b) \). The physical meaning of \( P_E(A) \) is that \emph{the outcome of the observable \( A \) belongs to \( E \)}. The elementary propositions of a quantum system obey a logic different from the classical one in view of the presence of physically incompatible statements (in the sense of Heisenberg principle). Von Neumann’s idea was that these elementary propositions are one-to-one described by the orthogonal projectors of a complex Hilbert space \( \mathcal{H} \) associated to the quantum system and a pair of such elementary propositions are physically incompatible if and only if the corresponding projectors do not commute. The set \( \mathcal{L}(\mathcal{H}) \) of orthogonal projectors over \( \mathcal{H} \) enjoys the structure of a non-Boolean lattice with respect to the ordering relation given by the set-theoretic inclusion of projection subspaces. In particular, the collection of elementary propositions associated to a physical quantity \( A \) as above are supposed to be pairwise compatible and the analysis of the mutual relations between these propositions (see, e.g., [Va68, Mo18]) shows that they give rise to the structure of a projection-valued measure over \( \mathbb{R} \). At this
point, the spectral machinery, integrating this projection-valued measure, permits one to associate the abstract observable \( A \) to a corresponding selfadjoint operator, motivating from a deeper viewpoint the standard assumption by physicists that observables are described by selfadjoint operators. It is also possible to assume an even more abstract viewpoint, where the elementary propositions are treated as elements of an abstract lattice with specific mutual properties reflecting the general phenomenology of quantum systems, without explicitly referring to the orthogonal projectors (the notion of commutativity of pairs of elements can be generally defined for abstract lattices exploiting the notion of Boolean sublattice). Adopting this quite abstract point of view, it turns out that, in addition to the lattice of orthogonal projectors over a complex Hilbert space, also the lattices of orthogonal projectors over either a real and a quaternionic Hilbert space fulfill the requirements which, in principle, may be justified by the quantum phenomenology (irreducibility, orthomodularity, \( \sigma \)-completeness, separability, atomicity and validity of the so-called covering property, see, e.g., [BeCa81]). The proof of the fact that these are the only three possibilities has a long history (see, e.g., [BeCa81, AeSt00, EGL09, Mo18] for technical discussions of the various intermediate results), starting from several remarkable results by Piron and other few authors like F. Maeda, S. Maeda and Araki in the ‘60s, and ending with the theorem by Solé in 1995 [So95]. Assuming that the abstract lattice of quantum elementary propositions is irreducible, orthomodular, \( \sigma \)-complete, separable, atomic, that it satisfies the covering property and that it includes an infinite set of orthogonal atoms, then the lattice is necessarily isomorphic to the lattice \( \mathcal{L}(H) \) of orthogonal projectors in a separable Hilbert space \( H \). There, the set of scalars may only be \( \mathbb{R}, \mathbb{C} \) or the real algebra of quaternions \( \mathbb{H} \).

We shall not address here the problem of the apparent absence of physical systems described in real Hilbert spaces [St60, StGu61, MoOp17] and the possibility (or impossibility) of quaternionic formulations [FJSS62, Ad95, Gan17, MoOp17b]. Instead, we restrict attention to a celebrated result, provided by Gleason’s theorem [Gl57], regarding the notion of quantum state. The lattice structure of \( \mathcal{L}(H) \) suggests a physically natural definition of a quantum state (at a given instant of time \( t \)) defined as a map associating every elementary proposition \( P \in \mathcal{L}(H) \) with the probability that it results to be true if measured. It should be evident that this is a good notion of state, as it includes all information necessary to compute probabilities, expectation values, standard deviations and so on for every observable (selfadjoint operator) viewed as the family of its spectral projectors (elementary propositions). Formally, a state is a generalized \( \sigma \)-additive probability measure over \( \mathcal{L}(H) \), that is a map \( \mu : \mathcal{L}(H) \rightarrow [0,1] \) satisfying both \( \mu(I) = 1 \), and \( \sum_{k \in K} \mu(P_k) = \mu(\bigvee_{k \in K} P_k) \) for every set \( \{P_k\}_{k \in K} \subset \mathcal{L}(H) \) with \( K \) finite or countably infinite, such that \( P_k \) and \( P_h \) are orthogonal when \( k \neq h \). That this definition is in agreement with the standard notion of state familiar to physicists was established by Gleason with his celebrated theorem.

**Theorem 1.1 (Gleason’s theorem)** Let \( H \) be a real or complex separable Hilbert space with dimension greater than 2. The class of \( \sigma \)-additive probability measures \( \mu \) over \( \mathcal{L}(H) \) is one-to-one with the class of self-adjoint, positive, trace-class operators \( T \) with unit trace.
This correspondence is defined by the requirement
\[ \mu(P) = tr(PT), \]
for every choice of \( P \in \mathcal{L}(H) \).

This way, the notion of statistical operator \( T \) (self-adjoint, positive, trace-class operators with unit trace) enters the mathematical formulation of Quantum Theories. However, also states represented by normalized vectors of \( H \), show up. To this end, first observe that a convex combination \( p\mu + q\mu' \) of probability measures (i.e., of statistical operators \( pT + qT' \)) – where \( p + q = 1 \) and \( p, q \in [0, 1] \) by definition of convex combination – is still a probability measure (a statistical operator). Among the class probability measures over \( \mathcal{L}(H) \) there are elements, defining the so-called pure states, which cannot be decomposed into non-trivial convex combinations of probability measures. They are by definition the extremal elements of the convex body of the said measures. It is easy to prove that they are the usual vector states of the standard formulation familiar to physicists, i.e., the associated statistical operators have the form \( \psi \langle \psi | \cdot \rangle \), where \( ||\psi|| = 1 \). The spectral theorem finally proves that a generic statistical operator is always a (generally infinite) convex combinations of pure states, i.e., a quantum mixture or an incoherent superposition of pure states familiar to physicists.

Remark 1.2

(a) It is worth stressing that pure states are unit vectors up to signs in the real formulation, whereas they are unit vectors up to phases in the complex formulation. So that, in particular, real quantum mechanics does not coincide with decomplexified complex quantum mechanics where pure states would be unit vectors up to \( SO(2) \) rotations.

(b) The statement of Gleason’s theorem is also trivially true for \( \text{dim}(H) = 1 \), whereas the constraint \( \text{dim}(H) \neq 2 \) cannot be removed, since well-known counterexamples can be constructed (see, e.g., [Mo18]). There exist several extensions of Gleason’s result towards various directions, e.g., dealing with a lattice of orthogonal projectors of a von Neumann algebra in \( H \) instead of the whole \( \mathcal{L}(H) \), or relaxing \( \sigma \)-additivity, or positivity of \( \mu \) or separability of \( H \). Remarkably, the requirement \( \text{dim}(H) \neq 2 \) or a corresponding constraint in terms of von Neuman algebras of definite type (type-\( I_2 \) is forbidden) survives all extensions. An exhaustive survey on the subject is [Dv93]. In the rest of the paper, we stick to the elementary version represented by the statement given above of Gleason’s theorem.

(c) The most important physical consequence of Gleason theorem is proving that the idea of states viewed from scratch as probability measures over the lattice of elementary is in agreement with the picture, more familiar to physicists, where the building blocks are pure (vector) states and statistical operators are secondary objects. Failure of Gleason’s theorem would imply that there are two inequivalent notions of state in the elementary formulation of Quantum Theories in Hilbert space.
(d) In physically important situations where not all the selfadjoint operators represent observables, as in the presence of superselection rules or gauge groups, the idea of states as probability measures over the restricted sublattice of physically meaningful projectors reveals to be more useful than the standard approach based on vectors states and their quantum mixtures. As is known, in the presence of superselection rules, many different state vectors and statistical operators may contain the same information. Conversely there is exactly one probability measure over the projectors representing the elementary propositions permitted by the superselection rules associated to a class of equivalent statistical operators and vector states. It is possible to extend Gleason’s theorem to encompass these cases in order to classify all possible statistical operators associated to a given probability measure when, for instance, superselection rules affect the theory (e.g., see [Mo18]).

(e) Within the framework based on the notion of lattice $\mathcal{L}(H)$ of elementary propositions and on the idea of quantum states viewed as probability measures, the notion of trace plays a crucial role. It is in fact a general mathematical tool useful to characterize the structure of the probability measures over $\mathcal{L}(H)$ somehow extending the notion of integral over non Boolean algebras. However, it has not an a priori direct physical meaning which would have if reversing the construction as in the formulation more familiar to physicists, where the fundamental objects are pure states and the operation of trace is often introduced with physical motivations to describe loss of coherence or (partial trace) loss of information on subsystems.

All the discussion above concerned the real and the complex Hilbert space cases. Let us pass to focus on the case of a quaternionic Hilbert space which is also permitted as a model of elementary propositions. Quantum states can be defined as probability measures over $\mathcal{L}(H)$ as well. Are these probability measures one-to-one represented by statistical operators? Differently from what erroneously asserted in [Va68, Va07], the formulation of Gleason’s result cannot directly encompass the quaternionic Hilbert space case in the form stated in Theorem [1.1]. This is because,

(A) in quaternionic Hilbert spaces the notion of trace is generally basis dependent as noticed in [Tor95] and [CGJ16], unless its argument is selfadjoint, so that Eq. (1), where $PT$ is generally not selfadjoint, needs further specification. (As a trivial example of basis dependence, using the notation of Sect. 1.2 below, consider the algebra of quaternions as one-dimensional quaternionic Hilbert space with scalar product $\langle q|q' \rangle := \bar{q}q'$. Both 1 and the first imaginary units $i$ are Hilbert bases. The operator $Jq := jq$ for $q \in H$ have different traces: $\langle 1|J1 \rangle = j$ and $\langle i|Ji \rangle = -j$.)

(B) $\text{tr}(PT)$ in Eq. (1) would produce quaternions rather than reals in $[0,1]$ as it should be since $\text{tr}(PT)$ has the meaning of a probability. This is because the cyclic property of the trace $\text{tr}(PT) = \text{tr}(TP)$ generally fails for quaternionic Hilbert spaces so that selfadjointness of $T$ and $P$ is not enough to guarantee $\text{tr}(PT) \in \mathbb{R}$ (whereas it is sufficient in the complex case: $\text{tr}(PT) = \text{tr}(T^*P^*) = \text{tr}(TP) = \text{tr}(PT)$).
(C) A minor issue, already noticed in [MoOp17], also affects the real-Hilbert space case in Varadarajan’s formulation, since the very definition of trace-class (bounded) operator $A : \mathcal{H} \rightarrow \mathcal{H}$ adopted in [Va07] for the three types of Hilbert spaces:

$$\sum_{x \in N} |\langle x |Ax \rangle| < +\infty \quad \text{for every Hilbert basis } N \subset \mathcal{H} \quad (2)$$

is ineffective in the real case. In infinite-dimensional real Hilbert spaces, there are nonvanishing bounded operators $A : \mathcal{H} \rightarrow \mathcal{H}$ such that $\langle z|Az \rangle = 0$ for every $z \in \mathcal{H}$ so that (2) holds, but it is false that

$$\sum_{x \in N} \langle x|A|x \rangle < +\infty \quad \text{for every Hilbert basis } N \subset \mathcal{H}. \quad (3)$$

(It is sufficient to construct $A$ such that $A^* = -A$ and $AA = -I$ which entails that $|A| = I$ so that the left-hand side of (2) diverges for every choice of $N$.) Even if (2) implies that some notion of trace can be defined, we cannot drop requirement (3) since it guarantees that the standard properties of trace-class operators are valid as we shall discuss in Proposition 2.2.

If $\mathcal{H}$ is complex or quaternionic, the counter-example in (C) does not work because (see Exercise 3.21 in [Mo18] for $\mathbb{D} = \mathbb{C}$ and Proposition 2.17 (a) [GMP13] for $\mathbb{D} = \mathbb{H}$) $\langle z|Az \rangle = 0$ for every $z \in \mathcal{H}$ imply $A = 0$ and the problem disappears. In real Hilbert spaces, this result is achieved only if $A = A^*$, it easily follows from $2\langle x|Ay \rangle + 2\langle y|Ax \rangle = \langle x+y|A(x+y) \rangle - \langle x-y|A(x-y) \rangle$. In complex and quaternionic Hilbert spaces actually (3) and (2) are equivalent, as we shall discuss shortly.

In spite of these technical difficulties, the statement in [Va68, Va07] seems physically safe since everything can be re-arranged in physical applications [Va68, Va07] to obtain the expected results. Physical soundness together with the absence of a sufficiently advanced rigorous formulation of quaternionic Hilbert space functional analysis when the book was written (1968) are, in our view, the reasons why these problems did not produce consequences into the physical literature. Furthermore, the central and hard part of Varadarajan’s proof is correct, and the final statements can be fixed into a formulation of Gleason’s theorem, presented as Theorem 3.3, that is valid for the three kinds of Hilbert spaces. This is the main goal of this paper. The key tool to formulate this common version of the theorem is the notion of real trace that will be introduced shortly. After the formulation of the theorem we shall show how this notion is in complete agreement with the standard formalism used by physicists.

1.2 Known results and notations

The real associative division algebra of **quaternions** is $\mathbb{H} = \{a+bi+cj+dk|a, b, c, d \in \mathbb{R}\}$. The three imaginary units $i, j, k$ pairwise anticommute, satisfy the well-known relations $i^2 = j^2 = k^2 = -1$, and $ij = k$ and cyclic permutations of it. The **quaternionic**
conjugation is defined as $a + bi + cj + dk = a - bi - cj - dk$ for $a, b, c, d \in \mathbb{R}$. Notice that $pq = qp$ if $p, q \in \mathbb{H}$. The real part of a quaternion is therefore $Re(q) = \frac{1}{2}(q + \overline{q})$. Notice that $Re(qq') = Re(q'q)$ for $q, q' \in \mathbb{H}$ as it immediately arises from a direct computation. $\mathbb{H}$ is a real Banach space if referring to the natural norm $|a + bi + cj + dk| := \sqrt{a^2 + b^2 + c^2 + d^2}$ that also satisfies $|pq| = |p| |q|$ and $|q| = \sqrt{\overline{q}q} = |\overline{q}|$ if $p, q \in \mathbb{H}$. In particular $\mathbb{H}$ is a unital real $C^*$-algebra, the * operation being the quaternionic conjugation.

We henceforth use the notion of quaternionic vector space and its properties as defined in [GMP13]. It is worth stressing that, in [GMP13, GMP14, GMP17], the multiplication of a vector of a quaternionic Hilbert space $x \in \mathbb{H}$ and a quaternion $q \in \mathbb{H}$ was supposed to act on the right,

$$H \times \mathbb{H} \ni (x, q) \mapsto xq \in H.$$  

In view of non-commutativity of quaternions, this choice is compulsory as soon as one requires that the inner product of a Hilbert space is linear in the right entry, as generally assumed in mathematical physics literature.

