



Existence of minimizers for causal variational principles on compact subsets of momentum space in the homogeneous setting

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Abstract

We prove the existence of minimizers for the causal action in the class of negative definite measures on compact subsets of momentum space in the homogeneous setting under several side conditions (constraints). The method is to employ Prohorov's theorem. Given a minimizing sequence of negative definite measures, we show that, under suitable side conditions, a unitarily equivalent subsequence thereof is bounded. By restricting attention to compact subsets, from Prohorov's theorem we deduce the existence of minimizers in the class of negative definite measures.

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1 Introduction

In the physical theory of causal fermion systems, spacetime and the structures therein are described by a minimizer of the so-called causal action principle (for an introduction to the physical background and the mathematical context, we refer the interested reader to the textbook [14], the survey articles [16, 17, 26] as well as the web platform [1]). Given a causal fermion system $(\mathcal{H}, \mathcal{F}, d\rho)$ together with a non-negative function $\mathcal{L} : \mathcal{F} \times \mathcal{F} \rightarrow \mathbb{R}_0^+ := [0, \infty)$ (the *Lagrangian*), the causal action principle is to minimize the *action* \mathcal{S} defined as the double integral over the Lagrangian

$$\mathcal{S}(\rho) = \int_{\mathcal{F}} d\rho(x) \int_{\mathcal{F}} d\rho(y) y \mathcal{L}(x, y)$$

under variations of the measure $d\rho$ within the class of regular Borel measures on \mathcal{F} under suitable side conditions. The corresponding minimizer is of crucial physical interest; for instance, its support is interpreted as physical spacetime (readers interested in the physical background of causal fermion systems are referred to the introductory paper [26]). In order to work out the existence theory for minimizers, *causal variational principles* evolved as a mathematical generalization of the causal action principle [12, 18]. The aim of the present paper is to prove the existence of minimizers for causal variational principles restricted to compact subsets in the *homogeneous* setting with respect to different side conditions.

Let us put the present paper into the mathematical context. In [10] it was proposed to formulate physics in terms of a new type of variational principle in spacetime. The suggestion in [10, Section 3.5] led to the causal action principle in discrete spacetime, which was first analyzed mathematically in [11]. A more general and systematic inquiry of causal variational principles on measure spaces was carried out in [12]. In [12, Section 3] the existence of minimizers for variational principles in indefinite inner product spaces is proven in the special case that the total spacetime volume as well as the number of particles therein are finite. Under the additional assumption that the kernel of the fermionic projector is *homogeneous* in the sense that it only depends on the difference of two spacetime points, variational principles for homogeneous systems were considered in [12, Section 4] in order to deal with an infinite number of particles in an infinite spacetime volume. More precisely, the main advantage in the homogeneous setting is that it allows for Fourier methods, thus giving rise to a natural correspondence between position and momentum space. As a consequence, one is led to minimize the causal action by varying in the class of negative definite measures, and the existence of minimizers on bounded subsets of momentum space is proven in [12, Theorem 4.2]. The aim of this paper is to prove the existence of minimizers on compact subsets with respect to additional side conditions (see Sect. 4) which were not considered in [12], thus partially generalizing [12, Theorem 4.2] concerning compact subsets.

The paper is organized as follows. In Sect. 2 we first outline some mathematical preliminaries (Sect. 2.1) and afterwards recall causal variational principles in infinite spacetime volume (Sect. 2.2). In order to put the causal variational principles into the context of the calculus of variations, in Sect. 3 we first introduce so-called operator-valued measures (Sect. 3.1); afterwards, we consider variational principles on compact subsets of momentum space in the homogeneous setting (Sect. 3.2). In Sect. 4, we prove the existence of minimizers for the causal variational principle on compact subsets in the class of negative definite measures (Theorem 4.1). To this end we first show that, under appropriate side conditions, minimizing sequences of negative definite measures are bounded with respect to the total variation (Sect. 4.1). We then state a preparatory result which ensures the existence of weakly convergent subsequences (Sect. 4.2). This allows us to prove our main result (Sect. 4.3). Afterwards

we show that the main result also holds in the case that a boundedness constraint is imposed (Sect. 4.4). Finally, we give a short discussion of the main results (Theorem 4.1 and Theorem 4.11) and compare them with [12, Theorem 4.2] (Sect. 4.5). In the appendix we motivate and justify our choice of side conditions (Appendix A).

2 Mathematical preliminaries

2.1 Mathematical preliminaries and notation

To begin with, let us compile some fundamental definitions being of central relevance throughout this paper. For details we refer the interested reader to Bognár [4], Gohberg et al. [21] and Langer [27]. Unless specified otherwise, we always let $n \geq 1$ be a given integer.

Definition 2.1 A mapping $\langle \cdot | \cdot \rangle : \mathbb{C}^n \times \mathbb{C}^n \rightarrow \mathbb{C}$ is called an *indefinite inner product* if the following conditions hold (cf. [21, Definition 2.1]):

- (i) $\langle y | \alpha x_1 + \beta x_2 \rangle = \alpha \langle y | x_1 \rangle + \beta \langle y | x_2 \rangle$ for all $x_1, x_2, y \in \mathbb{C}^n, \alpha, \beta \in \mathbb{C}$.
- (ii) $\langle x | y \rangle = \overline{\langle y | x \rangle}$ for all $x, y \in \mathbb{C}^n$.
- (iii) $\langle x | y \rangle = 0$ for all $y \in \mathbb{C}^n \implies x = 0$.

Definition 2.2 Let V be a finite-dimensional complex vector space, endowed with an indefinite inner product $\langle \cdot | \cdot \rangle$. Then $(V, \langle \cdot | \cdot \rangle)$ is called an *indefinite inner product space*.

As usual, by $L(V)$ we denote the set of (bounded) linear operators on a complex (finite-dimensional) vector space V of dimension $n \in \mathbb{N}$. The adjoint of $A \in L(V)$ with respect to the Euclidean inner product $\langle \cdot | \cdot \rangle$ on $V \simeq \mathbb{C}^n$ is denoted by A^\dagger . On the other hand, whenever $(V, \langle \cdot | \cdot \rangle)$ is an indefinite inner product space, unitary matrices and the adjoint A^* (with respect to $\langle \cdot | \cdot \rangle$) are defined as follows.

Definition 2.3 Let $\langle \cdot | \cdot \rangle$ be an indefinite inner product on $V \simeq \mathbb{C}^n$, and let S be the associated invertible hermitian matrix determined by Gohberg et al. [21, Eq. (2.1.1)],

$$\langle x | y \rangle = \langle Sx | y \rangle \quad \text{for all } x, y \in \mathbb{C}^n .$$

Then for every $A \in L(V)$, the adjoint of A (with respect to $\langle \cdot | \cdot \rangle$) is the unique matrix $A^* \in L(V)$ which satisfies

$$\langle Ax | y \rangle = \langle x | A^*y \rangle \quad \text{for all } x, y \in V .$$

A matrix $A \in L(V)$ is called *self-adjoint* (with respect to $\langle \cdot | \cdot \rangle$) if and only if $A = A^*$. In a similar fashion, an operator $U \in L(V)$ is said to be *unitary* (with respect to $\langle \cdot | \cdot \rangle$) if it is invertible and $U^{-1} = U^*$ (see [21, Section 4.1]).

We remark that every non-negative matrix (with respect to $\langle \cdot | \cdot \rangle$) is self-adjoint (with respect to $\langle \cdot | \cdot \rangle$) and has a real spectrum (cf. [21, Theorem 5.7.2]). Moreover, the adjoint A^* of $A \in L(V)$ satisfies the relation

$$A^* = S^{-1} A^\dagger S$$

in view of Gohberg et al. [21, Eq. (4.1.3)] (where A^\dagger denotes the adjoint with respect to $\langle \cdot | \cdot \rangle$) and A^* the adjoint with respect to $\langle \cdot | \cdot \rangle$). For details concerning self-adjoint operators (with respect to $\langle \cdot | \cdot \rangle$) we refer to Langer [27] and the textbook [4]. In the remainder of this

paper we will restrict attention exclusively to indefinite inner product spaces $(V, \langle \cdot | \cdot \rangle)$ with $V \simeq \mathbb{C}^{2n}$ for some $n \in \mathbb{N}$. It is convenient to work with a fixed pseudo-orthonormal basis $(\epsilon_i)_{i=1, \dots, 2n}$ of V in which the inner product has the standard representation with a signature matrix S ,

$$\langle u | v \rangle = \langle u | Sv \rangle_{\mathbb{C}^{2n}} \quad \text{with} \quad S = \text{diag}(\underbrace{1, \dots, 1}_{n \text{ times}}, \underbrace{-1, \dots, -1}_{n \text{ times}}), \tag{2.1}$$

where $\langle \cdot | \cdot \rangle_{\mathbb{C}^{2n}}$ denotes the standard inner product on \mathbb{C}^{2n} . The signature matrix can be regarded as an operator on V ,

$$S = \begin{pmatrix} \mathbb{1} & 0 \\ 0 & -\mathbb{1} \end{pmatrix} \in \text{Symm } V, \tag{2.2}$$

where $\text{Symm } V$ denotes the set of symmetric matrices on V with respect to the ‘‘spin scalar product’’ $\langle \cdot | \cdot \rangle$ (also cf. [12, proof of Lemma 3.4]). Without loss of generality we may assume that $\epsilon_i = (0, \dots, 0, 1, 0, \dots, 0)^T$ for all $i = 1, \dots, 2n$.

In what follows, we denote Minkowski space by $\mathcal{M} \simeq \mathbb{R}^4$ and momentum space by $\hat{\mathcal{M}} \simeq \mathbb{R}^4$. Identifying $\hat{\mathcal{M}}$ with Minkowski space \mathcal{M} , the Minkowski inner product (of signature $(+, -, -, -)$) can be considered as a mapping

$$\langle \cdot, \cdot \rangle : \hat{\mathcal{M}} \times \mathcal{M} \rightarrow \mathbb{R}, \quad (p, \xi) \mapsto \langle p, \xi \rangle = \eta_{\mu\nu} p^\mu \xi^\nu = p^0 \xi^0 - \sum_{i=1}^3 p^i \xi^i$$

for all $\xi = (\xi^0, \xi^1, \xi^2, \xi^3) \in \mathcal{M}$ and $p = (p^0, p^1, p^2, p^3) \in \hat{\mathcal{M}}$ (with Minkowski metric η , where we employed Einstein’s summation convention, cf. [19, Chapter 1]).

In the remainder of this paper, let $\hat{K} \subset \hat{\mathcal{M}}$ be a compact subset. By $\mathcal{B}(\hat{K})$ we denote the Borel σ -algebra on \hat{K} . The class of finite complex measures on \hat{K} is denoted by $\mathbf{M}_\mathbb{C}(\hat{K})$. By $C_c(\hat{\mathcal{M}})$ we denote the set of continuous functions on $\hat{\mathcal{M}}$ with compact support, whereas $C_b(\hat{\mathcal{M}})$ and $C_0(\hat{\mathcal{M}})$ indicate the sets of continuous functions on $\hat{\mathcal{M}}$ which are bounded or vanishing at infinity, respectively. Since \hat{K} is compact, the sets $C_c(\hat{K})$ and $C_b(\hat{K})$ coincide. By $L^1_{\text{loc}}(\mathcal{M})$ we denote the set of locally integrable functions on \mathcal{M} with respect to Lebesgue measure, denoted by $d\mu$. Unless otherwise specified, we always refer to locally finite measures on the Borel σ -algebra as Borel measures in the sense of Gardner and Pfeffer [20]. A Borel measure is said to be regular if it is inner and outer regular [7].

