

MINIMAX PRINCIPLES, HARDY-DIRAC INEQUALITIES,  
AND OPERATOR CORES FOR TWO AND THREE DIMENSIONAL  
COULOMB-DIRAC OPERATORS

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ABSTRACT. For  $n \in \{2, 3\}$  we prove minimax characterisations of eigenvalues in the gap of the  $n$  dimensional Dirac operator with an potential, which may have a Coulomb singularity with a coupling constant up to the critical value  $1/(4 - n)$ . This result implies a so-called Hardy-Dirac inequality, which can be used to define a distinguished self-adjoint extension of the Coulomb-Dirac operator defined on  $C_0^\infty(\mathbb{R}^n \setminus \{0\}; \mathbb{C}^{2(n-1)})$ , as long as the coupling constant does not exceed  $1/(4 - n)$ . We also find an explicit description of an operator core of this operator.

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## 1 INTRODUCTION

The relativistic dynamics of an electron moving in an atomic field is described by a Dirac operator with potential  $V$  having a Coulomb singularity. Since we want to consider such Dirac operators in two and three dimensions simultaneously, we assume throughout the text that  $n \in \{2, 3\}$ . In  $n$  dimensions the relativistic electron corresponds to a  $2(n - 1)$  component spinor and  $V$  is a  $2(n - 1) \times 2(n - 1)$  hermitian matrix function on  $\mathbb{R}^n$ . We say that  $V$  belongs to  $\mathfrak{P}_n$  if for some  $\nu \in [0, 1/(4 - n))$  the inequality  $0 \geq V \geq -\nu/|\cdot| \otimes \mathbb{I}_{\mathbb{C}^{2(n-1)}}$  holds.

This motivates the following question. Does the Dirac operator with potential  $V \in \mathfrak{P}_n \cup \{-1/((4-n)|\cdot|) \otimes \mathbb{I}_{\mathbb{C}^{2(n-1)}}\}$

$$\tilde{D}_n(V) := \begin{cases} -i\boldsymbol{\sigma} \cdot \nabla + \sigma_3 + V & \text{if } n = 2 \\ -i\boldsymbol{\alpha} \cdot \nabla + \beta + V & \text{if } n = 3 \end{cases} \text{ defined on } C_0^\infty(\mathbb{R}^n \setminus \{0\}; \mathbb{C}^{2(n-1)}), \tag{1}$$

have a unique self-adjoint extension? In (1) are  $\boldsymbol{\sigma} = (\sigma_1, \sigma_2)$ ,  $\boldsymbol{\alpha} = (\alpha_1, \alpha_2, \alpha_3)$  vectors;  $\sigma_1, \sigma_2, \sigma_3$  the standard Pauli matrices;  $\alpha_i = \begin{pmatrix} 0_{\mathbb{C}^2} & \sigma_i \\ \sigma_i & 0_{\mathbb{C}^2} \end{pmatrix}$  for  $i \in \{1, 2, 3\}$  and  $\beta = \begin{pmatrix} \mathbb{I}_{\mathbb{C}^2} & 0_{\mathbb{C}^2} \\ 0_{\mathbb{C}^2} & -\mathbb{I}_{\mathbb{C}^2} \end{pmatrix}$ . It is the uniqueness not the existence of a self-adjoint extension that is doubtful. For example the Coulomb-Dirac operator  $\tilde{D}_n(-\nu/|\cdot| \otimes \mathbb{I}_{\mathbb{C}^{2(n-1)}})$  is essentially self-adjoint if  $n = 2, \nu = 0$  or  $n = 3, \nu \in [0, \sqrt{3}/2]$  but for  $n = 2, \nu \in (0, 1/2]$  or  $n = 3, \nu \in (\sqrt{3}/2, 1]$  there are infinitely many self-adjoint extensions (see Lemma 14). Thus it is also natural to ask, whether there is a physically distinguished self-adjoint extension? In fact for  $V \in \mathfrak{P}_n$  there is a unique self-adjoint extension  $D_n(V)$  of  $\tilde{D}_n(V)$ , for which the wave functions in its domain possess finite mean kinetic energy, i.e.  $\mathfrak{D}(D_n(V)) \subset H^{1/2}(\mathbb{R}^n; \mathbb{C}^{2(n-1)})$ . The existence of this distinguished self-adjoint extension is proven in Section 3. There we apply some general results developed in [15]. Note that for  $\nu \in [0, 1/(4-n))$  the domain of the Coulomb-Dirac operator  $D_n(-\nu/|\cdot| \otimes \mathbb{I}_{\mathbb{C}^{2(n-1)}})$  is contained in  $H^{1/2}(\mathbb{R}^n; \mathbb{C}^{2(n-1)})$  and for  $\tilde{D}_n(((n-4)|\cdot|)^{-1} \otimes \mathbb{I}_{\mathbb{C}^{2(n-1)}})$  there is no self-adjoint extension with this property. In this sense  $1/(4-n)$  is a critical constant. At this point we want to mention that in the context of Theorem 5 we define a distinguished self-adjoint extension of  $\tilde{D}_n(-\nu/|\cdot| \otimes \mathbb{I}_{\mathbb{C}^{2(n-1)}})$  for  $\nu \in [0, 1/(4-n)]$ , i.e. the case of a Coulomb potential with the critical coupling constant  $1/(4-n)$  is in particular included here.

Let  $V \in \mathfrak{P}_n$ . As in Proposition 1 in [4] one can prove that there is a gap in the essential spectrum of  $D_n(V)$ . To be more precise

$$\sigma_{\text{ess}}(D_n(V)) = (-\infty, -1] \cup [1, \infty).$$

In 1986 James D. Talman proposed in [17] a formal minimax characterisation of the lowest eigenvalue in the gap of the essential spectrum of the operator  $D_3(V)$ . In this work we prove a minimax characterisation of eigenvalues in the gap of  $D_3(V)$  in the spirit of Talman and an analogous result for  $D_2(V)$ . The exact result is:

**THEOREM 1** (Talman minimax principle). *Let  $V \in \mathfrak{P}_n$ . If the  $k^{\text{th}}$  eigenvalue  $\mu_k$  of  $D_n(V)$  in  $(-1, 1)$ , counted from below with multiplicity, exists, then it is given by*

$$\mu_k = \inf_{\substack{\mathfrak{M} \subset H^{1/2}(\mathbb{R}^n; \mathbb{C}^{n-1}) \\ \dim \mathfrak{M} = k}} \sup_{\psi \in (\mathfrak{M} \oplus H^{1/2}(\mathbb{R}^n; \mathbb{C}^{n-1})) \setminus \{0\}} \frac{\mathbf{d}_n[\psi] + \mathbf{v}[\psi]}{\|\psi\|^2}.$$

Here  $\mathfrak{d}_n$  and  $\mathfrak{v}$  are the quadratic forms associated to the operators  $D_n(0)$  and  $V$ .

About Theorem 1 we want to remark that for  $n = 3$  there is an historical overview of results of the same type in [13] and that for  $n = 2$  there is no comparable result known. Moreover, Theorem 1 improves in the three dimensional case Theorem 3 in [13], which is the best known result for a Dirac operator with an electrostatic potential having strong Coulomb singularity. Furthermore, we give a different proof of the Esteban-Séré minimax principle (see Theorem 2 in [13] and [9]) and prove an analogous result for two dimensional Dirac operators:

**THEOREM 2** (Esteban-Séré minimax principle). *Let  $V \in \mathfrak{P}_n$ . If the  $k^{\text{th}}$  eigenvalue  $\mu_k$  of  $D_n(V)$  in  $(-1, 1)$ , counted from below with multiplicity, exists, then it is given by*

$$\mu_k = \inf_{\substack{\mathfrak{M} \subset P_n^+ H^{1/2}(\mathbb{R}^n; \mathbb{C}^{2(n-1)}) \\ \dim \mathfrak{M} = k}} \sup_{\psi \in (\mathfrak{M} \oplus P_n^- H^{1/2}(\mathbb{R}^n; \mathbb{C}^{2(n-1)})) \setminus \{0\}} \frac{\mathfrak{d}_n[\psi] + \mathfrak{v}[\psi]}{\|\psi\|^2}.$$

Here  $P_n^+$  is the projector on the non-negative spectral subspace of  $D_n(0)$  and  $P_n^- := \mathbb{I} - P_n^+$ .