**Remark 1.3** Since we want to use a common notation for real, complex, and quaternionic Hilbert spaces, we therefore adopt a right multiplication of scalars and vectors instead of the standard left one, also in real and complex Hilbert spaces. With this choice, linearity is written, if $H$ is a real, complex or quaternionic vector space, as

$$A(xq) = (Ax)q, \quad x \in H, \; q \in \mathbb{R} \text{ or } \mathbb{C} \text{ respectively.}$$

Linear combination of operators are defined according to this convention. However, in case of quaternionic Hilbert spaces, linear combinations of operators are therefore permitted only if the coefficients of the combination are real, since only reals in $\mathbb{H}$ commute with all quaternions. We stick to the standard notation $aA + bB$ to denote linear combination of (right-linear) operators $A, B : H \to H$

$$(aA + bB)(x) := Axa + Bxb,$$

where $x \in H$ and $a, b \in \mathbb{R}$ for both real or quaternionic Hilbert spaces, whereas $a, b \in \mathbb{C}$ for complex Hilbert spaces.

A scalar product $\langle \cdot | \cdot \rangle : H \times H \to \mathbb{H}$ over a quaternionic vector space $H$ is, by definition, $\mathbb{H}$-linear in the right entry $\langle x|yq +zp \rangle = \langle x|y \rangle q + \langle x|z \rangle p$, positive $\langle x|x \rangle \geq 0$ with $x = 0$ if $\langle x|x \rangle = 0$, and Hermitian $\langle x|y \rangle = \overline{\langle y|x \rangle}$, for every $x, y \in H$, $p, q \in \mathbb{H}$. $H$ is a Hilbert space if the norm $||x|| := \sqrt{\langle x|x \rangle}$ makes $H$ complete as a metric space.

We henceforth denote by $\mathbb{D}$ the set of scalars of the Hilbert space $H$ which can be real ($\mathbb{D} = \mathbb{R}$), complex ($\mathbb{D} = \mathbb{C}$), or quaternionic ($\mathbb{D} = \mathbb{H}$). The scalar product is henceforth denoted by $\langle \cdot | \cdot \rangle$ in the three types of Hilbert space.
The maps $\mathbb{D} \ni r \mapsto Re(r) \in \mathbb{R}$ and $\mathbb{D} \ni q \mapsto \overline{q} \in \mathbb{D}$ are therefore defined in the three cases with the standard meaning for $\mathbb{D} = \mathbb{C}$ and reducing to the identity map if $\mathbb{D} = \mathbb{R}$.

If $\mathbb{D} = \mathbb{H}$, the notion of Hilbert basis, orthogonal complement, bounded operator, adjoint operator, operator norm, uniform, strong and weak operator topologies are defined exactly as in real or complex Hilbert spaces and enjoy identical properties with the obvious trivial changes \cite{GMP13}. The (Hermitian) adjoint of an operator $A$ is always denoted by $A^\ast$.

A unitary operator $U : \mathcal{H} \to \mathcal{H}$ is a norm-preserving surjective $\mathbb{D}$-linear operator. In the special case of $\mathbb{D} = \mathbb{C}$, an anti-unitary operator $U : \mathcal{H} \to \mathcal{H}$ is an anti-$\mathbb{C}$-linear norm-preserving surjective operator. By polarization, unitary operators preserve the scalar product, while anti-unitary operators do it up to an overall complex conjugation of the scalar product ($\langle Cx|Cy \rangle = \langle \overline{x}y \rangle$).

If $\mathbb{D} = \mathbb{H}$, the set $\mathfrak{B}(\mathcal{H})$ of bounded operators $T : \mathcal{H} \to \mathcal{H}$ has the natural structure of a unital real $C^\ast$-algebra (the $^\ast$-operation being the Hermitian adjoint), since only real linear combinations of ($\mathbb{H}$-linear) operators are defined. Evidently $\mathfrak{B}(\mathcal{H})$ is also a unital real $C^\ast$-algebra if $\mathbb{D} = \mathbb{R}$, whereas $\mathfrak{B}(\mathcal{H})$ is a proper unital complex $C^\ast$-algebra when $\mathbb{D} = \mathbb{C}$.

$\mathcal{L}(\mathcal{H})$ denotes the set of orthogonal projectors over $\mathcal{H}$, i.e., of operators $P \in \mathfrak{B}(\mathcal{H})$ such that $PP = P$ and $P^\ast = P$. Orthogonal projectors $P$ are one-to-one with the class of all closed subspaces $\mathcal{M}$ of $\mathcal{H}$ through the requirement $P(\mathcal{M}) = \mathcal{M}$. $\mathcal{L}(\mathcal{H})$ is an orthomodular lattice with respect to the partial ordering relation $P \geq Q$ iff $P(\mathcal{H}) \supseteq Q(\mathcal{H})$ and the orthocomplement $P^\perp = I - P$ ($P, Q \in \mathcal{L}(\mathcal{H})$). $\mathcal{L}(\mathcal{H})$ is complete: the supremum $\vee_{j \in J} P_j$ exists in $\mathcal{L}(\mathcal{H})$ for every $\{P_j\}_{j \in J} \subset \mathcal{L}(\mathcal{H})$ (see, e.g., Theorems 7.22 and 7.56 in \cite{Mo18}, the proofs being identical in the three cases).

$T \in \mathfrak{B}(\mathcal{H})$ is said to be positive when $\langle u|Tu \rangle \geq 0$ if $u \in \mathcal{H}$. The polar decomposition $A = U|A|$ of an operator $A \in \mathfrak{B}(\mathcal{H})$, its absolute value $|A| := \sqrt{A^\ast A}$ and the positive squared root $\sqrt{B}$ of a selfadjoint positive operator $B \in \mathfrak{B}(\mathcal{H})$ are defined for quaternionic Hilbert spaces \cite{GMP13} exactly as for real (see, e.g., \cite{MoOp17}) and complex Hilbert spaces (see, e.g., Sect.3.5.2 in \cite{Mo18}): $\sqrt{B}$ is the unique selfadjoint positive operator in $\mathfrak{B}(\mathcal{H})$ such that $\sqrt{B}/\sqrt{B} = B$. In complex and quaternionic Hilbert spaces, positive operators are also selfadjoint, since positivity implies in particular that $\langle x|Ax \rangle = \langle Ax|x \rangle$ and so $\langle x|(A^\ast - A)x \rangle = 0$ for every $x \in \mathcal{H}$ that, in turn, leads to $A^\ast = A$ (see the comment just below (C) above). In real Hilbert spaces, positivity generally does not imply selfadjointness. If $S,T$ are bounded operator, we say that $T \geq S$ whenever $T - S \geq 0$. It should be noticed that in the real Hilbert space case $\geq$ is not a partial ordering relation in $\mathfrak{B}(\mathcal{H})$, the antisymmetry axiom not being satisfied (unless restricting to the set of selfadjoint operators in $\mathfrak{B}(\mathcal{H})$).

The notion of compact quaternionic-linear operator is identical to that in real or complex Hilbert spaces since it uses the real structure only. The set of compact operators will be denoted\footnote{That set was denoted by $\mathfrak{B}_0(\mathcal{H})$ in \cite{GMP14}.} by $\mathfrak{B}_\infty(\mathcal{H})$. With the same proof as for complex Hilbert spaces (Proposition 4.14 and Theorem 4.15 in \cite{Mo18}), it is easy to prove the following proposition irrespective of the nature of $\mathbb{D}$.
Proposition 1.4 Let $H$ be a real, complex or quaternionic Hilbert space. $\mathcal{B}_\infty(H)$ is a subspace and a closed two-sided $^\ast$-ideal of $\mathcal{B}(H)$. Furthermore $T \in \mathcal{B}_\infty(H)$ if and only if $|T| \in \mathcal{B}_\infty(H)$.

The Hilbert decomposition theorem over a Hilbert basis of eigenvectors, for compact self-adjoint operators (Theorem 4.20 in [Mo18] for $\mathbb{D} = \mathbb{C}$ which is also valid for $\mathbb{D} = \mathbb{R}$ as is easy to check) holds true also for $\mathbb{D} = \mathbb{H}$ as proved in Theorem 1.2 in [GMP14].

We specialize the statement to the case of a selfadjoint compact operator:

Proposition 1.5 Let $T^* = T \in \mathcal{B}_\infty(H)$ where $H$ is a real, complex or quaternionic Hilbert space (not necessarily separable). There is a Hilbert basis $N$ of $\text{Ker}(T)^\perp$ made of eigenvectors of $T$ such that $N$ is finite or countably infinite and

$$Tx = \sum_{u \in N} us(u)\langle u|x \rangle \quad \forall x \in H,$$

$s(u) \in \mathbb{R}$ is the eigenvalue of $u$ whose multiplicity (the dimension of the associated eigenspace of $T$) is finite and $\{s(u)\}_{u \in N} = \sigma(T) \setminus \{0\} = \sigma_p(T) \setminus \{0\}$ where, in case of $\mathbb{D} = \mathbb{H}$ the said spectra are interpreted as spherical spectra (see remark below). The series (4) can be re-ordered arbitrarily preserving the sum.

Remark 1.6 Some remarks on the $\mathbb{D} = \mathbb{H}$ case are listed below.

(a) Since quaternions are not commutative, the standard notion of spectrum must be changed in quaternionic Hilbert spaces. As first discussed in [CGSS07, CSS11] and later adapted to the Hilbert space theory in [GMP13] with a finer definition, the relevant quaternionic notion of spectrum does not concern the properties of the operator $T - \lambda I$ when $\lambda$ varies in $\mathbb{R}$ or $\mathbb{C}$, but those of the quadratic operator $\Delta_q(T) := T^2 - T(q + \overline{q}) + |q|^2$ when $q$ varies in $\mathbb{H}$. The spherical spectrum $\sigma_S(T)$ of $T \in \mathcal{B}(H)$ is made of all quaternions $q \in \mathbb{H}$ such that the spherical resolvent $\Delta_q(T)^{-1}$ does not exist in $\mathcal{B}(H)$. It is easy to see that, if $q \in \sigma_S(T)$, then $pqp^{-1} \in \sigma_S(T)$ for every $p \in H$ with $|p| = 1$. This property can be described as the action of $SO(3)$ on the non-real part of $q$ interpreted as an element of $\mathbb{R}^3$ and this is the reason why $\sigma_S(T)$ is called spherical. The generalization to unbounded and not everywhere defined operators is straightforward [GMP13] and the finer decomposition into point spherical spectrum, continuous spherical spectrum and residual spherical spectrum is strictly analogous to the corresponding decomposition of the standard spectrum in real and complex Hilbert spaces. As a matter of fact, it consists of systematically replacing $(T - \lambda I)$ for $\Delta_q(T)$ in every standard definition in real and complex Hilbert spaces (see [GMP13] for details). In particular, the point spherical spectrum $\sigma_{ps}(T)$ is made of the quaternions $q$ such that $\Delta_q(T)$ is not injective. Though apparently unrelated with the analogous definition in real and complex Hilbert spaces, the notion of spherical spectrum is that exploited in the extensions of all results of spectral theory (e.g., see [GMP13, ACK16, GMP17]).
(b) In spite of the different definition of (spherical) spectrum of $T$, the elements $q$ of the spherical point spectrum $\sigma_{Sp}(T)$ turn out to be the eigenvalues of $T$ (Proposition 4.5 in [GMP13]): there exists $u \in \mathbb{H}\setminus \{0\}$ such that $Tu = uq$ and $u$ is called eigenvector of $q$ as usual. As a consequence, the spherical point spectrum of an operator in a quaternionic Hilbert space coincides with the set of eigenvectors of that operator as it happens for the standard point spectrum in real or complex Hilbert spaces.

(c) Though this is generally false in view of the fact that quaternions does not commute, the eigenspaces of a selfadjoint operator $T$ are $\mathbb{H}$-subspaces of $\mathbb{H}$, since the eigenvalues are real (Theorem 4.8 in [GMP13]) and quaternionic combinations of eigenvectors of a given (real) eigenvalue are still eigenvectors with the same eigenvalue.

(d) When $\mathbb{D} = \mathbb{H}$, the spectrum of the operator $T$ in Proposition 1.5 has to be interpreted as its spherical spectrum which is however a subset of $\mathbb{R}$ since $T = T^*$. Furthermore, possibly up to $0$ which however does not play any role in the considered decomposition of $T$, the spherical spectrum of $T$ in Proposition 1.4 coincides with the set of eigenvalues of $T$, in accordance with the remark (b) above.

The spectral theorem for (densely defined, closed, and generally unbounded) normal operators and in particular selfadjoint operators can be formulated in complex Hilbert spaces (see e.g., [Ru91, Sc12, Mo18]) and quaternionic Hilbert spaces [ACK16, GMP17]. For the real case a corresponding theorem is suitable for selfadjoint (generally unbounded) operators (e.g., see [Li03, MoOp17]).

Remark 1.7 The spectral machinery developed in [GMP17] permits one to integrate quaternionic-valued functions with respect to a projector-valued measure (PVM). The problem pointed out in Remark 1.3 (apparently, only $\mathbb{R}$-linear combination of projectors are permitted in quaternionic Hilbert spaces) is not an obstruction. A PVM in a quaternionic Hilbert space is in fact “intertwining” [GMP17], i.e., equipped with a so-called left multiplications $L : \mathbb{H} \ni q \mapsto L_q \in \mathcal{B}(\mathbb{H})$ commuting with the PVM itself and giving rise to an operatorial representation of $\mathbb{H}$. $L$ permits to define quaternionic-linear combinations of projectors of the associated PVM as $L_q P_E + L_q' P_F$ (coinciding to $qP_E + q'P_F$ for $q,q' \in \mathbb{R}$). These $\mathbb{H}$-linear combinations of projectors are the building-blocks used to define the spectral integral of quaternionic-valued functions within the precise formulation of the quaternionic spectral theory established in [GMP17]. This work uses integrals of real-valued functions only, so that the left-multiplication of a PVM does not affect any result (though it exists). However, the full machinery of quantum mechanics exploits the general integration procedure, for instance dealing with continuous representations of Lie groups of symmetries [MoOp17b].

1.3 Structure of the work

The next section is devoted to extend the basic theory of trace-class operators to the case of quaternionic (generally non-separable) Hilbert spaces, focussing in particular on
the notion of real trace that naturally arises when requiring basis independence. The subsequent section proposes a common formulation of Gleason’s theorem relying to the correct part of Varadarajan formulation and replacing the notion of trace with that of real trace. The last section studies how our general formulation of Gleason’s theorem and the notion of real trace is compatible with the standard formalism handled by physicists. We find a total agreement. A final appendix includes the proofs of some technical propositions spread throughout the work.