2.2 Variational principles in infinite spacetime volume

This subsection is intended to give a motivating example, thereby illustrating the underlying physical ideas. More precisely, before entering variational principles in infinite spacetime volume as introduced in [12], let us briefly recall the concept of a Dirac sea as introduced by Paul Dirac in his paper [5]. In this article, he assumes that

‘‘(...) all the states of negative energy are occupied except perhaps a few of small velocity. (...) Only the small departure from exact uniformity, brought about by some of the negative-energy states being unoccupied, can we hope to observe. (...) We are therefore led to the assumption that the holes in the distribution of negative-energy electrons are the [positrons].’’

Dirac made this picture precise in his paper [6] by introducing a relativistic density matrix $R(t, \vec{x}; t', \vec{x}')$ with $(t, \vec{x}), (t', \vec{x}') \in \mathbb{R} \times \mathbb{R}^3$ defined by

$$R(t, \vec{x}; t', \vec{x}') = \sum_{l \text{ occupied}} \Psi_l(t, \vec{x}) \overline{\Psi_l(t', \vec{x}')}.$$

In analogy to Dirac’s original idea, in [9] the kernel of the fermionic projector is introduced as the sum over all occupied wave functions

$$P(x, y) = - \sum_{l \text{ occupied}} \Psi_l(x) \overline{\Psi_l(y)}$$

for spacetime points $x, y \in \mathcal{M}$ as outlined in [13]. A straightforward calculation shows that (see e.g. [15, Sect. 4.1]) the kernel of the fermionic projector takes the form

$$P(x, y) = \int_{\hat{\mathcal{M}}} \frac{d^4 p}{(2\pi)^4} (\not{p} + m) \delta(p^2 - m^2) \Theta(-p^0) e^{-ip(x-y)} \tag{2.3}$$

(where δ denotes Dirac’s delta distribution and Θ is the Heaviside function; moreover, by $\not{p} \equiv \gamma^\mu p_\mu$ we denote the so-called *Feynman slash*, where γ^μ ($\mu = 0, \dots, 3$) denote the *Dirac matrices*). We refer to $P(x, y)$ as the (unregularized) *kernel of the fermionic projector of the vacuum* (cf. [14, Eqs. (1.2.20) and (1.2.23)] as well as [10, Eq. (4.1.1)]; this object already appears in [8]). We also refer to (2.3) as a *completely filled Dirac sea*. The kernel of the fermionic projector (2.3) is the starting point for the analysis in [12, Section 4]. In order to deal with systems containing an infinite number of particles in an infinite spacetime volume, the main simplification in [12] is to assume that the kernel of the fermionic projector (2.3) is *homogeneous* in the sense that $P(x, y)$ only depends on the difference vector $y - x$ for all spacetime points $x, y \in \mathcal{M}$. The underlying homogeneity assumption $P(x, y) = P(y - x)$ for all $x, y \in \mathcal{M}$ is referred to as “homogeneous regularization of the vacuum” (cf. [10, Eq. (4.1.2)] and the explanations thereafter; also see [14, Assumption 3.3.1]). Introducing $\xi = \xi(x, y) := y - x$ for all $x, y \in \mathcal{M}$ and

$$\hat{P}(p) = (\not{p} + m) \delta(p^2 - m^2) \Theta(-p^0)$$

for all $p \in \hat{\mathcal{M}}$, the fermionic projector (2.3) can be written as a Fourier transform,

$$P(x, y) = \int_{\hat{\mathcal{M}}} \frac{d^4 p}{(2\pi)^4} \hat{P}(p) e^{i(p, \xi)}$$

(for details concerning the Fourier transform we refer to Folland [19]). In order to arrive at a measure-theoretic framework, it is convenient to regard $\hat{P}(p) d^4 p / (2\pi)^4$ as a Borel measure $d\nu$ on $\hat{\mathcal{M}}$, taking values in $L(V)$. In particular, the measure

$$d\nu(p) = (\not{p} + m) \delta(p^2 - m^2) \Theta(-p^0) d^4 p \tag{2.4}$$

has the remarkable property that $-d\nu$ is positive in the sense that

$$\langle u | -\nu(\Omega) u \rangle \geq 0 \quad \text{for all } \Omega \in \mathcal{B}(\hat{\mathcal{M}}) \text{ and for all } u \in \mathbb{C}^4 \tag{2.5}$$

with respect to the “spin scalar product” $\langle \cdot | \cdot \rangle$ on \mathbb{C}^4 introduced in Sect. 2.1.¹

In order to avoid ultraviolet problems (where “ultraviolet” refers to regions of high energy in momentum space), caused by measures of the form (2.4), one is led to restrict attention to

¹ In order to see this, we make use of the fact that the Dirac matrices anti-commute, i.e.

$$\gamma^\mu \gamma^\nu = -\gamma^\nu \gamma^\mu \quad \text{whenever } \mu \neq \nu.$$

compact subsets of momentum space [12]. Moreover, generalizing the "physical" indefinite inner product space $(\mathbb{C}^4, \langle \cdot | \cdot \rangle)$ to some abstract indefinite inner product space $(V, \langle \cdot | \cdot \rangle)$ of dimension $2n$, the above observations motivate the following definition (see [12, Definition 4.1]).

Definition 2.4 A vector-valued Borel measure $d\nu$ on a compact set $\hat{K} \subset \hat{\mathcal{M}}$ taking values in $L(V)$ is called a *negative definite measure* on \hat{K} with values in $L(V)$ whenever $d \prec u | -\nu u \succ$ is a positive finite measure for all $u \in V$. By \mathfrak{Ndm} we denote the class of negative definite measures on \hat{K} taking values in $L(V)$.

In terms of a negative definite measure $d\nu$, the *kernel of the fermionic projector* is then introduced by

$$P(\xi) := \int_{\hat{K}} e^{i\langle p, \xi \rangle} d\nu(p) \quad \text{for all } \xi \in \mathcal{M}. \tag{2.6}$$

In order to clarify the dependence on $d\nu$, we also write $P[\nu]$. For every $\xi \in \mathcal{M}$, the *closed chain* is defined by $A(\xi) := P(\xi)P(-\xi)$. In order to emphasize that the closed chain depends on $d\nu$, we also write $A[\nu]$. According to Finster [12, Eq. (3.7)], the *spectral weight* $|A|$ of an operator $A \in L(V)$ is given by the sum of the absolute values of the eigenvalues of A ,

$$|A| = \sum_{i=1}^{2n} |\lambda_i|,$$

where by λ_i we denote the eigenvalues of A , counted with algebraic multiplicities. In analogy to Finster [12, Eq. (3.8)], for every $\xi \in \mathcal{M}$ the Lagrangian is introduced via

$$\mathcal{L}[A(\xi)] := |A(\xi)|^2 - \frac{1}{2n}|A(\xi)|^2.$$

Thus for every $p \in \hat{\mathcal{M}}$ with $p = (p^0, \vec{p})$, the operators $p_{\pm}(\vec{p}) := (1/2p^0)(\not{p} + m)\gamma^0|_{p^0=\pm\omega(\vec{p})}$ as given by [29, Eq. (2.13)] satisfy

$$p_{\pm}(\vec{p})\gamma^0 = \frac{\not{p} + m}{2p^0}\gamma^0|_{p^0=\pm\omega(\vec{p})} = \gamma^0 p_{\pm}(-\vec{p})$$

with $\omega(\vec{p}) = \sqrt{\vec{p}^2 + m^2}$. Applying the fact that $p_{\pm}(\vec{p})$ is idempotent and symmetric with respect to the Euclidean scalar product $(\cdot | \cdot)_{\mathbb{C}^4}$ on \mathbb{C}^4 (cf. [29, Proposition 2.14]), the calculation

$$\begin{aligned} \langle u | (\not{p} + m) u \rangle &= 2p^0 \langle u | \gamma^0 p_{\pm}(-\vec{p}) u \rangle = 2p^0 \langle u | p_{\pm}(-\vec{p})^2 u \rangle_{\mathbb{C}^4} \\ &= 2p^0 \langle p_{\pm}(-\vec{p}) u | p_{\pm}(-\vec{p}) u \rangle_{\mathbb{C}^4} \end{aligned}$$

for any $u \in \mathbb{C}^4$ implies that

$$\langle \cdot | (\not{p} + m) \cdot \rangle \text{ is } \begin{cases} \text{positive semidefinite if } p^0 > 0 \\ \text{negative semidefinite if } p^0 < 0. \end{cases}$$

Introducing the set $\Omega^- = \Omega \cap \{p^0 < 0 : p = (p^0, \vec{p}) \in \hat{\mathcal{M}}\}$ for any $\Omega \in \mathcal{B}(\hat{\mathcal{M}})$, for all $u \in V$ we obtain

$$\begin{aligned} \langle u | -\nu(\Omega) u \rangle &= \langle u | - \int_{\Omega} (p_j \gamma^j + m) \delta(\langle p, p \rangle - m^2) \Theta(-p^0) d^4 p u \rangle \\ &= \int_{\Omega^-} \underbrace{\langle u | (\not{p} + m) u \rangle}_{\geq 0} \delta(\langle p, p \rangle - m^2) d^4 p \geq 0. \end{aligned}$$

Therefore, positivity (2.5) is a consequence of the corresponding behavior of the operator $(\not{p} + m)$.

Defining the *action* \mathcal{S} according to Finster [12, Eq. (4.5)] by

$$\mathcal{S} : \mathfrak{M} \rightarrow [0, +\infty] , \quad \mathcal{S}(v) := \int_{\mathfrak{M}} \mathcal{L}[A(\xi)] d\mu(\xi) ,$$

the *causal variational principle in the homogeneous setting* is to

minimize $\mathcal{S}(v)$ by suitably varying dv in \mathfrak{M} .

Given a negative definite measure dv , the complex measure $d \langle u | v v \rangle \in \mathbf{M}_{\mathbb{C}}(\hat{K})$ is defined by polarization for all $u, v \in V$,

$$d \langle u | v v \rangle := \frac{1}{4} \{ d \langle u + v | v (u + v) \rangle + i d \langle u + iv | v (u + iv) \rangle - d \langle u - v | v (u - v) \rangle - i d \langle u - iv | v (u - iv) \rangle \} \quad (2.7)$$

(see e.g. [21, Eq. (2.2.6)], also cf. [30, Section VIII.3]). Following Langer [25, Definition A.16], we define integration with respect to negative definite measures as follows.

Definition 2.5 Let $(V, \langle \cdot | \cdot \rangle)$ be an indefinite inner product space and let dv be a negative definite measure. Moreover, let $\phi : \hat{K} \rightarrow \mathbb{C}$ be a bounded Borel measurable function. For all $u, v \in V$, integration with respect to dv is defined by

$$\langle u | \left(\int_{\hat{K}} \phi(p) dv(p) \right) v \rangle := \int_{\hat{K}} \phi(p) d \langle u | v(p) v \rangle .$$

A similar definition in terms of operator-valued measures is stated below (see Definition 3.6). For a connection to spectral theory we refer to [28, Chapter 31].

3 Causal variational principles in the homogeneous setting

3.1 Operator-valued measures

in order to deal with causal variational principles in the homogeneous setting in sufficient generality, this subsection is devoted to put the definition of negative definite measures into the context of the calculus of variations. More precisely, as explained in Sect. 2.2, the variational principle as introduced in [12, Section 4] is to minimize the causal action \mathcal{S} in the class of negative definite measures. Unfortunately, in view of (2.5), the set of negative definite measures does not form a vector space, whereas in calculus of variations one usually considers functionals on a real, locally convex vector space (for details we refer to Zeidler [33, Section 43.2]). Hence in order to obtain a suitable framework, we first introduce operator-valued measures, which can be regarded as a generalization of negative definite measures, thus providing the basic structures required for the calculus of variations (see Lemma 3.3 below).

