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## On boundedness of isomerization paths for non- and semirelativistic molecules

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# ON BOUNDEDNESS OF ISOMERIZATION PATHS FOR NON- AND SEMIRELATIVISTIC MOLECULES 

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#### Abstract

This article focuses on isomerizations of molecules, i.e. chemical reactions during which a molecule is transformed into another one with the same atoms but in a different spatial configuration. We consider the special case in which the system breaks into two submolecules whose internal geometry is solid during the whole procedure. We prove, under some conditions, that the distance between the two submolecules stays bounded during the entire reaction. This paper generalizes [6] in two directions. The first one is that we relax the assumptions that the ground state eigenspaces of the submolecules have to fulfill. The second one is that we allow semirelativistic kinetic energy as well.


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An isomerization is a chemical reaction with the property that the reactant is a single molecule with the same atoms as the product, but in a different spatial configuration. One simple example is the isomerization HCN $\rightarrow$ CNH. Isomerizations play, especially in organic and organometallic chemistry, a very important role, see e.g. [24]. The question of how much energy is needed for an isomerization to take place is very fundamental. The study of it occupies a large amount of the numerical computations in quantum chemistry, see for example [32] and references therein.

As explained in [39, Chapter 17], real world chemical systems, involving a large number of particles, are very complicated and some approximations have to be made. We use here the Born-Oppenheimer approximation: we assume that electronic and nuclear motions occur at different time scales, so that one can first solve the Schrödinger equation for electrons with fixed nuclei, and then study the dynamics of the nuclei with a potential corresponding to the electronic energy. In this second step, we treat nuclei as classical particles. The Born-Oppenheimer approximation relies on the fact that nuclei are much heavier than electrons and it is customary to use it for slow chemical reactions. We assume as well that the temperature is zero, so that the only relevant electronic energy is the ground state energy. The approximation should be reasonable for small temperatures. A discussion about the Born-Oppenheimer approximation and its limits, from a chemistry point of view, can be found in [39, Section 3.1]. For a discussion from a mathematical point of view, we refer e.g. to [68].

Another usual approximation is to neglect the relativistic effects. In the non-relativistic setting, the kinetic energy of a particle with mass $m$ is given by the differential operator $-\left(\hbar^{2} / 2 m\right) \Delta$, where $\Delta$ denotes the Laplace operator. Nevertheless, the relativistic effects can be important, becoming stronger and stronger when the atomic numbers of involved atoms increase. Thus, if relativistic correction terms are needed only for very precise calculations for light atoms such as hydrogen, for atoms as heavy as silver or copper, relativistic effects cannot be neglected. For elements of the sixth row of the periodic table such as gold, lead or mercury, relativistic effects alter the chemical and physical behavior qualitatively, see [9]. Examples of relativistic effects are on molecular geometry, especially the length of covalent bonds, but they also give better values for dissociation energies of bonds and ionization potentials and explain common-life phenomena such as the yellow color of gold or the low melting point of mercury. Surveys about relativistic effects in structural chemistry can be found in [57] and [9]. The specific case of isomerizations has been studied for molecules containing heavy atoms such as platinum [61], gold [56] or silver [19] and more generally for organometallic molecules [24].

In order to be able to easily use variational formulations for relativistic systems, we will use the semirelativistic Hamiltonian, where the kinetic energy is obtained by applying to the classical relativistic kinetic energy $E=\sqrt{c^{2} p^{2}+m^{2} c^{4}}-m c^{2}$ the standard quantization $p \rightarrow-\mathrm{i} \hbar \nabla$, which gives the pseudodifferential operator

$$
T=\sqrt{-c^{2} \hbar^{2} \Delta+m^{2} c^{4}}-m c^{2} .
$$

As it is well-known, the operator $T$ is non local, which makes it more difficult to deal with than the standard Laplacian. Hamiltonians associated with this kinetic energy have been studied as early as in the 70s, for example by Herbst in the case of an hydrogenoid ion [34]. Lieb and Yau in [48] and Fefferman and de la Llave [28] discussed its stability in many-particle systems.

The general mathematical framework to study isomerizations is the same in the non-relativistic and in the semirelativistic settings. Two isomers of a molecule correspond to two stable configurations of the entire system. In these configurations, the ground state energy of the system has local minima with respect to the positions of the nuclei. In the Born-Oppenheimer approximation, the isomerization reaction is described by a path of the nuclei connecting these two local minima. The difference between the maximum and initial ground state energies along one path corresponds to the amount of energy needed to bring the system from one state to the other. The activation energy is the lowest energy needed for the reaction to happen, and it is obtained by minimizing the maximal energy over all possible paths linking the two stable configurations. While chemists work hard to compute the activation energy, from a mathematics point of view even its existence, namely the existence of an optimal path, is an open problem.

From the point of view of critical point theory, an isomerization corresponds to a mountain pass problem [3, 65, 38]. The difficulties in this situation are, however, very different from the ones present in most cases of the mathematical literature where mountain pass problems are considered. Usually, the mountain pass theorem is used to prove existence of a critical point of a functional, implying existence of a solution to a PDE. In this situation, the problem is infinite dimensional. In our case, the mountain pass problem is, in its simpler formulation, finite dimensional. However, as we will explain, proving that a min-maxing sequence stays on a bounded set is in our situation often very difficult, and it is the main issue that we deal with in this paper. This boundedness would physically mean that isomerizations take place without breaking the whole system infinitely apart. This was conjectured and proven in some specific cases in $[42,43,44]$. Another particularity of the mountain pass problem that we consider is that it would be important to have an optimal path. This would physically mean that the isomerization could happen in a way such that the energy that is needed can be optimized. The mountain pass theorem, however, does not imply the existence of such a path, even in finite dimensions.

One of the forces that play an important role in our analysis is the van der Waals force. It is a purely quantum and universal force that occurs when two molecules are placed at a distance $L$ from each other, where $L$ is not too small. It is originated by an attractive potential of the order $1 / L^{6}$ and it is thus an attractive force of the order of $1 / L^{7}$. It is induced as a result of quantum correlations between the two molecules. Lieb and Thirring [47] proved, using an intricate test function, an upper bound of the ground state energy of the system, which implies existence of an attractive potential that is, in the order of $L$, at least as strong as van der Waals potential for some orientations. The exact expression in the leading order has been more recently derived for individual non-relativistic atoms by Sigal and the first author of this paper in [8], assuming that the ground state energies of the atoms are irreducible. Irreducible roughly means that the ground states of the atoms are unique up to spin permutations, see e.g. Definition 1 and the related discussion in [6]. In [4], the entire remainder was analyzed in the spinless case in order to
extract information on its dependence on the number of atoms. The long-range asymptotics was studied coupled with the number of atoms going to infinity. Progress in this direction would be important for the study of gases. For the case of two atoms, the results of [8] were improved in [7] in the following way: the monotonicity and derivatives of the interaction energy were studied and the irreducibility assumption was dropped for one of the two atoms. In [5], the interaction energy of a molecule with a half infinite plate was analyzed. In [14], the van der Waals force was derived for semirelativistic atoms, and the results also included an expansion with higher order terms including e.g. the Axilrod-Teller-Muto three-body potential for atoms, see also [33].

In [6], the upper bound of [47] was improved and the results of [8] were extended to the case of molecules, still in a non-relativistic model. The van der Waals force implies that all neutral molecules can bind in the Born-Oppenheimer approximation, that is, the energy always has a global minimum with respect to the nuclear positions. This is because the energy of a molecule which splits into pieces is necessarily higher than the lowest possible energy, due to the van der Waals force. This was already observed in [47].

However, studying isomerizations has additional difficulties, even in the simplest case that the system splits into two solid submolecules, which is the case that we will consider. This is because if two molecules are far apart, there usually exist some orientations for which they repel, for example if they have dipole moments oriented in opposite directions. Deforming a whole path involving such states requires more information of the system. It is necessary to know the number of directions in which one can decrease the energy. At least two such directions are necessary to deform a one-dimensional path and decrease its maximal energy ${ }^{1}$ and to eliminate the loss of boundedness of optimizing sequences of paths. This was one of the main difficulties in the study of isomerizations in $[42,44,6]$.

In [42], it was proved that the critical points of the dipole-dipole interaction which have a positive energy all have a Morse index greater than or equal to 2, i.e. at least two directions of decrease. This was used to prove the boundedness of all reaction paths, for a system containing two rigid molecules, each having a non-degenerate ground state with a non-vanishing dipole moment. In [44], the completely different situation of a molecule with only one moving atom (like for the reaction $\mathrm{HCN} \rightarrow \mathrm{CNH}$ ) was treated. If the single atom escapes to infinity, then only the van der Waals force pertains since, by symmetry, the atomic ground state has no multipole in average. All the other cases were left open in [42, 44].

In [6], the results of [42, 44] were extended in many directions, and important progress was made for the case of a molecule composed of two rigid submolecules. The most important results were in the case when the two molecules have low order multipole moments, that is, such that the leading energy of the system is of the order $1 / L^{p}$ with $p<6$. In this case, the van der Waals force does not dominate in the long range asymptotics and the molecules may repel each other depending on their orientations. That was the main difficulty in [6]. On the other hand, the authors proved that if sufficiently many of the multipole moments of the two molecules vanish, then the van der Waals force dominates in the long range asymptotics, as physically expected. In that case, the boundedness of reaction paths followed quite immediately from the dominance of the van der Waals force.

Here, we extend the results of [6] in two directions. The first one is that we relax the assumption of irreducibility of the ground state eigenspaces of the molecules. The second one is that we consider molecules with semirelativistic kinetic energy as well. As we explained above, isomerizations can involve molecules with heavy atoms, for which the relativistic effects of the electrons become important. Our approach requires to extend results of [6] and [14] to the case of semirelativistic molecules on the van der Waals force, and relax the irreducibility assumption made in [6] and [14]. This direction is of interest on its own. Further difficulties in extending the results of [6] arise on computing derivatives with respect to rotations of a variant of the ground state energy of the system, not only due to the nonlocality of the semirelativistic Laplacian, but also because, in the semirelativistic case, ground states have less regularity than in the non-relativistic case. In particular, in Section 3, the approach of [6] does not seem to work in the semirelativistic setting. We use ideas of [37], which Hunziker used to prove smoothness of the ground state energy for bosonic non-relativistic systems. We rework these ideas so that we can quantify the derivatives and so that the analysis is applicable to the semirelativistic case as well. In the nonrelativistic case, our strategy is simpler than the one of [6] and quantifies higher order derivatives as well.

[^0]On the other hand, the careful study of the multipole-multipole interactions in section 4 of [6] does not depend on the kind of kinetic energy and is thus immediately applicable to our setting.

In the next section, we define the system properly and describe the main open questions. Then we state our new results. The rest of the paper is devoted to the proof of our theorems.

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## 1. Model and main results

1.1. Schrödinger Hamiltonian for molecules. We consider a molecule composed of

- $N$ electrons of charge $e$ with position variables $x:=\left(x_{1}, \ldots, x_{N}\right) \in \mathbb{R}^{3 N}$;
- $M$ nuclei of atomic numbers $\mathcal{Z}:=\left(Z_{1}, \ldots, Z_{M}\right) \in \mathbb{N}^{M}$, and positions $Y:=\left(y_{1}, \ldots, y_{M}\right) \in \mathbb{R}^{3 M}$ with $y_{j} \neq y_{k}$ for $j \neq k$.
The total nuclear charge will be denoted by

$$
\begin{equation*}
|\mathcal{Z}|:=\sum_{i=1}^{M} Z_{i} \tag{1.1}
\end{equation*}
$$

As explained in the introduction, we work in the Born-Oppenheimer approximation. In a first step, the nuclei are fixed and pointwise, and we study the spectrum of a Hamiltonian $H_{N}(Y, \mathcal{Z})$ only describing electrons. In a second step, we consider the ground state energy of $H_{N}(Y, \mathcal{Z})$ as an effective potential for the nuclei in order to investigate properties of reaction paths. To define $H_{N}(Y, \mathcal{Z})$, we need the Hilbert space associated with the electrons, which is

$$
\mathcal{H}:=\bigotimes_{j=1}^{N} L^{2}\left(\mathbb{R}^{3} \times\left\{-\frac{1}{2}, \frac{1}{2}\right\}\right)
$$

with the inner product corresponding to the quadratic form given by

$$
\langle\psi, \psi\rangle:=\sum_{s_{1} \in\{ \pm 1 / 2\}} \ldots \sum_{s_{N} \in\{ \pm 1 / 2\}} \int_{\mathbb{R}^{3}} \ldots \int_{\mathbb{R}^{3}}\left|\psi\left(x_{1}, s_{1}, x_{2}, s_{2}, \ldots, x_{N}, s_{N}\right)\right|^{2} \mathrm{~d} x_{1} \mathrm{~d} x_{2} \ldots \mathrm{~d} x_{N}
$$