2 Notions of trace

We introduce here a common notion of trace-class operators including real, complex and quaternionic Hilbert spaces and corresponding notions of traces.

2.1 Trace class operators

The notion of trace-class operator we introduce here is the direct extension to \( \mathbb{R} \) and \( \mathbb{H} \) of the standard notion adopted in complex Hilbert spaces (e.g., see Sect.4.4 of [Mo18]). For the quaternionic case we make use of the result established in [GMP14] in particular Proposition 1.5 and the remark below it.

**Definition 2.1** Let \( H \) be a real, complex or quaternionic Hilbert space. An operator \( T \in \mathcal{B}(H) \) is said to be of trace class if

\[
\sum_{u \in N} (u| |T| u) < +\infty,
\]

for some Hilbert basis \( N \subset H \). \( \mathcal{B}_1(H) \subset \mathcal{B}(H) \) denotes the set of trace-class operators.

Some of the relevant properties of trace-class operators are summarized below. We explicitly omit to discuss the interplay of trace-class operators and Hilbert-Schmidt ones for the sake of simplicity. Furthermore, no reference to the theory of Schatten class operators will be mentioned. We address the reader to [CGJ16] for a recent extension of the theory to quaternionic Hilbert spaces, where a different notion of trace class operators and a full definition of Schatten class of operators in quaternionic Hilbert spaces are proposed referring to a preferred, arbitrarily fixed, anti-selfadjoint unitary operator. No choice of such an operator is made within our approach.

**Proposition 2.2** Let \( H \) be a real, complex or quaternionic Hilbert space. The set \( \mathcal{B}_1(H) \) of trace-class operators enjoys the following properties.

(a) If \( T \in \mathcal{B}_1(H) \), then (5) is valid for every Hilbert basis \( M \subset H \) and \( \sum_{u \in M} (u| |T| u) \) does not depend on \( M \).

(b) \( T \in \mathcal{B}_1(H) \) if and only if both

(i) \( T \in \mathcal{B}_\infty(H) \) and
(ii) the following fact is true
\[ ||T||_1 := \sum_{s \in \sigma_p(|T|)} sd_s < +\infty, \]
where \( \sigma_p(|T|) \subset [0, +\infty) \) is the point-spectrum of the selfadjoint compact operator \( |T| := \sqrt{T^*T} \) (for \( D = \mathbb{H}, \sigma_p(|T|) \) is the spherical point-spectrum of \( |T| \)), and \( d_s = 1, 2, \ldots < +\infty \) is the dimension of the \( s \)-eigenspace of \( |T| \).

(c) If \( T \in \mathfrak{B}_1(H) \), then for every Hilbert basis \( M \subset H \) it holds that
\[ ||T||_1 = \sum_{u \in M} \langle u | T | u \rangle \]

(d) The set \( \mathfrak{B}_1(H) \) has the structure of

(i) \( \mathbb{R} \)-linear subspace of \( \mathfrak{B}(H) \) and also \( \mathbb{C} \)-linear if \( D = \mathbb{C} \),
(ii) two-sided \( * \)-ideal of \( \mathfrak{B}(H) \),
(iii) real Banach space with respect to the norm || \cdot ||_1.

(e) The following facts are true for \( A \in \mathfrak{B}_1(H) \) and \( B \in \mathfrak{B}(H) \).

(i) \( ||AB||_1 \leq ||A||_1||B|| \) and \( ||BA||_1 \leq ||A||_1||B|| \),
(ii) \( ||A||_1 = ||A^*||_1 \),
(iii) \( ||A|| \leq ||A||_1 \).

The set \( \mathfrak{B}_1(H) \) is therefore a real (complex if \( D = \mathbb{C} \)) Banach algebra which is also a \( * \)-algebra and the norm is \( * \)-invariant.

Proof. See Appendix \[A]\.

2.2 Real trace

We can now pass to discuss the notion of trace stating and proving some properties relevant for our final goal. We also include some of the results established in [Tor95] (here extended also to the non-separable case).

**Proposition 2.3** Let \( H \) be a Hilbert space over \( D = \mathbb{R}, \mathbb{C} \), or \( \mathbb{H} \) not necessarily separable. If \( N \subset H \) is a Hilbert basis and \( T \in \mathfrak{B}_1(H) \), then the \( N \)-trace of \( T \),
\[ \text{tr}_N(T) := \sum_{x \in N} \langle x | T | x \rangle \]
is a well defined element of \( D \) satisfying
\[ |\text{tr}_N(T)| \leq ||T||_1. \]
The right-hand side of (6) consists of a sum over a set at most countable of non-vanishing elements. That series absolutely converges and, for the given $N$, can be re-ordered arbitrarily preserving the sum.

The further facts are true.

(a) If $a, b \in \mathbb{R}$ (or $\mathbb{C}$ if $D = \mathbb{C}$) and $A, B \in \mathcal{B}_1(H)$,

(i) $\text{tr}_N(A^*) = \text{tr}_N(A)$,

(ii) $\text{tr}_N(aA + bB) = a\text{tr}_N(A) + b\text{tr}_N(B)$,

(iii) $\text{tr}_N(A) \geq 0$ whenever $A \geq 0$.

(b) If $D = \mathbb{R}$ or $\mathbb{C}$ and $A \in \mathcal{B}_1(H), B \in \mathcal{B}(H)$ then

(i) $\text{tr}_N(A) = \text{tr}_{N'}(A)$, for every pair of Hilbert basis $N, N' \subset H$,

(ii) $\text{tr}_N(AB) = \text{tr}_N(BA)$.

(c) If $D = \mathbb{H}$ and $A \in \mathcal{B}_1(H)$, then the following facts are equivalent.

(i) $\text{tr}_N(A) = \text{tr}_{N'}(A)$ for every pair of Hilbert basis $N, N' \subset H$.

(ii) $A = A^*$.

(d) If $A \in \mathcal{B}_1(H), B \in \mathcal{B}(H)$, the following facts hold,

(i) $\text{Re}(\text{tr}_N(A)) = \text{Re}(\text{tr}_{N'}(A))$ for every pair of Hilbert basis $N, N' \subset H$,

(ii) $\text{Re}(\text{tr}_N(AB)) = \text{Re}(\text{tr}_N(BA))$.

(e) If $A = A^* \in \mathcal{B}_1(H)$ and $N_A$ is a Hilbert basis of $H$ obtained by adding a Hilbert basis of $\text{Ker}(A)$ to a Hilbert basis of $\text{Ker}(A)^\perp$ made of eigenvectors of $A$ as in Proposition 1.5, then

$$\text{tr}_{N_A}(BA) = \text{tr}_{N_A}(AB)$$  \hspace{1cm} (8)

for any $B \in \mathcal{B}_1(H)$. Furthermore, if $B = B^*$, then $\text{tr}_N(AB) \in \mathbb{R}$.

(f) If $A, B \in \mathcal{B}_1(H)$ satisfy $A \geq B$, then

$$\text{Re}(\text{tr}_N(A)) \geq \text{Re}(\text{tr}_N(B))$$

for every Hilbert basis $N \subset H$.

**Remark 2.4** Notice that, although $\text{tr}$ does not satisfy the cyclic property $\text{tr}(AB) = \text{tr}(BA)$ in quaternionic Hilbert spaces, its real part does as stated in (d)(ii).

**Proof.** See Appendix A. \hfill \Box

Proposition 2.3 leads to the following definition.

**Definition 2.5** Let $H$ be a real, complex or quaternionic Hilbert space and $A \in \mathcal{B}_1(H)$. \hfill 13
(a) The real trace of $A$ is defined as 
\[ \text{tr}^R(A) := \text{Re}(\text{tr}_N(A)) \]
where $N \subset H$ is a Hilbert basis.

(b) If $\text{tr}_N(A)$ does not depend on the Hilbert basis $N \subset H$ we call it trace of $A$ and denote it by $\text{tr}(A)$.

Remark 2.6 It should be clear that, according to the given definitions, noticing that $|(|T|)| = |T|$, we have
\[ ||T||_1 = || |T| ||_1 = \text{tr}(|T|) \tag{9} \]
for every $T \in \mathfrak{B}_1(H)$ irrespective of the nature of $D$. This follows immediately from (c) of Proposition 2.2 and selfadjointness of $|T|$.

With the said definition, an immediate corollary of Proposition 2.3 follows.

Corollary 2.7 Let $H$ be a real, complex or quaternionic Hilbert space. Then the following statements hold.

(a) $\text{tr}^R : \mathfrak{B}_1(H) \rightarrow \mathbb{R}$ is $\mathbb{R}$-linear and $\ast$-invariant.

(b) $\text{tr}^R(A) \geq 0$ for $A \in \mathfrak{B}_1(H)$ positive.

(c) $\text{tr}^R(A) \geq \text{tr}^R(B)$ for $A, B \in \mathfrak{B}_1(H)$ satisfying $A \geq B$.

(d) $\text{tr}^R(AB) = \text{tr}^R(BA)$ for $A \in \mathfrak{B}_1(H)$ and $B \in \mathfrak{B}(H)$.

(e) $\text{tr}^R(A) = \text{tr}(A)$ for $A^* = A \in \mathfrak{B}_1(H)$.

(f) If $A^* = A \in \mathfrak{B}_1(H)$ and $B^* = B \in \mathfrak{B}(H)$ then
\begin{enumerate}
  
  \item[(i)] if $D = \mathbb{R}, \mathbb{C}$ then $\text{tr}^R(AB) = \text{tr}(AB)$ for $D = \mathbb{R}, \mathbb{C}$,
  
  \item[(ii)] if $D = \mathbb{H}$, then $\text{tr}^R(AB) = \text{tr}_{N_A}(AB)$ where $N_A$ is the completion of a Hilbert basis of $\text{Ker}(A) \perp$ made of eigenvectors of $A$.
\end{enumerate}

(g) If $A^* = A \in \mathfrak{B}_1(H)$ and $P \in \mathcal{L}(H)$, then $\text{tr}^R(PA) = \text{tr}^R(PAP) = \text{tr}(PAP)$.

Proof. The proofs of (a)-(f) follow easily from Proposition 2.3. To prove (g) just notice that $PAP$ is selfadjoint and that $\text{tr}^R(PA) = \text{tr}^R(PPA) = \text{tr}^R(PAP)$, which follows from $P = PP$. \hfill \qed

To conclude this list of properties of trace-class operators, we stress that, for complex Hilbert spaces, the following well known fact holds (e.g. see Proposition 4.41 in [Mo18]) which we extend here to the quaternionic case.

Proposition 2.8 If $H$ is a complex or quaternionic Hilbert space, the following asserts are equivalent for $A \in \mathfrak{B}(H)$.
(i) $A \in \mathcal{B}_1(H)$,

(ii) $\sum_{x \in N} |\langle x | Ax \rangle| < +\infty$ for every Hilbert basis of $H$.

If $H$ is real, then (i) implies (ii), but the converse implication is generally false in infinite dimensional spaces.

Proof. See Appendix A. \hfill \square

In [Va07] (ii) was used as definition of trace-class operator. We point out that this definition is equivalent to ours in the complex and quaternionic case only.

2.3 A basis-independence property of $tr(A)$ for $\mathbb{D} = \mathbb{H}$.

Contrarily to the real and complex case, in quaternionic Hilbert spaces $tr_N(A)$ may depend on the chosen Hilbert basis $N$ for a fixed $A \in \mathcal{B}_1(H)$. However, an invariance property remains and, once more, it can be stated in terms of real trace. Although this subject is not in the main stream of this work we spend some words about it. The reader interested in the application of previous results to Gleason’s theorem can safely skip this section.

We remind the reader that (Lemma 3.9 in [GMP13]), if $J$ is an anti selfadjoint unitary operator in a quaternionic Hilbert space and we take $i \in \mathbb{H}$ with $|i| = 1$, then the complex vector space $H_{Ji} := \{z \in H | Jz = z i\}$ – using $i$ as imaginary unit – inherits the structure of complex Hilbert space when restricting the scalar product of $H$ to $H_{Ji}$. Furthermore, if $N \subset H_{Ji}$ is a Hilbert basis of $H_{Ji}$, it is also a Hilbert basis of $H$ (Proposition 3.8 in [GMP13]). A Hilbert basis $N$ of $H$ is evidently also a Hilbert basis of $H_{Ji}$, whenever $N \subset H_{Ji}$, considering only $i$-complex combinations of elements of $N$. If $N$ is a Hilbert basis of $H_{Ji}$ and we change $i$ to $i'$, then $N' = \{us | u \in N\}$ is a Hilbert basis of $H_{Ji'}$ if $s^{-1} is = i'$ and $|s| = 1$.

Consider $A \in \mathcal{B}(H)$ for a quaternionic Hilbert space $H$. The anti selfadjoint part $A - A^*$ is normal so, according to Theorem 5.9 in [GMP13] (applied to $T = A - A^*$), its polar decomposition can be improved as follows. There exists $J \in \mathcal{B}(H)$ with

1. $J = -J^*$, $JJ = -I$,$^2$

2. $A - A^* = J|A - A^*|$, 

3. $J$ commutes with both $A - A^*$ and $|A - A^*|$ (not necessarily with $A$).

The polar decomposition theorem implies that $J$ is uniquely defined on $Ker(|A - A^*|)^{\perp} = Ker(A - A^*)^{\perp}$. $^2$

We can now state and prove a proposition regarding a basis-invariance property in the quaternionic case.

$^2$By definition of $|T|$, $||Tx|| = \langle x | T^*Tx \rangle = \langle x | T \rangle^2 x = ||Tx||^2$, so that $Ker(T) = Ker(|T|)$ if $T \in \mathcal{B}(H)$.
Proposition 2.9  Let $H$ be a quaternionic Hilbert space and $A \in \mathcal{B}_1(H)$, let $J$ satisfy (1)-(2) above with respect to $A - A^*$ and take $i \in \mathbb{H}$ with $|i| = 1$. For every Hilbert basis $N$ of $H$ such that $N \subset H_{Ji}$ the following identity holds

$$\text{tr}_N(A) = \text{tr}^R(A) + \frac{i}{2} \text{tr}(|A - A^*|).$$  \hspace{1cm} (10)$$

If $N' \subset H_{Ji}$ is another Hilbert basis of $H$, where $Jx = J'x$ for $x \in \text{Ker}(|A - A^*|)^\perp$, then (10) still holds for $N'$ in place of $N$.