Definition 3.1 A (vector-valued) measure $d\omega$ on $\mathcal{B}(\hat{K})$ taking values in $L(V)$ is called an *operator-valued measure* on \hat{K} with values in $L(V)$ whenever $d \langle u | \omega v \rangle$ is a finite complex measure in $\mathbf{M}_{\mathbb{C}}(\hat{K})$ for all $u, v \in V$.

Whenever \hat{K} and V are understood, the class of operator-valued measures on \hat{K} with values in $L(V)$ shall be denoted by $\mathfrak{O} \mathfrak{M}$.

In what follows, the variation of an operator-valued measure plays a central role:

Definition 3.2 Given an operator-valued measure $d\omega \in \mathfrak{D}\mathfrak{v}\mathfrak{m}$, the *variation* of $d\omega$, denoted by $d|\omega|$, is defined by

$$d|\omega| := \sum_{i,j=1}^{2n} d|\langle \epsilon_i | \omega \epsilon_j \rangle|,$$

where $d|\cdot|$ denotes the variation of a complex measure. Moreover, the *total variation* of $d\omega$, denoted by $d\|\omega\|$, is given by

$$d\|\omega\| := d|\omega|(\hat{K}) = \sum_{i,j=1}^{2n} d|\langle \epsilon_i | \omega \epsilon_j \rangle|(\hat{K}). \tag{3.1}$$

We point out that the variation as given by Definition 3.2 crucially depends on the pseudo-orthogonal $(\epsilon_i)_{i=1,\dots,2n}$ basis of V . Nevertheless, the set of operator-valued measures $\mathfrak{D}\mathfrak{v}\mathfrak{m}$ is a Banach space with respect to the total variation:

Lemma 3.3 *The total variation $d\|\cdot\|$ given by (3.1) defines a norm on $\mathfrak{D}\mathfrak{v}\mathfrak{m}$ in such a way that $(\mathfrak{D}\mathfrak{v}\mathfrak{m}, d\|\cdot\|)$ is a complex Banach space. In particular, $(\mathfrak{D}\mathfrak{v}\mathfrak{m}, d\|\cdot\|)$ is a real, locally convex vector space.*

Proof For the first part of the statement see the proof of Langer [25, Corollary 5.3]. In order to show that $\mathfrak{D}\mathfrak{v}\mathfrak{m}$ is a Banach space, let us consider a Cauchy sequence of operator-valued measures $(d\omega_k)_{k \in \mathbb{N}}$ with respect to the norm (3.1), that is, $d\|\omega_k - \omega_m\| \rightarrow 0$ as $k, m \rightarrow \infty$. Our task is to prove that its limit, denoted by $d\omega$, exists and that $d\omega$ is contained in $\mathfrak{D}\mathfrak{v}\mathfrak{m}$. Assuming that $(\epsilon_i)_{i=1,\dots,2n}$ is a pseudo-orthonormal basis of V satisfying (2.1), from (3.1) we deduce that

$$\lim_{k,m \rightarrow \infty} d|\langle \epsilon_i | (\omega_k - \omega_m) \epsilon_j \rangle| = 0 \quad \text{for all } i, j = 1, \dots, 2n.$$

Consequently, each sequence $(d|\langle \epsilon_i | \omega_k \epsilon_j \rangle|)_{k \in \mathbb{N}}$ is a Cauchy sequence of complex measures in $\mathbf{M}_{\mathbb{C}}(\hat{K})$ for all $i, j \in \{1, \dots, 2n\}$. Since $\mathbf{M}_{\mathbb{C}}(\hat{K})$ is a complex Banach space with respect to the total variation $d\|\cdot\|$ in virtue of Elstrodt [7, Aufgabe VII.1.7], there is a complex measure $d\omega_{i,j} \in \mathbf{M}_{\mathbb{C}}(\hat{K})$, being the unique limit of $(d|\langle \epsilon_i | \omega_k \epsilon_j \rangle|)_{k \in \mathbb{N}}$ for all $i, j \in \{1, \dots, 2n\}$.

Next, for all $i, j \in \{1, \dots, 2n\}$, the complex measures $d\omega_{i,j}$ in $\mathbf{M}_{\mathbb{C}}(\hat{K})$ give rise to an operator-valued measure $d\omega$ on \hat{K} with values in $L(V)$ in such a way that, for all $i, j = 1, \dots, 2n$, we are given $d|\langle \epsilon_i | \omega \epsilon_j \rangle| = d\omega_{i,j}$. More precisely, defining the operator $\omega(\Omega) \in L(V)$ for any $\Omega \in \mathcal{B}(\hat{K})$ by

$$\omega(\Omega) := \begin{pmatrix} \omega_{1,1}(\Omega) & \cdots & \omega_{1,2n}(\Omega) \\ \vdots & \ddots & \vdots \\ \omega_{n,1}(\Omega) & \cdots & \omega_{n,2n}(\Omega) \\ -\omega_{n+1,1}(\Omega) & \cdots & -\omega_{n+1,2n}(\Omega) \\ \vdots & \ddots & \vdots \\ -\omega_{1,2n}(\Omega) & \cdots & -\omega_{2n,2n}(\Omega) \end{pmatrix} \in L(V),$$

we obtain a mapping $d\omega: \mathcal{B}(\hat{K}) \rightarrow L(V)$ such that $d|\langle \epsilon_i | \omega \epsilon_j \rangle| = d\omega_{i,j} \in \mathbf{M}_{\mathbb{C}}(\hat{K})$ for all $i, j \in \{1, \dots, 2n\}$. Since $(\epsilon_i)_{i=1,\dots,2n}$ is a basis of V , for any $\Omega \in \mathcal{B}(\hat{K})$ and arbitrary

elements $u = \sum_{i=1}^{2n} \alpha_i \epsilon_i, v = \sum_{j=1}^{2n} \beta_j \epsilon_j \in V$ we arrive at

$$\langle u \mid \omega(\Omega) v \rangle = \sum_{i,j=1}^{2n} \bar{\alpha}_i \beta_j \langle \epsilon_i \mid \omega(\Omega) \epsilon_j \rangle = \sum_{i,j=1}^{2n} \bar{\alpha}_i \beta_j \omega_{i,j}(\Omega).$$

The fact that $\mathbf{M}_{\mathbb{C}}(\hat{K})$ is a complex Banach space implies that

$$d \langle u \mid \omega v \rangle = \sum_{i,j=1}^{2n} \bar{\alpha}_i \beta_j d\omega_{i,j} \in \mathbf{M}_{\mathbb{C}}(\hat{K}) \quad \text{for all } u, v \in V.$$

This shows that $d\omega \in \mathfrak{Ovm}$ is an operator-valued measure in view of Definition 3.1. Thus $(\mathfrak{Ovm}, d\|\cdot\|)$ is a complex Banach space with respect to the norm $d\|\cdot\|$ defined by (3.1). Since each norm induces a corresponding Fréchet metric, $(\mathfrak{Ovm}, d\|\cdot\|)$ can be regarded as a metric space. In particular, each complex vector space is a real one, and each Banach space is locally convex. This completes the proof. \square

Remark 3.4 The set of negative definite measures \mathfrak{Ndm} clearly is a subset of the vector space \mathfrak{Ovm} . However, \mathfrak{Ndm} itself is *not* a vector space (see [25, Remark 5.6]), but a *cone* (for the definition we refer to Schaefer and Wolff [32]).

Next, let us introduce the support of operator-valued measures as follows:

Definition 3.5 We define the *support* of an operator-valued measure $d\omega$ in \mathfrak{Ovm} as the support of its variation measure $d|\omega|$, i.e.

$$\text{supp } d\omega := \text{supp } d|\omega| = \hat{K} \setminus \bigcup \left\{ U \subset \hat{K} : U \text{ open and } d|\omega|(U) = 0 \right\}.$$

Since $d|\omega|$ is a locally finite measure on a locally compact Polish space, we conclude that $d|\omega|$ is regular and has support, $d|\omega|(\hat{K} \setminus \text{supp } d|\omega|) = 0$.

In a similar fashion, following Bogachev [3, Definition 7.1.5], an operator-valued measure $d\omega$ is called *regular* if and only if $d|\omega|$ is regular. Moreover, the measure $d\omega$ is said to be *tight* if for every $\varepsilon > 0$ there is a compact set $K_\varepsilon \subset \hat{K}$ such that $d|\omega|(\hat{K} \setminus K_\varepsilon) < \varepsilon$ (cf. [3, Definition 7.1.4]). Clearly, whenever $\hat{K} \subset \hat{M}$ is compact, every operator-valued measure on \hat{K} is tight.

Definition 3.6 In analogy to negative definite measures (see Definition 2.5), for any bounded Borel measurable function $\phi : \hat{K} \rightarrow \mathbb{C}$ we define integration with respect to operator-valued measures $d\omega$ by

$$\langle u \mid \left(\int_{\hat{K}} \phi(k) d\omega(k) \right) v \rangle := \int_{\hat{K}} \phi(k) d \langle u \mid \omega(k) v \rangle \quad \text{for all } u, v \in V.$$

Let us finally state the definition of weak convergence of operator-valued measures, which will be required later on (see Sect. 4.3 below).

Definition 3.7 We shall say that a sequence of operator-valued measures $(d\omega_k)_{k \in \mathbb{N}}$ in \mathfrak{Ovm} *converges weakly* to some operator-valued measure $d\omega$ if and only if

$$\lim_{k \rightarrow \infty} \int_{\hat{K}} \phi d \langle u \mid \omega_k v \rangle = \int_{\hat{K}} \phi d \langle u \mid \omega v \rangle \quad \text{for all } u, v \in V \text{ and } \phi \in C_b(\hat{K}).$$

We write symbolically $d\omega_k \rightarrow d\omega$.

Whenever $d\nu \in \mathfrak{N}\mathfrak{d}\mathfrak{m}$ is a negative definite measure, we recall that, for all $u, v \in V$, the complex measure $d \langle u | \nu v \rangle$ in $\mathbf{M}_{\mathbb{C}}(\hat{K})$ is defined by polarization (2.7). Thus a sequence of negative definite measures $(d\nu)_{k \in \mathbb{N}}$ converges weakly to some negative definite measure $d\nu \in \mathfrak{N}\mathfrak{d}\mathfrak{m}$ if and only if

$$\lim_{k \rightarrow \infty} \int_{\hat{K}} \phi d \langle u | \nu_k u \rangle = \int_{\hat{K}} \phi d \langle u | \nu u \rangle \quad \text{for all } u \in V \text{ and } \phi \in C_b(\hat{K}).$$

By polarization (2.7) we then conclude that

$$\lim_{k \rightarrow \infty} \int_{\hat{K}} \phi d \langle u | \nu_k v \rangle = \int_{\hat{K}} \phi d \langle u | \nu v \rangle \quad \text{for all } u, v \in V \text{ and } \phi \in C_b(\hat{K})$$

in accordance with Definition 3.7.

Note that, with the very same reasoning, the definitions and results stated in this section can be generalized to operator-valued measures on whole momentum space.