The introduction of the set $\{-1 / 2,1 / 2\}$ is due to the fact that an electron has 2 spin states.
The usual models of quantum mechanics do not consider relativistic effects. However, as explained in the introduction, when the molecule contains heavy atoms, we have to take these effects into account. Since we want to consider as well molecules without heavy atoms, we introduce both the non-relativistic and the semirelativistic kinetic energy operators, which are, in physical units and for the $j$-th electron,

$$
T_{j}:= \begin{cases}\sqrt{-\hbar^{2} c^{2} \Delta_{j}+m^{2} c^{4}}-m c^{2}, & \text { semirelativistic case (SR) }  \tag{1.2}\\ -\frac{\hbar^{2}}{2 m} \Delta_{j}, & \text { non-relativistic case (NR) }\end{cases}
$$

where $\Delta_{j}$ is the usual Laplace operator in the variable $x_{j} \in \mathbb{R}^{3}, \hbar$ is the reduced Planck's constant, $c$ the speed of light and $m$ the mass of the electron. The square root of the differential operator is defined through Fourier transformation: for $\psi \in H^{1}\left(\mathbb{R}^{3}\right)$ and almost every $x \in \mathbb{R}^{3}$, we have in the (SR) case

$$
\begin{equation*}
T \psi(x):=\left(\mathcal{F}^{-1}\left[\left(\sqrt{c^{2} \hbar^{2}|\cdot|^{2}+m^{2} c^{4}}-m c^{2}\right) \mathcal{F} \psi\right]\right)(x) \tag{1.3}
\end{equation*}
$$

where $\mathcal{F}$ denotes the Fourier transformation on $L^{2}\left(\mathbb{R}^{3}\right)$

$$
\begin{equation*}
(\mathcal{F} \psi)(p):=\frac{1}{(2 \pi)^{3 / 2}} \int_{\mathbb{R}^{3}} \mathrm{e}^{-\mathrm{i} p \cdot x} \psi(x) \mathrm{d} x \tag{1.4}
\end{equation*}
$$

Of course, $T_{j}$ can also be defined with the Fourier transformation as a sesquilinear form on $H^{\frac{1}{2}}\left(\mathbb{R}^{3}\right)$. Note that, due to the fermionic nature of electrons, we consider the cases where all kinetic energies are semirelativistic or all are non-relativistic. Since we consider a system with $N$ electrons, we introduce the $N$-particle Hamiltonian

$$
\begin{equation*}
H_{N}(Y, \mathcal{Z}):=\sum_{j=1}^{N} T_{j}+\hbar c \alpha I_{N}(Y, \mathcal{Z}) \tag{1.5}
\end{equation*}
$$

where $\alpha:=\frac{e^{2}}{4 \pi \epsilon_{0} \hbar c}$ is a dimensionless physical constant, which is called the fine structure constant $(\alpha \approx$ $1 / 137)$. The electric potential is

$$
\begin{equation*}
I_{N}(Y, \mathcal{Z}):=-\sum_{i=1}^{M} \sum_{j=1}^{N} \frac{Z_{i}}{\left|x_{j}-y_{i}\right|}+\sum_{1 \leqslant i<j \leqslant N} \frac{1}{\left|x_{i}-x_{j}\right|}+\sum_{1 \leqslant i<j \leqslant M} \frac{Z_{i} Z_{j}}{\left|y_{i}-y_{j}\right|} \tag{1.6}
\end{equation*}
$$

where the summands, respectively, represent attraction between the electrons and the nuclei, repulsion between the electrons and repulsion between the nuclei.

In the non-relativistic case, $H_{N}(Y, \mathcal{Z})$ is always self-adjoint with domain $H^{2}\left(\mathbb{R}^{3 N}\right)$ and bounded from below: see for example [45, Theorem 6.2], or [67, Section 11.1], where the proof is given only in the case of atoms but can be easily adapted to molecules. In the semirelativistic case, it was proved by Lieb and Yau in [48] and Fefferman and de la Llave in [28] that, if $\alpha$ is small enough, including the physical case $\alpha \approx 1 / 137$, and, for all $k \in\{1, \ldots, M\}$,

$$
\begin{equation*}
Z_{k}<\frac{2}{\pi \alpha} \tag{1.7}
\end{equation*}
$$

then there exists $b<1$ such that, for all $\psi \in H^{1 / 2}\left(R^{3 N}\right)$,

$$
\begin{equation*}
\left\langle\psi, I_{N}(Y, \mathcal{Z}) \psi\right\rangle \leqslant b\left\langle\psi, \sum_{i=1}^{N} T_{i} \psi\right\rangle \tag{1.8}
\end{equation*}
$$

uniformly in $Y$. As a consequence, $H_{N}(Y, \mathcal{Z})$ is bounded from below and self-adjoint with form domain $H^{\frac{1}{2}}\left(\mathbb{R}^{3 N}\right)$. We note here that the inequality (1.8) does not seem to provide information on the operator domain of $H_{N}(Y, \mathcal{Z})$. In particular, the operator domain may contain functions that are not in $H^{1}$. The $\sum_{i=1}^{N} T_{i}$ operator boundedness of $I_{N}(Y, \mathcal{Z})$ with relative bound less than 1 would guarantee that the domain of $H_{N}(Y, Z)$ is $H^{1}\left(\mathbb{R}^{3 N}\right)$. Such operator boundedness follows from the Hardy inequality if $\alpha$ is small enough but it is not clear to us how small $\alpha$ should be for large molecules. In particular, it is not clear if the physical value $\alpha \approx \frac{1}{137}$ would be included. For this reason, we avoid this assumption.

In the rest of the paper, we will assume that we have chosen units such that

$$
\hbar=c=m=1
$$

Moreover, we assume that $\alpha$ is small enough and that (1.7) holds, so that $H_{N}(Y, \mathcal{Z})$ is bounded from below and we are not going to mention these assumptions explicitly.

For a more precise analysis, one needs to consider the fermionic nature of the electrons. To this end, we consider the permutation group $S_{N}$ of $\{1, \ldots, N\}$. We denote by $\bigwedge_{j=1}^{N} L^{2}\left(\mathbb{R}^{3} \times\left\{-\frac{1}{2}, \frac{1}{2}\right\}\right)$ the Hilbert space of antisymmetric square-integrable functions $\Psi\left(x_{1}, s_{1}, \ldots, x_{N}, s_{N}\right)$ with spin, that is, such that

$$
\begin{equation*}
\pi \cdot \Psi\left(x_{1}, s_{1}, \ldots, x_{N}, s_{N}\right):=\Psi\left(x_{\pi(1)}, s_{\pi(1)}, \ldots, x_{\pi(N)}, s_{\pi(N)}\right)=(-1)^{\pi} \Psi\left(x_{1}, s_{1}, \ldots, x_{N}, s_{N}\right) \tag{1.9}
\end{equation*}
$$

for any permutation $\pi \in S_{N}$, and we make use of the orthogonal projection on this space

$$
Q_{N}: \bigotimes_{j=1}^{N} L^{2}\left(\mathbb{R}^{3} \times\left\{-\frac{1}{2}, \frac{1}{2}\right\}\right) \rightarrow \bigwedge_{j=1}^{N} L^{2}\left(\mathbb{R}^{3} \times\left\{-\frac{1}{2}, \frac{1}{2}\right\}\right)
$$

where

$$
\begin{equation*}
Q_{N} \Psi=\frac{1}{N!} \sum_{\pi \in S_{N}}(-1)^{\pi} \pi \cdot \Psi \tag{1.10}
\end{equation*}
$$

with $\pi \cdot \Psi$ defined in (1.9). Then, since the Hamiltonian $H_{N}(Y, \mathcal{Z})$, defined in (1.5), commutes with $Q_{N}$, we can consider the fermionic Hamiltonian defined by

$$
\begin{equation*}
\hat{H}_{N}(Y, \mathcal{Z}):=\left.H_{N}(Y, \mathcal{Z})\right|_{\operatorname{Ran} Q_{N}} \tag{1.11}
\end{equation*}
$$

which acts on the Hilbert space

$$
\begin{equation*}
Q_{N} \mathcal{H}=\bigwedge_{j=1}^{N} L^{2}\left(\mathbb{R}^{3} \times\left\{-\frac{1}{2}, \frac{1}{2}\right\}\right) \tag{1.12}
\end{equation*}
$$

The operator $\hat{H}_{N}(Y, \mathcal{Z})$ is essentially self-adjoint on $C_{c}^{\infty}\left(\mathbb{R}^{3 N}\right)$ and its form domain is

$$
Q\left(\hat{H}_{N}(Y, \mathcal{Z})\right)=\left\{\begin{array}{l}
H^{\frac{1}{2}}\left(\left(\mathbb{R}^{3} \times\{ \pm 1 / 2\}\right)^{N}\right) \cap Q_{N} \mathcal{H},  \tag{1.13}\\
H^{1}\left(\left(\mathbb{R}^{3} \times\{ \pm 1 / 2\}\right)^{N}\right) \cap Q_{N} \mathcal{H},
\end{array}\right.
$$

where we recall that (NR), (SR) respectively stand for non-relativistic and semirelativistic case, as defined by (1.2).

We denote by

$$
\begin{equation*}
E_{N}(Y, \mathcal{Z}):=\min \sigma\left(\hat{H}_{N}(Y, \mathcal{Z})\right) \tag{1.14}
\end{equation*}
$$

the bottom of the spectrum of $\hat{H}_{N}(Y, \mathcal{Z})$. Recall that $|\mathcal{Z}|$ was defined in (1.1). In the (NR) case, when $N<|\mathcal{Z}|+1$ (neutral or positively charged molecules), the HVZ [36, 70, 71] and Zhislin/Zhislin-Sigalov theorems [71, 72] imply that $E_{N}(Y, \mathcal{Z})$ is an eigenvalue of $\hat{H}_{N}(Y, \mathcal{Z})$, lying strictly below its essential spectrum:

$$
E_{N}(Y, \mathcal{Z})<\min \sigma_{\mathrm{ess}}\left(\hat{H}_{N}(Y, \mathcal{Z})\right)=E_{N-1}(Y, \mathcal{Z})
$$

and that the ground state eigenspace of $\hat{H}_{N}(Y, \mathcal{Z})$ is finite-dimensional. Moreover, exponential decay of eigenfunctions associated with eigenvalues below the bottom of the essential spectrum is known, see e.g. [23, 63, 35, 29, 31]. In the (SR) case, analogous statements were proven in [14] but only for atoms. Previous results on exponential decay of eigenfunctions of semirelativistic Hamiltonians had already been proven by Carmona, Masters and Simon in [18] and by Nardini in [54] and [55]. Nevertheless, [54] applied only in the 2-body case, [18] does not include the case of Coulomb potentials and [55] does not take fermionic statistics into account and was not applicable in the setting of [14]. Closely following [14], we extend their results for molecules in the simpler setting where we do not need to work with irreducible representations of $S_{N}$. Most of the parts of these proofs are simple adaptations of the analysis of [14] and earlier works but for convenience of the reader we shall include standard arguments.

As the theorems are known in the non-relativistic case, we state them here independently of the case. We begin with the extension of the HVZ theorem to the (SR) case.

Theorem 1.1. We have that

$$
\sigma_{\text {ess }}\left(\hat{H}_{N}(Y, \mathcal{Z})\right)=\left[\inf \sigma\left(\hat{H}_{N-1}(Y, \mathcal{Z})\right), \infty\right)
$$

The next theorem and corollary provide exponential decay of (approximate) eigenfunctions to eigenvalues below the essential spectrum. For $x \in \mathbb{R}^{d}$, we denote by $|x|$ the Euclidean norm of $x$. The ball of radius $R>0$ and centered in $y \in \mathbb{R}^{d}$ is denoted by

$$
B_{R}(y):=\left\{x \in \mathbb{R}^{d}| | x-y \mid \leqslant R\right\} .
$$

Theorem 1.2 (Exponential decay of approximate eigenfunctions). We define the ionization threshold of $\hat{H}_{N}(Y, \mathcal{Z})$ as

$$
\begin{equation*}
\widetilde{\Sigma}:=\lim _{R \rightarrow \infty} \inf _{\substack{\psi \in D\left(\hat{H}_{N}(Y, \mathcal{Z})\right) \backslash\{0\}, \operatorname{supp} \psi \cap B_{R}(0)=\emptyset}} \frac{\left\langle\psi, \hat{H}_{N}(Y, \mathcal{Z}) \psi\right\rangle}{\langle\psi, \psi\rangle} . \tag{1.15}
\end{equation*}
$$

For any fixed $\mu<\widetilde{\Sigma}$, let $\left(\gamma_{s}\right)_{s \in \mathcal{I}}$ be a family of functions in $\mathcal{D}\left(\hat{H}_{N}(Y, \mathcal{Z})\right)$ satisfying for all $s \in \mathcal{I}$,

$$
\begin{equation*}
\left\|\gamma_{s}\right\|_{L^{2}} \leqslant C, \quad\left(H_{N}(Y, \mathcal{Z})-\mu\right) \gamma_{s}=\Gamma_{s} \tag{1.16}
\end{equation*}
$$

Here $C>0$ is a constant and $\left(\Gamma_{s}\right)_{s \in \mathcal{I}}$ is such that there exists $C_{1}, a>0$ such that $\left\|e^{a|\cdot|} \Gamma_{s}\right\|_{L^{2}\left(\mathbb{R}^{3 N}\right)} \leqslant C_{1}$, for all $s \in \mathcal{I}$. Then $\exists b, D>0$ such that, for all $s \in \mathcal{I}$

$$
\begin{equation*}
\left\|e^{b|\cdot|} \gamma_{s}\right\|_{L^{2}\left(\mathbb{R}^{3 N}\right)} \leqslant D \tag{1.17}
\end{equation*}
$$

Note that the fact that $\gamma_{s}$ is not assumed to be an exact eigenfunction of $\hat{H}_{N}(Y, \mathcal{Z})$ can be useful: one may prove with Theorem 1.2 that an exponentially decaying function remains exponentially decaying after multiplication with a resolvent of $H_{N}(Y, \mathcal{Z})$.