Proof. See Appendix A. \hfill \Box

We stress that other basis-independence properties of the trace exist. Within the different approach of [CGJ16], for every $J = -J^*$ with $JJ = -I$ commuting with $A \in \mathcal{B}_1(H)$, $\text{tr}_N(A)$ is fixed when the Hilbert basis $N$ varies in $H_{Ji}$. (Such a $J$ does exist at least if $A$ is normal for Theorem 5.9 in [GMP13].) Our special choice of $J$ in Proposition 2.9 is always feasible for every $A \in \mathcal{B}_1(H)$ though, in general, $J$ does not commute with $A$ and the invariant trace in $H_{Ji}$ can be written in terms of real traces.

3  A common statement of Gleason’s theorem

We are in a position to establish a common form of Gleason’s theorem, valid for real, complex and quaternionic Hilbert spaces, which uses the real trace instead of the trace appearing in Theorem 1.1.

A convex body is a subset $K \neq \emptyset$ of a real linear space $X$ such that $(1 - \lambda)x + \lambda y \in K$ for every $x, y \in K$ and $\lambda \in [0, 1]$. A point $\omega \in K$ said to be extremal if it cannot be decomposed as $\omega = (1 - \lambda)x + \lambda y$ with $\lambda \in (0, 1)$ and $x, y \in K \setminus \{\omega\}$.

3.1 Probability measures over $\mathcal{L}(H)$

Definition 3.1  Let $H$ be a real, complex or quaternionic Hilbert space. A $\sigma$-additive probability measure over the lattice $\mathcal{L}(H)$ of orthogonal projectors over $H$ is a map $\mu : \mathcal{L}(H) \to [0, 1]$ satisfying both

(i) $\mu(I) = 1$,

(ii) $\sum_{k \in K} \mu(P_k) = \mu(\bigvee_{k \in K} P_k),$

for $\{P_k\}_{k \in K} \subset \mathcal{L}(H)$ with $K$ finite or countably infinite and $P_k \perp P_h$ for $k \neq h$.

We denote by $\mathcal{M}(H)$ the convex body of $\sigma$-additive probability measures over $\mathcal{L}(H)$.

Remark 3.2

(a) Since $K$ is countable and $P_k P_h = 0$ if $k \neq h$, the supremum $\bigvee_{k \in K} P_k$ can always be computed as $\bigvee_{k \in K} P_k x = \sum_{k \in K} P_k x$ for every $x \in H$, the sum converging in the topology of $H$ if $K$ is infinite. The proof is identical in the three Hilbert space cases (see, e.g., [Mo18], Theorem 7.22).
(b) \( \mathcal{M}(H) \) is to be understood as a convex subset of the set of functions \( f : \mathcal{L}(H) \to \mathbb{R} \) which is clearly a real linear space with respect to the operations \( (\alpha f + \beta g)(P) := \alpha f(P) + \beta g(P) \) for all \( P \in \mathcal{L}(H) \).

3.2 Fixing Varadarajan’s statement of Gleason’s theorem for \( D = \mathbb{H} \).

We present and prove the statement of Gleason’s theorem completing the assertion with remarks about extremal measures. A further statement is added on the fact that, in spite of an apparent lack of information embodied in the real trace if compared with the standard trace, probability measures are still able to distinguish between elements of \( \mathcal{L}(H) \).

**Theorem 3.3** Let \( H \) be a real, complex or quaternionic separable Hilbert space with dimension greater than 2. The following facts are true.

(a) The class \( \mathcal{M}(H) \) of \( \sigma \)-additive probability measures \( \mu \) over \( \mathcal{L}(H) \) is one-to-one with the class \( \mathcal{S}(H) \) of self-adjoint, positive, operators \( T \in \mathcal{B}_1(H) \) with unit trace. This correspondence is defined by the requirement

\[
\mu(P) = \text{tr}^\mathbb{R}(PT) \quad \text{for all } P \in \mathcal{L}(H). \tag{11}
\]

That correspondence preserves the real convex structures of \( \mathcal{S}(H) \) and \( \mathcal{M}(H) \).

(b) All orthogonal projectors onto one-dimensional subspaces of \( H \) belong to \( \mathcal{S}(H) \). They are precisely the extremal element of \( \mathcal{S}(H) \) and are one-to-one through (11) with the extremal elements of \( \mathcal{M}(H) \).

(c) As a consequence of (11), probability measures separate the elements of \( \mathcal{L}(H) \) (and evidently the elements of \( \mathcal{S}(H) \) separate the probability measures over \( \mathcal{L}(H) \)).

**Proof.** (a) (i) **Uniqueness of** \( T \). For a given \( \mu \), the associated \( T \in \mathcal{S}(H) \) is unique if exists. Indeed, if \( \text{tr}^\mathbb{R}(PT) = \mu(P) = \text{tr}^\mathbb{R}(PT') \) for all \( P \in \mathcal{L}(H) \), restricting to one-dimensional projectors \( P = \psi\langle \psi| \cdot \rangle \) we have \( \langle \psi|(T - T')\psi\rangle = 0 \) for every vector \( \psi \in H \). Since \( T - T' \) is bounded and selfadjoint we get \( T - T' = 0 \) (see Section 1.1 just below Point (C)).

(a)(ii) **Each** \( T \in \mathcal{S}(H) \) **defines a** \( \sigma \)-**additive probability measure** \( \mu \in \mathcal{M}(H) \). Let us now prove that selfadjoint positive operators \( T \in \mathcal{B}_1(H) \) with unit trace define probability measures (Definition 3.1) through (11). First focus on \( \sigma \)-additivity. Taking Remark 3.2 into account, this amounts to establish that

\[
\text{tr}^\mathbb{R} \left( \left( s- \sum_{k \in K} P_k \right) T \right) = \sum_{k \in K} \text{tr}^\mathbb{R}(P_k T)
\]

for every class of orthogonal projectors \( \{P_k\}_{k \in K} \) with \( K \) finite or countably infinite and \( P_kP_h = 0 \) for \( h \neq k \) and where \( s- \) denotes the strong limit. Assuming \( K = \mathbb{N} \), define the orthogonal projectors \( Q_n := \sum_{k=0}^n P_k \) and \( Q := s-\lim_{n \to +\infty} Q_n = s-\sum_{k \in K} P_k \). What
we have to prove is that \( \text{tr}^R(QT) = \lim_{n \to +\infty} \text{tr}^R(Q_nT) \). Since the real trace is basis independent we can evaluate it on a (countable in our hypotheses) Hilbert basis \( N \) of \( H \) obtained by completing a Hilbert basis of \( Q(H) \) that, in turn, is made of the union of Hilbert bases of each \( P_n(H) \). Exploiting Corollary 2.7(g), we get \( \text{tr}^R(Q_nT) = \text{tr}_N(Q_nTQ_n) \) and similarly \( \text{tr}^R(QT) = \text{tr}_N(QTQ) \). What we have to prove is that

\[
\sum_{u \in N} \langle u|QTQu \rangle = \lim_{n \to +\infty} \sum_{u \in N} \langle u|Q_nTQ_nu \rangle.
\]

(12)

Observe that, for every fixed \( u \in N \), we have \( 0 \leq \langle u|Q_nTQ_nu \rangle \leq \langle u|QTQu \rangle \) (indeed, if \( u \in P_n(H) \) then \( \langle u|Q_nTQ_nu \rangle = 0 \) for \( n < m \) and \( \langle u|Q_nTQ_nu \rangle = \langle u|QTQu \rangle \) for \( n \geq m \). If \( u \in Q(H)^\perp \) both sides vanish). Since \( 0 \leq \langle u|Q_nTQ_nu \rangle \leq \langle u|QTQu \rangle \) and \( \sum_u \langle u|QTQu \rangle < +\infty \), the dominated convergence theorem proves (12) completing the proof of \( \sigma \)-additivity. Let us prove that the constructed measure ranges in \([0,1]\). So, take any \( P \in \mathscr{L}(H) \) and consider a Hilbert basis \( N \) of \( H \) completing a Hilbert basis \( N_P \) of \( P(H) \). Corollary 2.7(g) guarantees that \( \text{tr}^R(PT) = \text{tr}_N(PTP) \). Now, exploiting the positivity of \( T \) and Proposition 2.3(a)-(iii) we get

\[
0 \leq \text{tr}_N(PTP) = \sum_{z \in N} \langle z|PTPz \rangle = \sum_{z \in N_P} \langle z|Tz \rangle \leq \sum_{z \in N} \langle z|Tz \rangle = \text{tr}_N(T) = 1,
\]

so that \( \text{tr}^R(PT) \in [0,1] \). Evidently \( \text{tr}^R(IT) = 1 \). We proved that \( T \) defines a \( \sigma \)-additive probability measure over \( \mathscr{L}(H) \) through (11).

(iii) Every \( \mu \in \mathcal{M}(H) \) can be written as in (11) for some \( T \in \mathscr{L}(H) \). This is the hard part of the proof and we exploit the results established in \([Va07]\) extending the original procedure by Gleason \([Gl57]\). Take a probability measure \( \mu : \mathscr{L}(H) \to [0,1] \). In view of its definition, for every Hilbert basis \( N \), which is at most countable by hypothesis, we have \( 1 = \mu(I) = \sum_{u \in N} \mu(\langle u|\rangle) \) independently from the choice of \( N \). Therefore, the map \( f_\mu : \{x \in H \mid ||x|| = 1\} \to [0,1] \) defined by \( f_\mu(x) := \mu(\langle x|\rangle) \) satisfies \( \sum_{u \in N} f_\mu(u) = 1 \) independently from the choice of \( N \). Such a map is called frame-function with weight 1. As a consequence of Lemma 4.22 in \([Va07]\), since \( f_\mu(x) \geq 0 \) for every \( x \), there is a \( T \in \mathfrak{B}(H) \) such that \( T = T^* \) and \( f_\mu(x) = \langle x|Tx \rangle \). In particular \( T \geq 0 \) because \( f_\mu \geq 0 \). Since \( 1 = \sum_{u \in N} f_\mu(u) \) for every Hilbert basis \( N \), we also have \( 1 = \sum_{u \in N} \langle u|Tu \rangle \) which, from \( T \geq 0 \) can be re-written \( 1 = \sum_{u \in N} \langle u||u| \rangle \). According to Definition 2.1 we have proved that \( T \in \mathfrak{B}(H) \) and also \( tr(T) = 1 \) in view of (b) in Definition 2.5. It remains to prove that (11) holds true. The identity \( f_\mu(x) = \langle x|Tx \rangle \) can be rephrased to \( \mu(\langle x|\rangle x) = \langle x|Tx \rangle \) for every \( x \in H \) with \( ||x|| = 1 \). If \( P \in \mathscr{L}(H) \), let \( M \) be a Hilbert basis of \( P(H) \) so that \( Pz = \sum_{x \in M} x\langle z| x \rangle \) for every \( z \in H \). \( \sigma \)-additivity of \( \mu \) implies \( \mu(P) = \sum_{x \in M} \langle x|Tx \rangle \). Completing \( M \) to a Hilbert basis \( N \) of \( H \), the found identity can be re-written \( \mu(P) = tr_M(PT) \). Taking the real part of both sides, we eventually get \( \mu(P) = \text{Re} \left( \text{tr}_M(PT) \right) = \text{tr}^R(PT) \), where the last term no longer depends on \( M \). The proof of (a) ends: the fact that the correspondence between operators \( T \) and measures \( \mu \) preserves the real convex structures of \( \mathscr{L}(H) \) and \( \mathcal{M}(H) \) follows trivially from (a) in Corollary 2.7.
(b) We identify every $\mu \in \mathcal{M}(H)$ with the corresponding $T \in \mathcal{F}(H)$ according to (a). Consider such a measure, that is $T \in \mathcal{B}_1(H)$ which is selfadjoint positive and $tr(T) = 1$ and suppose that $T$ is not a one-dimensional orthogonal projector. Now consider the spectral decomposition in Proposition 1.5 $T = \sum_{u \in N} us(u)\langle u|\cdot \rangle$. There the finite or countable orthonormal system of vectors $N$ (which is a Hilbert basis of $Ker(T)^\perp$) can be completed to a Hilbert basis of $H$ (by adding a Hilbert space of $Ker(T)$), the positive reals $s(u) \in (0, 1]$ form the point spectrum of $T$ except for, possibly, the zero eigenvalue. Finally they satisfy $\sum_{u \in N} s(u) = 1$. If $T$ is not a one-dimensional orthogonal projector there are at least two different $u$, say $u_1$ and $u_2$ with $s(u_1) > 0$ and $1 - s(u_1) \geq s(u_2) > 0$. As a consequence, $T$ decomposes into the convex decomposition $T = s(u_1)T_1 + (1 - s(u_1))T_2$ for

$$ T_1 = u_1\langle u_1|\cdot \rangle \quad \text{and} \quad T_2 := \sum_{u \neq u_1} \frac{s(u)}{1 - s(u_1)} u\langle u|\cdot \rangle. $$

Notice that (i) $T_1 \neq T_2$, (ii) $T_1, T_2 \neq 0$, (iii) $T_1, T_2 \in \mathcal{B}_1(H)$ by construction, (iv) they are selfadjoint, (v) $T_1, T_2 \geq 0$ and (vi) $tr(T_1) = tr(T_2) = 1$, so $T_1$ and $T_2$ represent two different probability measures over $\mathcal{L}(H)$. We conclude that $T$ cannot be extremal in the convex body of probability measures, because it admits a non-trivial convex decomposition.

To conclude the proof, we prove that a one-dimensional orthogonal projector $P$ does not admit a non-trivial convex decomposition into a pair of different operators $T_1$ and $T_2$ representing probability measures over $\mathcal{L}(H)$. So, suppose that $P = \lambda_1T_1 + \lambda_2T_2$ with $T_1, T_2 \in \mathcal{B}_1(H)$ and $\lambda_1, \lambda_2 \in (0, 1)$ such that $\lambda_1 + \lambda_2 = 1$. Thus $P = P^4P = \lambda_1PT_1P + \lambda_2PT_2P$. As a consequence of (d)(i) Proposition 2.2 $PT_1P, PT_2P \in \mathcal{B}_1(H)$. In particular, if $P = \psi\langle \psi|\cdot \rangle$, it must be

$$ PT_rP = \psi\langle \psi|T_r\psi\langle \psi|\cdot \rangle = \psi\langle \psi|T_r\psi\rangle\langle \psi|\cdot \rangle = \psi q_r\langle \psi|\cdot \rangle, $$

for $q_1, q_2 \geq 0$, the operators $PT_rP$ being selfadjoint and positive. Furthermore, completing $\psi$ to a Hilbert basis of $H$ and exploiting $tr(T_r) = 1$ and the positivity of $T_r$ we see that $q_r = tr(PT_rP) \leq 1$. Finally, taking the trace of both sides of $P = \lambda_1PT_1P + \lambda_2PT_2P$ we have $\sum_i \lambda_iq_i = 1$. Summing up we have four reals $\lambda_1, \lambda_2, q_1, q_2$ such that

$$ \lambda_1, \lambda_2 \in (0, 1), \quad q_1, q_2 \in [0, 1], \quad \sum_{i=1}^2 \lambda_i = 1 \quad \text{and} \quad \sum_{i=1}^2 \lambda_iq_i = 1. $$

At this point we need the following result

**Lemma 3.4** Let $2 \leq N \leq \infty$ and $(p_n)_{n=0}^N \subset (0, 1)$ and $(q_n)_{n=0}^N \subset [0, 1]$ such that

$$ \sum_{n=0}^N p_n = \sum_{n=0}^N p_nq_n = 1, $$

then it holds that $q_n = 1$ for any choice of $n$.  