3.2 Causal variational principles on compact subsets

After these technical preliminaries, let us now return to causal variational principles in the homogeneous setting. Motivated by (2.3), the fermionic projector $P(x, y)$ in the homogeneous setting takes the form (2.6) for all $x, y \in \mathcal{M}$, where the measure $d\nu$ is given by (2.4). Generalizing $d\nu$ according to Sects. 2.2.3.1 to operator-valued measures, for a given operator-valued measure $d\omega$ on \hat{K} with values in $L(V)$ and all $x, y \in \mathcal{M}$ we introduce the *kernel of the fermionic projector* by

$$P(x, y) : V \rightarrow V, \quad P(x, y) := \int_{\hat{K}} e^{ip(y-x)} d\omega(p).$$

In order to emphasize the dependence on the operator-valued measure $d\omega$, we also write $P[\omega](x, y)$. As $P(x, y)$ is supposed to be *homogeneous*, only the difference of two spacetime points $x, y \in \mathcal{M}$ matters; denoting the difference vector by $\xi = y - x \in \mathcal{M}$, the kernel of the fermionic projector reads

$$P(\xi) : V \rightarrow V, \quad P(\xi) = \int_{\hat{K}} e^{ip\xi} d\omega(p). \tag{3.2}$$

The first step in order to set up the variational principle is to form the *closed chain*, which (as motivated by Finster [10, Sect. 3.5]) for any $\xi \in \mathcal{M}$ is defined as the mapping

$$A(\xi) : V \rightarrow V, \quad A(\xi) := P(\xi) P(-\xi).$$

We also write $A[\omega](\xi)$ in order to clarify the dependence of the closed chain on the operator-valued measure $d\omega$. Next, given a linear operator $A : V \rightarrow V$, we define the *spectral weight* by

$$|A| := \sum_{i=1}^{2n} |\lambda_i|,$$

where by $(\lambda_i)_{i=1, \dots, 2n}$ we denote the eigenvalues of the operator A , counted with algebraic multiplicities. In this way, the spectral weight furnishes a connection between endomorphisms and scalar functionals.

In order to set up a real-valued variational principle on the set of operator-valued measures, for every $d\omega \in \mathfrak{D}\mathfrak{v}\mathfrak{m}$ we introduce the *Lagrangian*

$$\mathcal{L}[\omega] : \mathcal{M} \rightarrow \mathbb{R}_0^+, \quad \mathcal{L}[\omega](\xi) := |A[\omega](\xi)|^2 - \frac{1}{2n} |A[\omega](\xi)|^2 .$$

Defining the *causal action* $\mathcal{S} : \mathfrak{D}\mathfrak{v}\mathfrak{m} \rightarrow \mathbb{R}_0^+ \cup \{+\infty\}$ by

$$\mathcal{S}(\omega) := \int_{\mathcal{M}} \mathcal{L}[\omega](\xi) d\mu(\xi)$$

(where $d\mu$ denotes the Lebesgue measure on $\mathcal{M} \simeq \mathbb{R}^4$), the *causal variational principle in the homogeneous setting* is to

minimize $\mathcal{S}(v)$ by suitably varying dv in $\mathfrak{N}\mathfrak{d}\mathfrak{m}$.

(3.3)

In order to exclude trivial minimizers, for given parameters $c, f > 0$ we either introduce the constraints

$$\text{Tr}_V (v(\hat{K})) = c \quad \text{and} \quad \text{Tr}_V (-Sv(\hat{K})) \leq f \tag{3.4}$$

(where S denotes the signature matrix (2.2)) or the side conditions

$$\text{Tr}_V (v(\hat{K})) = c \quad \text{and} \quad |v(\hat{K})| \leq f \tag{3.5}$$

(where $|\cdot|$ denotes the spectral weight), respectively. The side condition

$$\text{Tr}_V (v(\hat{K})) = c \tag{3.6}$$

is referred to as *trace constraint*. For later convenience, we also label the second conditions in (3.4) and (3.5) separately,

$$\text{Tr}_V (-Sv(\hat{K})) \leq f, \tag{3.7}$$

$$|v(\hat{K})| \leq f. \tag{3.8}$$

A motivation for the constraints (3.4)–(3.5) can be found in Appendix A. For the connection to the boundedness constraint as considered in [12, Section 4] we refer to Sect. 4.4 below.

Definition 3.8 Given a subset $N \subset \mathfrak{N}\mathfrak{d}\mathfrak{m}$, the *causal variational principle in the homogeneous setting* is to

$$\text{minimize } \mathcal{S}(v) \text{ by varying } dv \text{ in } N \subset \mathfrak{N}\mathfrak{d}\mathfrak{m}. \tag{3.9}$$

Concerning the side conditions (3.4), (3.5), the subset N takes either the form

$$N = \{dv \in \mathfrak{N}\mathfrak{d}\mathfrak{m} : dv \text{ satisfies conditions (3.4)}\} \quad \text{or} \\ N = \{dv \in \mathfrak{N}\mathfrak{d}\mathfrak{m} : dv \text{ satisfies conditions (3.5)}\} .$$

In agreement with [33, Definition 43.4], we define a minimizer for \mathcal{S} as follows:

Definition 3.9 A negative definite measure $dv \in N$ is said to be a *minimizer* for the causal variational principle (3.9) if and only if

$$\mathcal{S}(\tilde{v}) \geq \mathcal{S}(v) \quad \text{for all } d\tilde{v} \in N .$$

4 Existence of minimizers on compact subsets

This section is devoted to developing the existence theory for minimizers of the causal action principle (3.3) for given parameters $c, f > 0$ either with respect to the constraints (3.4) or (3.5), respectively.

The main result of this section can be stated as follows:

Theorem 4.1 *Let $(d\nu^{(j)})_{j \in \mathbb{N}}$ be a minimizing sequence of negative definite measures in \mathfrak{Ndm} of the causal variational principle (3.3) with respect to the constraints (3.4) or (3.5) for given constants $c, f > 0$, respectively. Then there exists a sequence of unitary operators $(U_j)_{j \in \mathbb{N}}$ on V (with respect to $\langle \cdot | \cdot \rangle$) and a subsequence $(d\nu^{(j_k)})_{k \in \mathbb{N}}$ such that $(U_{j_k} d\nu^{(j_k)} U_{j_k}^{-1})_{k \in \mathbb{N}}$ converges weakly to some negative definite measure $d\nu \neq 0$. Moreover,*

$$S(\nu) \leq \liminf_{k \rightarrow \infty} S(\nu^{(j_k)}),$$

and the limit measure $d\nu \in \mathfrak{Ndm}$ satisfies the side conditions (3.4) or (3.5), respectively. In particular, the limit measure $d\nu$ is a non-trivial minimizer of the causal variational principle (3.3) with respect to the side conditions (3.4) or (3.5), respectively. A fortiori, the above statements remain true in case that “ \leq ” in (3.4), (3.5) is replaced by “ $=$ ”.

The remainder of this section is devoted to the proof of Theorem 4.1. The key idea for proving Theorem 4.1 is essentially to apply Prohorov’s theorem (see e.g. [3, Section 8.6]). To this end, we proceed in several steps. Given a minimizing sequence of negative definite measures which satisfies the side conditions (3.4) or (3.5), we first prove boundedness of a unitarily equivalent subsequence thereof (Sect. 4.1). The proof of Theorem 4.1 is completed afterwards (Sect. 4.3). Once this is accomplished, we show that Theorem 4.1 also applies in the case that a boundedness constraint is imposed (Sect. 4.4). Afterwards the obtained results will be discussed (Sect. 4.5).

4.1 Boundedness of minimizing sequences

Let us assume that $(d\nu^{(k)})_{k \in \mathbb{N}}$ is a sequence of negative definite measures in \mathfrak{Ndm} , either satisfying

$$\text{Tr}_V(-S\nu^{(k)}(\hat{K})) \leq f \quad \text{or} \quad |\nu^{(k)}(\hat{K})| \leq f$$

for all $k \in \mathbb{N}$ and some positive constant $f > 0$. The aim of this subsection is to show that in both cases, there exists a sequence of unitary matrices $(U_k)_{k \in \mathbb{N}}$ in $L(V)$ (with respect to $\langle \cdot | \cdot \rangle$) such that the resulting sequence $(U_k d\nu^{(k)} U_k^{-1})_{k \in \mathbb{N}}$ is bounded in \mathfrak{Ndm} (with respect to the norm (3.1)). In particular, whenever the first condition is imposed, it eventually turns out that one can choose $U_k = \mathbb{1}_V$ for all $k \in \mathbb{N}$. In preparation, let us state the following results:

Proposition 4.2 *For all $B, C \in L(V)$, the operator products BC and CB have the same spectrum.*

Proof Follow the arguments in [11, Section 3] or cf. [10, Eq. (3.5.6)]. □

Lemma 4.3 *Assume that $U \in L(V)$ is unitary (with respect to $\langle \cdot | \cdot \rangle$), and let $d\nu$ in \mathfrak{Ndm} . Then the operators $\nu(\hat{K})$ and $U \nu(\hat{K}) U^{-1}$ have the same spectrum.*

Proof Applying Proposition 4.2, we infer that the operators $\nu(\hat{K}) = (\nu(\hat{K})U^{-1})U$ and $U\nu(\hat{K})U^{-1} = U(\nu(\hat{K})U^{-1})$ have the same spectrum for any unitary matrix U in $L(V)$. \square

Corollary 4.4 For any negative definite measure $d\nu \in \mathfrak{N}\mathfrak{d}\mathfrak{m}$ and arbitrary unitary transformations U on V (with respect to $\langle \cdot | \cdot \rangle$),

$$\mathcal{L}[U\nu U^{-1}] = \mathcal{L}[\nu] \quad \text{and} \quad \mathcal{S}(U\nu U^{-1}) = \mathcal{S}(\nu). \tag{4.1}$$

Proof Introducing the kernel of the fermionic projector by (3.2) and making use of Definition 2.5, for all $u, w \in V$ and $\xi \in \mathcal{M}$ we obtain

$$\begin{aligned} & \langle u | P[U\nu U^{-1}](\xi) w \rangle \\ &= \langle u | \int_{\hat{K}} e^{ip\xi} d(U\nu U^{-1})(p) w \rangle \\ &= \int_{\hat{K}} e^{ip\xi} d \langle u | U\nu(p)U^{-1} w \rangle = \int_{\hat{K}} e^{ip\xi} d \langle U^{-1}u | \nu(p)U^{-1}w \rangle \\ &= \langle U^{-1}u | \int_{\hat{K}} e^{ip\xi} d\nu(p)U^{-1}w \rangle = \langle u | U \int_{\hat{K}} e^{ip\xi} d\nu(p)U^{-1}w \rangle \\ &= \langle u | U P[\nu](\xi)U^{-1}w \rangle \end{aligned}$$

for any negative definite measure $d\nu \in \mathfrak{N}\mathfrak{d}\mathfrak{m}$ and any unitary matrix U (with respect to $\langle \cdot | \cdot \rangle$). Thus non-degeneracy of the indefinite inner product implies that

$$P[U\nu U^{-1}] = U P[\nu]U^{-1}.$$

Henceforth, employing Lemma 4.3, we deduce that the spectral weight of the closed chain A is unaffected by unitary similarity, i.e.

$$|A[U\nu U^{-1}](\xi)| = |U A[\nu](\xi)U^{-1}| = |A[\nu](\xi)| \quad \text{for all } \xi \in \mathcal{M}.$$

Analogously, for every $\xi \in \mathcal{M}$ we obtain

$$|A[U\nu U^{-1}](\xi)^2| = |(U A[\nu](\xi)U^{-1})^2| = |A[\nu](\xi)^2|,$$

thus implying that

$$\mathcal{L}[U\nu U^{-1}](\xi) = \mathcal{L}[\nu](\xi) \quad \text{for all } \xi \in \mathcal{M}$$

as well as $\mathcal{S}(U\nu U^{-1}) = \mathcal{S}(\nu)$. This completes the proof. \square

Lemma 4.5 Let $f > 0$ and assume that $(d\nu^{(k)})_{k \in \mathbb{N}}$ is a sequence in $\mathfrak{N}\mathfrak{d}\mathfrak{m}$ such that

$$\text{Tr}_V(-S\nu^{(k)}(\hat{K})) \leq f \quad \text{for all } k \in \mathbb{N}$$

(where S denotes the signature matrix). Then there exists a positive constant $C > 0$ in such a way that $d\|\nu^{(k)}\| \leq C$ for all $k \in \mathbb{N}$, where $d\|\cdot\|$ denotes the total variation according to Definition 3.2.