Corollary 1.3 (Exponential decay of eigenfunctions). Choosing $\Gamma=0$, the above theorem implies that any eigenfunction $\gamma$ of $\hat{H}_{N}(Y, \mathcal{Z})$ with associated eigenvalue below its essential spectrum is exponentially decaying, in the sense that (1.17) holds.

The following theorem is an extension of Zhislin's theorem to the (SR) case.
Theorem 1.4. If $|\mathcal{Z}|>N-1$, then $\hat{H}_{N}(Y, \mathcal{Z})$ has infinitely many eigenvalues below the infimum of its essential spectrum. In particular, it has a ground state and the ground state eigenspace is finitedimensional.

We prove Theorems 1.1, 1.2 and 1.4 in Appendix B.
Remark 1.5. Note that, in [14], a variant of Theorem 1.4 was proven for (SR) atoms, but only the existence of at least one eigenvalue below the bottom of the essential spectrum was shown. Due to the nonlocality of the (SR) kinetic energy operator, we have to modify their argument in order to prove that there are infinitely many eigenvalues below the bottom of the essential spectrum.

From Theorems 1.1 and 1.4 , it follows that, for every $Y=\left(y_{1}, \ldots, y_{M}\right) \in \mathbb{R}^{3 M}$ with $y_{j} \neq y_{k}$ for $j \neq k$, there exists at least one eigenfunction $\Psi \in D\left(\hat{H}_{N}(Y, \mathcal{Z})\right)$, such that

$$
\hat{H}_{N}(Y, \mathcal{Z}) \Psi=E_{N}(Y, \mathcal{Z}) \Psi
$$

Moreover, the eigenspace associated with $E_{N}(Y, \mathcal{Z})$ has finite dimension.
1.2. Isomerizations in the case of two rigid submolecules: the model. From now on, we consider only globally neutral systems, i.e. the case where

$$
N=|\mathcal{Z}|
$$

Moreover, we assume that the molecule is composed of two rigid submolecules, which can only be rotated and translated with respect to each other. This restriction is unsatisfactory as this has no physical justification and, in many isomerizations, there is splitting to more than two solid submolecules. However, removing the assumption seems to be very hard. In some particular cases, e.g. reaction HCN $\rightarrow \mathrm{CNH}$, the assumption is reasonable. Some bigger organic molecules are made of 2 parts, which, if isomerization happens, should approximately remain rigid. This should be the case, for example, for the cis-trans isomerization of azobenzene, see e.g. [66]. We refer to the introduction and appendix C of [6] for some discussion of the general case of several submolecules.

Let $Y_{1}:=\left(y_{1}, \ldots, y_{M_{1}}\right) \in\left(\mathbb{R}^{3}\right)^{M_{1}}, Y_{2}:=\left(y_{M_{1}+1}, \ldots, y_{M}\right) \in\left(\mathbb{R}^{3}\right)^{M_{2}}$ be the nuclear positions of the two solid submolecules, where $M_{1}+M_{2}=M$. The collections of their atomic numbers are denoted by $\mathcal{Z}_{1}:=\left(Z_{1}, \ldots, Z_{M_{1}}\right) \in \mathbb{N}^{M_{1}}, \mathcal{Z}_{2}:=\left(Z_{M_{1}+1}, \ldots, Z_{M}\right) \in \mathbb{N}^{M_{2}}$.

We assume that each of these molecules is neutral, which means that the $i$-th molecule has $N_{i}=\left|\mathcal{Z}_{i}\right|$ electrons, where $\left|\mathcal{Z}_{i}\right|$ is defined as in (1.1). For $i=1,2$, we can define the Hamiltonian for molecule $i$ as in (1.5):

$$
\begin{equation*}
H_{i}:=H_{\left|\mathcal{Z}_{i}\right|}\left(Y_{i}, \mathcal{Z}_{i}\right) \tag{1.18}
\end{equation*}
$$

as well as its version with statistics

$$
\begin{equation*}
\hat{H}_{i}:=\hat{H}_{\left|\mathcal{Z}_{i}\right|}\left(Y_{i}, \mathcal{Z}_{i}\right) \tag{1.19}
\end{equation*}
$$

as in (1.11) and its ground state energy:

$$
\begin{equation*}
E_{i}:=E_{\left|\mathcal{Z}_{i}\right|}\left(Y_{i}, \mathcal{Z}_{i}\right) \tag{1.20}
\end{equation*}
$$

in the same way as in (1.14). We then consider all possible ways of placing these two molecules in space. Without loss of generality, we can place the first molecule at the origin and the second one at a distance $L$ in the direction $e_{1}:=(1,0,0)$, and simply rotate the two molecules using $U, V \in \mathrm{SO}(3)$ (see Figure 1). Our sole variables are therefore $(L, U, V) \in(0, \infty) \times \mathrm{SO}(3) \times \mathrm{SO}(3)$. For shortness, we introduce the new variable

$$
\tau:=(L, U, V) \in(0, \infty) \times \mathrm{SO}(3) \times \mathrm{SO}(3)
$$

and denote by

$$
\begin{equation*}
Y(\tau):=\left(U y_{1}, \ldots, U y_{M_{1}}, V y_{M_{1}+1}+L e_{1}, \ldots, V y_{M}+L e_{1}\right) \tag{1.21}
\end{equation*}
$$

the nuclear positions, as well as by

$$
\begin{equation*}
\mathcal{E}_{\tau}:=E_{N}(Y(\tau), \mathcal{Z})=\inf \sigma\left(\hat{H}_{N}(Y(\tau), \mathcal{Z})\right) \tag{1.22}
\end{equation*}
$$



Figure 1. The case of two rigid molecules (the picture was taken from [6] and was created by Mathieu Lewin).
the corresponding ground state energy. The Hamiltonian has now the explicit expression

$$
\begin{equation*}
H_{\tau}:=H_{N}(Y(\tau), \mathcal{Z})=H_{1, \tau}+H_{2, \tau}+I_{\tau} \tag{1.23}
\end{equation*}
$$

where $H_{1, \tau} H_{2, \tau}$ and $I_{\tau}$ are, respectively, the Hamiltonians for the first and the second molecule and the interaction among them. Explicitly, they are given by

$$
\begin{align*}
& H_{1, \tau}:=\sum_{j=1}^{N_{1}} T_{j}-\sum_{j=1}^{N_{1}} \sum_{k=1}^{M_{1}} \frac{Z_{k} \alpha}{\left|x_{j}-U y_{k}\right|}+\sum_{1 \leqslant i<j \leqslant N_{1}} \frac{\alpha}{\left|x_{j}-x_{i}\right|}+\sum_{1 \leqslant k<l \leqslant M_{1}} \frac{Z_{k} Z_{l} \alpha}{\left|U y_{k}-U y_{l}\right|},  \tag{1.24}\\
& H_{2, \tau}:=\sum_{j=N_{1}+1}^{N} T_{j}-\sum_{j=N_{1}+1}^{N} \sum_{k=M_{1}+1}^{M} \frac{Z_{k} \alpha}{\left|x_{j}-V y_{k}-L e_{1}\right|}+\sum_{N_{1}+1 \leqslant i<j \leqslant N} \frac{\alpha}{\left|x_{j}-x_{i}\right|} \\
& +\sum_{M_{1}+1 \leqslant k<l \leqslant M} \frac{Z_{k} Z_{l} \alpha}{\left|V y_{k}-V y_{l}\right|},  \tag{1.25}\\
& I_{\tau}:=-\sum_{j=1}^{N_{1}} \sum_{k=M_{1}+1}^{M} \frac{Z_{k} \alpha}{\left|x_{j}-V y_{k}-L e_{1}\right|}-\sum_{j=N_{1}+1}^{N} \sum_{k=1}^{M_{1}} \frac{Z_{k} \alpha}{\left|x_{j}-U y_{k}\right|} \\
& \\
& \quad+\sum_{k=1}^{M_{1}} \sum_{l=M_{1}+1}^{M} \frac{Z_{k} Z_{l} \alpha}{\left|U y_{k}-V y_{l}-L e_{1}\right|}+\sum_{i=1}^{N_{1}} \sum_{j=N_{1}+1}^{N} \frac{\alpha}{\left|x_{i}-x_{j}\right|} .
\end{align*}
$$

Note that there is here a small abuse of notation, since $H_{j, \tau}$ can act either on $L^{2}\left(\mathbb{R}^{3 N_{j}}\right)$ or on the whole $L^{2}\left(\mathbb{R}^{3 N}\right)$. We define as well the versions with statistics of the Hamiltonians by:

$$
\begin{equation*}
\hat{H}_{\tau}:=\left.H_{\tau}\right|_{\operatorname{Ran} Q_{N}}, \quad \hat{H}_{1, \tau}:=\left.H_{1, \tau}\right|_{\operatorname{Ran} Q_{N_{1}}}, \quad \hat{H}_{2, \tau}:=\left.H_{2, \tau}\right|_{\operatorname{Ran} Q_{N_{2}}}, \tag{1.26}
\end{equation*}
$$

where $Q_{N_{i}}$ is the antisymmetrizer for $N_{i}$ particles, as defined in (1.10).
Since we are interested in isomerizations, we have to assume that there exist different isomers for the considered molecule. This is done with the following condition.

Condition. (Min) There exist two local minima $\tau_{0}, \tau_{1} \in(0, \infty) \times \mathrm{SO}(3) \times \mathrm{SO}(3)$ of $\tau \mapsto \mathcal{\mathcal { E } _ { \tau }}$ that are strict with respect to the relative displacements $\tau$ that do not leave $Y(\tau)$ invariant.

This condition is natural, as it means that the reactant and the product of the isomerization are stable molecules, in the sense that energy is needed to change their structure.

We consider the mountain pass level

$$
\begin{equation*}
c:=\inf _{\substack{\tau(0)=\tau_{0}, \tau(1)=\tau_{1}}} \max _{t \in[0,1]} \mathcal{E}_{\tau(t)} \tag{1.27}
\end{equation*}
$$

where, as usual, it is understood that $\tau:[0,1] \rightarrow(0, \infty) \times \mathrm{SO}(3) \times \mathrm{SO}(3)$ is continuous. If $\tau_{0}$ gives the shape of the reactant of the isomerization, the difference $c-\mathcal{E}_{\tau_{0}}$ can be understood as the least energy needed for the isomerization to happen. The existence of a continuous path $\widetilde{\tau}$ with $\widetilde{\tau}(0)=\tau_{0}, \widetilde{\tau}(1)=\tau_{1}$ and $c=\max _{t \in[0,1]} \mathcal{E}_{\widetilde{\tau}(t)}$ would imply that the needed energy for the interaction can be indeed optimized. As we mentioned above, this is, to our best of knowledge, an open problem.

In our setting, we consider paths where the system is in a ground state at each time. However, it makes sense to let the electrons evolve along the path as well. For this reason, considering two local
$\operatorname{minima} \tau_{0}$ and $\tau_{1}$ as in Condition (Min), we choose two ground states $\bar{\Psi}_{0}$ and $\bar{\Psi}_{1}$ corresponding to the ground state energies $E_{\tau_{0}}$ and $E_{\tau_{1}}$, respectively, and define a new mountain pass level by

$$
\begin{equation*}
c^{\prime}:=\inf _{\substack{(\tau, \Psi)(0)=\left(\tau_{0}, \bar{\Psi}_{0}\right) \\(\tau, \Psi)(1)=\left(\tau_{1}, \Psi_{1}\right)}} \max _{t \in[0,1]}\left\langle\Psi(t), H_{\tau(t)} \Psi(t)\right\rangle \tag{1.28}
\end{equation*}
$$

which is independent of the choices of $\bar{\Psi}_{0}, \bar{\Psi}_{1}$ and obviously larger than or equal to $c$. Recall that $H_{\tau(t)}$ is defined by (1.23). It is understood here that $\tau(t)$ is a path as before and that $\Psi(t)$ is a continuous path of wavefunctions in the form domain of $\hat{H}_{\tau(t)}$ (which is in fact independent of $t$, as stated in (1.13)), such that $\|\Psi(t)\|_{L^{2}}=1$ for all $t \in[0,1]$.