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Proof. See Appendix A.

Exploiting Lemma 3.4 we get \( q_1 = q_2 = 1 \). Since the operators \( T_r \) are selfadjoint, positive, compact, trace-class and unit-trace, Proposition 1.5 yields \( T_r = \sum_{u \in N_r} s_r(u) \langle u | \psi | \cdot \rangle \) where \( N_r \) is a Hilbert basis of \( Ker(T_r) \perp \), \( s_r(u) > 0 \) and \( \sum_{u \in N_r} s_r(u) = 1 \). Since also \( PT_r P \) is unit-trace, completing \( N_r \) to a Hilbert basis \( N'_r \) of \( H \) and exploiting Corollary 2.7 (g) we get

\[
1 = tr(PT_r P) = tr^R(PT_r) = \sum_{z \in N'_r} Re \langle z | PT_r z \rangle = \sum_{z \in N_r} s_r(u) \langle u | Pu \rangle = \sum_{u \in N_r} s_r(u) \| Pu \|^2.
\]

Since \( 0 \leq \| Pu \|^2 \leq \| u \|^2 = 1 \) and \( \sum_{u \in N_r} s(u) = 1 \) we argue that \( \| Pu \| = 1 \) for all \( u \in N_r \) which is equivalent to \( u \in P(H) \) for all \( u \in N_r \). Since \( P(H) \) is one-dimensional and generated by \( \psi \) and the elements of \( N_r \) are orthogonal to each other, it must be \( N_r = \{ u = \psi p_r \} \) for some \( p_r \in \mathbb{H} \) with \( |p_r| = 1 \). At this point the proof would be complete, since \( T_r = u(u|\cdot) = \psi \langle \psi | \cdot \rangle = P \). So, let us prove that \( \| Pu \| = 1 \) for all \( u \in N_r \).

First, notice that if there exists \( u_0 \in N_r \) such that \( s_r(u_0) = 1 \), then \( \sum_{u \in N_r} s_r(u) = 1 \) and \( s_r(u) > 0 \) forces \( N_r = \{ u_0 \} \) and so \( 1 = s_r(u_0) \| Pu_0 \| = \| Pu_0 \| \). If, instead, all of the elements \( u \in N_r \) satisfy \( 0 < s_r(u) < 1 \), then the thesis follows from Lemma 3.4 (notice that \( N_r \) is at most countable).

(c) (Obviously, \( \mathcal{L}(H) \) separates the convex body of probability measures because they admit \( \mathcal{L}(H) \) as domain.) To prove that probability measures separate orthogonal projectors, suppose that \( \mu(P) = \mu(P') \) for every probability measure \( \mu \). As a consequence of (a), \( tr^R(PT) = tr^R(P'T) \) for every self-adjoint positive \( T \in \mathfrak{B}_1(H) \) with unit trace. Restricting to one-dimensional projectors \( T = \psi \langle \psi | \cdot \rangle \), the written identity specialised to \( \langle \psi | (P - P') \psi \rangle = 0 \) for every vector \( \psi \in \mathbb{H} \). Since \( P - P' \) is bounded, polarization identity proves \( P - P' = 0 \) if \( \mathbb{D} = \mathbb{C}, \mathbb{H} \). The same result arises from polarization identity, for \( \mathbb{D} = \mathbb{R} \), using the fact that \( P - P' \) is also self-adjoint.

\[\text{Remark 3.5}\]

(a) In case \( \mathbb{D} = \mathbb{R} \) or \( \mathbb{C} \), the real trace in (11) can be replaced by the standard trace without affecting the result because of (f) in Corollary 2.7 recovering the known statement of Gleason’s theorem as in Theorem 1.1.

(b) As alternate possibilities in stating the theorem above, we observe that right-hand side of (11) satisfies

\[\text{tr}^R(PT) = \text{tr}^R(PTP) = \text{tr}^R(PTP) = \text{tr}(PTP) = (13)\]

which follows from Corollary 2.7 (f) (notice that \( \text{tr}^R(PT) = \text{tr}^R(TPP) = \text{tr}^R(PTP) \) since \( P = PP \) and the real trace is cyclic).

(c) The sequence of identities (13) proves that the standard trace can be used in place of the real trace in some elementary physical computations involving orthogonal
projectors and statistical operators. A price to pay is that the formulas have to be re-arranged in order to always deal with selfadjoint arguments of the trace as 
\( \mu(P) = \text{tr}(PTP) \) (the other possibility \( \mu(P) = \text{tr}\left[\frac{1}{2}(PT + PT)\right] \) would be a trivial rephrasing of \( \mu(P) = \text{tr}^\mathbb{R}(PT) \)). However, passing to physical objects more complicated than probabilities as the expectation value or the standard deviation of (generally unbounded) observables (see identities (i) and (ii) in (a) and (b) of Proposition 4.1 below), this route would become technically very complicated and unnatural. As an example, replacing \( P \) for a (generally unbounded) selfadjoint operator \( A \), the sequence of identities \( \text{tr}\mathbb{R}(AT) = \text{tr}\mathbb{R}(TA) = \text{tr}\mathbb{R}(ATA) = \text{tr}(ATA) \) would result trivially false. Regardless subtle problems related with domains, the crucial obstruction is that \( AA \neq A \). Furthermore, even starting from \( \mu(P_E^{(A)}) = \text{tr}(P_E^{(A)}TP_E^{(A)}) \) where \( P_E^{(A)} \) is the projection-valued measure of the selfadjoint operator \( A \), the quadratic \( P_E^{(A)} \)-dependence of \( \text{tr}(P_E^{(A)}TP_E^{(A)}) \) makes impossible to take advantage of the spectral integration procedure to achieve results similar to those asserted in Proposition 4.1 below regarding the expectation value and the standard deviation of \( A \).

4 Compatibility with the standard physical formalism

In physical applications in complex Hilbert spaces the standard notion of \textit{trace}, instead of \textit{real trace} is exploited. This section proves that the notion of real trace is completely enough to deal with physical formalism and its use provides a common mathematical tool valid for real, complex and quaternionic formulations.

With the hypotheses of Theorem 3.3, extremal probability measures over \( \mathcal{L}(H) \) are called \textbf{pure states} in the language of physicists. According to the statement of Gleason’s theorem, they are represented by all operators of the form \( T = \psi\langle \psi | \) for any unit vector \( \psi \in H \) fixed up to a scalar \( q \in \mathbb{D} \) with \( |q| = 1 \). This is because, as the reader can immediately prove, that is the general form of orthogonal projectors onto one-dimensional subspaces spanned by unit vectors \( \psi \). The remaining operators \( T \in \mathcal{S}(H) \) describing generic probability measures over \( \mathcal{L}(H) \) according to the theorem above, are called \textbf{mixed states} or also \textbf{statistical operators}. Pure and mixed states are called \textbf{quantum states}.

4.1 Observables

According to the standard terminology of physicists, an \textbf{observable} is a selfadjoint (generally unbounded) operator \( A : D(A) \rightarrow H \), where \( D(A) \) is a dense subspace in the Hilbert space describing a given quantum system. We henceforth assume it to be real, complex or quaternionic. \textbf{Elementary observables} (also called \textbf{elementary propositions}) are in particular \( P \in \mathcal{L}(H) \). \( A \) can be spectrally decomposed

\[
A = \int_{\sigma(A)} \lambda dP^{(A)}(\lambda),
\]
where we have introduced the projector-valued measure (PVM)
\[
\left\{ P_E^{(A)} \right\}_{E \in \mathcal{B}(\sigma(A))} \subset \mathcal{L}(H)
\]
of \(A\) over the Borel \(\sigma\)-algebra \(\mathcal{B}(\sigma(A))\) defined over the spectrum \(\sigma(A)\) of \(A\). That is a (closed) subset of \(\mathbb{R}\) because \(A = A^*\). See, e.g., [Mo18] for the complex case, [MoOp17] for the real case, and [GMP17] for the quaternionic case, noticing that \(\sigma(A)\) should be interpreted as the intersection of the spherical spectrum and a complex slice \(\sigma(A) = \sigma_S(A) \cap \mathbb{C}_i^+\). However as \(A = A^*\), the spherical spectrum is completely included in \(\mathbb{R}\) and this intersection turns out to be a subset of \(\mathbb{R}\) independent of the chosen complex slice \(\mathbb{C}_i \subset \mathbb{H}\). The PVM is uniquely associated with \(A\) and satisfies
\[
P_{\sigma(A)}^{(A)} = I, \quad P_{E \cap F}^{(A)} = P_E^{(A)} P_F^{(A)} \quad \text{and} \quad P_{\cup_{n=1}^{\infty} E_n}^{(A)} x = \sum_{n=1}^{\infty} P_{E_n}^{(A)} x
\]
for any \(x \in H\) and for any \(E, F, E_n \in \mathcal{B}(\sigma(A))\) with the subsets \(E_n\) pairwise disjoint.

If \(A : D(A) \to H\) is selfadjoint and \(f : \mathbb{R} \to \mathbb{R}\) is a Borel-measurable function, the selfadjoint operator
\[
f(A) := \int_{\sigma(A)} f(\lambda) dP^{(A)}(\lambda)
\]
is defined with domain
\[
D(f(A)) = \left\{ x \in H \left| \int_{\sigma(A)} |f(\lambda)|^2 d\mu^{(A)}_x(\lambda) < +\infty \right. \right\},
\]
where we introduced the finite positive Borel measure associated with \(x\) and \(A\):
\[
\mu^{(A)}_x : \mathcal{B}(\sigma(A)) \ni E \mapsto ||P_E^{(A)} x||^2 \in [0, ||x||^2].
\]
Moreover, for any \(x \in D(f(A))\) it holds that
\[
\langle x | f(A) x \rangle = \int_{\sigma(A)} f(\lambda) d\mu^{(A)}_x(\lambda) \quad \text{and} \quad ||f(A) x||^2 = \int_{\sigma(A)} |f(\lambda)|^2 d\mu^{(A)}_x(\lambda).
\]
For \(D = \mathbb{C}\) and \(D = \mathbb{H}\) a larger class of measureable functions \(f\) can be used ([Mo18, GMP17]), but we stick here to the real-valued ones because they are completely sufficient for the rest of the work.

Physically speaking, the Borel sets \(E \in \mathcal{B}(\sigma(A))\) are the outcomes of measurement procedures of \(A\), and \(P_E^{(A)}\) is the elementary propositions corresponding to the statement “the measurement of \(A\) belongs to \(E\)”. More precisely, if \(T\) is a quantum state over the Hilbert space \(H\) satisfying the hypotheses of Theorem 3.3, the natural interpretation of
\[
\mu^{(A)}_T : \mathcal{B}(\sigma(A)) \ni E \mapsto tr^\mathbb{R}(P_E^{(A)} T) \in [0, 1]
\]
is the probability to obtain $E$ after a measurement of $A$ in the quantum state $T \in \mathcal{S}(\mathcal{H})$. In particular, if $T$ is pure, so that $T = \psi\langle\psi|\cdot\rangle$ for some unit vector $x \in \mathcal{H}$, we have
\[
\mu_T^{(A)}(E) = ||P_E^{(A)}\psi||^2 = \mu_\psi^{(A)}(E).
\] (19)

The proof is trivial, just complete $\{\psi\}$ as a Hilbert basis of $\mathcal{H}$ and compute the real trace along that basis. This way, the expectation value of $A$ with respect to the state $T$ can be defined
\[
\langle A \rangle_T := \int_{\sigma(A)} \lambda \, d\mu_T^{(A)}(\lambda),
\]
provided the function $\sigma(A) \ni \lambda \rightarrow \lambda \in \mathbb{R}$ is $\mathcal{L}^1(\sigma(A), \mu_T^{(A)})$. Similarly, the standard deviation is defined as
\[
\Delta A_T := \sqrt{\int_{\sigma(A)} (\lambda - \langle A \rangle_T)^2 \, d\mu_T^{(A)}(\lambda)} = \sqrt{\int_{\sigma(A)} \lambda^2 \, d\mu_T^{(A)}(\lambda) - \langle A \rangle_T^2},
\]
provided $\sigma(A) \ni \lambda \rightarrow \lambda \in \mathbb{R}$ is $\mathcal{L}^2(\sigma(A), \mu_T^{(A)})$. (Notice $\mathcal{L}^2(\sigma(A), \mu_T^{(A)}) \subset \mathcal{L}^1(\sigma(A), \mu_T^{(A)})$ since the measure is finite.)

The next proposition establishes that, with our statement of Gleason’s theorem, the usual formal results handled by physicists (see formulas in (b)-(d) below) are however valid when $\mathbb{D} = \mathbb{R}, \mathbb{C}, \mathbb{H}$. It happens if systematically replacing the standard trace with $tr^\mathbb{R}$ and assuming natural conditions on the states.\(^3\) Referring to domain issues in (b) and (c) below we observe that from (16) we have $D(A^2) \subset D(A) = D(|A|)$.

**Proposition 4.1** Let $\mathcal{H}$ be a real, complex or quaternionic Hilbert space satisfying the hypotheses of Theorem 3.3, $T \in \mathcal{S}(\mathcal{H})$ a quantum state and $A : D(A) \rightarrow \mathcal{H}$, densely defined, an observable (i.e., $A = A^\dagger$). The following facts hold.

(a) $\mu_T^{(A)}$ as in (18) is a well-defined probability measure over $\mathcal{B}(\sigma(A))$.