Proof For convenience, we fix an arbitrary integer $k \in \mathbb{N}$ and let $d\nu = d\nu^{(k)}$. Next, we let $(\mathbf{e}_i)_{i=1, \dots, 2n}$ be a pseudo-orthonormal basis of V with signature matrix S such that (2.1) is satisfied. Then $d \langle \mathbf{e}_i | \nu \mathbf{e}_j \rangle$ is a finite complex measure in $\mathbf{M}_{\mathbb{C}}(\hat{K})$ for every $i, j \in \{1, \dots, 2n\}$ according to Definition 3.1, i.e.

$$d\|\langle \mathbf{e}_i | \nu \mathbf{e}_j \rangle\| = d\|\langle \mathbf{e}_i | \nu \mathbf{e}_j \rangle(\hat{K})\| < \infty \quad \text{for all } i, j = 1, \dots, 2n.$$

Employing the definition of the total variation of complex measures and applying the Schwarz inequality (see e.g. [25, Lemma A.13] or [21, ineq. (2.3.9)]), we obtain

$$\begin{aligned}
 d\|\langle e_i \mid \nu e_j \rangle\| &= \sup_{n \in \mathbb{N}} \sum |\langle e_i \mid \nu(E_n) e_j \rangle| = \sup_{n \in \mathbb{N}} \sum |\langle e_i \mid -\nu(E_n) e_j \rangle| \\
 &\leq \sup_{n \in \mathbb{N}} \sum \sqrt{|\langle e_i \mid -\nu(E_n) e_i \rangle|} \sqrt{|\langle e_j \mid -\nu(E_n) e_j \rangle|},
 \end{aligned}$$

where the supremum is taken over all partitions $(E_n)_{n \in \mathbb{N}}$ of \hat{K} (cf. [31, Chapter 6]). Applying Young’s inequality (see e.g. [2, Sect. 1]), for all $i, j \in \{1, \dots, 2n\}$ we arrive at

$$\begin{aligned}
 d\|\langle e_i \mid \nu e_j \rangle\| &\leq \frac{1}{2} \sup_{n \in \mathbb{N}} \sum (|\langle e_i \mid -\nu(E_n) e_i \rangle| + |\langle e_j \mid -\nu(E_n) e_j \rangle|) \\
 &\leq \frac{1}{2} \left[\sup_{n \in \mathbb{N}} \sum |\langle e_i \mid -\nu(E_n) e_i \rangle| + \sup_{n \in \mathbb{N}} \sum |\langle e_j \mid -\nu(E_n) e_j \rangle| \right] \\
 &= \frac{1}{2} (d\|\langle e_i \mid \nu e_i \rangle\| + d\|\langle e_j \mid \nu e_j \rangle\|).
 \end{aligned}$$

Due to the fact that $d\|\langle e_i \mid -\nu e_i \rangle\|$ is a positive measure for each $i \in \{1, \dots, 2n\}$, the total variation $d\|\langle e_i \mid \nu e_j \rangle\|$ is bounded by

$$d\|\langle e_i \mid \nu e_j \rangle\| \leq \sum_{i=1}^{2n} d\|\langle e_i \mid \nu e_i \rangle\| = \sum_{i=1}^{2n} \langle e_i \mid -\nu(\hat{K}) e_i \rangle$$

for all $i, j \in \{1, \dots, 2n\}$. The last expression can be estimated by

$$\sum_{i=1}^{2n} \langle e_i \mid -\nu(\hat{K}) e_i \rangle = \sum_{i=1}^{2n} \langle e_i \mid -S\nu(\hat{K}) e_i \rangle = \text{Tr}_V (-S\nu(\hat{K})) \leq f,$$

thus completing the proof. □

In the case that the spectral weight is bounded (in analogy to Finster [11, Theorem 6.1]), we obtain the following result:

Lemma 4.6 *Let $f > 0$ and assume that $(d\nu^{(k)})_{k \in \mathbb{N}}$ is a sequence in \mathfrak{M} such that*

$$|\nu^{(k)}(\hat{K})| \leq f \quad \text{for all } k \in \mathbb{N}$$

(where $|\cdot|$ denotes the spectral weight). Then there is a sequence $(U_k)_{k \in \mathbb{N}}$ of unitary operators on V (with respect to $\langle \cdot \mid \cdot \rangle$) as well as a positive constant $C > 0$ such that $d\|U_k \nu^{(k)} U_k^{-1}\| \leq C$ for all $k \in \mathbb{N}$ (where $d\|\cdot\|$ denotes the total variation according to Definition 3.2).

For the proof of this result we make use of the next lemma:

Lemma 4.7 *Let W be a finite-dimensional vector space and let $T \in L(W)$. Then for any sequence $(T_n)_{n \in \mathbb{N}}$ of operators in $L(W)$ with $\|T_n - T\| \rightarrow 0$ as $n \rightarrow \infty$ (where $\|\cdot\|$ denotes any norm on $L(W)$), the eigenvalues of T_n converge to those of T .*

Proof See [24, Chapter II, Sect. 5-1]. □

Proof of Lemma 4.6 The basic idea is to make use of Finster [12, Lemma 4.4]. For convenience, we fix an arbitrary integer $k \in \mathbb{N}$ and let $d\nu = d\nu^{(k)}$. Moreover, let $(\epsilon_i)_{i=1, \dots, 2n}$ be a pseudo-orthonormal basis of V with signature matrix S such that (2.1) is satisfied (see for instance [21, Sect. 2.3] or [25, Sect. 3.3]). Since V is a finite-dimensional vector space, all norms on $L(V)$ are equivalent, and one of these norms is given by

$$\|A\|_1 = \max_{j=1, \dots, 2n} \sum_{i=1}^{2n} |\langle \epsilon_i \mid A \epsilon_j \rangle| \tag{4.2}$$

for any $A \in L(V)$, where $|\cdot|$ denotes the absolute value. Moreover, for any unitary matrix U in $L(V)$ (with respect to $\langle \cdot \mid \cdot \rangle$), we may introduce another pseudo-orthonormal basis $(f_j)_{j=1, \dots, 2n}$ by

$$f_i := U^{-1} \epsilon_i \quad \text{for all } i = 1, \dots, 2n. \tag{4.3}$$

Making use of $U^* = U^{-1}$, for all $i, j = 1, \dots, 2n$ we obtain

$$d \langle \epsilon_i \mid U \nu U^{-1} \epsilon_j \rangle = d \langle U^* \epsilon_i \mid \nu U^{-1} \epsilon_j \rangle = d \langle f_i \mid \nu f_j \rangle. \tag{4.4}$$

Since $d\nu$ is a negative definite measure, the operator $-\nu(\hat{K})$ is positive (2.5). Thus in view of Finster [12, Lemma 4.4], for any $\epsilon > 0$ there is a unitary matrix $U = U(\epsilon)$ in $L(V)$ (with respect to $\langle \cdot \mid \cdot \rangle$) so that $U \nu(\hat{K}) U^{-1}$ is diagonal, up to an arbitrarily small error term $\Delta\nu(\hat{K})$ with $\|\Delta\nu(\hat{K})\|_1 < \epsilon$. Since $k \in \mathbb{N}$ is arbitrary, we thus obtain a sequence of negative definite measures $(U_k d\nu^{(k)} U_k^{-1})_{k \in \mathbb{N}}$.

Next, in order to prove that $(U_k d\nu^{(k)} U_k^{-1})_{k \in \mathbb{N}}$ is bounded with respect to the total variation defined by (3.1), for each $k \in \mathbb{N}$ we consider the basis $(f_i)_{i=1, \dots, 2n}$ given by (4.3) with respect to the unitary matrix $U = U_k$. Accordingly, each $d \langle f_i \mid \nu f_j \rangle$ is a finite complex measure in $\mathbf{M}_{\mathbb{C}}(\hat{K})$ in view of Definition 3.1,

$$d\| \langle f_i \mid \nu f_j \rangle \| = d\| \langle f_i \mid \nu f_j \rangle \|(\hat{K}) < \infty \quad \text{for all } i, j = 1, \dots, 2n.$$

Employing the definition of the total variation of complex measures and applying the Schwarz inequality and Young’s inequality in analogy to the proof of Lemma 4.5, we arrive at

$$d\| \langle f_i \mid \nu f_j \rangle \| \leq \frac{1}{2} (d\| \langle f_i \mid \nu f_i \rangle \| + d\| \langle f_j \mid \nu f_j \rangle \|)$$

for all $i, j \in \{1, \dots, 2n\}$. Since $S = S^{-1}$ and $U^* = S^{-1} U^\dagger S$ in view of Gohberg et al. [21, Eq. (4.1.3)] (where U^\dagger denotes the adjoint with respect to $\langle \cdot \mid \cdot \rangle$) and U^* the adjoint with respect to $\langle \cdot \mid \cdot \rangle$), for all $i = 1, \dots, 2n$ we obtain

$$\begin{aligned} & d\| \langle f_i \mid -\nu f_i \rangle \| \\ &= \langle f_i \mid -\nu(\hat{K}) f_i \rangle \leq \sum_{i,j=1}^{2n} \left| \langle U^{-1} \epsilon_i \mid \nu(\hat{K}) U^{-1} \epsilon_j \rangle \right| \\ &\leq \sum_{i,j=1}^{2n} \left| \langle S U^* S S \epsilon_i \mid S \nu(\hat{K}) U^{-1} \epsilon_j \rangle \right| = \sum_{i,j=1}^{2n} \left| \langle U^\dagger \epsilon_i \mid \nu(\hat{K}) U^{-1} \epsilon_j \rangle \right| |s_i| \\ &= \sum_{i,j=1}^{2n} \left| \langle \epsilon_i \mid U \nu(\hat{K}) U^{-1} \epsilon_j \rangle \right| \stackrel{(4.2)}{\leq} 2n \|U \nu(\hat{K}) U^{-1}\|_1, \end{aligned}$$

where we made use of $S\mathbf{e}_i = s_i \mathbf{e}_i$ with $|s_i| = |\langle \mathbf{e}_i | S\mathbf{e}_i \rangle| = 1$ for all $i = 1, \dots, 2n$ and employed the fact that $d \prec f_i | -\nu f_i \succ$ is a positive measure for any $i \in \{1, \dots, 2n\}$.