A direct application of Theorem 1 is:

**THEOREM 3** (Hardy-Dirac inequality). *Let  $v$  be a scalar function on  $\mathbb{R}^n$  such that  $v \otimes \mathbb{I}_{\mathbb{C}^{2(n-1)}} \in \mathfrak{P}_n$ . Moreover, we define the operator:*

$$K_n := \begin{cases} -i\partial_1 - \partial_2 & \text{if } n = 2, \\ -i\boldsymbol{\sigma} \cdot \nabla & \text{if } n = 3, \end{cases}$$

and denote by  $\lambda(v)$  the smallest eigenvalue of  $D_n(v \otimes \mathbb{I}_{\mathbb{C}^{2(n-1)}})$  in the gap  $(-1, 1)$ . Then for all  $\varphi \in H^1(\mathbb{R}^n; \mathbb{C}^{n-1})$  the inequality

$$0 \leq \int_{\mathbb{R}^n} \frac{|K_n \varphi(\mathbf{x})|^2}{1 + \lambda(v) - v(\mathbf{x})} dx + \int_{\mathbb{R}^n} (1 - \lambda(v) + v(\mathbf{x})) |\varphi(\mathbf{x})|^2 dx \tag{2}$$

holds.

We follow the tradition of [5] and call these type of inequality Hardy-Dirac inequality. In [6] it is demonstrated, how one can prove Hardy-Dirac inequalities for  $n = 3$  with the help of the Talman minimax principle.

We know that the lowest eigenvalue of  $D_n(-\nu/|\cdot| \otimes \mathbb{I}_{\mathbb{C}^{2(n-1)}})$  in  $(-1, 1)$  is  $\sqrt{1 - ((4 - n)\nu)^2}$  for  $\nu \in (0, 1/(4 - n))$  (see [7] and [19]). Thus Theorem 3 implies with a simple limiting argument

COROLLARY 4. *Let  $\nu \in [0, 1/(4 - n)]$ . Then*

$$0 \leq \int_{\mathbb{R}^n} \left( \frac{|K_n \varphi|^2}{1 + \sqrt{1 - ((4 - n)\nu)^2} + \frac{\nu}{|x|}} + \left( 1 - \sqrt{1 - ((4 - n)\nu)^2} - \frac{\nu}{|x|} \right) |\varphi|^2 \right) dx$$

*holds for all  $\varphi \in H^1(\mathbb{R}^n; \mathbb{C}^{n-1})$ .*

Let  $\nu \in [0, 1/(4 - n)]$ . With the help of Corollary 4 and Theorem 1 in [8] ( $\tilde{D}_n(-\nu/|\cdot| \otimes \mathbb{I}_{\mathbb{C}^{2(n-1)})}$  corresponds to  $H$  there) we know that there is only one self-adjoint extension of  $\tilde{D}_n(-\nu/|\cdot| \otimes \mathbb{I}_{\mathbb{C}^{2(n-1)})}$  with a positive Schur complement. We denote this distinguished self-adjoint extension by  $D_n^\nu$ . Now we want to give an explicit description of an operator core of  $D_n^\nu$ . For this purpose we introduce polar and spherical coordinates. We denote by the coordinate pair  $(\rho, \vartheta) \in [0, \infty) \times [0, 2\pi)$  the radial and angular polar coordinates in  $\mathbb{R}^2$  and by the coordinate triplet  $(r, \theta, \phi) \in [0, \infty) \times [0, \pi) \times [0, 2\pi)$  the radial, inclination and azimuthal spherical coordinates in  $\mathbb{R}^3$ . For  $m \in \{-1/2, 1/2\}^{n-1}$  we define the function  $\zeta_{n,m}^\nu$  in polar coordinates for  $n = 2$

$$\zeta_{2,m}^\nu(\rho, \vartheta) := \xi(\rho) \rho^{\sqrt{1/4 - \nu^2} - 1/2} \begin{pmatrix} \frac{\nu e^{-i(1/2+m)\vartheta}}{\sqrt{2\pi}} \\ -i(\sqrt{1/4 - \nu^2} + (-1)^{1/2-m}/2) \frac{e^{i(1/2-m)\vartheta}}{\sqrt{2\pi}} \end{pmatrix}; \tag{3}$$

and in spherical coordinates for  $n = 3$

$$\zeta_{3,m}^\nu(r, \theta, \phi) := \xi(r) r^{\sqrt{1 - \nu^2} - 1} \begin{pmatrix} \nu \Omega_{\frac{1}{2} + m_2, m_1, -m_2}(\theta, \phi) \\ -i(\sqrt{1 - \nu^2} + (-1)^{\frac{1}{2} - m_2}) \Omega_{\frac{1}{2} - m_2, m_1, m_2}(\theta, \phi) \end{pmatrix}; \tag{4}$$

with the spherical spinor  $\Omega_{l,m,s}$  (see Relation (7) in [10]) and the smooth cut-off function  $\xi$  (i.e.,  $\xi \in C^\infty(\mathbb{R}_+; \mathbb{R}_+)$ ,  $\xi(t) = 1$  for  $t \in (0, 1)$  and  $\xi(t) = 0$  for  $t > 2$ ). In the next theorem we give a characterisation of an operator core of  $D_n^\nu$  with the help of the functions  $\zeta_{n,m}^\nu$  introduced in (3) and (4).

THEOREM 5 (Operator core). *Let  $\nu \in [0, 1/(4 - n)]$ . The set*

$$\mathfrak{C}_n^\nu := C_0^\infty(\mathbb{R}^n \setminus \{0\}; \mathbb{C}^{2(n-1)}) \dot{+} \begin{cases} \{0\}, & \text{if } n = 2, \nu = 0 \text{ or } n = 3, \nu \in [0, \frac{\sqrt{3}}{2}]; \\ \text{span}\{\zeta_{n,m}^\nu : m \in \{-1/2, 1/2\}^{n-1}\}, & \text{else;} \end{cases} \tag{5}$$

*is an operator core for  $D_n^\nu$ .*

The knowledge of the operator core of  $D_n^\nu$  is important for the proof of estimates on the square of the operator, see e.g. [14]. In Remark 15 we show that for  $\nu \in [0, 1/(4 - n))$  the set  $\mathfrak{C}_n^\nu$  is an operator core for  $D_n(-\nu/|\cdot| \otimes \mathbb{I}_{\mathbb{C}^{2(n-1)})}$ . A direct consequence is:

COROLLARY 6. *Let  $\nu \in [0, 1/(4-n))$ . The distinguished self-adjoint extensions of  $\tilde{D}_n(-\nu/|\cdot| \otimes \mathbb{I}_{\mathbb{C}^{2(n-1)}})$  in the sense of [15] and [8] coincide, i.e.,*

$$D_n^\nu = D_n(-\nu/|\cdot| \otimes \mathbb{I}_{\mathbb{C}^{2(n-1)}}).$$

The proofs of the minimax characterisations rely on the angular momentum channel decomposition of the Coulomb-Dirac operator in the momentum space. This representation and the corresponding unitary transformations are introduced in the next section. In the remaining sections we prove in the order of enumeration: Theorems 1, 2, 3 and 5.

2 ANGULAR MOMENTUM CHANNEL DECOMPOSITION IN THE MOMENTUM SPACE

The Fourier transform connects the quantum mechanical descriptions of a particle in the configuration and momentum space. We use the standard unitary Fourier transform  $\mathcal{F}_n$  in  $L^2(\mathbb{R}^n)$  given for  $\varphi \in L^1(\mathbb{R}^n) \cap L^2(\mathbb{R}^n)$  by

$$\mathcal{F}_n \varphi := \frac{1}{(2\pi)^{n/2}} \int_{\mathbb{R}^n} e^{-i\langle \cdot, \mathbf{x} \rangle} \varphi(\mathbf{x}) d\mathbf{x}. \tag{6}$$

For the angular momentum channel decomposition in  $n$  dimensions we use an orthonormal basis in  $L^2(\mathbb{S}^{n-1}; \mathbb{C}^{n-1})$ . For  $n = 2$  this orthonormal basis is  $((2\pi)^{-1/2} e^{im(\cdot)})_{m \in \mathbb{Z}}$ . In three dimensions we use spherical spinors  $\Omega_{l,m,s}$ , which are defined in Relation (7) in [10], with  $l \in \mathbb{N}_0$ ,  $m \in \{-l - 1/2, \dots, l + 1/2\}$  and  $s \in \{-1/2, 1/2\}$ . The corresponding index sets are denoted by

$$\mathfrak{T}_2 := \mathbb{Z}; \tag{7}$$

and

$$\mathfrak{T}_3 := \left\{ (l, m, s) : l \in \mathbb{N}_0, m \in \left\{ -l - \frac{1}{2}, \dots, l + \frac{1}{2} \right\}, s = \pm \frac{1}{2}, \Omega_{l,m,s} \neq 0 \right\}. \tag{8}$$