We note that the mountain pass problem (1.28) is an infinite dimensional mountain pass problem, and even if existence of an optimizer for the problem (1.27) were known, existence of an optimizer for (1.28) would not immediately follow. Nevertheless, it is proven in [42, Theorem 4] (see also [6, Appendix A]) that

$$
c^{\prime}=c
$$

Moreover, existence of a bounded min-maxing sequence $\left(\tau_{n}\right)$ for (1.27) implies existence of a bounded min-maxing sequence $\left(\tau_{n}, \Psi_{n}\right)$ for (1.28). If an optimal path $\left(\tau_{n}, \Psi_{n}\right)$ exists, the position of its maximum is called a transition state in chemistry. In the (NR) setting, Theorem 4 of [42] provides, under condition (Min) and boundedness of a min-maxing sequence, existence of such a transition state in a more general setting: it is not necessary that the system splits into only two rigid submolecules. We conjecture that this should hold in the (SR) case too.
1.3. Case of two rigid molecules: new results. In this section, we state our results, which generalize those in [6, 42, 44].

We begin by giving an asymptotic expansion of the energy of the 2 interacting submolecules, before stating our main theorem.
1.3.1. Expansion of the energy: multipoles and the van der Waals force. In this section, we expand the energy $\mathcal{E}_{\tau}=\mathcal{E}_{(L, U, V)}$ for large $L$ up to order $L^{-6}$ and get the van der Waals energy as well as all the lower order multipolar energies. In particular, we prove that there always exists a non-vanishing term of order $L^{-6}$. In the analysis, error terms depending on $L, U, V$ are going to appear. Since we very often need to bound them uniformly in $U, V$ with functions of $L$, we introduce the following notation: Let $g=g(L, U, V)$ be a function with values on a normed vector space $W$ and $f: \mathbb{R}^{+} \rightarrow \mathbb{R}^{+}$. We say that

$$
g=O_{\left(L^{\infty}, W\right)}(f(L))
$$

if there exist constants $C, D>0$ such that if $L \geqslant D$ then

$$
\|g(L, U, V)\|_{W} \leqslant C f(L)
$$

uniformly in $U, V$.
Let us now introduce the set of matrices

$$
\begin{equation*}
\Gamma:=\left\{M \in \mathfrak{s o}(3)| | M_{i j} \mid \leqslant 1, \forall i, j \in\{1,2,3\}\right\} \tag{1.29}
\end{equation*}
$$

where recall that $\mathfrak{s o}(3)$, the Lie algebra of $\mathrm{SO}(3)$, is the set of 3 -by- 3 real antisymmetric matrices. We say that

$$
\begin{equation*}
g=O_{\left(C^{\infty}, W\right)}(f(L)) \tag{1.30}
\end{equation*}
$$

if there exists a constant $D>0$ such that, for all $n \in \mathbb{N} \cup\{0\}$, there exists a constant $C_{n}$ such that, for all $L \geqslant D,(U, V) \in \operatorname{SO}(3)^{2}$, all $A, B \in \Gamma$ and all $t$ in a neighborhood of 0 , we have

$$
\begin{equation*}
\left\|\frac{\mathrm{d}^{n}}{\mathrm{~d} t^{n}} g\left(L, e^{t A} U, e^{t B} V\right)\right\|_{W} \leqslant C_{n} f(L) \tag{1.31}
\end{equation*}
$$

If $W=\mathbb{R}$ or $W=\mathbb{C}$, then $W$ is going to be omitted from the notation and we will write $O_{L^{\infty}}(f(L))$ or $O_{C^{\infty}}(f(L))$. When we write $O_{\left(L^{\infty}, W\right)}\left(e^{-d L}\right)$ or $O_{\left(C^{\infty}, W\right)}\left(e^{-d L}\right)$, we will mean, even if it is not mentioned, that there exists a $d>0$ such that we have the respective estimates. The constant $d$ might change from one equation to the other.

In order to be able to formulate our theorems, we need to introduce the multipolar interactions. We recall here the expansions introduced in [6, Section 2.1], to which it can be referred for a more detailed explanation. The reader already aware of it can directly go to Definition 1.8.

For a bounded measure $\rho$ decaying faster than any polynomial at infinity, we can define the $2^{n}$-pole moment associated with this measure as a tensor of order $n$ on $\mathbb{R}^{3}$ from the Taylor expansion of $\rho * 1 / \| \cdot \mid$ :

$$
\begin{equation*}
\mathcal{M}_{\rho}^{(n)}\left(h_{1}, \ldots, h_{n}\right):=\frac{(-1)^{n}}{n!} \int_{\mathbb{R}^{3}}|x|^{2 n+1}\left[\left(h_{1} \cdot \nabla_{x}\right) \ldots\left(h_{n} \cdot \nabla_{x}\right) \frac{1}{|x|}\right] \mathrm{d} \rho(x) \tag{1.32}
\end{equation*}
$$

These multipoles are used in the following expansion:
Lemma 1.6 (Lemma 2.1, [6]). For any $N \geqslant 1$, there exists $c_{N}>0$ such that for all $L>1, x \in \mathbb{R}^{3}$ such that $|x|<L / 2$, and any multi-index $\alpha \in \mathbb{N}^{N}$ we have:

$$
\left|\partial_{x}^{\alpha}\left(\frac{1}{\left|x-L e_{1}\right|}-\sum_{n=0}^{N} \frac{\mathcal{M}_{\delta_{x}}^{(n)}\left(e_{1}, \ldots, e_{1}\right)}{L^{n+1}}\right)\right| \leqslant c_{N} \frac{1+|x|^{N+1}}{L^{N+1}}
$$

where $\delta_{x}$ is the Dirac mass at $x$.
Here, we note that [ 6 , Lemma 2.1] was stated with the restriction $0 \leqslant|\alpha| \leqslant 2$ because it was needed only for such $\alpha$. This restriction, however, was not necessary for the proof.

We will use multipole moments defined by (1.32) in the case where the measures are charge densities, which we define now. Given a state $\psi \in Q_{N} \mathcal{H}$, where recall that $Q_{N} \mathcal{H}$ has been defined in (1.12), we define the electronic density on $\mathbb{R}^{3}$ by

$$
\begin{equation*}
\rho_{\psi}(x):=N \sum_{s_{1} \in\left\{ \pm \frac{1}{2}\right\}} \ldots \sum_{s_{N} \in\left\{ \pm \frac{1}{2}\right\}} \int_{\mathbb{R}^{3}} \ldots \int_{\mathbb{R}^{3}}\left|\psi\left(x, s_{1}, x_{2}, s_{2}, \ldots, x_{N}, s_{N}\right)\right|^{2} \mathrm{~d} x_{2} \ldots \mathrm{~d} x_{N} \tag{1.33}
\end{equation*}
$$

where we multiplied the first marginal by $N$ instead of summing over all marginals, which we can do because electrons are fermions.

Given normalized ground states $\Psi_{j}$ of $\hat{H}_{j}, j=1,2$, as defined in (1.19), we can define the molecular charge densities $\rho_{j}^{\Psi_{j}}, j=1,2$ by

$$
\begin{equation*}
\rho_{1}^{\Psi_{1}}(x):=\sum_{i=1}^{M_{1}} Z_{i} \delta_{Y_{i}}(x)-\rho_{\Psi_{1}}(x) ; \quad \rho_{2}^{\Psi_{2}}(x):=\sum_{i=M_{1}+1}^{M} Z_{i} \delta_{Y_{i}}(x)-\rho_{\Psi_{2}}(x) . \tag{1.34}
\end{equation*}
$$

Note that they depend, in general, on the choice of the ground state.
We introduce the multipolar interaction associated with two bounded measures $\rho_{1}$ and $\rho_{2}$ on $\mathbb{R}^{3}$ decaying faster than any polynomial at infinity to be

$$
\begin{align*}
& \mathcal{F}^{(n, m)}\left(\rho_{1}, \rho_{2}\right):=K_{m, n} \sum_{j_{1}, \ldots, j_{n}} \sum_{k_{1}, \ldots, k_{m}} \mathcal{M}_{\rho_{1}}^{(n)}\left(e_{j_{1}}, \ldots, e_{j_{n}}\right) \mathcal{M}_{\rho_{2}}^{(m)}\left(e_{k_{1}}, \ldots, e_{k_{m}}\right) \times \\
& \times\left.\left(\partial_{z_{j_{1}}} \ldots \partial_{z_{k_{m}}} \frac{1}{\left|z+e_{1}\right|}\right)\right|_{z=0} \tag{1.35}
\end{align*}
$$

where $K_{m, n}=(-1)^{m} / \Pi_{i=1}^{m}(2 i-1) \Pi_{j=1}^{n}(2 j-1)$, and define then the multipolar interactions between the two molecules in respective states $\Psi_{1}$ and $\Psi_{2}$ to be $\mathcal{F}^{(n, m)}\left(\rho_{1}^{\Psi_{1}}, \rho_{2}^{\Psi_{2}}\right)$.

The following Lemma is key in order to apply the expansion of the Coulomb interaction and exploit the dependence of the energy on the multipole moments. As we will see, the $2^{n}$-pole- $2^{m}$-pole interaction decays as $1 / L^{n+m+1}$. More explicitly, the charge-charge, dipole-charge, dipole-dipole and dipole-quadrupole interactions are proportional to $L^{-1}, L^{-2}, L^{-3}$ and $L^{-4}$, respectively. The quadrupole-quadrupole and dipole-octopole interactions are proportional to $L^{-5}$.

Lemma 1.7 (Lemma 2.2, [6]). Let $\rho_{1}$ and $\rho_{2}$ be two bounded measures with support in $B_{L / 3}(0)$ such that $\rho_{1}\left(\mathbb{R}^{3}\right)=\rho_{2}\left(\mathbb{R}^{3}\right)=0$. For any $U \in \mathrm{SO}(3)$, we define the rotated measure $U \rho_{j}(j=1,2)$ by $U \rho_{j}(D):=\rho\left(\left\{x, U^{-1} x \in D\right\}\right)$ for any Borel set $D$. For any $N \geqslant 1$ and $U, V \in \mathrm{SO}(3)$,

$$
\int_{\mathbb{R}^{3} \times \mathbb{R}^{3}} \frac{\mathrm{~d} \rho_{1}(x) \mathrm{d} \rho_{2}(y)}{\left|U x-V y-L e_{1}\right|}=\sum_{1 \leqslant n+m \leqslant N} \frac{\mathcal{F}^{(n, m)}\left(U \rho_{1}, V \rho_{2}\right)}{L^{n+m+1}}+K O_{C^{\infty}}\left(\frac{1}{L^{N+2}}\right)
$$

where

$$
K=\prod_{j=1,2} \int_{\mathbb{R}^{3}} \mathrm{~d}\left|\rho_{j}\right|(x)\left(1+|x|^{N+1}\right)
$$

Here again, we claim $C^{\infty}$ regularity while the statement in [6] only gives $C^{2}$ regularity. However, looking at the proof in that paper, we see that this restriction was only justified by the context and that, in fact, we can get information on derivatives of any order.

Lemma 1.7 implies that the leading term of the expansion of the Coulomb interaction is given by the first multipole-multipole interaction which does not vanish. It is therefore important to know the first non-vanishing $2^{n}$-pole moment of each submolecule.

Definition 1.8 (First non-vanishing multipole moment). For $k \in\{1,2\}$, let $\Psi_{k}$ be a normalized ground state of $\hat{H}_{k}$ and let $\rho_{k}^{\Psi_{k}}$ denote the corresponding molecular density. We call $n_{k}$ the smallest integer $n \geqslant 1$ such that $\mathcal{M}_{\rho_{\Psi_{k}}}^{(n)} \neq 0$, see (1.32). If all the multipoles vanish, we let $n_{k}:=+\infty$.

Note that, in principle, $n_{k}$ depends on the choice of a normalized ground state $\Psi_{k}$. Nevertheless, we will often assume that it is in fact not the case, at least for the leading term.

Condition. (Poles) We assume that one of the following condition holds
(Poles-all) : The numbers $n_{1}, n_{2}$ are independent of the choice of normalized ground states $\Psi_{1}, \Psi_{2}$ of $\hat{H}_{1}, \hat{H}_{2}$, respectively, and so are all their multipole moments $\mathcal{M}_{\rho_{\Psi_{1}}}^{(n)}, \mathcal{M}_{\rho_{\Psi_{2}}}^{(n)}$ for all $n \leqslant 4$.
(Poles-lead) : The numbers $n_{1}, n_{2}$ are independent of the choice of normalized ground states $\Psi_{1}, \Psi_{2}$ of $\hat{H}_{1}, \hat{H}_{2}$, respectively, and so are their leading multipole moments $\mathcal{M}_{\rho_{\Psi_{1}}}^{\left(n_{1}\right)}, \mathcal{M}_{\rho_{\Psi_{2}}}^{\left(n_{2}\right)}$.

Remark 1.9. Condition (Poles) follows from the assumption of the irreducibility of the ground state eigenspaces that was assumed in previous works like [6]. Condition (Poles) is weaker, at least we do not see a reason why Condition (Poles-all) or (Poles-lead) should imply irreducibility. However, proving even Condition (Poles-lead) is in the general case, to our best of knowledge, an open problem and it is known only for the hydrogen atom and molecule and for the helium atom, at least in the non-relativistic case, see e.g. [60, Section XIII.12]. In these cases, irreducibility is also known. However, we find Condition (Poles) much more natural from a physical perspective, as chemistry handbooks speak about the dipole or quadrupole etc. moment of a molecule, see for example [62]. If we ignore the fermionic statistics, the ground state eigenspaces of the molecules are, at least in the non-relativistic case, non-degenerate, see again [60, Section XIII.12]. As a consequence, Condition (Poles) is automatically satisfied.