Taken the previous results together, by (4.4) we obtain the inequality

$$d\|\langle \mathbf{e}_i | U \nu U^{-1} \mathbf{e}_i \rangle\| = d\|\langle f_i | -\nu f_i \rangle\| \leq 2n \|U \nu(\hat{K}) U^{-1}\|_1 \tag{4.5}$$

for all $i = 1, \dots, 2n$. Thus it only remains to find an upper bound for $\|U \nu(\hat{K}) U^{-1}\|_1$ in terms of f by establishing a connection to the spectral weight $|\nu(\hat{K})|$. To this end we exploit the fact that $U \nu(\hat{K}) U^{-1}$ is diagonal according to Finster [12, Lemma 4.4], up to an arbitrarily small error term $\Delta\nu(\hat{K})$,

$$U \nu(\hat{K}) U^{-1} = \text{diag}(\tilde{\lambda}_1(U), \dots, \tilde{\lambda}_{2n}(U)) + \Delta\nu(\hat{K}).$$

Denoting the eigenvalues of $U \nu(\hat{K}) U^{-1}$ by $\lambda_i(U)$ for all $i = 1, \dots, 2n$, by choosing the error term $\Delta\nu(\hat{K})$ sufficiently small, in virtue of Lemma 4.7 we can arrange that

$$\sum_{i=1}^{2n} |\tilde{\lambda}_i(U) - \lambda_i(U)| < \varepsilon \quad \text{for any } \varepsilon > 0.$$

Since the off-diagonal elements $\|\Delta\nu(\hat{K})\|_1 < \varepsilon$ are arbitrarily small, we thus obtain

$$\|U \nu(\hat{K}) U^{-1}\|_1 \leq \|\text{diag}(\tilde{\lambda}_1(U), \dots, \tilde{\lambda}_{2n}(U))\|_1 + \|\Delta\nu(\hat{K})\|_1 \leq \sum_{i=1}^{2n} |\lambda_i(U)| + 2\varepsilon.$$

Applying Lemma 4.3, we infer that $|\nu(\hat{K})| = |U \nu(\hat{K}) U^{-1}|$ (where $|\cdot|$ again denotes the spectral weight). Choosing $\varepsilon < 1/2$, we arrive at

$$\|U \nu(\hat{K}) U^{-1}\|_1 \leq |\nu(\hat{K})| + 1 \leq f + 1.$$

Hence in view of Definition 3.2 and (4.5), we finally obtain

$$d\|U_k \nu^{(k)} U_k^{-1}\| = \sum_{i,j=1}^{2n} d\|\langle \mathbf{e}_i | U_k \nu^{(k)} U_k^{-1} \mathbf{e}_j \rangle\| \leq (2n)^3 (f + 1) =: C.$$

This completes the proof. □

The major simplification when restricting attention to compact subsets is that any minimizing sequence is uniformly tight a priori. As a consequence, we may apply Prohorov’s theorem to each component, thereby obtaining the desired minimizer.

4.2 Preparatory result

Given a sequence of negative definite measures which is bounded and uniformly tight, we employ Prohorov’s theorem to prove that a subsequence thereof converges weakly (see Definition 3.7) to a negative definite measure:

Lemma 4.8 *Let $(d\nu_k)_{k \in \mathbb{N}}$ be a sequence of negative definite measures in \mathfrak{M}^{dm} with the following properties:*

- (a) *There is a constant $C > 0$ such that $d|\nu_k|(\hat{K}) \leq C$ for all $k \in \mathbb{N}$.*
- (b) *The sequence $(d\nu_k)_{k \in \mathbb{N}}$ is uniformly tight in the sense that, for every $\varepsilon > 0$, there is a compact subset $K_\varepsilon \subset \hat{K}$ such that $d|\nu_k|(\hat{K} \setminus K_\varepsilon) < \varepsilon$ for all $k \in \mathbb{N}$.*

Then a subsequence of $(d\nu_k)_{k \in \mathbb{N}}$ converges weakly to some negative definite measure $d\nu$.

Proof The main idea is to apply Prohorov’s theorem. More precisely, let $(\mathbf{e}_i)_{i=1, \dots, 2n}$ be a pseudo-orthonormal basis of V satisfying (2.1), and for every $k \in \mathbb{N}$ we denote by $d|\nu_k|$ the corresponding variation of $d\nu_k$ according to Definition 3.2. Decomposing the complex measure $d \langle \mathbf{e}_i | -\nu_k \mathbf{e}_j \rangle$ into its real and imaginary part,

$$d \langle \mathbf{e}_i | -\nu_k \mathbf{e}_j \rangle = \operatorname{Re} d \langle \mathbf{e}_i | -\nu_k \mathbf{e}_j \rangle + i \operatorname{Im} d \langle \mathbf{e}_i | -\nu_k \mathbf{e}_j \rangle,$$

and introducing the (positive) measures

$$d\mathfrak{R}_{[i,j],k}^\pm := \operatorname{Re} d \langle \mathbf{e}_i | -\nu_k \mathbf{e}_j \rangle^\pm \quad \text{and} \quad d\mathfrak{I}_{[i,j],k}^\pm := \operatorname{Im} d \langle \mathbf{e}_i | -\nu_k \mathbf{e}_j \rangle^\pm$$

by applying the Jordan decomposition [22, Sect. 29], we arrive at

$$d \langle \mathbf{e}_i | -\nu_k \mathbf{e}_j \rangle = d\mathfrak{R}_{[i,j],k}^+ - d\mathfrak{R}_{[i,j],k}^- + i d\mathfrak{I}_{[i,j],k}^+ - i d\mathfrak{I}_{[i,j],k}^-$$

for all $i, j \in \{1, \dots, 2n\}$ and each $k \in \mathbb{N}$. Then the conditions (a) and (b) imply that the sequences $(d\mathfrak{R}_{[i,j],k}^\pm)_{k \in \mathbb{N}}$ and $(d\mathfrak{I}_{[i,j],k}^\pm)_{k \in \mathbb{N}}$ are bounded and uniformly tight for all $i, j = 1, \dots, 2n$. Iteratively applying Prohorov’s theorem, we deduce that $(d\nu_k)_{k \in \mathbb{N}}$ contains a subsequence (which for convenience we again denote by $(d\nu_k)_{k \in \mathbb{N}}$) such that the corresponding sequences $(d\mathfrak{R}_{[i,j],k}^\pm)_{k \in \mathbb{N}}$ and $(d\mathfrak{I}_{[i,j],k}^\pm)_{k \in \mathbb{N}}$ weakly converge to (positive) measures $d\mathfrak{R}_{[i,j]}^\pm$ and $d\mathfrak{I}_{[i,j]}^\pm$, respectively, i.e.

$$d\mathfrak{R}_{[i,j],k}^\pm \rightharpoonup d\mathfrak{R}_{[i,j]}^\pm \quad \text{and} \quad d\mathfrak{I}_{[i,j],k}^\pm \rightharpoonup d\mathfrak{I}_{[i,j]}^\pm$$

for all $i, j \in \{1, \dots, 2n\}$ as $k \rightarrow \infty$. Introducing the measures

$$d\nu_{i,j} := d\mathfrak{R}_{[i,j]}^+ - d\mathfrak{R}_{[i,j]}^- + i d\mathfrak{I}_{[i,j]}^+ - i d\mathfrak{I}_{[i,j]}^- \quad \text{for all } i, j \in \{1, \dots, 2n\},$$

for every $\phi \in C_b(\hat{K})$ we obtain

$$\lim_{k \rightarrow \infty} \int_{\hat{K}} \phi d \langle \mathbf{e}_i | -\nu_k \mathbf{e}_j \rangle = \int_{\hat{K}} \phi d\nu_{i,j} \quad \text{for all } i, j \in \{1, \dots, 2n\}.$$

Following the proof of Lemma 3.3, we introduce the operator-valued measure $d\nu$ by

$$\nu(\Omega) := \begin{pmatrix} \nu_{1,1}(\Omega) & \cdots & \nu_{1,2n}(\Omega) \\ \vdots & \ddots & \vdots \\ \nu_{n,1}(\Omega) & \cdots & \nu_{n,2n}(\Omega) \\ -\nu_{n+1,1}(\Omega) & \cdots & -\nu_{n+1,2n}(\Omega) \\ \vdots & \ddots & \vdots \\ -\nu_{1,2n}(\Omega) & \cdots & -\nu_{2n,2n}(\Omega) \end{pmatrix} \in L(V)$$

for every $\Omega \in \mathcal{B}(\hat{K})$. The measure $d\nu$ has the property that, for all $i, j \in \{1, \dots, 2n\}$,

$$d \langle \mathbf{e}_i | \nu \mathbf{e}_j \rangle = d \langle \mathbf{e}_i | S \nu \mathbf{e}_j \rangle = d\nu_{i,j} \in \mathbf{M}_{\mathbb{C}}(\hat{K})$$

is a complex measure. For elements $u = \sum_{j=1}^{2n} \alpha_j(u) \mathbf{e}_j$ and $v = \sum_{j=1}^{2n} \alpha_j(v) \mathbf{e}_j$ in V , by linearity we conclude that $d \langle u | \nu v \rangle \in \mathbf{M}_{\mathbb{C}}(\hat{K})$ for all $u, v \in V$. Hence $d\nu$ is an

operator-valued measure in the sense of Definition 3.1, and by linearity we arrive at

$$\begin{aligned} \lim_{k \rightarrow \infty} \int_{\hat{K}} \phi \, d \langle u \mid -v_k v \rangle &= \lim_{k \rightarrow \infty} \sum_{\ell, m=1}^{2n} \overline{\alpha_\ell(u)} \alpha_m(v) \int_{\hat{K}} \phi \, d \langle \epsilon_\ell \mid -v_k \epsilon_m \rangle \\ &= \sum_{\ell, m=1}^{2n} \overline{\alpha_\ell(u)} \alpha_m(v) \int_{\hat{K}} \phi \, d \langle \epsilon_\ell \mid -v \epsilon_m \rangle = \int_{\hat{K}} \phi \, d \langle u \mid -v v \rangle \end{aligned}$$

for all $\phi \in C_b(\hat{K})$ and $u, v \in V$. This yields weak convergence $d v_k \rightarrow d v$ of operator-valued measures in the sense of Definition 3.7. In particular, $d \|v\| < \infty$.

It remains to show that $d v$ is indeed negative definite. To this end, we need to prove that $d \langle u \mid -v u \rangle$ is a positive measure for all $u \in V$. We point out that, by assumption, the measures $d \langle u \mid -v_k u \rangle$ are positive for each $u \in V$ and all $k \in \mathbb{N}$. Assume now, for some $u \in V$, that $d \mu_u := d \langle u \mid -v u \rangle$ is a signed measure with

$$d \mu_u = d \mu_u^+ - d \mu_u^-$$

such that $d \mu_u^-$ is non-zero. In this case, there is $\Omega \in \mathcal{B}(\hat{K})$ with the property that

$$\mu_u^+(\Omega) < \mu_u^-(\Omega)$$

(assuming conversely that $\mu_u^+(\Omega) \geq \mu_u^-(\Omega)$ for all $\Omega \in \mathcal{B}(\hat{K})$, then the measure $d \mu_u$ is non-negative, implying that $d \mu_u^- = 0$). In virtue of Ulam’s theorem we know that both measures $d \mu_u^\pm$ are regular on \hat{K} . As a consequence, there is an open set $U \supset \Omega$ and a compact set $K \subset \Omega$ such that $\mu_u^+(U) < \mu_u^-(K)$. Hence a partition of unity yields a function $\phi \in C_c(U; [0, 1])$ with $\text{supp } \phi \subset U$ and $\phi|_K \equiv 1$, thus giving rise to the contradiction

$$0 \leq \lim_{k \rightarrow \infty} \int_{\hat{K}} \phi \, d \langle u \mid -v_k u \rangle = \int_{\hat{K}} \phi \, d \langle u \mid -v u \rangle \leq \mu_u^+(U) - \mu_u^-(K) < 0.$$

This completes the proof. □

4.3 Proof of the existence theorem

In order for proving Theorem 4.1, we require some more preparatory results. The proof of Theorem 4.1 will be completed towards the end of this subsection. To begin with, let us state the following proposition.