(b) If $\text{Ran}(T) \subset D(A)$ and $|A|T \in \mathcal{B}_1(\mathcal{H})$ (always valid if $A \in \mathcal{B}(\mathcal{H})$), then
\[(i) \langle A \rangle_T \text{ is defined},
(ii) \langle A \rangle_T = tr^\mathbb{R}(AT).
\]

(c) If $\text{Ran}(T) \subset D(A^2)$ and $|A|T, A^2T \in \mathcal{B}_1(\mathcal{H})$ (always valid if $A \in \mathcal{B}(\mathcal{H})$), then
\[(a) \Delta A_T \text{ is defined},
(b) \Delta A_T = \sqrt{tr^\mathbb{R}(A^2T) - (tr^\mathbb{R}(AT))^2}.
\]

\(^3\)Weaker necessary and sufficient conditions assuring that these formulas are valid can be found in [Mo18] in the complex case, referring to the Hilbert-Schmidt class of the operators we do not consider here.
(d) Assume that \( T = \psi\langle \psi| \) with \(|\psi| = 1 \)

(i) If \( \psi \in D(A) \) then (b) is valid and \( \langle A\rangle_T = \langle \psi|A\psi \rangle \),

(ii) If \( \psi \in D(A^2) \) then (c) is valid and \( \Delta A_T = \sqrt{\langle \psi|A^2\psi \rangle} - \langle \psi|A\psi \rangle^2 \).

**Proof.** (a) Taking the definition of projector valued measure into account, the proof is a trivial re-adaptation of the part (a)(ii) of the proof of Theorem 3.3.

(b)(i) We prove the thesis using the weak hypotheses (1) and (2), since they are automatically true if \( A \in \mathcal{B}(H) \). As already stressed, \( D(|A|) = D(A) \) so \( \text{Ran}(T) \subset D(A) = D(|A|) \) is true and both \( AT, |A|T \) are well defined with the said hypotheses. Next the polar decomposition theorem for (generally unbounded) self-adjoint operators gives \( A = U|A| \) with \( |A| \) and \( U := \text{sign}(A) \in \mathcal{B}(H) \) defined as in [13]). As a consequence, \( AT = U|A|T \in \mathcal{B}_1(H) \) because \( U \in \mathcal{B}(H) \) and \( \mathcal{B}_1(H) \) is two-sided ideal. Now, referring to the Borel \( \sigma \)-algebra over \( \sigma(A) \subset \mathbb{R} \) we can construct [Mo18] a sequence of real simple functions

\[
s_n = \sum_{i_n \in I_n} c_{i_n}^{(n)} \chi_{E_{i_n}^{(n)}} : \sigma(A) \to \mathbb{R} \quad \text{with } c_{i_n}^{(n)} \in \mathbb{R}, \text{ and } I_n \text{ finite}
\]

which satisfies

\[
0 \leq |s_n| \leq |s_{n+1}| \leq |id|, \quad s_n \to id \quad \text{point-wise for } n \to +\infty,
\]

where \( id : \sigma(A) \ni \lambda \mapsto \lambda \in \mathbb{R} \). By direct application of the given definitions, if

\[
A_n := \int_{\sigma(A)} s_n dP^{(A)} = \sum_{i_n \in I_n} c_{i_n}^{(n)} P_{E_{i_n}^{(n)}}^{(A)} \in \mathcal{B}(H),
\]

exploiting [17], the monotone and the dominated convergence theorems, we have both

\[
\langle \psi|A_n\psi \rangle \to \langle \psi|A\psi \rangle, \quad \langle \psi||A_n|\psi \rangle \to \langle \psi||A|\psi \rangle \quad \forall \psi \in D(A) \quad \text{as } n \to +\infty
\]

and also

\[
|\langle \psi|A_n\psi \rangle| \leq \langle \psi||A_n|\psi \rangle \leq \langle \psi||A|\psi \rangle.
\]

On the other hand, if \( M \) is a Hilbert basis of \( H \) obtained by completing a Hilbert basis \( N \) of \( \text{Ker}(T) \dagger \) made of eigenvectors of \( T \) according to Proposition 1.3 exploiting Corollary 2.7 (f)(ii) and Proposition 2.3 (e) we have both

\[
\text{tr}_M(A_n T) = \text{tr}_M \left( \sum_{i_n \in I_n} c_{i_n}^{(n)} P_{E_{i_n}^{(n)}}^{(A)} T \right) = \sum_{i_n \in I_n} c_{i_n}^{(n)} \text{tr}_M(T P_{E_{i_n}^{(n)}}^{(A)}) = \sum_{i_n \in I_n} c_{i_n}^{(n)} \text{tr}_M(T P_{E_{i_n}^{(n)}}^{(A)}) = \sum_{i_n \in I_n} c_{i_n}^{(n)} \text{tr}_M(T P_{E_{i_n}^{(n)}}^{(A)}) = \sum_{i_n \in I_n} c_{i_n}^{(n)} \mu_T(E_{i_n}^{(n)}) = \int_{\sigma(A)} s_n d\mu_T^{(A)}
\]

(23)
and similarly
\[
tr_M(|A_n|T) = \int_{\sigma(A)} |s_n| \, d\mu_T^{(A)}.
\] (24)

Looking at the identity (24), by monotone convergence theorem, for \(n \to +\infty\),
\[
tr_M(|A_n|T) = \int_{\sigma(A)} |s_n| \, d\mu_T^{(A)}(\lambda) \to \int_{\sigma(A)} |\lambda| \, d\mu_T^{(A)}(\lambda),
\]
and simultaneously we have
\[
tr_M(|A_n|T) = \sum_{u \in N} s(u)\langle u|A_n|u \rangle \to \sum_{u \in N} s(u)\langle u|Au \rangle = tr_M(|A|T),
\]
where \(s(u) \geq 0\) are the eigenvalues of \(T\), again by monotone convergence theorem and (21). Putting all together and taking the real part, we get
\[
tr^\mathbb{R}(|A|T) = \int_{\sigma(A)} |\lambda| \, d\mu_T^{(A)}(\lambda).
\]
We have in particular established that the integral in the right-hand side is finite (because the left-hand side exists by hypothesis) and thus \(\langle A \rangle_T\) is well defined.

(ii) Let us look at the identity in (23). From the dominated convergence theorem taking (20) into account, we obtain for \(n \to \infty\)
\[
tr_M(A_nT) = \int_{\sigma(A)} s_n(\lambda) \, d\mu_T^{(A)} \to \int_{\sigma(A)} \lambda \, d\mu_T^{(A)}.
\]
On the other hand,
\[
tr_M(A_nT) = \sum_{u \in N} \langle u|A_nu \rangle s(u) \to \sum_{u \in N} \langle u|Au \rangle s(u) = tr_M(AT),
\]
where we have once again applied the dominated convergence theorem as is permitted by (22). Putting all together and taking the real part we get
\[
tr^\mathbb{R}(AT) = \int_{\sigma(A)} \lambda \, d\mu_T^{(A)}(\lambda) =: \langle A \rangle_T,
\]
concluding the proof of (ii).

(c) The proof is strictly analogous to that of (b) also noticing that the hypotheses of (c) implies those of (b) and that \( L^2(\sigma(A), \mu_T^{(A)}) \subset L^1(\sigma(A), \mu_T^{(A)}) \) because \(\mu_T^{(A)}\) is finite.

(d) The thesis consists of trivial subcases of (b) and (c) in particular completing \(\{\psi\}\) to a Hilbert basis of \(\mathcal{H}\) to be used to explicitly compute the various traces. \(\square\)
4.2 Symmetries

Symmetries (including time evolution) of a quantum system described on a Hilbert space $\mathcal{H}$ can be represented in terms of various transformations of mathematical structures entering the game (see [Lan17, Mo18] for exhaustive surveys when $\mathbb{D} = \mathbb{C}$). We only say that, when the set of elementary observables consists of the whole $\mathcal{L}(\mathcal{H})$ (absence of superselection rules and gauge group), symmetries are represented by unitary operators $U : \mathcal{H} \to \mathcal{H}$ defined up to signs for $\mathbb{D} = \mathbb{R}$, $\mathcal{H}$; they are represented by either unitary or anti-unitary operators $U : \mathcal{H} \to \mathcal{H}$ defined up to phases for $\mathbb{D} = \mathbb{C}$ (see [Va07]). Symmetries act on quantum states in a standard way: $T \mapsto UTU^{-1}$, where it is easy to prove that $UTU^{-1}$ is still a quantum state if $T$ is. It is more strongly evident that, for a fixed $U$, the map $T \to UTU^{-1}$ defines an automorphism of the space of the quantum states $\mathcal{S}(\mathcal{H})$ viewed as a convex body in the real vector space $\mathfrak{B}_1(\mathcal{H})$.

For $\mathbb{D} = \mathbb{R}$ or $\mathbb{C}$ that action on states has a corresponding dual action on observables according to the real and complex formulation of Gleason’s theorem as presented in Theorem 1.1. There, the trace is cyclic so that

$$tr(P UTU^{-1}) = tr(U^{-1} PU T)$$

and, evidently, $U^{-1} PU \in \mathcal{L}(\mathcal{H})$ if $P \in \mathcal{L}(\mathcal{H})$. It is easy to see that, for a fixed $U$, this map defines an automorphism of the orthocomplemented lattice $\mathcal{L}(\mathcal{H})$ (all this surely holds for the unitary case and it can be proved to hold also in the antiunitary case for $\mathbb{D} = \mathbb{C}$).

Summing up, symmetries induced by (anti) unitaries $U : \mathcal{H} \to \mathcal{H}$ can be viewed as automorphisms of the real convex body of the states or, alternatively, of the orthomodular lattice of (elementary) observables:

$$\mathcal{S}(\mathcal{H}) \ni T \mapsto UTU^{-1} \in \mathcal{S}(\mathcal{H}), \quad \mathcal{L}(\mathcal{H}) \ni P \mapsto U^{-1} PU \in \mathcal{L}(\mathcal{H}).$$

and this duality interplay corresponds to the physical fact that, looking at measurements, the result of the action of a symmetry on a state can be obtained by a corresponding dual action on observables keeping fixed the state. Everything is encoded in (25), and that identity is responsible for several crucial theoretical tools in theoretical physics, like the duality between Schrödinger picture and Heisenberg picture dealing with time evolution. All that holds for $\mathbb{D} = \mathbb{R}, \mathbb{C}$. Due to (f)(i) in Corollary 2.7 with this choice of $\mathbb{D}$, (25) can be equivalently re-written without loss of information replacing $tr$ for $tr^\mathbb{R}$, since $P$, $UTU^{-1}$ and $T$, $U^{-1} PU$ are pairs of self-adjoint operators.

The next straightforward result states that (25) survives the extension to the quaternionic formulation, provided $tr$ is replaced for $tr^\mathbb{R}$ according to Theorem 3.3. More generally:

**Proposition 4.2** Let $\mathcal{H}$ be a real, complex or quaternionic Hilbert space and $A \in \mathfrak{B}_1(\mathcal{H})$, $B \in \mathfrak{B}(\mathcal{H})$ or vice versa. It holds that

$$tr^\mathbb{R}(A UBU^{-1}) = tr^\mathbb{R}(U^{-1} AU B)$$

(26)
for every operator $U : H \to H$ which is unitary if $D = \mathbb{R}, \mathbb{H}$ or indifferently unitary or anti-unitary if $D = \mathbb{C}$.

In particular, the thesis is true if $A = P \in \mathcal{L}(H)$ and $B = T \in \mathcal{F}(H)$.

Proof. This is an immediate consequence of (d)(ii) Proposition 2.2 and (d)(ii) Proposition 2.3.

As the last result, we prove that the action of a continuous symmetry makes the probabilities computed through $tr^R$ continuous as well, as it is expected from physics.

Proposition 4.3 Let $H$ be a real, complex or quaternionic Hilbert space, $A \in \mathfrak{B}(H)$ and $B^* = B \in \mathfrak{B}_1(H)$ with $B \geq 0$ and $G \ni g \mapsto U_g$ a strongly continuous unitary representation of the topological group $G$. Then the map

$$G \ni g \mapsto tr^R(A U_g BU_g^{-1})$$

is continuous. In particular the thesis is true if $A = P \in \mathcal{L}(H)$ and $B = T \in \mathcal{F}(H)$.

Proof. We know that $tr^R(A U_g BU_g^{-1}) = tr^R(U_g^{-1} A U_g B)$ so we prove the thesis in this second form. It is clear that, as $U_g A U_g^{-1} \in \mathfrak{B}(H)$ if $A \in \mathfrak{B}(H)$ and $U_{gg'} = U_g U_g'$, continuity at $g$ is equivalent to continuity at the neutral element $e$ of $G$. Let us prove it. If $N \subset H$ is a Hilbert basis, $tr^R(U_g^{-1} A U_g B) = Re\left(tr_N(U_g^{-1} A U_g B)\right)$. To conclude, it is enough proving that the right-hand side tends to $Re\left(tr_N(U_e^{-1} A U_e B)\right) = Re\left(tr_N(A B)\right)$ as $g \to e$ with a suitable choice of the basis $N$. It is convenient to chose $N$ as the basis of eigenvectors $u$ of $B$ that exists for Proposition 1.5 (completing the basis of $Ker(B)^\perp$ by adding a Hilbert basis of $Ker(B)$). Non-vanishing eigenvalues are real numbers $s(u) > 0$ (because $B \geq 0$) and $tr_N(U_g^{-1} A U_g B) = \sum_u (u|U_g^{-1} A U_g u)s_u = \sum u||AU_g u||^2$. Since $B = |B|$, we have $\sum s(u) = ||B||_1 < +\infty$ from (3) and $||AU_g u|| \leq ||A|| ||U_g|| ||u|| \leq ||A||$, the dominated convergence theorem proves that $tr_N(U_g^{-1} A U_g B) \to tr_N(U_e^{-1} A U_e B)$ as $g \to e$, concluding the proof using the fact that $Re : \mathbb{D} \to \mathbb{R}$ is continuous.

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A Proof of some propositions

Proof of Proposition 2.2 In the complex case all statements are established, e.g., in Theorems 4.31 and 4.34, including Remark 4.35, of [Mo18]. The proofs for the real case

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4(Continuous) groups of symmetries are more generally described in terms of unitary projective representations due to the "phase" ambiguity in associating unitaries to symmetries (see, e.g., Ch.12 of [Mo18] for the complex case and [Va07] for the general case), however we only stick here to the more elementary case.

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are identical. What we need to prove is that (a)-(e) are valid for the quaternionic case, too. The route we follow is a reduction procedure. Observe that, if \((H, \langle | \rangle)\) is a quaternionic Hilbert space, the set \(H\) becomes a real Hilbert space when considering only real linear combinations among the quaternionic ones and making use of the real symmetric scalar product \((x|y) := \text{Re}\langle x|y \rangle\) for \(x, y \in H\). In particular, completeness survives this change of viewpoint because \(|x| = \sqrt{\langle x|x \rangle} = \sqrt{|x|^2}\) if \(x \in H\). In the rest of the proof, the real afore-mentioned Hilbert space constructed out of the quaternionic Hilbert space \(H\) will be denoted by \(H_{\mathbb{R}}\). We stress that as sets \(H = H_{\mathbb{R}}\), the difference stays in the structures over these sets. The following lemma will be useful.