Another assumption, inspired by physical evidence, is that the configuration of lowest energy is the one where each of the two submolecules is neutral.
Condition. (Neutr) The system satisfies the neutral best configuration assumption: namely $\left(Y_{1}, \mathcal{Z}_{1}\right)$ and $\left(Y_{2}, \mathcal{Z}_{2}\right)$ are such that the only configuration of the system with the minimum energy possible is the one where the charges are distributed such that the molecules are neutral:

$$
\begin{equation*}
E_{\left|\mathcal{Z}_{1}\right|}\left(Y_{1}, \mathcal{Z}_{1}\right)+E_{\left|\mathcal{Z}_{2}\right|}\left(Y_{2}, \mathcal{Z}_{2}\right)<\min _{\substack{N_{1}^{\prime} \neq\left|\mathcal{Z}_{1}\right|, N_{1}^{\prime}+N_{2}^{\prime}=N}}\left(E_{N_{1}^{\prime}}\left(Y_{1}, \mathcal{Z}_{1}\right)+E_{N_{2}^{\prime}}\left(Y_{2}, \mathcal{Z}_{2}\right)\right) \tag{1.36}
\end{equation*}
$$

where the energies appearing in the inequality are defined in the same way as $E_{N}(Y, \mathcal{Z})$ in (1.14).
It is a famous conjecture that (1.36) holds, see e.g. the introduction of [8] for some discussions. If this is not the case, the mountain pass problem is actually rather easy, see [42, Thm. 4].

From the results of [51, 52, 42], we have

$$
\begin{equation*}
\lim _{L \rightarrow \infty} \mathcal{E}_{(L, U, V)}=E_{\infty}:=\min _{N_{1}^{\prime}+N_{2}^{\prime}=N}\left(E_{N_{1}^{\prime}}\left(Y_{1}, \mathcal{Z}_{1}\right)+E_{N_{2}^{\prime}}\left(Y_{2}, \mathcal{Z}_{2}\right)\right) \tag{1.37}
\end{equation*}
$$

uniformly in $U, V \in \operatorname{SO}(3)$, where $E_{N_{j}}\left(Y_{j}, \mathcal{Z}_{j}\right), j=1,2$ is defined in the same way as $E_{N}(Y, \mathcal{Z})$ in (1.14).
In order to state the theorem about the asymptotic expansion of the energy, we introduce the following objects.

- $H_{\infty}$ is the sum of the Hamiltonians of the two submolecules, without interaction:

$$
\begin{equation*}
H_{\infty}:=\hat{H}_{1} \otimes \mathbb{1}^{\otimes 3 N_{2}}+\mathbb{1}^{\otimes 3 N_{1}} \otimes \hat{H}_{2} \tag{1.38}
\end{equation*}
$$

where $\hat{H}_{1}, \hat{H}_{2}$ were defined in (1.19) and $\mathbb{1}$ is the identity operator on $L^{2}(\mathbb{R})$.

- $f_{(U, V)}$ is the dipole interaction function

$$
\begin{align*}
& f_{(U, V)}\left(x_{1}, \ldots, x_{N}\right):=U D_{1}\left(x_{1}, \ldots, x_{N_{1}}\right) \cdot V D_{2}\left(x_{N_{1}+1}, \ldots, x_{N}\right) \\
&  \tag{1.39}\\
& -3\left(e_{1} \cdot U D_{1}\left(x_{1}, \ldots, x_{N_{1}}\right)\right)\left(e_{1} \cdot V D_{2}\left(x_{N_{1}+1}, \ldots, x_{N}\right)\right)
\end{align*}
$$

where

$$
\begin{aligned}
D_{1}\left(x_{1}, \ldots, x_{N_{1}}\right) & :=\sum_{j=1}^{N_{1}} x_{j}-\sum_{m=1}^{M_{1}} Z_{m} y_{m} \\
D_{2}\left(x_{N_{1}+1}, \ldots, x_{N}\right) & :=\sum_{j=N_{1}+1}^{N} x_{j}-\sum_{m=M_{1}+1}^{M} Z_{m} y_{m}
\end{aligned}
$$

are the instantaneous dipoles of the submolecules and

$$
\begin{aligned}
U D_{1}\left(x_{1}, \ldots, x_{N_{1}}\right) & :=D_{1}\left(U^{-1} x_{1}, \ldots, U^{-1} x_{N_{1}}\right), \\
V D_{2}\left(x_{N_{1}+1}, \ldots, x_{N}\right) & :=D_{2}\left(V^{-1} x_{N_{1}+1}, \ldots, V^{-1} x_{N}\right) .
\end{aligned}
$$

- Let

$$
\begin{equation*}
\mathcal{G}_{j}:=\operatorname{Ker}\left(\hat{H}_{j}-E_{j}\right), j=1,2, \tag{1.40}
\end{equation*}
$$

denote the ground state eigenspace of the $j$-th molecule, where $\hat{H}_{j}$ and $E_{j}$ have been defined in (1.19) and (1.20). Furthermore, we denote by $\Pi$ the orthogonal projection onto $\mathcal{G}_{1} \otimes \mathcal{G}_{2}$. Let also

$$
\begin{equation*}
H_{\infty}^{\perp}:=\Pi^{\perp} H_{\infty} \Pi^{\perp} \tag{1.41}
\end{equation*}
$$

where $\Pi^{\perp}:=1-\Pi$ is the orthogonal projection on $\left(\mathcal{G}_{1} \otimes \mathcal{G}_{2}\right)^{\perp}$.

- For $\phi \in \mathcal{G}_{1} \otimes \mathcal{G}_{2}$ with $\|\phi\|=1$, we define the van der Waals coefficient

$$
\begin{equation*}
C_{\mathrm{vdW}}(\phi, U, V):=\left\langle\Pi^{\perp} f_{(U, V)} \phi,\left(H_{\infty}^{\perp}-E_{\infty}\right)^{-1} \Pi^{\perp} f_{(U, V)} \phi\right\rangle \tag{1.42}
\end{equation*}
$$

with $E_{\infty}$ defined in (1.37), and

$$
\begin{equation*}
C_{\mathrm{vdW}}(U, V):=\max _{\phi \in \mathcal{G}_{1} \otimes \mathcal{G}_{2},\|\phi\|=1} C_{\mathrm{vdW}}(\phi, U, V) \tag{1.43}
\end{equation*}
$$

- For the multipole-multipole terms $\mathcal{F}^{(m, n)}\left(U \rho_{\Psi_{1}}, V \rho_{\Psi_{2}}\right)$, where $\Psi_{j} \in \mathcal{G}_{j}$, if there is independence of the choice of normalized ground states implied by (Poles), we will write $\mathcal{F}^{(m, n)}(U, V)$.
We remark that the functions $f_{(U, V)}$ and thus $C_{\mathrm{vdW}}(\phi, U, V)$ for each fixed $\phi$ depend continuously on $U, V$. Note that, if (Neutr) holds, then (1.20), (1.36) and (1.37) imply that

$$
\begin{equation*}
E_{\infty}=E_{1}+E_{2}=\min \sigma\left(H_{\infty}\right) \tag{1.44}
\end{equation*}
$$

Recalling that we have denoted $\tau=(L, U, V)$, we have the following results.
Theorem 1.10 (Multipolar/van der Waals expansion of the energy).
(a): Let $\Psi_{1}$ and $\Psi_{2}$ be any two normalized ground states of, respectively, $\hat{H}_{1}$ and $\hat{H}_{2}$. Then,

$$
\begin{equation*}
\mathcal{E}_{\tau} \leqslant E_{1}+E_{2}+\sum_{2 \leqslant n+m \leqslant 5} \frac{\mathcal{F}^{(n, m)}\left(U \rho_{\Psi_{1}}, V \rho_{\Psi_{2}}\right)}{L^{n+m+1}}-\frac{C_{\mathrm{vdW}}\left(\Psi_{1} \otimes \Psi_{2}, U, V\right)}{L^{6}}+O_{L^{\infty}}\left(\frac{1}{L^{7}}\right) . \tag{1.45}
\end{equation*}
$$

Furthermore, the function $C_{\mathrm{vdW}}\left(\Psi_{1} \otimes \Psi_{2}, U, V\right)$ is strictly positive for all $(U, V)$ in $\mathrm{SO}(3) \times \mathrm{SO}(3)$.
(b): If in addition Conditions (Neutr) and (Poles-all) hold, then

$$
\begin{equation*}
\mathcal{E}_{\tau}=E_{\infty}+\sum_{2 \leqslant n+m \leqslant 5} \frac{\mathcal{F}^{(n, m)}(U, V)}{L^{n+m+1}}-\frac{C_{\mathrm{vdW}}(U, V)}{L^{6}}+O_{L^{\infty}}\left(\frac{1}{L^{7}}\right) \tag{1.46}
\end{equation*}
$$

If, instead of Condition (Poles-all), we assume the weaker Condition (Poles-lead), we prove the following theorem providing a weaker form of equation (1.46), which intervenes in the proof of the main Theorem 1.13 below.
Theorem 1.11. Assume Conditions (Neutr) and (Poles-lead). If $n_{1}+n_{2}<5$, then

$$
\begin{equation*}
\mathcal{E}_{\tau}=E_{\infty}+\frac{\mathcal{F}^{\left(n_{1}, n_{2}\right)}(U, V)}{L^{n_{1}+n_{2}+1}}+O_{L^{\infty}}\left(\frac{1}{L^{n_{1}+n_{2}+2}}\right) . \tag{1.47}
\end{equation*}
$$

The fact that we assume only (Poles) and not irreducibility as in [6] for Theorems 1.10 (b) and 1.11 creates some additional difficulties in the proof, see Section 2.3 below for explanations.

Proof of Theorem 1.10 and 1.11. The proof of both the theorems can be found in Section 2. We derive an upper and lower bound in Subsections 2.1 and 2.2, respectively, and in Subsection 2.3 we collect the results to prove the theorems.


Figure 2. This is a hypothetical simplified graph of $\mathcal{E}=\mathcal{E}_{\tau}$ as an example in which the conclusion of Theorem 1.13 does not hold. The red paths here would be the first terms of a min-maxing sequence, and the distance of the submolecules along these sequence of paths is not uniformly bounded.
1.3.2. The main theorem. The most natural result we would like to have would be the existence of a bounded optimal path. However, we have not been able to obtain it. Nevertheless, we prove the existence of a bounded min-maxing sequence, as defined below.

Definition 1.12. A min-maxing sequence is a sequence of paths

$$
\left\{\tau_{n}(\cdot)\right\} \subset C^{0}\left([0,1] ; \mathbb{R}^{+} \times \mathrm{SO}(3) \times \mathrm{SO}(3)\right)
$$

such that

$$
\begin{equation*}
\lim _{n \rightarrow \infty} \max _{t \in[0,1]} \mathcal{E}_{\tau_{n}(t)}=c \tag{1.48}
\end{equation*}
$$

where $c$ was defined in (1.27).
Now we are ready to state our main theorem. Recall that $n_{1}, n_{2}$ indicate the first nonvanishing multipole moments, defined in Definition 1.8, and in general can depend on the ground state even if this is normalized. Under Condition (Poles), the leading multipole moments $\mathcal{M}_{\rho_{\Psi_{1}}}^{\left(n_{1}\right)}, \mathcal{M}_{\rho_{\Psi_{2}}}^{\left(n_{2}\right)}$ of the molecules are independent of the choice of the normalized ground state and we will write instead $\mathcal{M}_{1}^{\left(n_{1}\right)}, \mathcal{M}_{2}^{\left(n_{2}\right)}$.

Theorem 1.13. Let us assume Condition (Neutr). Moreover, let us assume that Condition (Poleslead) holds and $n_{1}+n_{2}<5$. In the case $n_{k}=3$ for $k=1$ or $k=2$, we also assume the implication

$$
\begin{equation*}
\mathcal{M}_{k}^{\left(n_{k}\right)}(v, \cdot, \cdot) \equiv 0 \Longrightarrow v=0 \tag{1.49}
\end{equation*}
$$

Then, one can find a bounded min-maxing sequence of paths for the min-max problem (1.27): there exists $L_{\text {cut }}>0$ and a min-maxing sequence of paths $\tau_{n}(t)=\left(L_{n}(t), U_{n}(t), V_{n}(t)\right)$ such that

$$
L_{n}(t) \leqslant L_{\mathrm{cut}}
$$

for all $n$ and all $t \in[0,1]$.

Remark 1.14. Theorem 1.13, aside from being an important step towards proving existence of an optimal path, also has a meaning on its own. To see this we refer to Figure 2 showing a hypothetical simplified graph of $\mathcal{E}=\mathcal{E}_{\tau}$ as an example in which a min-maxing sequence has to be unbounded. In this case, optimization of the energy would happen through infinite separation of the submolecules, a situation that would be unphysical. Moreover, a transition state would not exist. We recall, however, that Theorem 1.13 together with Theorem 4 in [42] implies existence of a transition state in the (NR) case. As we mentioned above, we conjecture that Theorem 4 in [42] holds in the (SR) case as well, and, if this is the case, then existence of a transition state follows similarly.