Furthermore, we define subsets  $\mathfrak{T}_n^\pm$  of  $\mathfrak{T}_n$ :

$$\mathfrak{T}_n^a := \begin{cases} 2\mathbb{Z} & \text{if } n = 2, a = +; \\ 2\mathbb{Z} + 1 & \text{if } n = 2, a = -; \\ \{(l, m, s) \in \mathfrak{T}_3 : s = \pm 1/2\} & \text{if } n = 3, a = \pm. \end{cases} \tag{9}$$

Note that if  $(l, m, -1/2) \in \mathfrak{T}_3^-$  then  $l \in \mathbb{N}$ . Moreover, we introduce bijective maps

$$T_2 : \mathfrak{T}_2 \rightarrow \mathfrak{T}_2, T_2 k := k + 1; \tag{10}$$

and

$$T_3 : \mathfrak{T}_3 \rightarrow \mathfrak{T}_3, T_3(l, m, s) := (l + 2s, m, -s). \tag{11}$$

We can represent any  $\varphi \in L^2(\mathbb{R}^2; \mathbb{C})$  in polar coordinates and  $\zeta \in L^2(\mathbb{R}^3; \mathbb{C}^2)$  in spherical coordinates as

$$\varphi(\rho, \vartheta) = \sum_{k \in \mathfrak{T}_2} (2\pi\rho)^{-1/2} \varphi_k(\rho) e^{ik\vartheta}; \tag{12}$$

$$\zeta(r, \theta, \phi) = \sum_{(l,m,s) \in \mathfrak{T}_3} r^{-1} \zeta_{(l,m,s)}(r) \Omega_{l,m,s}(\theta, \phi); \tag{13}$$

with

$$\varphi_k(\rho) := \sqrt{\frac{\rho}{2\pi}} \int_0^{2\pi} \varphi(\rho, \vartheta) e^{-ik\vartheta} d\vartheta; \tag{14}$$

$$\zeta_{(l,m,s)}(r) := r \int_0^{2\pi} \int_0^\pi \langle \Omega_{l,m,s}(\theta, \phi), \zeta(r, \theta, \phi) \rangle_{\mathbb{C}^2} \sin(\theta) d\theta d\phi. \tag{15}$$

With the help of (14) and (15) we define the unitary operator

$$\mathcal{U}_n : L^2(\mathbb{R}^n; \mathbb{C}^{n-1}) \rightarrow \bigoplus_{j \in \mathfrak{T}_n} L^2(\mathbb{R}_+); \quad \psi \mapsto \bigoplus_{j \in \mathfrak{T}_n} \psi_j. \tag{16}$$

For the proof of the following lemma see Theorem 2.2.5 in [1] (based on Lemmata 2.1, 2.2 of [2]) for  $n = 2$  and Section 2.2 in [1] for  $n = 3$ .

LEMMA 7. For  $j \in (\mathbb{N}_0/2 - 1/2)$  and  $z \in (1, \infty)$  let

$$Q_j(z) = 2^{-j-1} \int_{-1}^1 (1 - t^2)^j (z - t)^{-j-1} dt \tag{17}$$

be a Legendre function of the second kind (see Section 15.3 in [21]). Let the sesquilinear form  $\mathbf{q}_j$  be defined on  $L^2(\mathbb{R}_+, (1+p^2)^{1/2} dp) \times L^2(\mathbb{R}_+, (1+p^2)^{1/2} dp)$  by

$$\mathbf{q}_j[f, g] := \pi^{-1} \int_0^\infty \int_0^\infty \overline{f(p)} Q_j\left(\frac{1}{2}\left(\frac{q}{p} + \frac{p}{q}\right)\right) g(q) dq dp. \tag{18}$$

For the special case  $f = g$  we introduce  $\mathbf{q}_j[f] := \mathbf{q}_j[f, f]$ .

Then for every  $\zeta, \eta \in H^{1/2}(\mathbb{R}^n)$  the relation

$$\int_{\mathbb{R}^n} \frac{\overline{\zeta(\mathbf{x})} \cdot \eta(\mathbf{x})}{|\mathbf{x}|} d\mathbf{x} = \begin{cases} \sum_{k \in \mathfrak{T}_2} \mathbf{q}_{|k|-1/2} [(\mathcal{F}_2\zeta)_k, (\mathcal{F}_2\eta)_k] & \text{if } n = 2, \\ \sum_{(l,m,s) \in \mathfrak{T}_3} \mathbf{q}_l [(\mathcal{F}_3\zeta)_{(l,m,s)}, (\mathcal{F}_3\eta)_{(l,m,s)}] & \text{if } n = 3, \end{cases} \tag{19}$$

holds.

The operators  $-i\boldsymbol{\sigma} \cdot \nabla$  and  $-i\boldsymbol{\alpha} \cdot \nabla$  are partially diagonalised in the momentum space by the unitary transforms

$$\mathcal{W}_2 : L^2(\mathbb{R}^2; \mathbb{C}^2) \rightarrow \bigoplus_{k \in \mathfrak{I}_2} L^2(\mathbb{R}_+; \mathbb{C}^2); \quad \begin{pmatrix} \varphi \\ \psi \end{pmatrix} \mapsto \bigoplus_{k \in \mathfrak{I}_2} \begin{pmatrix} \varphi_k \\ \psi_{T_2 k} \end{pmatrix} \quad (20)$$

and

$$\mathcal{W}_3 : L^2(\mathbb{R}^3; \mathbb{C}^4) \rightarrow \bigoplus_{(l,m,s) \in \mathfrak{I}_3} L^2(\mathbb{R}_+; \mathbb{C}^2); \quad \begin{pmatrix} \psi_1 \\ \psi_2 \\ \psi_3 \\ \psi_4 \end{pmatrix} \mapsto \bigoplus_{(l,m,s) \in \mathfrak{I}_3} \begin{pmatrix} \psi_{(l,m,s)}^+ \\ \psi_{T_3(l,m,s)}^- \end{pmatrix} \quad (21)$$

with

$$\psi_{(l,m,s)}^+ := \begin{pmatrix} \psi_1 \\ \psi_2 \end{pmatrix}_{(l,m,s)} \quad \text{and} \quad \psi_{(l,m,s)}^- := \begin{pmatrix} \psi_3 \\ \psi_4 \end{pmatrix}_{(l,m,s)} \quad (22)$$

for  $(l, m, s) \in \mathfrak{I}_3$ . To be more precise:

LEMMA 8. *For the self-adjoint operators  $-i\boldsymbol{\sigma} \cdot \nabla$  and  $-i\boldsymbol{\alpha} \cdot \nabla$  the relations*

$$(\mathcal{W}_n \mathcal{F}_n)^* \left( \bigoplus_{j \in \mathfrak{I}_n} \begin{pmatrix} 0 & (\cdot) \\ (\cdot) & 0 \end{pmatrix} \right) (\mathcal{W}_n \mathcal{F}_n) = \begin{cases} -i\boldsymbol{\sigma} \cdot \nabla & \text{if } n = 2, \\ -i\boldsymbol{\alpha} \cdot \nabla & \text{if } n = 3, \end{cases} \quad (23)$$

hold.

*Proof.* By a straightforward calculation and Relation 2.1.28 in [1] the relations

$$\boldsymbol{\sigma} \cdot \mathbf{x} = \begin{pmatrix} 0 & e^{-i\vartheta} \rho \\ e^{i\vartheta} \rho & 0 \end{pmatrix} \quad \text{for } \mathbf{x} \in \mathbb{R}^2; \quad (24)$$

$$\boldsymbol{\sigma} \cdot \frac{\mathbf{x}}{|\mathbf{x}|} \Omega_{l,m,s} = \Omega_{l+2s,m,-s} \quad \text{for } \mathbf{x} \in \mathbb{R}^3 \text{ and } (l, m, s) \in \mathfrak{I}_3; \quad (25)$$

hold.