Proposition 4.9 *Let $(d v^{(k)})_{j \in \mathbb{N}}$ be a sequence of negative definite measures in \mathfrak{Ndm} which converges weakly to some negative definite measure $d v \in \mathfrak{Ndm}$. Then*

$$\lim_{j \rightarrow \infty} \mathcal{L}[v^{(j)}](\xi) = \mathcal{L}[v](\xi) \quad \text{for all } \xi \in \mathcal{M}$$

and

$$\mathcal{S}(v) \leq \liminf_{j \rightarrow \infty} \mathcal{S}(v^{(j)}).$$

Proof Let us first consider the behavior of the kernel of the fermionic projector and the closed chain. For convenience, we introduce the notation $P_j(\xi) := P[v^{(j)}](\xi)$ as well as $A_j(\xi) :=$

$A[v^{(j)}](\xi)$ for all $j \in \mathbb{N}$ and arbitrary $\xi \in \mathcal{M}$. Then weak convergence (see Definition 3.7 and the remark thereafter) implies that

$$\begin{aligned} \lim_{j \rightarrow \infty} \langle u \mid P_j(\xi) v \rangle &= \lim_{j \rightarrow \infty} \int_{\hat{K}} e^{ip\xi} d \langle u \mid v^{(j)}(p) v \rangle \\ &= \int_{\hat{K}} e^{ip\xi} d \langle u \mid v(p) v \rangle = \langle u \mid P[v](\xi) v \rangle \end{aligned}$$

for all $u, v \in V$ and arbitrary $\xi \in \mathcal{M}$. Given a pseudo-orthonormal basis $(\epsilon_i)_{i=1, \dots, 2n}$ of V satisfying (2.1), we thus obtain

$$\lim_{j \rightarrow \infty} \langle \epsilon_\alpha \mid P_j(\xi) \epsilon_\beta \rangle = \lim_{j \rightarrow \infty} \langle S\epsilon_\alpha \mid P_j(\xi) \epsilon_\beta \rangle = \langle S\epsilon_\alpha \mid P[v](\xi) \epsilon_\beta \rangle = \langle \epsilon_\alpha \mid P[v](\xi) \epsilon_\beta \rangle$$

for all $\alpha, \beta \in \{1, \dots, 2n\}$ and arbitrary $\xi \in \mathcal{M}$. From this we deduce that

$$\lim_{j \rightarrow \infty} (A_j(\xi))_{\alpha, \beta} = \lim_{j \rightarrow \infty} (P_j(\xi) P_j(-\xi))_{\alpha, \beta} = (P[v](\xi) P[v](-\xi))_{\alpha, \beta} = (A[v](\xi))_{\alpha, \beta}$$

for all $\alpha, \beta \in \{1, \dots, 2n\}$ and arbitrary $\xi \in \mathcal{M}$. By continuity of the spectral weight,

$$\lim_{j \rightarrow \infty} \mathcal{L}[v^{(j)}](\xi) = \mathcal{L}[v](\xi) \quad \text{for all } \xi \in \mathcal{M}.$$

The second statement follows from Fatou’s lemma (see e.g. [23, Theorem 16.4]),

$$S(v) = \int_{\mathcal{M}} \mathcal{L}[v](\xi) d\mu(\xi) = \int_{\mathcal{M}} \liminf_{j \rightarrow \infty} \mathcal{L}[v^{(j)}](\xi) d\mu(\xi) \leq \liminf_{j \rightarrow \infty} \int_{\mathcal{M}} \mathcal{L}[v^{(j)}](\xi) d\mu(\xi).$$

This completes the proof. □

Proposition 4.10 *Let $(dv^{(j)})_{j \in \mathbb{N}}$ be a sequence of negative definite measures in \mathfrak{Ndm} which converges weakly to some negative definite measure $dv \in \mathfrak{Ndm}$. Then*

$$\lim_{j \rightarrow \infty} \text{Tr}_V(v^{(j)}(\hat{K})) = \text{Tr}_V(v(\hat{K}))$$

as well as

$$\lim_{j \rightarrow \infty} \text{Tr}_V(-Sv^{(j)}(\hat{K})) = \text{Tr}_V(-Sv(\hat{K})) \quad \text{and} \quad \lim_{j \rightarrow \infty} |v^{(j)}(\hat{K})| = |v(\hat{K})|.$$

Proof By weak convergence, the first two equalities can be verified as follows:

$$\begin{aligned} \lim_{j \rightarrow \infty} \text{Tr}_V(v^{(j)}(\hat{K})) &= \lim_{j \rightarrow \infty} \sum_{\alpha=1}^{2n} \langle \epsilon_\alpha \mid v^{(j)}(\hat{K}) \epsilon_\alpha \rangle = \lim_{j \rightarrow \infty} \sum_{\alpha=1}^{2n} \int_{\hat{K}} d \langle S\epsilon_\alpha \mid v^{(j)}(p) \epsilon_\alpha \rangle \\ &= \sum_{\alpha=1}^{2n} \int_{\hat{K}} d \langle S\epsilon_\alpha \mid v(p) \epsilon_\alpha \rangle = \sum_{\alpha=1}^{2n} \langle \epsilon_\alpha \mid v(\hat{K}) \epsilon_\alpha \rangle = \text{Tr}_V(v(\hat{K})), \end{aligned}$$

and analogously

$$\lim_{j \rightarrow \infty} \text{Tr}_V(-Sv^{(j)}(\hat{K})) = \text{Tr}_V(-Sv(\hat{K})).$$

In order to prove the remaining equality, we essentially make use of the fact that the spectral weight is continuous. More precisely, by continuity of the absolute value and weak convergence we obtain

$$\begin{aligned} & \lim_{j \rightarrow \infty} \|v^{(j)}(\hat{K}) - v(\hat{K})\|_1 \\ & \leq \lim_{j \rightarrow \infty} \sum_{\alpha, \beta=1}^{2n} \left| \langle \epsilon_\alpha \mid (v^{(j)}(\hat{K}) - v(\hat{K})) \epsilon_\beta \rangle \right| \\ & = \lim_{j \rightarrow \infty} \sum_{\alpha, \beta=1}^{2n} \left| \int_{\hat{K}} d \prec S \epsilon_\alpha \mid v^{(j)}(p) \epsilon_\beta \succ - \int_{\hat{K}} d \prec S \epsilon_\alpha \mid v(p) \epsilon_\beta \succ \right| = 0 \end{aligned}$$

(where $\|\cdot\|_1$ is given by (4.2)). Denoting the eigenvalues of $v(\hat{K})$ by $(\lambda_i)_{i=1, \dots, 2n}$ and those of $v^{(j)}(\hat{K})$ for every $j \in \mathbb{N}$ by $(\lambda_i^{(j)})_{i=1, \dots, 2n}$, by applying Lemma 4.7 together with the inverse triangle inequality we thus arrive at

$$\lim_{j \rightarrow \infty} \left| |v^{(j)}(\hat{K})| - |v(\hat{K})| \right| \leq \lim_{j \rightarrow \infty} \sum_{i=1}^{2n} \left| |\lambda_i^{(j)}| - |\lambda_i| \right| \leq \lim_{j \rightarrow \infty} \sum_{i=1}^{2n} |\lambda_i^{(j)} - \lambda_i| = 0.$$

This completes the proof. □

After these preliminaries we are finally in the position to prove Theorem 4.1.

Proof of Theorem 4.1 Let us first assume that the side conditions (3.5) are satisfied. In this case, Lemma 4.6 yields a sequence of unitary operators $(U_j)_{j \in \mathbb{N}}$ in $L(V)$ (with respect to $\prec \cdot \mid \cdot \succ$) as well as a constant $C > 0$ such that

$$d \|U_j v^{(j)} U_j^{-1}\| \leq C \quad \text{for all } j \in \mathbb{N}.$$

Since $\hat{K} \subset \hat{\mathcal{M}}$ is compact, the sequence of measures $(dv^{(j)})_{j \in \mathbb{N}}$ is uniformly tight. As a consequence, we may apply Lemma 4.8 in order to conclude that a subsequence of $(U_j dv^{(j)} U_j^{-1})_{j \in \mathbb{N}}$ converges weakly to some negative definite measure $d\nu \in \mathfrak{Ndm}$,

$$d\tilde{v}^{(j_k)} := U_{j_k} dv^{(j_k)} U_{j_k}^{-1} \rightharpoonup d\nu \quad \text{weakly}.$$

Making use of (4.1), from Proposition 4.9 we deduce that

$$\mathcal{S}(v) \leq \lim_{k \rightarrow \infty} \mathcal{S}(\tilde{v}^{(j_k)}) = \lim_{k \rightarrow \infty} \mathcal{S}(v^{(j_k)}).$$

In the case that the constraints (3.4) are imposed, the above arguments remain valid by applying Lemma 4.5 instead of Lemma 4.6 and choosing $U_j = \mathbb{1}_V$ for all $j \in \mathbb{N}$.

Thus it only remains to prove that the measure $d\nu$ satisfies the conditions (3.4) or (3.5), respectively. In both cases, this follows readily from Proposition 4.10. In particular, the limit measure $d\nu$ is non-trivial, which completes the proof. □

4.4 Imposing a boundedness constraint

Let us finally establish a connection to the boundedness constraint as considered in [12, Section 4] (which originally was proposed in [10, Eq. (3.5.10)] as a constraint for the causal

action principle). In the homogeneous setting, for any operator-valued measure $d\omega \in \mathfrak{D}\mathfrak{v}\mathfrak{m}$ we introduce the mapping $\mathfrak{t}[\omega] : \mathcal{M} \rightarrow \mathbb{R}_0^+$ by

$$\mathfrak{t}[\omega](\xi) := |A[\omega](\xi)|^2 \quad \text{for all } \xi \in \mathcal{M} .$$

We then define the functional $\mathcal{T} : \mathfrak{D}\mathfrak{v}\mathfrak{m} \rightarrow \mathbb{R}_0^+ \cup \{+\infty\}$ by

$$\mathcal{T}(\omega) := \int_{\mathcal{M}} \mathfrak{t}[\omega](\xi) \, d\mu(\xi) = \int_{\mathcal{M}} |A[\omega](\xi)|^2 \, d\mu(\xi) .$$

Given $C > 0$, the corresponding *boundedness constraint* reads

$$\mathcal{T}(\omega) \leq C . \tag{4.6}$$

In analogy to Theorem 4.1 we then obtain the following existence result:

Theorem 4.11 *Assume that $(d\nu^{(j)})_{j \in \mathbb{N}}$ is a minimizing sequence of negative definite measures in $\mathfrak{N}\mathfrak{D}\mathfrak{m}$ for the causal variational principle (3.3) with respect to the side conditions (3.6) and (4.6) for some positive constants $c, C > 0$. Then there exists a sequence of unitary operators $(U_j)_{j \in \mathbb{N}}$ on V (with respect to $\langle \cdot | \cdot \rangle$) as well as a subsequence $(d\nu^{(j_k)})_{k \in \mathbb{N}}$ such that the sequence $(U_{j_k} d\nu^{(j_k)} U_{j_k}^{-1})_{k \in \mathbb{N}}$ converges weakly to some non-trivial negative definite measure $d\nu \neq 0$. Moreover,*

$$\mathcal{S}(v) \leq \liminf_{k \rightarrow \infty} \mathcal{S}(v^{(j_k)}) ,$$

and the limit measure $d\nu \in \mathfrak{N}\mathfrak{D}\mathfrak{m}$ satisfies the side conditions

$$\text{Tr}_V(v(\hat{K})) = c \quad \text{and} \quad \mathcal{T}(v) \leq C .$$

In particular, the limit measure $d\nu$ is a non-trivial minimizer of the causal variational principle (3.3) with respect to the side conditions (3.6) and (4.6).