**Lemma A.1** Let \((H, \langle | \rangle)\) be a quaternionic Hilbert space and \((H_{\mathbb{R}}, \langle | \rangle)\) the associated real Hilbert space as said above. The following facts are true.

(A.1a) \(\mathcal{B}(H) \subset \mathcal{B}(H_{\mathbb{R}})\) is exactly made of the \(\mathbb{R}\)-linear operators in \(\mathcal{B}(H_{\mathbb{R}})\) commuting with the three \(\mathbb{R}\)-linear maps \(H_{\mathbb{R}} \ni x \mapsto xi \in H_{\mathbb{R}}, H_{\mathbb{R}} \ni x \mapsto xj \in H_{\mathbb{R}}\) and \(H_{\mathbb{R}} \ni x \mapsto xk \in H_{\mathbb{R}}\).

(A.1b) if \(N \subset H\) is a quaternionic Hilbert basis, then the associated set

\[
N_{\mathbb{R}} := \{u, ui, uj, uk \mid u \in N\}
\]

is a real Hilbert basis of \(H_{\mathbb{R}}\).

(A.1c) If \(A \in \mathcal{B}(H)\), the adjoint \(A^* \in \mathcal{B}(H)\) is also the adjoint in \(\mathcal{B}(H_{\mathbb{R}})\).

(A.1d) If \(\mathcal{B}(H) \ni A \geq 0\) then \(A^* = A \geq 0\) as operator in \(\mathcal{B}(H_{\mathbb{R}})\) and the squared root \(\sqrt{A} \in \mathcal{B}(H)\) coincides with that computed in \(\mathcal{B}(H_{\mathbb{R}})\).

(A.1e) If \(A \in \mathcal{B}(H)\), \(|A| \in \mathcal{B}(H)\) coincides with the absolute value computed in \(\mathcal{B}(H_{\mathbb{R}})\).

(A.1f) If \(A \in \mathcal{B}(H)\), \(A \in \mathcal{B}_{\infty}(H)\) if and only if \(A \in \mathcal{B}_{\infty}(H_{\mathbb{R}})\).

**Proof.** (A.1a) is evident. (A.1b) is true because \(N_{\mathbb{R}}\) is \(|\cdot|\)-orthonormal and

\[
|\langle x | x \rangle|^2 = \sum_{u \in N} |\langle u | x \rangle|^2 = \sum_{u \in N} |(u| x)|^2 + |(u| x)|^2 + |(u| x)|^2 + |(u| x)|^2 \quad \text{if} \quad x \in H = H_{\mathbb{R}}.
\]

(A.1c) If \(A \in \mathcal{B}(H)\), the adjoint \(A^* \in \mathcal{B}(H)\) is also the adjoint in \(\mathcal{B}(H_{\mathbb{R}})\) because \(\langle A^* x | y \rangle = \langle x | Ay \rangle\) for \(x, y \in H\) implies \(\langle A^* x | y \rangle = \langle x | Ay \rangle\) for \(x, y \in H_{\mathbb{R}} = H\) and this identity completely defines the adjoint operators for elements of \(\mathcal{B}(H_{\mathbb{R}})\).

(A.1d) and (A.1e) If \(\mathcal{B}(H) \ni A \geq 0\) then \(A^* = A\). Moreover, since \(A \geq 0\), it holds in particular that \((x|Ax) \in \mathbb{R}\) and so \((x|Ax) = \text{Re}\langle x|Ax \rangle = \langle x|Ax \rangle \geq 0\) for all \(x \in H = H_{\mathbb{R}}\), so \(A^* = A \geq 0\) in \(\mathcal{B}(H_{\mathbb{R}})\) and, there, \(A\) admits positive squared root. By uniqueness of the positive squared root, \(\sqrt{A} \in \mathcal{B}(H)\) of (self-adjoint) positive operators coincides with that computed in \(\mathcal{B}(H_{\mathbb{R}})\). Consequently, \(|A| = \sqrt{A^* A} \in \mathcal{B}(H)\) coincides with the absolute value computed in \(\mathcal{B}(H_{\mathbb{R}})\).
from (27) and (c), the sequence is also Cauchy in $B$. We finally observe that for $A \in \mathfrak{B}(H)$, identity (27) and the validity of (a) in the real case prove $A \in \mathfrak{B}_1(H)$ if and only if $A \in \mathfrak{B}_1(H_{\mathbb{R}})$.

(a) (For $H$ quaternionic.) Suppose that (5) is true. Passing to the Hilbert basis $N_{\mathbb{R}}$ of $H_{\mathbb{R}}$ and using (A,1a) and (A,1b) we have

$$\sum_{u \in N} \langle u || T || u \rangle = \frac{1}{4} \sum_{v \in N_{\mathbb{R}}} \langle v || T || v \rangle .$$

Since (a) is valid in real Hilbert spaces, (5) must be in particular valid for any other basis $N'_{\mathbb{R}}$ constructed out of a Hilbert basis $N'$ of $H$ and

$$+ \infty > \sum_{u \in N} \langle u || T || u \rangle = \frac{1}{4} \sum_{v \in N_{\mathbb{R}}} \langle v || T || v \rangle = \frac{1}{4} \sum_{v' \in N'_{\mathbb{R}}} \langle v' || T || v' \rangle = \sum_{u' \in N} \langle u' || T || u' \rangle .$$

We finally observe that for $A \in \mathfrak{B}(H)$, identity (27) and the validity of (a) in the real case prove $A \in \mathfrak{B}_1(H)$ if and only if $A \in \mathfrak{B}_1(H_{\mathbb{R}})$.

(b) (For $H$ quaternionic.) Suppose that $A \in \mathfrak{B}_1(H)$, we prove that (i) and (ii) are true. The proof of (a) implies $A \in \mathfrak{B}_1(H_{\mathbb{R}})$, so that, since (b) is valid in the real case, $A \in \mathfrak{B}_\infty(H_{\mathbb{R}})$. Finally (A,1c) implies that $A \in \mathfrak{B}_\infty(H)$ and (i) holds. The validity of (ii) immediately follows from Proposition 1.5 since $|A|$ is compact because $A$ is (Proposition 1.4). To conclude, let us prove that (i),(ii) imply $A \in \mathfrak{B}_1(H)$. This immediately arises by applying Proposition 1.5 ($|A|$ is compact because $A$ is) and proving that (5) is true for a Hilbert basis of $H$ which completes the Hilbert basis of $Ker(A)^{\perp}$ existing for the said proposition. The proof of (b) is over.

(c) (For $H$ quaternionic.) Again, the identity arises by direct application of Proposition 1.5 and (a), (b).

(d) (For $H$ quaternionic.) (i) Suppose that $a,b \in \mathbb{R}$ and $A,B \in \mathfrak{B}_1(H)$. $A,B \in \mathfrak{B}_1(H_{\mathbb{R}})$ form the last assertion in the proof of (a) above. Since (i) holds true in the real case, $aA + bB \in \mathfrak{B}_1(H_{\mathbb{R}})$. The last statement in the proof of (a) proves that $aA + bB \in \mathfrak{B}_1(H)$. (ii) can be proved similarly exploiting (A,1d) and the validity of (ii) in the real case. (iii) We already know that $\mathfrak{B}_1(H_{\mathbb{R}})$ is Banach referring to the corresponding norm here denoted by $|| \ ||_{\mathbb{R}}$. Consider a Cauchy sequence $\{A\}_{n \in \mathbb{N}} \subset \mathfrak{B}_1(\mathbb{R}) \subset \mathfrak{B}_1(H_{\mathbb{R}})$. Since $||A||_{\mathbb{R}} = \frac{1}{4} ||A||_1^{\mathbb{R}}$ from (27) and (c), the sequence is also Cauchy in $\mathfrak{B}_1(H_{\mathbb{R}})$ and thus there is $A \in \mathfrak{B}_1(H_{\mathbb{R}})$ such that $A_n \to A$ in the norm $|| \ ||_{\mathbb{R}}$. However, from (e)(iii) whose proof does not depend on this argument, the convergence of the sequence is also valid referring to the operator norm $|| \ ||$ which is the same for the real and the quaternionic Hilbert space. In particular is also holds point-wise. Since every $A_n$ in $\mathfrak{B}(H_{\mathbb{R}})$ commutes with the three $\mathbb{R}$-linear maps $H_{\mathbb{R}} \ni x \mapsto xi \in H_{\mathbb{R}}$, $H_{\mathbb{R}} \ni x \mapsto xj \in H_{\mathbb{R}}$ and $H_{\mathbb{R}} \ni x \mapsto xk \in H_{\mathbb{R}}$, also $A$ does and thus $A \in \mathfrak{B}_1(H)$ for (A,1a) and the last assertion in the proof of (a) above. This concludes the proof of (d).
(e) (For $\mathbb{H}$ quaternionic.) The proof of (i) and (ii) immediately arises with the same argument exploited proving (d), using again $||A||_1 = \frac{1}{4}||A||$, the fact that (i) and (ii) are valid in the real Hilbert space case and the invariance of the norm $|||\cdot|||$ when passing from $\mathfrak{B}(H)$ to $\mathfrak{B}(H_\mathbb{R})$. Let us conclude the proof by establishing the validity of (iii). Take $A \in \mathfrak{B}_1(H)$. Since $||Ax||^2 = \langle x,A^*Ax \rangle = |||A|||^2$ so that $||A|| = |||A|||$, proving the thesis amounts to prove that $|||A||| \leq ||A||_1$. Since $|A|$ is compact selfadjoint there is a Hilbert basis $M$ of $H$ made of eigenvectors $u$ with eigenvalues $s(u) \geq 0$ (since $|A| \geq 0$) according to Proposition 1.5 (it is sufficient to complete the basis $N$ introduced there). If $x \in H$, $x = \sum_{u \in M} u x_u$ for quaternions $x_u \in \mathbb{H}$. Therefore

$$||A||^2 = ||A||^2 = \sup_{||x||^2} ||A|| x ||^2 = \sup_{||x||^2} \left( \sum_{v \in M} |v x_v| A |v x_v| \right) = \sup_{||x||^2} \sum_{u \in M} |x_u|^2 s(u)^2.$$ 

Since $1 = ||x||^2 = \sum_{u \in M} |x_u|^2$, it must hold $|x_u| \leq 1$ and we can write

$$||A||^2 \leq \sup_{||x||^2} \sum_{u} s(u)^2 = \sum_{u} s(u)^2 \leq \left( \sum_{u} s(u) \right)^2 = |||A|||^2,$$

where we exploited $s(u) \geq 0$ and we made use of (c) in evaluating $||A||_1$ with respect to the said basis $M$ of eigenvectors of $|A|$. \hfill \Box

**Proof of Proposition 2.3.** Consider a Hilbert basis $N \subset H$ and a finite subset $F \subset N$. We have, for $A \in \mathfrak{B}_1(H)$ using its polar decomposition $A = U |A|$, 

$$\sum_{u \in F} |\langle u |Au \rangle| = \sum_{u \in F} |\langle u |U |A| U^* u \rangle| = \sum_{u \in F} |\langle U |A| U^* u \rangle| \leq \sum_{u \in F} |||A||| |\sqrt{A}| u |||\sqrt{A}| u || |.$$ 

$$\leq \sqrt{\sum_{u \in F} |||A||| |\sqrt{A}| u |||} \sqrt{\sum_{u \in F} |||A||| |\sqrt{A}| u |||} = \sqrt{\sum_{u \in F} |\langle u |A| U^* u \rangle|} \sqrt{\sum_{u \in F} |\langle u |A| U^* u \rangle|}.$$ 

Notice that $U |A| U^*$ is selfadjoint and positive, furthermore it belongs to $\mathfrak{B}_1(H)$ in view of (d) Proposition 2.2. In summary, we have found that

$$\sum_{u \in F} |\langle u |Au \rangle| \leq \sqrt{\sum_{u \in F} |\langle u |A| U^* u \rangle|} \sqrt{\sum_{u \in F} |\langle u |A| U^* u \rangle|} \leq \sqrt{|||A|||} \sqrt{|||A|||}.$$ 

Taking advantage of (e) Proposition 2.2, noticing that $||U|| = ||U^*|| \leq 1$ since $U$ is a partial isometry, we have

$$\sum_{u \in F} |\langle u |Au \rangle| \leq |||A||| < +\infty.$$ 

Since this result holds true for every finite subset $F \subset N$, we conclude that

$$\sum_{u \in N} |\langle u |Au \rangle| := \sup \left\{ \sum_{u \in F} |\langle u |Au \rangle| \bigg| F \text{ finite } \subset N \right\} \leq |||A||| < +\infty,$$

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Thus only an countably at most number of elements $|\langle u | Au \rangle|$ do not vanish and the sum can computed as a standard series. As a consequence, the series $\sum_{u \in \mathcal{N}} \langle u | Au \rangle$ in $\mathbb{D}$ absolutely converges and therefore can be re-ordered arbitrarily without changing its sum.

(a) (i), (ii) and (iii) are easy consequences of the given definitions.

(b) (i) and (ii) are standard results. E.g., see Proposition 4.38 [Mo18] for the complex case and the real case can be identically proved.