The set  $C_0^\infty(\mathbb{R}^n; \mathbb{C}^{2(n-1)})$  is dense in  $H^1(\mathbb{R}^n; \mathbb{C}^{2(n-1)})$ . Thus it is enough to work with  $\psi \in C_0^\infty(\mathbb{R}^2; \mathbb{C}^2)$  and  $\zeta \in C_0^\infty(\mathbb{R}^3; \mathbb{C}^4)$ .

Moreover, the Fourier transform diagonalises differential operators:

$$\langle \psi, -i\boldsymbol{\sigma} \cdot \nabla \psi \rangle = \langle \mathcal{F}_2 \psi, \boldsymbol{\sigma} \cdot \mathbf{p} \mathcal{F}_2 \psi \rangle, \quad (26)$$

$$\langle \zeta, -i\boldsymbol{\alpha} \cdot \nabla \zeta \rangle = \langle \mathcal{F}_3 \zeta, \boldsymbol{\alpha} \cdot \mathbf{p} \mathcal{F}_3 \zeta \rangle. \quad (27)$$

Here we denote by  $\mathbf{p}$  the independent variable of multiplication in  $L^2(\mathbb{R}^n; d\mathbf{p})$ . Now we prove (23) for  $n = 3$ . We obtain by the representation (13) of the upper and lower bispinor of  $\mathcal{F}_3 \zeta$  and the notation introduced in (22) that the right hand side of (27) is equal to

$$2 \sum_{\substack{(l',m',s') \in \mathfrak{I}_3 \\ (l,m,s) \in \mathfrak{I}_3}} \Re \left( \langle |\mathbf{p}|^{-1} (\mathcal{F}_3 \zeta)_{(l',m',s')}^+ \Omega_{l',m',s'}, (\boldsymbol{\sigma} \cdot \mathbf{p}) |\mathbf{p}|^{-1} (\mathcal{F}_3 \zeta)_{(l,m,s)}^- \Omega_{l,m,s} \rangle \right). \quad (28)$$

The expression in (28) is equal to

$$\begin{aligned}
 & 2 \sum_{(l,m,s) \in \mathfrak{I}_3} \Re \left( \langle (\mathcal{F}_3\zeta)_{(l+2s,m,-s)}^+, (\cdot)(\mathcal{F}_3\zeta)_{(l,m,s)}^- \rangle \right) \\
 &= \sum_{(l,m,s) \in \mathfrak{I}_3} \left\langle \left( \begin{array}{c} (\mathcal{F}_3\zeta)_{(l,m,s)}^+ \\ (\mathcal{F}_3\zeta)_{T_3(l,m,s)}^- \end{array} \right), \left( \begin{array}{cc} 0 & (\cdot) \\ (\cdot) & 0 \end{array} \right) \left( \begin{array}{c} (\mathcal{F}_3\zeta)_{(l,m,s)}^+ \\ (\mathcal{F}_3\zeta)_{T_3(l,m,s)}^- \end{array} \right) \right\rangle \\
 &= \left\langle \mathcal{W}_3 \mathcal{F}_3 \zeta, \left( \bigoplus_{(l,m,s) \in \mathfrak{I}_3} \left( \begin{array}{cc} 0 & (\cdot) \\ (\cdot) & 0 \end{array} \right) \right) \mathcal{W}_3 \mathcal{F}_3 \zeta \right\rangle \tag{29}
 \end{aligned}$$

by the sequential application of (25), (21) and (6). Thus the claim of Lemma 8 is a consequence of (27), (28) and (29).

For  $n = 2$  we obtain (23) by an analogous procedure, i.e., we represent the upper and lower component of  $\mathcal{F}_2\psi$  by (12) in (26) and perform a calculation, which involves (24). □

### 3 PROOF OF THEOREM 1

Let  $V \in \mathfrak{P}_n$ . We use the abstract minimax principle Theorem 1 of [13] to prove the Talman minimax principle. We apply the theorem with  $q := \mathfrak{d}_n$  (quadratic form associated to  $D_n(0)$ ),  $B := D_n(V)$  and  $\Lambda_{\pm}$  as the projector  $T_n^{\pm}$  on the upper and lower  $(n - 1)$  components of a  $2(n - 1)$  spinor, i.e.,

$$T_n^+ \begin{pmatrix} \varphi \\ \psi \end{pmatrix} = \begin{pmatrix} \varphi \\ 0 \end{pmatrix}, \quad T_n^- \begin{pmatrix} \varphi \\ \psi \end{pmatrix} = \begin{pmatrix} 0 \\ \psi \end{pmatrix}, \text{ for } \varphi, \psi \in \mathbb{L}^2(\mathbb{R}^n; \mathbb{C}^{n-1}).$$

That  $D_n(V)$  plays the role of  $B$  in [13] is a consequence of Theorem 2.1 in [15] and the following lemma.

LEMMA 9. *Let  $V \in \mathfrak{P}_n$ . Then the quadratic form  $\mathfrak{v}$  associated to the operator  $V$  is a form perturbation of  $D_n(0)$  in the sense of Definition 2.1 in [15].*

*Proof.*  $V$  is  $D_n(0)$  form bounded by the Herbst inequality (see Theorem 2.5 in [11]). Moreover, the inequality

$$\|r^{-1/2} D_n(0)^{-1} r^{-1/2}\| \leq 4 - n$$

holds. This is proven in Section 2 in [12] for  $n = 3$ . The same arguments also apply for  $n = 2$  (see Step 1 in the proof of Theorem 1 in [4]). Thus

$$\|V^{1/2} D_n(0)^{-1} V^{1/2}\| \leq \|V^{1/2} r^{1/2}\|^2 \cdot \|r^{-1/2} D_n(0)^{-1} r^{-1/2}\| < 1.$$

Hence  $1 + V^{1/2} D_n(0)^{-1} V^{1/2}$  has a bounded inverse by the Neumann series. Now the claim follows from Theorem 2.2 in [15] with  $A := D_n(0)$  and  $t := 0$ . □

Since the assumptions (i) and (ii) of Theorem 1 in [13] are obviously fulfilled, it remains to check assumption (iii). Thus it is enough to find an operator  $L_n : \mathbf{H}^{1/2}(\mathbb{R}^n; \mathbb{C}^{n-1}) \rightarrow \mathbf{H}^{1/2}(\mathbb{R}^n; \mathbb{C}^{n-1})$  such that

$$\inf_{\varphi \in \mathbf{H}^{1/2}(\mathbb{R}^n; \mathbb{C}^{n-1}) \setminus \{0\}} \frac{\mathbf{d}_n \left[ \begin{pmatrix} \varphi \\ L_n \varphi \end{pmatrix} \right] + \mathbf{v} \left[ \begin{pmatrix} \varphi \\ L_n \varphi \end{pmatrix} \right]}{\left\| \begin{pmatrix} \varphi \\ L_n \varphi \end{pmatrix} \right\|^2} > -1.$$

Now we give in three steps an explicit construction of  $L_n$  and show that  $L_n$  satisfies the requirements. For  $k \in \mathfrak{T}_2$  and  $(l, m, s) \in \mathfrak{T}_3$  we define in the first step various constants:

$$c_n := 2(4 - n) \frac{\Gamma(\frac{n+1}{4})^2}{\Gamma(\frac{n-1}{4})^2}; \tag{30}$$

$$c_{2,k} := \begin{cases} c_2^{-1} & \text{if } k \in \mathfrak{T}_2^-, \\ c_2 & \text{if } k \in \mathfrak{T}_2^+; \end{cases} \tag{31}$$

$$c_{3,(l,m,s)} := c_3^{2s}. \tag{32}$$

In the second step we define the operator  $R_n$

$$R_n : \bigoplus_{j \in \mathfrak{T}_n} \mathbf{L}^2(\mathbb{R}_+) \rightarrow \bigoplus_{j \in \mathfrak{T}_n} \mathbf{L}^2(\mathbb{R}_+); \quad \bigoplus_{j \in \mathfrak{T}_n} \psi_j \mapsto \bigoplus_{j \in \mathfrak{T}_n} c_{n,j} \psi_{T_n^{-1}j}. \tag{33}$$

Finally we define

$$L_n := (\mathcal{U}_n \mathcal{F}_n)^* R_n (\mathcal{U}_n \mathcal{F}_n). \tag{34}$$

The desired properties of  $L_n$  are proven in the following lemma:

LEMMA 10. *Let  $\varphi \in \mathbf{H}^{1/2}(\mathbb{R}^n; \mathbb{C}^{n-1})$  then  $L_n \varphi \in \mathbf{H}^{1/2}(\mathbb{R}^n; \mathbb{C}^{n-1})$  and the following inequality*

$$\frac{c_n^2 - 1}{c_n^2 + 1} \left\| \begin{pmatrix} \varphi \\ L_n \varphi \end{pmatrix} \right\|^2 \leq \mathbf{d}_n \left[ \begin{pmatrix} \varphi \\ L_n \varphi \end{pmatrix} \right] - \frac{1}{4-n} \int_{\mathbb{R}^n} \frac{1}{|\mathbf{x}|} \left| \begin{pmatrix} \varphi(\mathbf{x}) \\ (L_n \varphi)(\mathbf{x}) \end{pmatrix} \right|^2 d\mathbf{x} \tag{35}$$

holds.