For the proof of Theorem 4.11 we make use of the following result:

Proposition 4.12 *Whenever $d\nu \in \mathfrak{N}\mathfrak{D}\mathfrak{m}$ is a negative definite measure satisfying the boundedness constraint (4.6), it satisfies condition (3.8) for some constant $f > 0$.*

Proof To begin with, from $\int_{\mathcal{M}} |A[v]|^2 \, d\mu \leq C$ we know that $|A[v]|^2 \in L^1(\mathcal{M})$. By continuity of $|A[v]|^2$ we infer that $\max_{\xi \in \mathcal{M}} |A[v](\xi)|^2 < \infty$, and that the maximum is attained. Let $\xi_{\max} \in \mathcal{M}$ such that

$$|A[v](\xi_{\max})|^2 = \max_{\xi \in \mathcal{M}} |A[v](\xi)|^2 =: c_{\max} .$$

We thus clearly deduce that

$$|A[v](0)|^2 \leq |A[v](\xi_{\max})|^2 \leq C + c_{\max} . \tag{4.7}$$

We now apply (4.7) in order to prove that $|v(\hat{K})| < f$ for some constant $f > 0$. To this end, we essentially employ [12, Lemma 4.4]. More precisely, for any negative definite measure $d\nu$ and arbitrary $\varepsilon > 0$, there is a unitary operator $U \in L(V)$ (with respect to $\langle \cdot | \cdot \rangle$) such that

$$U v(\hat{K}) U^{-1} = -\text{diag}(\tilde{\lambda}_1, \dots, \tilde{\lambda}_{2n}) + \Delta v(\hat{K}) ,$$

where the real parameters $\tilde{\lambda}_i$ ($i = 1, \dots, 2n$) are ordered according to Finster [12, Eq. (2.6)], and $\|\Delta v(\hat{K})\| < \varepsilon$. Denoting by $\{.,.\}$ the anti-commutator, i.e. $\{A, B\} = AB + BA$ for any $A, B \in L(V)$, we thus obtain

$$\begin{aligned}
 U A[v](0) U^{-1} &= (U v(\hat{K}) U^{-1})^2 \\
 &= \text{diag}(\tilde{\lambda}_1^2, \dots, \tilde{\lambda}_{2n}^2) - \left\{ \text{diag}(\tilde{\lambda}_1, \dots, \tilde{\lambda}_{2n}), \Delta v(\hat{K}) \right\} + \Delta v(\hat{K})^2.
 \end{aligned}$$

Since $\|v(\hat{K})\| < \infty$, the absolute values of $\tilde{\lambda}_i$ are bounded for all $i = 1, \dots, 2n$; from this we conclude that the spectrum of $\text{diag}(\tilde{\lambda}_1^2, \dots, \tilde{\lambda}_{2n}^2)$ coincides with the spectrum of $A[v](0)$, up to an arbitrarily small error term (where we applied the fact that the spectra of $A[v](0)$ and $U A[v](0) U^{-1}$ coincide according to Lemma 4.3). In a similar fashion, one can show that the spectra of $v(\hat{K})$ and $-\text{diag}(\tilde{\lambda}_1, \dots, \tilde{\lambda}_{2n})$ coincide, up to an arbitrarily small error term. Neglecting the error terms in what follows, we thus can arrange that

$$|v(\hat{K})| \leq 2 \sum_{i=1}^{2n} |\tilde{\lambda}_i| \quad \text{and} \quad \sum_{i=1}^{2n} \tilde{\lambda}_i^2 \leq 2|A[v](0)|.$$

Employing Jensen’s inequality, we conclude that

$$|v(\hat{K})|^2 \leq 4 \left(\sum_{i=1}^{2n} |\tilde{\lambda}_i| \right)^2 \leq 8n \sum_{i=1}^{2n} |\tilde{\lambda}_i|^2 \leq 16n|A[v](0)|.$$

Applying (4.7), the boundedness constraint gives rise to the desired estimate

$$|v(\hat{K})| < 4\sqrt{n(C + c_{\max})} =: f,$$

which completes the proof. □

This allows us to prove Theorem 4.11:

Proof of Theorem 4.11 We basically combine Proposition 4.12 and Theorem 4.1. To this end let $(d\nu^{(j)})_{j \in \mathbb{N}}$ be a minimizing sequence of negative definite measures which satisfies the side conditions (3.6) and (4.6) for some positive constants $c, C > 0$. Then by Proposition 4.12, there exists $f > 0$ in such a way that condition (3.8) is satisfied for every $j \in \mathbb{N}$. As a consequence, according to Theorem 4.1, there is a sequence of unitary operators $(U_j)_{j \in \mathbb{N}}$ in $L(V)$ (with respect to $\langle \cdot | \cdot \rangle$) such that the sequence $(U_j d\nu^{(j)} U_j^{-1})_{j \in \mathbb{N}}$ contains a subsequence (which for simplicity we again denote by $(U_j d\nu^{(j)} U_j^{-1})_{j \in \mathbb{N}}$) with the property that it converges weakly to some limit measure $d\nu \in \mathfrak{M}$. Applying Fatou’s lemma one can show that

$$S(v) \leq \liminf_{j \rightarrow \infty} S(v^{(j)}) \quad \text{and} \quad \mathcal{T}(v) \leq \liminf_{j \rightarrow \infty} \mathcal{T}(v^{(j)}).$$

By virtue of Proposition 4.10 we conclude that $d\nu$ satisfies condition (3.6), thus implying that $d\nu \neq 0$ is non-zero. This completes the proof. □

Thus for compact subsets of momentum space, Theorem 4.11 gives an alternative proof of [12, Theorem 4.2] considering an additional trace constraint.

4.5 Discussion of the results

This section is finally devoted to discuss the results obtained in Theorem 4.1 and Theorem 4.11. To this end, let us recall the first part of the main theorem obtained in [12, Sect. 4] (see [12, Theorem 4.2]):

Theorem 4.13 *Assume that $(dv_k)_{k \in \mathbb{N}}$ is a sequence of negative definite measures on the bounded set $\hat{K} \subset \hat{M}$ such that the functional \mathcal{T} is bounded by some constant $C > 0$, i.e.*

$$\mathcal{T}(v_k) \leq C \quad \text{for all } k \in \mathbb{N}.$$

Then there is a subsequence $(dv_{k_\ell})_{\ell \in \mathbb{N}}$ as well as a sequence of unitary transformations $(U_\ell)_{\ell \in \mathbb{N}}$ on V (with respect to $\langle \cdot | \cdot \rangle$) such that the measures $U_\ell dv_{k_\ell} U_\ell^{-1}$ converge weakly to a negative definite measure dv with the properties

$$\mathcal{T}(v) \leq \liminf_{k \rightarrow \infty} \mathcal{T}(v_k), \quad \mathcal{S}(v) \leq \liminf_{k \rightarrow \infty} \mathcal{S}(v_k).$$

Theorem 4.13 is stated as a compactness result. Applying it to a minimizing sequence yields statements similar to Finster [12, Theorems 2.2 and 2.3], asserting that the functional \mathcal{S} attains its minimum.

Restricting attention to compact subsets of momentum space, Theorem 4.11 shows that Theorem 4.13 can be extended to variational problems in which apart from the boundedness constraint also a trace constraint is imposed. Theorem 4.1, on the other hand, applies to situations in which the weaker side conditions (3.4) or (3.5) are of central interest. This might be the case for deriving the Euler-Lagrange equations in momentum space which shall be postponed to future projects.

The second part of Finster [12, Theorem 4.2] deals with the counting measure on lattices, thereby considering two different side conditions. It remains an open task to verify if all the above arguments also go through in case of a counting measure instead of Lebesgue measure with respect to the side conditions under consideration.

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Appendix A. Justifying the side conditions

This appendix is devoted to justify and explain the side conditions (3.4)–(3.5). Apart from excluding trivial minimizers in a quite simple way, the following reasoning provides a strong argument for imposing condition (3.6). Given a causal fermion system $(\mathcal{H}, \mathcal{F}, d\rho)$, the so-called local trace $\text{tr}(x)$ defined by

$$\text{tr}(x) = \text{Tr}_{\mathcal{S}_x} (P(x, x)) \quad \text{for all } x \in \text{supp } d\rho$$

is constant on $\text{supp } d\rho$ whenever the measure $d\rho$ is a minimizer of the causal action principle (for details see [14, Sect. 1.1.3, Proposition 1.4.1 and Section 2.5]). Considering homogeneous causal fermion systems, this suggests to impose that

$$\text{Tr}_V(\nu(\hat{K})) = \text{Tr}_V\left(\int_{\hat{K}} d\nu(k)\right) = \text{Tr}_V(P(0)) = \text{Tr}_V(P(x, x)) = c \quad \text{for all } x \in \mathcal{M},$$

thus motivating the side condition (3.6). Following the arguments in [14, Sect. 1.4.1], we shall always assume that $c \neq 0$, thereby excluding trivial minimizers. Let us briefly explain why the quantity $\text{Tr}_V(P(0)) = \text{Tr}_V(\nu(\hat{K}))$ in (3.6) is also referred to as *mass density*.² In order to see that $\text{Tr}_V(P(0))$ can indeed be regarded as a density, let us assume that $(\mathcal{H}, \mathcal{F}, d\rho)$ is a causal fermion system. Whenever $P^\varepsilon(x, y)$ is a regularization of the kernel of the fermionic projector of the vacuum $P(x, y)$ with regularization length ε (where $P(x, y)$ coincides with (2.3), cf. [14, Eq. (1.2.23)]), its trace is given by (see [14, Eq. (2.5.1)])

$$\text{Tr}_{S_x}(P^\varepsilon(x, x)) \sim \frac{m}{\varepsilon^2} \quad \text{for all } x \in \text{supp } d\rho.$$

Making use of the fact that the unit of mass equals one over length, we conclude that $\text{Tr}_{S_x}(P^\varepsilon(x, x))$ is a density, which apparently is proportional to the mass m . Carrying these observations over to $\text{Tr}_V(P(0))$ in the homogeneous case justifies the terminology of mass density.

A possible explanation for introducing the constraint (3.8) is that a similar side condition for the closed chain is imposed in the existence theorem [11, Theorem 6.1]. Since the fermionic projector $P(0) = \nu(\hat{K})$ can be diagonalized (up to an arbitrarily small error term) according to Finster [12, Lemma 4.4], in order to develop the existence theory of minimizers in the homogeneous setting it seems promising to demand that constraint (3.8) is satisfied. On the other hand, following the original ideas in [10] and its modifications in [12], it is natural to impose a boundedness constraint (4.6). The arguments in Sect. 4.4 show that (4.6) already implies condition (3.8).

Let us finally discuss the remaining side condition (3.7). Since working with the spectral weight as appearing in the constraint (3.8) may be awkward, it might seem preferable to work with a similar condition which is easier to handle. Bearing in mind that the operator $\nu(\hat{K})$ may be diagonalized (up to an arbitrarily small error term) in virtue of [12, Lemma 4.4] in such a way that its diagonal entries are ordered according to Finster [12, Eq. (2.6)], the specific form of the signature matrix S (see (2.2)) suggests to replace condition (3.8) by (3.7),

$$\text{Tr}_V(-S\nu(\hat{K})) = f.$$

The same arguments as before illustrate that $\text{Tr}_V(-S\nu(\hat{K}))$ is a density; we refer to this quantity as *particle density*.

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² Note that the quantity $\text{Tr}_V(P(0)) = \text{Tr}_V(\nu(\hat{K}))$ coincides with the *local particle density* f_{loc} as introduced in [12, Eq. (4.4)]. In order to avoid confusion, this notion will not be used in what follows.

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