Remark 1.15. We believe that the non-degeneracy assumption (1.49) on the octopole, in case this is the first non-vanishing multipole moment, is physical. We refer the interested reader to [6] Lemma 1.9 and the related discussion.

We have not been able to obtain a result if $n_{1}+n_{2}=5$, i.e. when the leading multipole-multipole interaction has the same order as the van der Waals attraction. In the case where $n_{1}+n_{2} \geqslant 6$, as soon as we have the upper bound (1.45), the following result of [6] can be adapted without any modification to the (SR) case, thus we will not give a proof of it here.

Theorem 1.16 ([6], Theorem 1.7). Let us assume Condition (Neutr) and that $\hat{H}_{1}, \hat{H}_{2}$ have ground states so that, for the leading multipole terms, we have $n_{1}+n_{2} \geqslant 6$.

Then, one can find a bounded min-maxing sequence of paths $\tau_{n}(t)$ for the min-max problem (1.27).
1.4. Sketch of the proof of Theorem 1.13. In this section, we give the proof of Theorem 1.13, but defer the proof of many intermediate results to the rest of the paper. The general strategy is the same as in [6]. In this sketch, we will also briefly explain additional difficulties coming in the (SR) case.

We will only consider the case where

$$
\begin{equation*}
E_{\infty}=c, \tag{1.50}
\end{equation*}
$$

with $c, E_{\infty}$ defined in (1.27), (1.37), respectively. If (1.50) is not fulfilled, then proving boundedness is easy and already known from previous works, see e.g. [42, Theorem 4]. Roughly speaking, the reason is that if $E_{\infty} \neq c$, then the energy of the system does not go to the mountain pass level $c$ when the molecules go infinitely apart. Therefore, the mountain pass level $c$ should be attained at a finite distance.

We want to prove that there exists a min-maxing sequence $\left(\tau_{n}\right)$ for the mountain pass problem (1.27), such that, for all $n$ and $t$,

$$
\begin{equation*}
L_{n}(t) \leqslant L_{\mathrm{cut}} \tag{1.51}
\end{equation*}
$$

for an appropriate $L_{\text {cut }}>0$. Definition 1.12 and (1.27) imply that there exists a min-maxing sequence $\left(\widetilde{\tau}_{n}(t)\right)_{n \in \mathbb{N}}=\left(\widetilde{L}_{n}(t), \widetilde{U}_{n}(t), \widetilde{V}_{n}(t)\right)_{n \in \mathbb{N}}$, for which $\widetilde{L}_{n}(t)$ might be unbounded. We will modify it to construct a min-maxing sequence $\left(\tau_{n}(t)\right)_{n \in \mathbb{N}}=\left(L_{n}(t), U_{n}(t), V_{n}(t)\right)_{n \in \mathbb{N}}$ such that (1.51) is fulfilled. We consider $L_{\text {cut }}$ as a parameter and, for a given $n$, we start by looking at the first time $t_{0}$ and the last time $t_{1}$ for which $\widetilde{L}_{n}(t)=L_{\text {cut }}$. These times depend on $L_{\text {cut }}$ and $n$, but we do not emphasize this in our notation.

We will construct a path $\tau_{n}(\cdot)$ which coincides with $\widetilde{\tau}_{n}(\cdot)$ outside $\left(t_{0}, t_{1}\right)$ and for which $L_{n}(t) \equiv L_{\text {cut }}$ for all $t \in\left[t_{0}, t_{1}\right]$. The difficulty is to link the two molecular orientations $\left(\widetilde{U}_{n}\left(t_{0}\right), \widetilde{V}_{n}\left(t_{0}\right)\right)$ and $\left(\widetilde{U}_{n}\left(t_{1}\right), \widetilde{V}_{n}\left(t_{1}\right)\right)$ by a continuous path in $\mathrm{SO}(3) \times \mathrm{SO}(3)$ on which the energy does not increase too much. More precisely, we want to find, for $L_{\text {cut }}$ large enough, a continuous path $(U(t), V(t))_{t \in\left[t_{0}, t_{1}\right]}$ of rotations such that

$$
\left(U\left(t_{0}\right), V\left(t_{0}\right)\right)=\left(\widetilde{U}_{n}\left(t_{0}\right), \widetilde{V}_{n}\left(t_{0}\right)\right), \quad\left(U\left(t_{1}\right), V\left(t_{1}\right)\right)=\left(\widetilde{U}_{n}\left(t_{1}\right), \widetilde{V}_{n}\left(t_{1}\right)\right)
$$

and

$$
\begin{equation*}
\max _{t \in\left[t_{0}, t_{1}\right]} \mathcal{E}_{\left(L_{\mathrm{cut}}, U(t), V(t)\right)} \leqslant \max \left\{\mathcal{E}_{\left(L_{\mathrm{cut}}, \widetilde{U}_{n}\left(t_{0}\right), \widetilde{V}_{n}\left(t_{0}\right)\right)}, \mathcal{E}_{\left(L_{\mathrm{cut}}, \widetilde{U}_{n}\left(t_{1}\right), \widetilde{V}_{n}\left(t_{1}\right)\right)}, E_{\infty}\right\} . \tag{1.52}
\end{equation*}
$$

An arbitrary path would not work here, since there are orientations for which the molecules repel each other. As a consequence, the maximal energy along the path could be bigger than $E_{\infty}$ and thus (1.52) would not hold. This is the case when $\mathcal{F}^{\left(n_{1}, n_{2}\right)}(U, V)>0$, see (1.47). If (1.52) holds, then, due to (1.50) and (1.48), $\tau_{n}$ is going to be a min-maxing sequence fulfilling (1.51), as desired.

We proceed as follows:
Step 1. Replacing the two points by local (pseudo-)minima. Our first step is to rotate the molecules at each point $t_{0}, t_{1}$, without changing the distance $L_{\text {cut }}$, so that we reach a kind of local minimum of the energy of the system with respect to rotations. In general, this is impossible if, in all directions, the energy takes values higher and lower than that of the current point, at arbitrarily small distances. Since we cannot exclude this pathological behavior, we introduce, as in [6], the weaker notion of local pseudo-minimum.

Definition 1.17 (Local pseudo-minimum). Let $\mathcal{W}$ be any continuous function on $\mathrm{SO}(3) \times \mathrm{SO}(3)$. A point $(\bar{U}, \bar{V}) \in \mathrm{SO}(3) \times \mathrm{SO}(3)$ is called a local pseudo-minimum of $\mathcal{W}$ if, for any $A, B \in \Gamma$, as defined in (1.29), there exists a sequence $t_{n} \rightarrow 0+$ with $t_{n} \neq 0$ such that $\mathcal{W}\left(e^{t_{n} A} \bar{U}, e^{t_{n} B} \bar{V}\right) \geqslant \mathcal{W}(\bar{U}, \bar{V})$.

In other words, a local pseudo-minimum is a point at which one cannot find a direction in $\mathrm{SO}(3)^{2}$ in which $\mathcal{W}$ strictly decreases. If $\mathcal{W}$ is $C^{2}$ and has a local pseudo-minimum at $(\bar{U}, \bar{V})$, one must then have

$$
\begin{equation*}
\left.\frac{\mathrm{d}}{\mathrm{~d} t} \mathcal{W}\left(e^{t A} \bar{U}, e^{t B} \bar{V}\right)\right|_{t=0}=0,\left.\quad \frac{\mathrm{~d}^{2}}{\mathrm{~d} t^{2}} \mathcal{W}\left(e^{t A} \bar{U}, e^{t B} \bar{V}\right)\right|_{t=0} \geqslant 0 \tag{1.53}
\end{equation*}
$$

for all $A, B \in \Gamma$, but $(\bar{U}, \bar{V})$ does not need to be a true local minimum. A counterexample can be found even for functions of one variable: if

$$
f(x)=x^{4} \sin (1 / x) \chi_{\{x<0\}}(x)+x^{4} \chi_{\{x>0\}}(x), \quad x \in \mathbb{R},
$$

where $\chi_{C}$ denotes the characteristic function of the set $C$, then $f$ has a local pseudo-minimum at zero but not a local minimum at 0 .

The fact that we can reach a local pseudo-minimum can be shown with the help of the following lemma, already proven in [6]. We state it for simplicity on $\mathrm{SO}(3) \times \mathrm{SO}(3)$ but it is much more general.

Lemma 1.18 ([6], Lemma 1.10: linking any point to a local pseudo-minimum). Let $\mathcal{W}$ be a continuous function on $\mathrm{SO}(3) \times \mathrm{SO}(3)$. Let $(U, V)$ be a point that is not a local pseudo-minimum of $\mathcal{W}$. Then there exists a continuous path on $\mathrm{SO}(3) \times \mathrm{SO}(3)$ linking $(U, V)$ to a local pseudo-minimum $(\bar{U}, \bar{V})$, such that the maximum value of $\mathcal{W}$ on this path is $\mathcal{W}(U, V)$.

We will apply these results to the function $\mathcal{E}_{L_{\text {cut }}}$, where, for $L>2 \max \left\{\left|y_{j}\right|: j \in\{1, \ldots, M\}\right\}$,

$$
\begin{equation*}
\mathcal{E}_{L}(U, V):=\mathcal{E}_{(L, U, V)} . \tag{1.54}
\end{equation*}
$$

Note that the last restriction on $L$ makes sure that no nuclei coincide. Using Lemma 1.18, we can connect the point $\left(\widetilde{U}_{n}\left(t_{0}\right), \widetilde{V}_{n}\left(t_{0}\right)\right)$ to a local pseudo-minimum $\left(\bar{U}_{0}, \bar{V}_{0}\right)$ of $\mathcal{E}_{L_{\text {cut }}}$, with a path along which the energy stays below the value $\mathcal{E}_{\left(L_{\text {cut }}, \widetilde{U}_{n}\left(t_{0}\right), \widetilde{V}_{n}\left(t_{0}\right)\right)}$. We proceed analogously for $t_{1}$ and connect $\left(\widetilde{U}_{n}\left(t_{1}\right), \widetilde{V}_{n}\left(t_{1}\right)\right)$ to a local pseudo-minimum $\left(\bar{U}_{1}, \bar{V}_{1}\right)$ of $\mathcal{E}_{L_{\text {cut }}}$, with a path along which the energy stays below $\mathcal{E}_{\left(L_{\mathrm{cut}}, \widetilde{U}_{n}\left(t_{1}\right), \tilde{V}_{n}\left(t_{1}\right)\right)}$. In the next steps, we will study the properties of these two points and finally connect them by a path, without increasing the energy too much.

Step 2. Properties of the leading multipolar interaction at the local pseudominima ( $\bar{U}_{0}, \bar{V}_{0}$ ) and $\left(\bar{U}_{1}, \bar{V}_{1}\right)$. In Step 1, we have managed to reach two points ( $\left.\bar{U}_{0}, \bar{V}_{0}\right)$ and ( $\bar{U}_{1}, \bar{V}_{1}$ ) which are local pseudo-minima of the energy $\mathcal{E}_{L_{\text {cut }}}$. In this step we prove that

$$
\begin{equation*}
\mathcal{E}_{\left(L_{\mathrm{cut}}, \bar{U}_{0}, \bar{V}_{0}\right)}<E_{\infty}, \quad \mathcal{E}_{\left(L_{\mathrm{cut}}, \bar{U}_{1}, \bar{V}_{1}\right)}<E_{\infty} \tag{1.55}
\end{equation*}
$$

independently of $n$, provided we had chosen $L_{\text {cut }}$ large enough. This means that the configurations $\left(L_{\mathrm{cut}}, \bar{U}_{0}, \bar{V}_{0}\right),\left(L_{\mathrm{cut}}, \bar{U}_{1}, \bar{V}_{1}\right)$ have an energy strictly below the dissociation threshold $c=E_{\infty}$ - recall (1.27) and (1.50). More precisely, we show that $\exists \delta>0$ such that

$$
\begin{equation*}
\mathcal{F}^{\left(n_{1}, n_{2}\right)}\left(\bar{U}_{0}, \bar{V}_{0}\right) \leqslant-\delta, \quad \mathcal{F}^{\left(n_{1}, n_{2}\right)}\left(\bar{U}_{1}, \bar{V}_{1}\right) \leqslant-\delta, \tag{1.56}
\end{equation*}
$$

where we recall that $n_{1}$ and $n_{2}$ are the indices for the first non-vanishing multipoles of the two neutral submolecules, defined in Definition 1.8 and the $\mathcal{F}^{\left(n_{1}, n_{2}\right)}$ are the multipolar interactions defined by (1.35). The estimate (1.55) follows from (1.56) and the expansion in Theorem 1.11. We now explain how to prove (1.56). We first need the following proposition.