*Proof.* We recall that

$$\mathbf{H}^{1/2}(\mathbb{R}^n) = \{ \psi \in \mathbf{L}^2(\mathbb{R}^n) : (1 + |\cdot|^2)^{1/4} \mathcal{F}_n \psi \in \mathbf{L}^2(\mathbb{R}^n) \}.$$

Thus the unitarity of  $\mathcal{U}_n$  implies

$$\mathbf{H}^{1/2}(\mathbb{R}^n) = \{ \psi \in \mathbf{L}^2(\mathbb{R}^n) : \bigoplus_{j \in \mathfrak{T}_n} (1 + (\cdot)^2)^{1/4} (\mathcal{F}_n \psi)_j \in \bigoplus_{j \in \mathfrak{T}_n} \mathbf{L}^2(\mathbb{R}_+) \}. \tag{36}$$

Moreover we observe that the operator  $R_n$  is bounded, which together with (36) and (34) implies that  $L_n\varphi \in \mathbf{H}^{1/2}(\mathbb{R}^n)$ .

Now we define the quadratic form  $\mathbf{p}$  on  $\mathbf{L}^2(\mathbb{R}_+, (1 + p^2)^{1/2} dp)$  by

$$\mathbf{p}[\chi] := \int_0^\infty p|\chi(p)|^2 dp.$$

For the proof of (35) we recall that the quadratic form (18) satisfy the inequalities

$$\begin{aligned} \mathbf{q}_{k+1/2}[\zeta] &\leq \mathbf{q}_{k-1/2}[\zeta]; \\ \mathbf{q}_{k+1}[\zeta] &\leq \mathbf{q}_k[\zeta]; \\ \mathbf{q}_0[\zeta] &\leq c_3^{-1}\mathbf{p}[\zeta], \quad \mathbf{q}_1[\zeta] \leq c_3\mathbf{p}[\zeta]; \\ \mathbf{q}_{-1/2}[\zeta] &\leq 2c_2^{-1}\mathbf{p}[\zeta], \quad \mathbf{q}_{1/2}[\zeta] \leq 2c_2\mathbf{p}[\zeta]; \end{aligned} \tag{37}$$

for  $k \in \mathbb{N}_0$  and  $\zeta \in \mathbf{L}^2(\mathbb{R}_+, (1 + p^2)^{1/2} dp)$  (see [2] and [10]).

By Lemma 7 we obtain

$$\int_{\mathbb{R}^n} \frac{|\varphi(\mathbf{x})|^2}{|\mathbf{x}|} d\mathbf{x} = \begin{cases} \sum_{k \in \mathfrak{X}_2} \mathbf{q}_{|k|-1/2}[(\mathcal{F}_2\varphi)_k] & \text{if } n = 2; \\ \sum_{(l,m,s) \in \mathfrak{X}_3} \mathbf{q}_l[(\mathcal{F}_3\varphi)_{(l,m,s)}] & \text{if } n = 3; \end{cases} \tag{38}$$

and by (31) - (34)

$$\begin{aligned} &\int_{\mathbb{R}^n} \frac{|(L_n\varphi)(\mathbf{x})|^2}{|\mathbf{x}|} d\mathbf{x} \\ &= \begin{cases} \sum_{k \in \mathfrak{X}_2^+} c_2^2 \mathbf{q}_{|k|-\frac{1}{2}}[(\mathcal{F}_2\varphi)_{k-1}] + \sum_{k \in \mathfrak{X}_2^-} c_2^{-2} \mathbf{q}_{|k|-\frac{1}{2}}[(\mathcal{F}_2\varphi)_{k-1}] & \text{if } n = 2; \\ \sum_{(l,m,\frac{1}{2}) \in \mathfrak{X}_3^+} c_3^2 \mathbf{q}_l[(\mathcal{F}_3\varphi)_{(l+1,m,-\frac{1}{2})}] + \sum_{(l,m,-\frac{1}{2}) \in \mathfrak{X}_3^-} c_3^{-2} \mathbf{q}_l[(\mathcal{F}_3\varphi)_{(l-1,m,\frac{1}{2})}] & \text{if } n = 3. \end{cases} \end{aligned} \tag{39}$$

Note that  $(l, m, s) \in \mathfrak{X}_3^-$  implies  $l \in \mathbb{N}$ . Hence (37) implies that the right hand sides of (38) can be estimated by

$$(4 - n) \left( \sum_{j \in \mathfrak{X}_n^+} c_n^{-1} \mathbf{p}[(\mathcal{F}_n\varphi)_j] + \sum_{j \in \mathfrak{X}_n^-} c_n \mathbf{p}[(\mathcal{F}_n\varphi)_j] \right); \tag{40}$$

and the right hand side of (39) by

$$(4 - n) \left( \sum_{j \in \mathfrak{X}_n^+} c_n \mathbf{p}[(\mathcal{F}_n\varphi)_{T_n^{-1}j}] + \sum_{j \in \mathfrak{X}_n^-} c_n^{-1} \mathbf{p}[(\mathcal{F}_n\varphi)_{T_n^{-1}j}] \right). \tag{41}$$

By  $T_n(\mathfrak{I}_n^\pm) = \mathfrak{I}_n^\mp$  we conclude that (41) is equal to (40). This together with the relation

$$(\mathcal{F}_n L_n \varphi)_{T_n j} = c_{n, T_n j} (\mathcal{F}_n \varphi)_j \text{ for all } j \in \mathfrak{I}_n,$$

implies

$$\begin{aligned} & \frac{1}{4-n} \int_{\mathbb{R}^n} \frac{1}{|\mathbf{x}|} \left| \begin{pmatrix} \varphi(\mathbf{x}) \\ (L_n \varphi)(\mathbf{x}) \end{pmatrix} \right|^2 d\mathbf{x} \leq \\ & \sum_{j \in \mathfrak{I}_n} \int_{\mathbb{R}_+} \left\langle \begin{pmatrix} (\mathcal{F}_n \varphi)_j(p) \\ (\mathcal{F}_n L_n \varphi)_{T_n j}(p) \end{pmatrix}, \begin{pmatrix} 0 & p \\ p & 0 \end{pmatrix} \begin{pmatrix} (\mathcal{F}_n \varphi)_j(p) \\ (\mathcal{F}_n L_n \varphi)_{T_n j}(p) \end{pmatrix} \right\rangle_{\mathbb{C}^2} dp. \end{aligned} \tag{42}$$

A straightforward calculation using (31) - (34) gives

$$\begin{aligned} & \left\langle \begin{pmatrix} \varphi \\ L_n \varphi \end{pmatrix}, \begin{pmatrix} \mathbb{I}_{\mathbb{C}^{n-1}} & 0 \\ 0 & \mp \mathbb{I}_{\mathbb{C}^{n-1}} \end{pmatrix} \begin{pmatrix} \varphi \\ L_n \varphi \end{pmatrix} \right\rangle \\ & = (1 \mp c_n^{-2}) \sum_{j \in \mathfrak{I}_n^+} \|(\mathcal{F}_n \varphi)_j\|^2 + (1 \mp c_n^2) \sum_{j \in \mathfrak{I}_n^-} \|(\mathcal{F}_n \varphi)_j\|^2. \end{aligned} \tag{43}$$

By Lemma 8 we know that the right hand side of Relation (42) plus the minus case of the left hand side of (43) is equal to  $d_n[(L_n^\varphi)]$ . Thus we obtain (35) by (42) and (43).  $\square$