(c) Suppose that (i) holds. Define $\mathcal{N}'$ out of $\mathcal{N}$ just replacing the element $u_0 \in \mathcal{N}$ for $u_0q \in \mathcal{N}'$ for a given $q \in \mathbb{H}$ with $|q| = 1$ and keeping all remaining elements. By direct computation $0 = tr_{\mathcal{N}'}(A) - tr_{\mathcal{N}}(A) = \overline{q}\langle u_0 | Au_0 \rangle q - \langle u_0 | Au_0 \rangle$. Namely, $\langle u_0 | Au_0 \rangle q = q\langle u_0 | Au_0 \rangle$. Since $q$ is arbitrary, we have $\langle u_0 | Au_0 \rangle \in \mathbb{R}$. Every unit vector $u_0 \in \mathbb{H}$ can be completed to a Hilbert basis and therefore $\langle x | Ax \rangle = \overline{\langle x | Ax \rangle} = \langle Ax | x \rangle$ for $x \in \mathbb{H}$. In summary $\langle x | (A - A^*)x \rangle = 0$ for every $x \in \mathbb{H}$. By polarization $A = A^*$. Suppose conversely that $A^* = A \in \mathcal{B}_{1}(\mathbb{H})$. Since $A \in \mathcal{B}_{\infty}(\mathbb{H})$ for (b) Proposition 2.2, we can decompose $A$ along a Hilbert basis of eigenvectors $M$ (completing the Hilbert basis $\mathcal{N} \subset \text{Ker}(A)^\perp$ in Proposition 1.5),

$$A = \sum_{u \in \mathcal{M}} us(u)\langle u |$$,  \hspace{1cm} (28)

and in particular using $M$ to compute the trace, with $s(u) = 0$ if $u \in \text{ker}(A)$:

$$tr_{\mathcal{M}}(A) = \sum_{u \in \mathcal{M}} \langle u | Au \rangle = \sum_{u \in \mathcal{M}} s(u) .$$

As shown in the first part of this very proof this sum is absolutely convergent. In particular, since $s(u) = \langle u | Au \rangle$ for every $u \in \mathcal{M}$, it holds that

$$\sum_{u \in \mathcal{M}} |s(u)| < +\infty .$$

Next consider another Hilbert basis $\mathcal{B} \in \mathbb{H}$, and from (28) we easily find

$$\langle v | Av \rangle = \sum_{u \in \mathcal{M}} s(u) |\langle v | u \rangle|^2 \hspace{0.5cm} \forall v \in \mathcal{B} ,$$

so that

$$\sum_{v \in \mathcal{B}} \langle v | Av \rangle = \sum_{v \in \mathcal{B}} \sum_{u \in \mathcal{M}} s(u) |\langle v | u \rangle|^2 \hspace{0.5cm} \forall v \in \mathcal{B} .$$  \hspace{1cm} (29)

If we were allowed to swap the two summations in the right-hand side of (29), we would obtain

$$tr_{\mathcal{B}}(A) = \sum_{v \in \mathcal{B}} \sum_{u \in \mathcal{M}} |\langle v | u \rangle|^2 = \sum_{u \in \mathcal{M}} s(u) ||u||^2 = \sum_{u \in \mathcal{M}} s(u) = tr_{\mathcal{M}}(A) ,$$  \hspace{1cm} (30)

proving that the trace does not depend on the used Hilbert basis and concluding the proof.

To prove that those summations can be in fact interchanged, first observe that, in double
summation in (29), $u$ varies in a set at most countable according to Proposition [13]: Only the elements with eigenvalue $s(u) \neq 0$ give contribution and every eigenspace but the kernel of $A$ have finite dimension. These $u$ are exactly the elements of the countable Hilbert basis $N$ of $\text{Ker}(A)^{\perp}$. Similarly, only a set at most countable of elements $v$ gives contribution to the second sum (for every $u$, the set of non-vanishing Fourier coefficients along $B$ is at most countable). Summing up, the double sum in (29) can be interpreted as an iterated integral with respect the product measure over $N \times B'$ of a pair of $\sigma$-finite counting measures. $B'$ is countable as it is defined as the union of the supports (including an at most countable set of elements) of the functions $f_u : B \ni v \mapsto |\langle v | u \rangle|^2$ for $u$ varying in the at most countable set $N$. The integrals can be interchanged (Fubini) provided the function $N \times B' \ni (u, v) \mapsto s(u)|\langle v | u \rangle|^2$ is absolutely integrable with respect the product measure. In turn, this is equivalent (Tonelli theorem) to requiring that one of the two iterated integrals, either

$$\sum_{u \in N} \sum_{v \in B'} |s(u)||\langle v | u \rangle|^2 < +\infty \quad \text{or} \quad \sum_{v \in B'} \sum_{u \in N} |s(u)||\langle v | u \rangle|^2 < +\infty.$$  

This is the case, indeed, $M \times B \supset N \times B'$ and

$$\sum_{u \in M} \sum_{v \in B} |s(u)||\langle v | u \rangle|^2 = \sum_{u \in M} |s(u)| \sum_{v \in B} |\langle v | u \rangle|^2 = \sum_{u \in M} |s(u)| < +\infty.$$  

Therefore (30) is valid, concluding the proof of (c).

(d) We have to consider the case $\mathbb{D} = \mathbb{H}$, since the validity of (i) and (ii) for $\mathbb{D} = \mathbb{R}, \mathbb{C}$ immediately follows from (b).

(i) (For $\mathbb{D} = \mathbb{H}$) It arises from (a) and (c): $2\Re(\text{tr}_N(A)) = \text{tr}_N(A) + \overline{\text{tr}_N(A)} = 2\Re(\text{tr}_N(A) + \text{tr}_N(A^*)) = \text{tr}_N(A + A^*)$ does not depend on $N$ since $A + A^*$ is selfadjoint.

(ii) (For $\mathbb{D} = \mathbb{H}$) We consider the real Hilbert space $\mathbb{H}_R$ as in Lemma [A.1]. According to (A.1b) in Lemma [A.1] if $N$ is a Hilbert basis of $\mathbb{H}$, $N_R$ is a Hilbert basis of the real Hilbert space $\mathbb{H}_R$ equipped with the real scalar product $(\cdot | \cdot) = \Re(\langle \cdot | \cdot \rangle)$. In view of the identity (27), if $AB \in \mathfrak{B}_1(\mathbb{H})$, then $AB \in \mathfrak{B}_1(\mathbb{H}_R)$ and the same fact holds true for $BA$, so we can compute the trace $\text{tr}_{N_R}(AB)$ and $\text{tr}_{N_R}(BA)$ in the Hilbert space $\mathbb{H}_R$ where the cyclic property is valid for (b). Using the real scalar product $(\cdot | \cdot) = \Re(\langle \cdot | \cdot \rangle)$ of $\mathbb{H}_R$ and continuity of the function $Re$, we therefore have

$$\sum_{z \in N} \Re(\langle z | ABz \rangle) + \sum_{z \in N} \Re(\langle zi | ABzi \rangle) + \sum_{z \in N} \Re(\langle zj | ABzj \rangle) + \sum_{z \in N} \Re(\langle zk | ABzk \rangle)$$

$$= \text{tr}_{N_R}(AB) = \text{tr}_{N_R}(BA) =$$

$$= \sum_{z \in N} \Re(\langle z | BAz \rangle) + \sum_{z \in N} \Re(\langle zi | BAzi \rangle) + \sum_{z \in N} \Re(\langle zj | BAzj \rangle) + \sum_{z \in N} \Re(\langle zk | BAzk \rangle).$$

Next observe that, since $\Re(q'_q) = \Re(q'q)$, it holds that $\Re(\langle zi | Tzi \rangle) = -\Re(i\langle zi | Tz \rangle i) = -\Re(i'i\langle zi | Tz \rangle) = \Re(\langle zi | Tz \rangle)$. Similarly $\Re(\langle jz | Tjz \rangle) = \Re(\langle kz | T kz \rangle) = \Re(\langle z | Tz \rangle)$ and
therefore, the long identity written above can be rephrased into
\[ 4 \sum_{z \in N} Re(\langle z|ABz \rangle) = tr_{N_e}(AB) = tr_{N_e}(BA) = 4 \sum_{z \in N} Re(\langle z|BAz \rangle). \]

In summary, \( Re(tr_N(AB)) = Re(tr_N(BA)) \).

Let us prove point (e). Let \( N \) a Hilbert basis as in the hypotheses obtained by completing a Hilbert basis \( N_e \) of \( Ker(A)^\perp \) made of eigenvectors of \( A \). Then we have
\[
tr_N(BA) = \sum_{u \in N} \langle u|BAu \rangle = \sum_{u \in N} \langle u|Bu \rangle s(u) = \sum_{u \in N} \langle Au|Bu \rangle = \sum_{u \in N} \langle Au|Bu \rangle = tr_N(AB),
\]
where \( s(u) \in \mathbb{R} \) is the eigenvalue associated with the eigenvector \( u \in N_e \). In particular, if \( B = B^* \), then \( tr_N(AB) = tr_N((AB)^*) = tr_N(BA) = tr_N(AB) \), and so \( tr_N(AB) \in \mathbb{R} \).

To conclude let us prove point (f). The condition \( A \geq B \) means that \( \langle x|Ax \rangle - \langle x|Bx \rangle \geq 0 \), thus in particular \( Re\langle x|Ax \rangle \geq Re\langle x|Bx \rangle \), for all \( x \in \mathcal{H} \). So, consider any Hilbert basis \( N \subset \mathcal{H} \), then
\[
Re(tr_N(A)) = \sum_{z \in N} Re(\langle z|Az \rangle) \geq \sum_{z \in N} Re(\langle z|Bz \rangle) = Re(tr_N(B)).
\]

\[ \square \]

**Proof of Proposition 2.8.** (i) implies (ii) as established with the first statement of Proposition 2.3. The fact that (ii) entails (i) if \( \mathbb{D} = \mathbb{C} \) is a known result (e.g., see Proposition 4.41 in [Mo18]). Failure of the implication (ii) \( \Rightarrow \) (i) in the real case is evident considering \( A : \mathcal{H} \rightarrow \mathcal{H} \) with \( AA = -I \) and \( A^* = A \) (such an operator can be constructed easily referring to a Hilbert basis of an infinite dimensional real Hilbert space) so that \( |A| = I \) and \( \sum_{u \in N} \langle u|A|u \rangle = +\infty \) falsifying (i), but \( \langle u|Au \rangle = -\langle Au|u \rangle = -\langle u|Au \rangle \) leads to \( \sum_{u \in N} \langle u|Au \rangle = 0 \), satisfying (ii). To conclude, we prove that (ii) \( \Rightarrow \) (i) for \( \mathbb{D} = \mathbb{H} \). Suppose that \( A \in \mathfrak{B}(\mathcal{H}) \) satisfies (ii). Define \( B := (A + A^*)/2 \) and \( C = (A - A^*)/2 \). Since \( |\langle x|(A \pm A^*)x \rangle| = |\langle x|Ax \rangle \pm \langle x|Ax \rangle| \leq 2|\langle x|Ax \rangle| \), both \( B \) and \( C \) satisfy (ii). Since \( \mathfrak{B}_1(\mathcal{H}) \) is a real vector space ((d) Proposition 2.2), it is sufficient to prove that \( B \) and \( C \) satisfy (i) to conclude. Regarding the selfadjoint operator \( B \), the fact that it satisfies (i) from (ii) can be proved passing to its integral spectral decomposition (using the fact that the spherical spectrum is completely included in \( \mathbb{R} \) since \( B = B^* \)) and following the same route as that used to prove Proposition 4.41 in [Mo18] for \( T = T^* \). Let us pass to \( C \). Due to Theorem 5.9 in [GMP13] (specialised to \( T = C \)), we can write \( C = J|C| \) for some \( J \in \mathfrak{B}(\mathcal{H}) \) with \( JJ = -I \) and \( J^* = -J \) such that \( JC = CJ \) and \( J|C| = |C|J \). As a consequence (GMP13 Sect.3.3), the complex Hilbert space \( \mathcal{H}_{J|} := \{ u \in \mathcal{H} | J|u = u \} \) equipped with the restriction of the scalar product of \( \mathcal{H} \) is invariant under \( C \). If \( N \subset \mathcal{H}_{J|} \) is a Hilbert basis, it is also a Hilbert basis of \( \mathcal{H} \) in view of (f) Proposition 3.8 in [GMP13]. So, condition (ii) for \( C \) specializes to \( \sum_{u \in N} |\langle u|J|C|u \rangle| < +\infty \) for the said simultaneous Hilbert
basis of $H_J$ and $H$. From the definition of $H_J$, $\sum_{u \in N} |z\langle u | C | u \rangle| < +\infty$, which can be written as $\sum_{u \in N} \langle u | C | u \rangle < +\infty$ so that $C \in \mathcal{B}_1(H)$ by definition, concluding the proof. □

**Proof of Proposition 2.9** Consider a Hilbert basis $N$ of $H$ such that $N \subset H_J$. We have

$$tr_N(A) = \frac{1}{2} tr_N(A + A^*) + \frac{1}{2} tr_N(A - A^*) = Re(tr_N(A)) + \frac{1}{2} tr_N(J(A - A^*)) .$$

We have proved that

$$tr_N(A) = tr^R(A) + \frac{1}{2} \sum_{z \in N} \langle z | J(A - A^*) | z \rangle = tr^R(A) - \frac{1}{2} \sum_{z \in N} \langle Jz | A - A^* | z \rangle .$$

Since $Jz = z^*$ and $|A - A^*|$ is selfadjoint so that $\langle z | A - A^* | z \rangle \in \mathbb{R}$,

$$tr_N(A) = tr^R(A) + \frac{i}{2} \sum_{z \in N} \langle z | A - A^* | z \rangle = tr^R(A) + \frac{i}{2} tr^R(|A - A^*|) = tr^R(A) + \frac{i}{2} tr(|A - A^*|),$$

where we have used (e) Corollary 2.7. Eventually, observe that any change of $J$ on $Ker(A - A^*) = Ker(|A - A^*|)$ does not affect the result because it would only change vanishing terms $\langle z | J(A - A^*) | z \rangle$ in view of $|A - A^*| z = 0$. □

**Proof of Lemma 3.4** Let $N$ be as in the hypotheses, and consider the decomposition

$$1 = \sum_{n \in N} p_n q_n = (1 - p_0) \left[ \sum_{n = 1}^N \frac{p_n}{1 - p_0} q_n \right] + p_0 q_0 .$$

Notice the following inequalities

$$0 < \frac{p_i}{1 - p_0} < \sum_{n = 1}^N \frac{p_n}{1 - p_0} = \frac{1}{1 - p_0} \sum_{n = 1}^N p_n = \frac{1}{1 - p_0} (1 - p_0) = 1 , \quad (32)$$

from which we also have

$$0 \leq q := \sum_{n = 1}^N \frac{p_n}{1 - p_0} q_n \leq \sum_{n = 1}^N \frac{p_n}{1 - p_0} = 1 . \quad (33)$$

Thus we reduce to $(1 - p_0)q + p_0 q_0 = 1$ with $p_0 \in (0, 1)$ and $q_0, q \in [0, 1]$. Assume by contradiction that $q \neq q_0$ and, without loss of generality, suppose that $q_0 > q$. The identity above can be rewritten as $q - p_0 q + p_0 q_0 = 1$, i.e. $p_0 (q_0 - q) = 1 - q$. Thus, from $(1 - q)(q_0 - q)^{-1} = p_0 < 1$ we get $1 - q < q_0 - q$ which is equivalent to $1 < q_0$, in turn impossible. Thus $q = q_0$ and more precisely $1 = (1 - p_0)q + p_0 q_0 = q_0 = q$. Repeating the argument on the sums $\sum_{n = 1}^N \frac{p_n}{1 - p_0} = \sum_{n = 1}^N \frac{p_n}{1 - p_0} q_n = 1$ we get $q_1 = 1$. By induction we get $q_n = 1$ for all $n$. □
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