Proposition 1.19 (Leading multipolar energy at a local pseudo-minimum). Assume that Condition (Poles-lead) holds. Then, for any $\delta>0$, there exists an $L_{m}>0$ such that, if $L>L_{m}$ and $(\bar{U}, \bar{V})$ is any local pseudo-minimum of $\mathcal{E}_{L}$, defined in (1.54), then we have

$$
\begin{equation*}
\left|\frac{\mathrm{d}}{\mathrm{~d} t} \mathcal{F}^{\left(n_{1}, n_{2}\right)}\left(e^{t A} \bar{U}, e^{t B} \bar{V}\right)\right|_{t=0}\left|\leqslant \delta, \quad \frac{\mathrm{~d}^{2}}{\mathrm{~d} t^{2}} \mathcal{F}^{\left(n_{1}, n_{2}\right)}\left(e^{t A} \bar{U}, e^{t B} \bar{V}\right)\right|_{t=0} \geqslant-\delta \tag{1.57}
\end{equation*}
$$

for all $A, B \in \Gamma$, recall (1.29).
Proof. The proof can be found in Section 3.
Intuitively, Proposition 1.19 says that a local minimum of $\mathcal{E}_{L}$ is close to being a local minimum of the leading multipole-multipole term $\mathcal{F}^{\left(n_{1}, n_{2}\right)}$. This was already proven in [6] in the (NR) case under an irreducibility assumption on the ground state eigenspaces $\mathcal{G}_{1}, \mathcal{G}_{2}$. Nevertheless, it has not been possible for us to adapt these arguments to the (SR) case. Indeed, in the (NR) case, the eigenstates are in $H^{2}$. In the (SR) case, it is not even clear that they are in $H^{1}$. As usual, the nonlocality of the kinetic energy operator makes things more complicated too. We also emphasize the fact that, for the results of [14], it was not necessary to differentiate the energy and therefore their methods are not applicable for the proof of Proposition 1.19. We overcome these difficulties by using ideas in [37], that Hunziker used to prove that the ground state energy of $H_{N}(Y, \mathcal{Z})$ depends analytically on $Y$ in the (NR) case and if the fermionic statistics is ignored. We rework these ideas, so that they provide quantitative information for


Figure 3. Graphical representation of the proof of Theorem 1.13 in Case 1 (the picture was taken from [6] and was created by Mathieu Lewin).
derivatives and are applicable in the (SR) case as well. Our proof of Proposition 1.19 in the (NR) case is simpler than the one in [6] and provides information for higher order derivatives as well.

The rest of the proof is the same as in [6], as Propositions 1.20 and 1.21 below were proven in [6, Section 4] independently of the kind of the kinetic energy and of the irreducibility assumption. Proving Propositions 1.20 and 1.21 was the hardest part of the analysis of [6] but, since there are no modifications needed in our setting, we will only state them and explain how we can conclude the proof with them. The first result is rather intuitive, since the average of $\mathcal{F}^{\left(n_{1}, n_{2}\right)}$ over all orientations is zero.

Proposition 1.20 ([6], Proposition 1.12). Let $m, n \in \mathbb{N}$ with $n+m \in\{2,3,4\}$. If $n$ or $m$ is 3 , assume that the condition on the octopole moment stated in Theorem 1.13 holds. Then there exists a $\delta>0$ such that, if $\mathcal{F}^{(n, m)}$ fulfills (1.57), then

$$
\mathcal{F}^{(n, m)}(\bar{U}, \bar{V}) \leqslant-\delta
$$

Using Propositions 1.19 and 1.20 , we obtain that, if $L_{\text {cut }}$ is large enough, there exists $\delta>0$ such that (1.56) holds and, in view of (1.47),

$$
\mathcal{E}_{\left(L_{\mathrm{cut}}, \bar{U}_{0}, \bar{V}_{0}\right)} \leqslant E_{\infty}-\frac{\delta}{2}, \quad \mathcal{E}_{\left(L_{\mathrm{cut}}, \bar{U}_{1}, \bar{V}_{1}\right)} \leqslant E_{\infty}-\frac{\delta}{2},
$$

provided that $L_{\text {cut }}$ is large enough.
Step 3. Linking the local pseudominima $\left(\bar{U}_{0}, \bar{V}_{0}\right)$ and ( $\left.\bar{U}_{1}, \bar{V}_{1}\right)$ with the path of low energy. We have so far managed to show that the leading multipole-multipole interaction of the two new points $\left(\bar{U}_{0}, \bar{V}_{0}\right)$ and $\left(\bar{U}_{1}, \bar{V}_{1}\right)$ is negative, away from zero. We connect $\left(\bar{U}_{0}, \bar{V}_{0}\right)$ and $\left(\bar{U}_{1}, \bar{V}_{1}\right)$ with a path along which $\mathcal{F}^{\left(n_{1}, n_{2}\right)}$ stays negative, which implies, if $L_{\text {cut }}$ is large enough, that $\mathcal{E}_{\left(L_{\text {cut }}, U(t), V(t)\right)}<E_{\infty}$ for all $t$, by Theorem 1.11. Note that, for this part, the non-degeneracy assumption (1.49) for the octopole moment is not needed. The existence of such a path is given in the following proposition.
Proposition 1.21 ([6] Proposition 1.13, Connectedness of $\left.\left\{\mathcal{F}^{(n, m)} \leqslant-\delta\right\}\right)$. Let $n, m \in \mathbb{N}$ with $n+m \in$ $\{2,3,4\}$. Then there exists $\delta_{0}>0$ such that, for all $0<\delta<\delta_{0}$, the set $\left\{(U, V) \in \mathrm{SO}(3): \mathcal{F}^{(n, m)}(U, V)<\right.$ $-\delta\}$ is nonempty and pathwise connected.

Thus, we can connect $\left(\bar{U}_{0}, \bar{V}_{0}\right)$ and $\left(\bar{U}_{1}, \bar{V}_{1}\right)$ with a path along which the leading term of the interaction energy is negative. If $L_{\text {cut }}$ is large enough, then the difference $\mathcal{E}_{\left(L_{\text {cut }}, U(t), V(t)\right)}-E_{\infty}$ remains negative along the path connecting $\left(\bar{U}_{0}, \bar{V}_{0}\right)$ and $\left(\bar{U}_{1}, \bar{V}_{1}\right)$. We have thus found a new sequence $\left\{\tau_{n}\right\}_{n \in \mathbb{N}}, \tau_{n}=$ ( $L_{n}, U_{n}, V_{n}$ ) where $L_{n} \leqslant L_{\text {cut }}$ for any $n \in \mathbb{N}$ and for which

$$
\begin{equation*}
\max _{t \in[0,1]} \mathcal{E}_{\tau_{n}(t)} \leqslant \max \left(E_{\infty}, \max _{t \in[0,1]} \mathcal{E}_{\widetilde{\tau}_{n}(t)}\right) \tag{1.58}
\end{equation*}
$$

But, since we have assumed that $\left(\widetilde{\tau}_{n}\right)_{n \in \mathbb{N}}$ is a min-maxing sequence, we know that

$$
\lim _{n \rightarrow \infty} \max _{t \in[0,1]} \mathcal{E}_{\widetilde{\tau}_{n}(t)}=c .
$$

We have assumed too, in (1.50), that $E_{\infty}=c$. We arrive thus from (1.58) at

$$
\lim \sup \max _{t \in[0,1]} \mathcal{E}_{\tau_{n}(t)} \leqslant c
$$

which implies that $\left(\tau_{n}\right)_{n \in \mathbb{N}}$ is a min-maxing sequence as well.

## 2. Proof of Theorems 1.10 and 1.11

2.1. Upper bound: Proof of Theorem 1.10(a). In this section, we prove the first part of Theorem 1.10 in the (SR) case. The analogous result for the (NR) case has already been proved in Section 2.2 of [6]. We estimate the interaction energy $\mathcal{E}_{\tau}-E_{1}-E_{2}$ from above by the associated quadratic form calculated on a suitable trial function.

Let us introduce

$$
H_{\infty, \tau}:=H_{1, \tau} \otimes \mathbb{1}^{\otimes 3 N_{2}}+\mathbb{1}^{\otimes 3 N_{1}} \otimes H_{2, \tau}
$$

with $H_{j, \tau}$, defined in (1.24), (1.25): $H_{\infty, \tau}$ is the Hamiltonian describing the 2 submolecules without interaction between each other. It will often be useful to work in coordinates in which $H_{\infty, \tau}$ does not depend on $\tau$. For that reason, we introduce the unitary transformation $K_{\tau}: L^{2}\left(\mathbb{R}^{3 N}\right) \rightarrow L^{2}\left(\mathbb{R}^{3 N}\right)$ defined by

$$
\begin{align*}
& \left(K_{\tau} \Psi\right)\left(x_{1}, \ldots, x_{N_{1}}, x_{N_{1}+1}, \ldots, x_{N}\right) \\
& \quad:=\Psi\left(U^{-1} x_{1}, \ldots, U^{-1} x_{N_{1}}, V^{-1}\left(x_{N_{1}+1}-L e_{1}\right), \ldots, V^{-1}\left(x_{N}-L e_{1}\right)\right) \tag{2.1}
\end{align*}
$$

In words, $K_{\tau}$ applies to the first $N_{1}$ electrons the transformation applied by $\tau$ to the nuclei of the first molecule and to the last $N_{2}$ electrons the transformation applied by $\tau$ to the nuclei of the second molecule. We observe that

$$
\begin{equation*}
H_{\infty}=K_{\tau}^{-1} H_{\infty, \tau} K_{\tau} \tag{2.2}
\end{equation*}
$$

where we recall that $H_{\infty}$ was defined in (1.38). In light of (2.2) and (1.23), it makes sense to introduce

$$
\begin{equation*}
\tilde{I}_{\tau}:=K_{\tau}^{-1} I_{\tau} K_{\tau} \tag{2.3}
\end{equation*}
$$

which depends on $\tau$ even though the left-hand side of (2.2) does not. Explicitly, we have

$$
\begin{align*}
\tilde{I}_{\tau}=-\sum_{i=1}^{N_{1}} \sum_{l=M_{1}+1}^{M} \frac{Z_{l} \alpha}{\left|U x_{i}-V y_{l}-L e_{1}\right|} & -\sum_{j=N_{1}+1}^{N} \sum_{k=1}^{M_{1}} \frac{Z_{k} \alpha}{\left|U y_{k}-V x_{j}-L e_{1}\right|} \\
& +\sum_{k=1}^{M_{1}} \sum_{l=M_{1}+1}^{M} \frac{Z_{l} Z_{k} \alpha}{\left|U y_{k}-V y_{l}-L e_{1}\right|}+\sum_{i=1}^{N_{1}} \sum_{j=N_{1}+1}^{N} \frac{\alpha}{\left|U x_{i}-V x_{j}-L e_{1}\right|} . \tag{2.4}
\end{align*}
$$

We also define a localizing function on $\mathbb{R}^{3}$ by

$$
\chi_{L}(x):=\chi\left(\frac{|x|}{L}\right), \quad \text { where } \chi(r)= \begin{cases}1, & \text { if } r \leqslant 1 / 10  \tag{2.5}\\ 0, & \text { if } r \geqslant 1 / 8\end{cases}
$$

and $\chi$ is a smoothed out characteristic function: $\chi \in C^{\infty}(\mathbb{R})$, $\operatorname{Ran} \chi=[0,1]$. We define on $L^{2}\left(\mathbb{R}^{3 N}\right)$ the operators

$$
\begin{equation*}
\xi_{\tau}:=\left(H_{\infty}^{\perp}-E_{1}-E_{2}\right)^{-1} \Pi^{\perp} \tilde{I}_{\tau} \chi_{L}^{\otimes N} \tag{2.6}
\end{equation*}
$$

and

$$
\begin{equation*}
W_{\tau}:=\chi_{L}^{\otimes N}-\chi_{4 L / 3}^{\otimes N} \xi_{\tau}=\left(1-\chi_{4 L / 3}^{\otimes N}\left(H_{\infty}^{\perp}-E_{1}-E_{2}\right)^{-1} \Pi^{\perp} \tilde{I}_{\tau}\right) \chi_{L}^{\otimes N} \tag{2.7}
\end{equation*}
$$

where $H_{\infty}^{\perp}, \Pi^{\perp}$ were defined in (1.41) and below (1.41), respectively.
We consider a normalized function $\Psi \in \mathcal{G}_{1} \otimes \mathcal{G}_{2}$, where we recall that $\mathcal{G}_{1}, \mathcal{G}_{2}$ were defined in (1.40). For this section, the interesting case is $\Psi=\Psi_{1} \otimes \Psi_{2}$, where $\Psi_{1}$ and $\Psi_{2}$ are in $\mathcal{G}_{1}$ and $\mathcal{G}_{2}$, respectively. However, as most of our calculations are useful for the next sections for general normalized $\Psi \in \mathcal{G}_{1} \otimes \mathcal{G}_{2}$, we will work in this section for general $\Psi$ as well, and we will specify the parts of the calculations for which the factorized form $\Psi=\Psi_{1} \otimes \Psi_{2}$ is necessary. As we will prove below, in (2.26), for $\Psi \in \mathcal{G}_{1} \otimes \mathcal{G}_{2}$, $W_{\tau} \Psi$ is a small perturbation of $\Psi$. More precisely, it is an approximation of the ground state given by the Feshbach map, see the brief sketch of the proof of Theorem 1.5 in [4]. The cut-off functions make some calculations easier.