#### 4 PROOF OF THEOREM 2

We proceed analogously to the proof of Theorem 1. Thus it is enough to find an operator  $G_n : P_n^+ H^{1/2}(\mathbb{R}^n; \mathbb{C}^{2(n-1)}) \rightarrow P_n^- H^{1/2}(\mathbb{R}^n; \mathbb{C}^{2(n-1)})$  such that

$$\inf_{\varphi \in P_n^+ H^{1/2}(\mathbb{R}^n; \mathbb{C}^{2(n-1)}) \setminus \{0\}} \frac{d_n[\varphi + G_n \varphi] + v[\varphi + G_n \varphi]}{\|\varphi + G_n \varphi\|^2} > -1 \tag{44}$$

holds. In the following lemma we prove that a possible choice of  $G_n$  is

$$G_n := (\mathcal{W}_n \mathcal{F}_n)^* E_n (\mathcal{W}_n \mathcal{F}_n), \tag{45}$$

with

$$E_n : \bigoplus_{j \in \mathfrak{I}_n} L^2(\mathbb{R}_+; \mathbb{C}^2) \rightarrow \bigoplus_{j \in \mathfrak{I}_n} L^2(\mathbb{R}_+; \mathbb{C}^2); \tag{46}$$

$$\bigoplus_{j \in \mathfrak{I}_n} \Psi_j \mapsto \bigoplus_{j \in \mathfrak{I}_n} \frac{1 - c_{n,j}(\cdot) + \sqrt{1 + (\cdot)^2}}{c_{n,j} + (\cdot) + c_{n,j} \sqrt{1 + (\cdot)^2}} \begin{pmatrix} 0 & -1 \\ 1 & 0 \end{pmatrix} \Psi_j. \tag{47}$$

LEMMA 11. Let  $\varphi \in P_n^+ H^{1/2}(\mathbb{R}^n; \mathbb{C}^{2(n-1)})$  then  $G_n \varphi \in P_n^- H^{1/2}(\mathbb{R}^n; \mathbb{C}^{2(n-1)})$  and the relation

$$L_n(\varphi + G_n \varphi)_1 = (\varphi + G_n \varphi)_2 \tag{48}$$

holds.

REMARK 12. *By Lemma 10 and Relation (48) we conclude (44).*

*Proof of Lemma 11.* By Lemma 8 we deduce that  $\psi \in P_n^\pm H^{1/2}(\mathbb{R}^n; \mathbb{C}^{2(n-1)})$  if and only if there exists  $\bigoplus_{j \in \mathfrak{I}_n} \zeta_j \in \bigoplus_{j \in \mathfrak{I}_n} L^2(\mathbb{R}_+; (1+p^2)^{1/2} dp)$  such that

$$(\mathcal{W}_n \mathcal{F}_n \psi)_j(p) = \begin{cases} \zeta_j(p) \left( \frac{1}{1+\sqrt{1+p^2}} \right) & (\text{"+" case}); \\ \zeta_j(p) \left( \frac{-p}{1+\sqrt{1+p^2}} \right) & (\text{"-" case}); \end{cases} \tag{49}$$

holds for every  $j \in \mathfrak{I}_n$  and  $p \in \mathbb{R}_+$ . Hence we get  $G_n \varphi \in P_n^- H^{1/2}(\mathbb{R}^n; \mathbb{C}^{2(n-1)})$ . By (49),(46) we obtain that there exists  $\bigoplus_{j \in \mathfrak{I}_n} \chi_j \in \bigoplus_{j \in \mathfrak{I}_n} L^2(\mathbb{R}_+; (1+p^2)^{1/2} dp)$  such that

$$(\mathcal{W}_n \mathcal{F}_n \varphi)_j(p) = \chi_j(p) \left( \frac{1}{1+\sqrt{1+p^2}} \right)$$

and

$$\begin{aligned} ((\mathbb{I} + E_n) \mathcal{W}_n \mathcal{F}_n \varphi)_j &= \begin{pmatrix} \tilde{\chi}_j \\ c_{n,T_n j} \tilde{\chi}_j \end{pmatrix} \text{ with} \\ \tilde{\chi}_j(p) &:= \frac{c_{n,j} (p^2 + (1 + \sqrt{1+p^2})^2)}{(1 + \sqrt{1+p^2})(c_{n,j} + p + c_{n,j} \sqrt{1+p^2})} \chi_j(p) \text{ for } p \in \mathbb{R}_+, \end{aligned} \tag{50}$$

hold for every  $j \in \mathfrak{I}_n$ . Hence we get by (45),(33) and (34) the relation

$$\varphi + G_n \varphi = (\mathcal{W}_n \mathcal{F}_n)^* \bigoplus_{j \in \mathfrak{I}_n} \begin{pmatrix} \tilde{\chi}_j \\ c_{n,T_n j} \tilde{\chi}_j \end{pmatrix} = \begin{pmatrix} (\mathcal{U}_n \mathcal{F}_n)^* \bigoplus_{j \in \mathfrak{I}_n} \tilde{\chi}_j \\ L_n (\mathcal{U}_n \mathcal{F}_n)^* \bigoplus_{j \in \mathfrak{I}_n} \tilde{\chi}_j \end{pmatrix}.$$

Thus we have proven Relation (48). □

### 5 PROOF OF THEOREM 3

Since the right hand side of (2) is continuous in the  $H^1(\mathbb{R}^n; \mathbb{C}^{n-1})$  norm (see Theorem 2.5 in [11]), we can assume that  $\varphi \in C_0^\infty(\mathbb{R}^n \setminus \{0\}; \mathbb{C}^{n-1}) \setminus \{0\}$  by the density of  $C_0^\infty(\mathbb{R}^n \setminus \{0\}; \mathbb{C}^{n-1})$  in  $H^1(\mathbb{R}^n; \mathbb{C}^{n-1})$ .

By the application of Theorem 1 we obtain

$$\lambda(v) \leq \sup_{\psi \in H^1(\mathbb{R}^n \setminus \{0\}; \mathbb{C}^{n-1})} I_{n,v,\varphi}(\psi) \text{ with} \tag{51}$$

$$I_{n,v,\varphi} : H^1(\mathbb{R}^n \setminus \{0\}; \mathbb{C}^{n-1}) \rightarrow \mathbb{R}; \tag{52}$$

$$I_{n,v,\varphi}(\psi) := \frac{\left\langle \begin{pmatrix} \varphi \\ \psi \end{pmatrix}, \begin{pmatrix} (1+v) \otimes \mathbb{I}_{\mathbb{C}^{n-1}} & K_n \\ K_n & (-1+v) \otimes \mathbb{I}_{\mathbb{C}^{n-1}} \end{pmatrix} \begin{pmatrix} \varphi \\ \psi \end{pmatrix} \right\rangle}{\left\| \begin{pmatrix} \varphi \\ \psi \end{pmatrix} \right\|^2}. \tag{53}$$

Note that we calculate the suprema in (51) over  $H^1(\mathbb{R}^n \setminus \{0\}; \mathbb{C}^{n-1})$  instead of  $H^{1/2}(\mathbb{R}^n; \mathbb{C}^{n-1})$ . This is justified by a density argument, which makes use of the form boundedness of  $v \otimes \mathbb{I}_{\mathbb{C}^{2(n-1)}}$  with respect to  $D_n(0)$  (see Lemma 9) and the density of  $H^1(\mathbb{R}^n \setminus \{0\}; \mathbb{C}^{n-1})$  in  $H^{1/2}(\mathbb{R}^n; \mathbb{C}^{n-1})$ .

Thus the proof of Theorem 3 basically follows from the following lemma.

LEMMA 13. *We define*

$$J_{n,v,\varphi} : (-1, \infty) \rightarrow \mathbb{R};$$

$$J_{n,v,\varphi}(\lambda) := \int_{\mathbb{R}^n} \left( \frac{|K_n \varphi(\mathbf{x})|^2}{1 + \lambda - v(\mathbf{x})} + (1 - \lambda + v(\mathbf{x})) |\varphi(\mathbf{x})|^2 \right) dx.$$

For  $\lambda \in (-1, \infty)$ ,  $J_{n,v,\varphi}(\lambda) \leq 0$  implies

$$\sup_{\psi \in H^1(\mathbb{R}^n \setminus \{0\}; \mathbb{C}^{n-1})} I_{n,v,\varphi}(\psi) \leq \lambda.$$

*Proof.* We introduce

$$\psi_{n,v,\varphi} : (-1, \infty) \rightarrow H^1(\mathbb{R}^n \setminus \{0\}; \mathbb{C}^{n-1}); \quad \psi_{n,v,\varphi}(\lambda) := \frac{K_n \varphi}{1 + \lambda - v}. \tag{54}$$

For every  $\zeta \in H^1(\mathbb{R}^n \setminus \{0\}; \mathbb{C}^{n-1})$  the inequality

$$\begin{aligned} & (I_{n,v,\varphi}(\psi_{n,v,\varphi}(\lambda) + \zeta) - \lambda) (\|\varphi\|^2 + \|\psi_{n,v,\varphi}(\lambda) + \zeta\|^2) \\ &= J_{n,v,\varphi}(\lambda) + 2\Re \langle \zeta, K_n \varphi - (1 + \lambda - v) \psi_{n,v,\varphi}(\lambda) \rangle + \\ & \langle K_n \varphi - (1 + \lambda - v) \psi_{n,v,\varphi}(\lambda), \psi_{n,v,\varphi}(\lambda) \rangle - \langle \zeta, (1 + \lambda - v) \zeta \rangle \leq J_{n,v,\varphi}(\lambda) \end{aligned}$$

holds, and thus we conclude the claim. □

By Lemma 13 and (51) we obtain

$$J_{n,v,\varphi}(\lambda(v) - \varepsilon) > 0 \text{ for } \varepsilon \in (0, 1 + \lambda(v)). \tag{55}$$

Letting  $\varepsilon \searrow 0$  in (55) we obtain Theorem 3.

6 PROOF OF THEOREM 5

The proof is based on:

LEMMA 14. *Let  $\nu \in [0, 1/(4 - n)]$ . The restriction of  $(\tilde{D}_n(-\nu/|\cdot| \otimes \mathbb{I}_{\mathbb{C}^{2(n-1)}}))^*$  to  $\mathfrak{C}_n^\nu$  is essentially self-adjoint.*

*Proof.* For  $m \in \mathfrak{T}_2$  and  $(l, m, s) \in \mathfrak{T}_3$  we define

$$\begin{aligned} \kappa_m &:= m + 1/2; \\ \kappa_{(l,m,s)} &:= 2sl + s + 1/2. \end{aligned}$$

Furthermore we introduce for every  $j \in \mathfrak{T}_n$  the operator  $D^{j,\nu}$  in  $L^2(\mathbb{R}_+; \mathbb{C}^2)$  by the differential expression

$$d^{j,\nu} := \begin{pmatrix} -\frac{\nu}{r} & -\frac{d}{dr} - \frac{\kappa_j}{r} \\ \frac{d}{dr} - \frac{\kappa_j}{r} & -\frac{\nu}{r} \end{pmatrix}$$

on  $C_0^\infty(\mathbb{R}_+; \mathbb{C}^2)$ . Now we observe that any solution of the equation  $d^{j,\nu}\varphi = 0$  in  $\mathbb{R}_+$  is a linear combination of the two functions

$$\varphi_{j,1}^\nu(r) := \begin{cases} \begin{pmatrix} 1 \\ 0 \end{pmatrix} r^{\kappa_j} & \text{if } \nu = 0, \\ \begin{pmatrix} \nu \\ \sqrt{\kappa_j^2 - \nu^2} - \kappa_j \end{pmatrix} r^{\sqrt{\kappa_j^2 - \nu^2}} & \text{else,} \end{cases}$$

and

$$\varphi_{j,2}^\nu(r) := \begin{cases} \begin{pmatrix} 0 \\ 1 \end{pmatrix} r^{-\kappa_j} & \text{if } \nu = 0, \\ \begin{pmatrix} \nu \\ -\sqrt{\kappa_j^2 - \nu^2} - \kappa_j \end{pmatrix} r^{-\sqrt{\kappa_j^2 - \nu^2}} & \text{if } 0 < \nu^2 < \kappa_j^2, \\ \begin{pmatrix} \nu \ln(r) \\ 1 - \kappa_j \ln(r) \end{pmatrix} & \text{if } \nu^2 = \kappa_j^2. \end{cases}$$

Through the application of the results of [20] as in Section 2 in [14] we obtain that the closure  $D_{\text{ex}}^{j,\nu}$  of the restriction of  $(D^{j,\nu})^*$  to  $\mathfrak{C}^{j,\nu}$  is self-adjoint with

$$\mathfrak{C}^{j,\nu} := \begin{cases} C_0^\infty(\mathbb{R}_+; \mathbb{C}^2) \dot{+} \text{span}\{\xi\varphi_{j,1}^\nu\} & \text{if } \kappa_j^2 - \nu^2 < 1/4; \\ C_0^\infty(\mathbb{R}_+; \mathbb{C}^2) & \text{else.} \end{cases}$$

Here  $\xi$  is a smooth cut-off function with  $\xi \in C^\infty(\mathbb{R}_+; \mathbb{R}_+)$ ,  $\xi(t) = 1$  for  $t \in (0, 1)$

and  $\xi(t) = 0$  for  $t > 2$ . Thus we conclude the claim by

$$(\tilde{D}_n(-\nu/|\cdot| \otimes \mathbb{I}_{\mathbb{C}^{2(n-1)}}))^* = (\mathcal{W}_n \mathcal{M}_n)^* \left( \bigoplus_{j \in \mathfrak{I}_n} (D^{j,\nu} + \sigma_3)^* \right) \mathcal{W}_n \mathcal{M}_n \text{ with,} \tag{56}$$

$$\mathcal{M}_n := \text{diag}(1, i) \otimes \mathbb{I}_{\mathbb{C}^{n-1}}$$

(see Section 7.3.3 in [19] for  $n = 2$  and Section 2.1 in [1] for  $n = 3$ ) and the fact that  $\sigma_3$  is a bounded operator in  $L^2(\mathbb{R}_+; \mathbb{C}^2)$ .  $\square$

REMARK 15. Let  $\nu \in [0, 1/(4 - n)]$  and  $j \in \mathfrak{I}_n$ . By the embedding

$$H^{1/2}(\mathbb{R}^n) \subset L^2(\mathbb{R}^n, (1 + |\mathbf{x}|^{-1})d\mathbf{x})$$

and (56) we obtain that the domain of  $(\mathcal{W}_n M_n D_n(-\nu/|\cdot| \otimes \mathbb{I}_{\mathbb{C}^{2(n-1)}}) (\mathcal{W}_n M_n)^*)_j$  is in  $L^2(\mathbb{R}_+, (1 + r^{-1})dr)$ . Hence there is a self-adjoint extension of  $D^{j,\nu}$  with domain in  $L^2(\mathbb{R}_+, (1 + r^{-1})dr)$ . By  $\xi\varphi_{j,2}^\nu \notin L^2(\mathbb{R}_+, (1 + r^{-1})dr)$  for  $\nu > 0$  and Theorem 1.5 in [20] we get that  $D_{\text{ex}}^{j,\nu}$  is the unique self-adjoint extension of  $D^{j,\nu}$  with domain in  $L^2(\mathbb{R}_+, (1 + r^{-1})dr)$ . Therefore, we obtain

$$(\mathcal{W}_n M_n D_n(-\nu/|\cdot| \otimes \mathbb{I}_{\mathbb{C}^{2(n-1)}}) (\mathcal{W}_n M_n)^*)_j = D_{\text{ex}}^{j,\nu}.$$

We conclude that the closure of  $(\tilde{D}_n(-\nu/|\cdot| \otimes \mathbb{I}_{\mathbb{C}^{2(n-1)}}))^*$  restricted to  $\mathfrak{C}_n^\nu$  is  $D_n(-\nu/|\cdot| \otimes \mathbb{I}_{\mathbb{C}^{2(n-1)}})$ .

As a consequence of Lemma 14 it remains to prove that  $\zeta_{n,m}^\nu \in \mathfrak{D}(D_n^\nu)$  for  $m \in \{-1/2, 1/2\}^{n-1}$  and  $(n, \nu) \in (\{2\} \times (0, 1/2]) \cup (\{3\} \times (\sqrt{3}/2, 1])$ . We introduce the symmetric and non-negative (by Corollary 4) quadratic form  $q_n^\nu$  on  $C_0^\infty(\mathbb{R}^n \setminus \{0\}; \mathbb{C}^{n-1})$  by

$$q_n^\nu[\varphi] := \int_{\mathbb{R}^n} \left( \frac{|K_n \varphi|^2}{1 + \sqrt{1 - ((4 - n)\nu)^2} + \frac{\nu}{|\mathbf{x}|}} + \left( 1 - \sqrt{1 - ((4 - n)\nu)^2} - \frac{\nu}{|\mathbf{x}|} \right) |\varphi|^2 \right) d\mathbf{x}.$$

Note that  $q_n^\nu$  is closable by Theorem X.23 in [16]. We denote the domain of the closure of  $q_n^\nu$  by  $\Omega_n^\nu$ .

By the characterisation of  $\mathfrak{D}(D_n^\nu)$  in Theorem 1 in [8], it is enough to show that for all  $m \in \{-1/2, 1/2\}^{n-1}$  the upper  $(n - 1)$  spinor of  $\zeta_{n,m}^\nu$  is in  $\Omega_n^\nu$ , i.e.,  $\varsigma_{n,m}^\nu \in \Omega_n^\nu$  with  $\varsigma_{2,m}^\nu$  given in polar coordinates by

$$\varsigma_{2,m}^\nu(\rho, \vartheta) := \xi(\rho)\rho^{\sqrt{1/4 - \nu^2} - 1/2} e^{-i(m+1/2)\vartheta};$$

and  $\varsigma_{3,m}^\nu$  in spherical coordinates by

$$\varsigma_{3,m}^\nu(r, \theta, \phi) := \xi(r)r^{\sqrt{1 - \nu^2} - 1} \Omega_{1/2+m_2, m_1, -m_2}(\theta, \phi).$$

We achieve this goal by the application of the following abstract lemma

LEMMA 16. *Let  $q$  be a closable and non-negative quadratic form on a dense linear subspace  $\mathfrak{Q}$  of the Hilbert space  $\mathfrak{H}$  and  $\psi \in \mathfrak{H}$ . If there is a sequence  $(\psi_n)_{n \in \mathbb{N}} \subset \mathfrak{Q}$  with  $\sup_{n \in \mathbb{N}} q[\psi_n] < \infty$  which converges weakly in  $\mathfrak{H}$  to  $\psi$ , then  $\psi$  is in the domain of the closure of  $q$ .*

*Proof.* We denote by  $\bar{q}$  the closure of  $q$  and by  $\bar{\mathfrak{Q}}$  the domain of  $\bar{q}$ . There is a unique self-adjoint operator  $B : \bar{\mathfrak{Q}} \rightarrow \mathfrak{H}$  with

$$\bar{q}[\varphi] = \|B\varphi\|^2 \text{ for all } \varphi \in \bar{\mathfrak{Q}}$$

by Theorem 2.13 in [18] ( $B^2$  corresponds to  $A$  there). Thus we know that

$$\sup_{n \in \mathbb{N}} \|B\psi_n\|^2 < \infty.$$

Hence there is a  $\Psi \in \mathfrak{H}$  and a subsequence  $(B\psi_{n_k})_{n_k \in \mathbb{N}}$  of  $(B\psi_n)_{n \in \mathbb{N}} \subset \mathfrak{H}$  that converges weakly to  $\Psi$  by the Banach-Alaoglu Theorem. This implies that  $((\psi_{n_k}, B\psi_{n_k}))_{n_k \in \mathbb{N}}$  converges weakly to  $(\psi, \Psi) \in \mathfrak{H} \oplus \mathfrak{H}$ . By the closedness of the graph of  $B$  and Theorem 8 in Chapter 1 of [3] we deduce the claim.  $\square$

Now we pick  $v \in C_0^\infty(\mathbb{R}_+)$  such that  $v(r) = \xi(r)$  for all  $r \in [1, \infty)$  and  $0 \leq v(r) \leq 1$  for  $r \in (0, 1)$ . Let  $k \in \mathbb{N}$ . We define

$$v_k(r) := \begin{cases} v(kr) & \text{if } r \in (0, 1/k]; \\ 1 & \text{if } r \in (1/k, 1]; \\ \xi(r) & \text{else ;} \end{cases}$$

and the function  $\varsigma_{2,m,k}^\nu$  in the polar coordinates by

$$\varsigma_{2,m,k}^\nu(\rho, \vartheta) := v_k(\rho)\rho^{\sqrt{1/4-\nu^2}-1/2}e^{-i(m+1/2)\vartheta},$$

and  $\varsigma_{3,m,k}^\nu$  in the spherical coordinates by

$$\varsigma_{3,m,k}^\nu(r, \theta, \phi) := v_k(r)r^{\sqrt{1-\nu^2}-1}\Omega_{1/2+m_2, m_1, -m_2}(\theta, \phi).$$

The sequence  $(\varsigma_{n,m,k}^\nu)_{k \in \mathbb{N}}$  converges to  $\varsigma_{n,m}^\nu$  in  $L^2(\mathbb{R}^n; \mathbb{C}^{n-1})$ . By Lemma 16 it is thus enough to prove that

$$\sup_{k \in \mathbb{N}} q_n^\nu[\varsigma_{n,m,k}^\nu] < \infty. \tag{57}$$

Let  $\varphi \in C_0^\infty(\mathbb{R}^n \setminus \{0\}; \mathbb{C}^{n-1})$ . At first we observe that

$$q_n^\nu[\varphi] \leq \int_{\mathbb{R}^n} \left( \frac{|x|}{\nu} |K_n \varphi|^2 - \frac{\nu}{|x|} |\varphi|^2 + |\varphi|^2 \right) dx. \tag{58}$$

A tedious calculation shows

$$K_n = \begin{cases} -ie^{i\vartheta}(\partial_\varrho - \frac{1}{\rho}A_2) \text{ with } A_2 := -i\partial_\vartheta \text{ if } n = 2; \\ -i\left(\sigma \cdot \frac{x}{|x|}\right)\left(\partial_r - \frac{1}{r}A_3\right) \text{ with } A_3 := \sigma \cdot (-ix \wedge \nabla) \text{ if } n = 3. \end{cases} \tag{59}$$

Using (59) and integration by parts we obtain that the right hand side of (58) is equal to

$$\int_{\mathbb{R}^n} \left( \frac{|\mathbf{x}|}{\nu} |\partial_{|\mathbf{x}|} \varphi|^2 + \frac{1}{\nu|\mathbf{x}|} |(1/(4-n) + A_n)\varphi|^2 - \frac{\left(\nu + \frac{1}{(4-n)^2\nu}\right)}{|\mathbf{x}|} |\varphi|^2 + |\varphi|^2 \right) d\mathbf{x}. \tag{60}$$

By (60) and Relation 2.1.37 in [1] we obtain

$$\begin{aligned} & \int_{\mathbb{R}^n} \left( \frac{|x|}{\nu} |K_n \varsigma_{n,m,k}^\nu|^2 - \frac{\nu}{|x|} |\varsigma_{n,m,k}^\nu|^2 + |\varsigma_{n,m,k}^\nu|^2 \right) d\mathbf{x} \\ &= \int_0^\infty \left( \frac{t^n}{\nu} \left| \partial_t v_k(t) t^{\sqrt{(4-n)^{-2}-\nu^2}-(4-n)^{-1}} \right|^2 \right. \\ & \quad \left. - \nu v_k(t)^2 t^2 \sqrt{(4-n)^{-2}-\nu^2}-1 + v_k(t)^2 t^2 \sqrt{(4-n)^{-2}-\nu^2} \right) dt. \end{aligned} \tag{61}$$

A straightforward calculation shows that (61) is equal to

$$\begin{aligned} & \int_0^\infty \left( \nu^{-1} v_k'(t)^2 t^2 \sqrt{(4-n)^{-2}-\nu^2}+1 + v_k(t)^2 t^2 \sqrt{(4-n)^{-2}-\nu^2} \right) dt \\ &= \int_0^{1/k} \nu^{-1} k^2 v'(kt)^2 t^2 \sqrt{(4-n)^{-2}-\nu^2}+1 dt + \int_1^\infty \nu^{-1} v'(t)^2 t^2 \sqrt{(4-n)^{-2}-\nu^2}+1 dt \tag{62} \\ &+ \int_0^\infty v_k(t)^2 t^2 \sqrt{(4-n)^{-2}-\nu^2} dt. \end{aligned}$$

An upper bound for the expression in (62) is

$$\int_0^\infty \nu^{-1} v'(t)^2 t^2 \sqrt{(4-n)^{-2}-\nu^2}+1 dt + \int_0^\infty \xi(t)^2 t^2 \sqrt{(4-n)^{-2}-\nu^2} dt. \tag{63}$$

The combination of (63), (62) (61) and (58) implies (57).

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