Since ground states of $\hat{H}_{1}$ and $\hat{H}_{2}$ have at least $H^{1 / 2}$ regularity, $\Psi \in \mathcal{G}_{1} \otimes \mathcal{G}_{2}$ has $H^{1 / 2}$ regularity as well and, by [46, Theorem 7.18], this is the case for $\chi_{L}^{\otimes N} \Psi$ too. Since the function $\tilde{I}_{\tau}$ as well as all its derivatives are bounded on the support of $\chi_{L}^{\otimes N} \Psi$, we see by (2.7) that $W_{\tau} \Psi$ is the sum of 2 terms, each of them being in $H^{1 / 2}\left(\mathbb{R}^{3 N}\right)$. Since $K_{\tau}$, defined in (2.1), is unitary in $H^{1 / 2}\left(\mathbb{R}^{3 N}\right), K_{\tau} W_{\tau} \Psi$ is in $H^{1 / 2}\left(\mathbb{R}^{3 N}\right)$ as well, so it is in the form domain of $H_{\tau}$, see (1.13). We consider the following test function:

$$
\frac{Q_{N} K_{\tau} W_{\tau} \Psi}{\left\|Q_{N} K_{\tau} W_{\tau} \Psi\right\|}
$$

where $Q_{N}$ is defined by (1.10). A similar test function was used in [4] and [6] with different notation. By definition of the ground state energy, see (1.22)-(1.23), we have that

$$
\begin{equation*}
\mathcal{E}_{\tau}-E_{1}-E_{2} \leqslant \frac{\left\langle Q_{N} K_{\tau} W_{\tau} \Psi,\left(H_{\tau}-E_{1}-E_{2}\right) Q_{N} K_{\tau} W_{\tau} \Psi\right\rangle}{\left\|Q_{N} K_{\tau} W_{\tau} \Psi\right\|^{2}} \tag{2.8}
\end{equation*}
$$

Note that $Q_{N}$ does not commute with $K_{\tau}$.
We first get rid of the antisymmetrizer $Q_{N}$ in the right-hand side of (2.8). How to do this is not as obvious as in the (NR) case and, for this reason, we will explain it in detail. First, since $Q_{N}$ is a projection and it commutes with $H_{\tau}-E_{1}-E_{2}$,

$$
\left\langle Q_{N} K_{\tau} W_{\tau} \Psi,\left(H_{\tau}-E_{1}-E_{2}\right) Q_{N} K_{\tau} W_{\tau} \Psi\right\rangle=\left\langle Q_{N} K_{\tau} W_{\tau} \Psi,\left(H_{\tau}-E_{1}-E_{2}\right) K_{\tau} W_{\tau} \Psi\right\rangle
$$

In view of (1.10), we can write

$$
\begin{equation*}
\left\langle Q_{N} K_{\tau} W_{\tau} \Psi,\left(H_{\tau}-E_{1}-E_{2}\right) K_{\tau} W_{\tau} \Psi\right\rangle=\sum_{\sigma \in S_{N}} \frac{(-1)^{\sigma}}{N!}\left\langle\sigma \cdot K_{\tau} W_{\tau} \Psi,\left(H_{\tau}-E_{1}-E_{2}\right) K_{\tau} W_{\tau} \Psi\right\rangle \tag{2.9}
\end{equation*}
$$

where we recall that $\sigma \cdot K_{\tau} W_{\tau} \Psi$ has been defined in (1.9).
We now have to distinguish between two cases: when $\sigma\left(\left\{1, \ldots, N_{1}\right\}\right)=\left\{1, \ldots, N_{1}\right\}$, i.e. when the permutation $\sigma$ leaves invariant the set of electrons associated with each of the submolecules, and when $\sigma\left(\left\{1, \ldots, N_{1}\right\}\right) \neq\left\{1, \ldots, N_{1}\right\}$.

In the first case, we have by construction, for such a $\sigma$, that $\sigma \cdot K_{\tau} W_{\tau} \Psi=(-1)^{\sigma} K_{\tau} W_{\tau} \Psi$. Indeed, the ground states $\Psi_{1}$ and $\Psi_{2}$ are antisymmetric and $H_{\infty}^{\perp}, K_{\tau}$ and $W_{\tau}$ are invariant under these permutations of variables. As a consequence

$$
\begin{align*}
& \sum_{\sigma\left(\left\{1, \ldots, N_{1}\right\}\right)=\left\{1, \ldots, N_{1}\right\}}(-1)^{\sigma}\left\langle\sigma \cdot K_{\tau} W_{\tau} \Psi,\left(H_{\tau}-E_{1}-E_{2}\right) K_{\tau} W_{\tau} \Psi\right\rangle \\
&=N_{1}!N_{2}!\left\langle K_{\tau} W_{\tau} \Psi,\left(H_{\tau}-E_{1}-E_{2}\right) K_{\tau} W_{\tau} \Psi\right\rangle \tag{2.10}
\end{align*}
$$

On the other hand, if $\sigma\left(\left\{1, \ldots, N_{1}\right\}\right) \neq\left\{1, \ldots, N_{1}\right\}$, then, due to the cut-offs $\chi_{L}^{\otimes N}, \chi_{4 L / 3}^{\otimes N}$ and the translation in the definition of $K_{\tau}$ in (2.1), the functions $K_{\tau} W_{\tau} \Psi$ and $\sigma \cdot K_{\tau} W_{\tau} \Psi$ have disjoint supports. Observing additionally that the only non-local terms in $H_{\tau}-E_{\infty}$ are the kinetic energy operators, we find that

$$
\begin{equation*}
\left\langle\sigma \cdot K_{\tau} W_{\tau} \Psi,\left(H_{\tau}-E_{1}-E_{2}\right) K_{\tau} W_{\tau} \Psi\right\rangle=\sum_{j=1}^{N}\left\langle\sigma \cdot K_{\tau} W_{\tau} \Psi, T_{j} K_{\tau} W_{\tau} \Psi\right\rangle \tag{2.11}
\end{equation*}
$$

But, for each $\sigma$ and $j$, we can find $k \neq j$ such that, either $k \in\left\{1, \ldots, N_{1}\right\}$ and $\sigma(k) \in\left\{N_{1}+1, \ldots, N\right\}$, or $k \in\left\{N_{1}+1, \ldots, N\right\}$ and $\sigma(k) \in\left\{1, \ldots, N_{1}\right\}$. Thus, $K_{\tau} W_{\tau} \Psi$ and $\sigma \cdot K_{\tau} W_{\tau} \Psi$ have disjoint supports in the $k$-th variable and, as a consequence,

$$
\begin{equation*}
\left\langle\sigma \cdot K_{\tau} W_{\tau} \Psi, T_{j} K_{\tau} W_{\tau} \Psi\right\rangle=0, \quad \forall j \in\{1, \ldots, N\} \tag{2.12}
\end{equation*}
$$

From (2.9), (2.10), (2.11) and (2.12), we find that

$$
\begin{equation*}
\left\langle Q_{N} K_{\tau} W_{\tau} \Psi,\left(H_{\tau}-E_{1}-E_{2}\right) Q_{N} K_{\tau} W_{\tau} \Psi\right\rangle=\frac{\left\langle K_{\tau} W_{\tau} \Psi,\left(H_{\tau}-E_{1}-E_{2}\right) K_{\tau} W_{\tau} \Psi\right\rangle}{\binom{N}{N_{1}}} \tag{2.13}
\end{equation*}
$$

A similar but easier calculation gives

$$
\begin{equation*}
\left\|Q_{N} K_{\tau} W_{\tau} \Psi\right\|^{2}=\frac{\left\|K_{\tau} W_{\tau} \Psi\right\|^{2}}{\binom{N}{N_{1}}} \tag{2.14}
\end{equation*}
$$

With (2.13) and (2.14), we find

$$
\begin{equation*}
\frac{\left\langle Q_{N} K_{\tau} W_{\tau} \Psi,\left(H_{\tau}-E_{1}-E_{2}\right) Q_{N} K_{\tau} W_{\tau} \Psi\right\rangle}{\left\|Q_{N} K_{\tau} W_{\tau} \Psi\right\|^{2}}=\frac{\left\langle K_{\tau} W_{\tau} \Psi,\left(H_{\tau}-E_{1}-E_{2}\right) K_{\tau} W_{\tau} \Psi\right\rangle}{\left\|K_{\tau} W_{\tau} \Psi\right\|^{2}} \tag{2.15}
\end{equation*}
$$

which, using the fact that $K_{\tau}$ is unitary together with (2.2), (2.3) and (1.23), gives

$$
\begin{equation*}
\frac{\left\langle Q_{N} K_{\tau} W_{\tau} \Psi,\left(H_{\tau}-E_{1}-E_{2}\right) Q_{N} K_{\tau} W_{\tau} \Psi\right\rangle}{\left\|Q_{N} K_{\tau} W_{\tau} \Psi\right\|^{2}}=\frac{\left\langle W_{\tau} \Psi,\left(H_{\infty}+\tilde{I}_{\tau}-E_{1}-E_{2}\right) W_{\tau} \Psi\right\rangle}{\left\|W_{\tau} \Psi\right\|^{2}} \tag{2.16}
\end{equation*}
$$

In the following, we compute the right-hand side of (2.16). Moreover, we prove that it depends smoothly on $U$ and $V$, which will be very important for Section 3 .

We remark that, in $W_{\tau} \Psi$, the rotations $U$ and $V$ are present only in the interaction $\tilde{I}_{\tau}$ in $\xi_{\tau} \Psi$.

Recall that $\Pi$ is the orthogonal projection on $\mathcal{G}_{1} \otimes \mathcal{G}_{2}$. Since all functions on Ran $\Pi$ are exponentially decaying in the sense of Corollary 1.3, we have

$$
\begin{equation*}
\left[\Pi, \chi_{4 L / 3}^{\otimes N}\right]=O_{\left(C^{\infty}, \mathcal{B}\left(L^{2}\right)\right)}\left(e^{-d L}\right), \quad\left[\Pi, \chi_{L}^{\otimes N}\right]=O_{\left(C^{\infty}, \mathcal{B}\left(L^{2}\right)\right)}\left(e^{-d L}\right) \tag{2.17}
\end{equation*}
$$

where $\mathcal{B}\left(L^{2}\right)$ denotes the bounded operators on $L^{2}\left(\mathbb{R}^{3 N}\right)$ and we recall the notation (1.30). The $C^{\infty}$ estimate is here totally for free, since the left-hand sides do not depend on $U$ and $V$. Since by (2.6)

$$
\begin{equation*}
\Pi^{\perp} \xi_{\tau}=\xi_{\tau} \tag{2.18}
\end{equation*}
$$

we can use (2.17) to find that

$$
\begin{equation*}
\left\langle\chi_{L}^{\otimes N} \Psi, \chi_{4 L / 3}^{\otimes N} \xi_{\tau} \Psi\right\rangle=\left\langle\chi_{L}^{\otimes N} \Pi^{\perp} \Psi, \chi_{4 L / 3}^{\otimes N} \xi_{\tau} \Psi\right\rangle+O_{C^{\infty}}\left(e^{-d L}\right)=O_{C^{\infty}}\left(e^{-d L}\right) \tag{2.19}
\end{equation*}
$$

where in the last step we have used that $\Psi \in \operatorname{Ran} \Pi$. From (2.7) and (2.19),

$$
\begin{equation*}
\left\|W_{\tau} \Psi\right\|^{2}=\left\|\chi_{L}^{\otimes N} \Psi\right\|^{2}+\left\|\chi_{4 L / 3}^{\otimes N} \xi_{\tau} \Psi\right\|^{2}+O_{C^{\infty}}\left(e^{-d L}\right) \tag{2.20}
\end{equation*}
$$

By Corollary 1.3, there exists some $c>0$ such that

$$
\begin{equation*}
e^{c|\cdot|} \Psi \in L^{2}\left(\mathbb{R}^{3 N}\right) \tag{2.21}
\end{equation*}
$$

Consequently, as $\Psi$ is normalized,

$$
\begin{equation*}
\chi_{L}^{\otimes N} \Psi=\Psi+O_{\left(C^{\infty}, L^{2}\left(\mathbb{R}^{3 N}\right)\right)}\left(e^{-d L}\right) \tag{2.22}
\end{equation*}
$$

The $C^{\infty}$ estimate follows again, since the left-hand side does not depend on $U$ and $V$.
As e.g. in Equation (70) of [6], we use a multipolar expansion from Lemma 1.6 for $\tilde{I}_{\tau}$ to find that

$$
\begin{equation*}
\tilde{I}_{\tau} \chi_{L}^{\otimes N} \Psi=\frac{f_{(U, V)}}{L^{3}} \chi_{L}^{\otimes N} \Psi+O_{\left(C^{\infty}, L^{2}\left(\mathbb{R}^{3 N}\right)\right)}\left(\frac{1}{L^{4}}\right) \tag{2.23}
\end{equation*}
$$

where $f_{(U, V)}$ is defined in (1.39). To achieve a $C^{\infty}$ control of the error in (2.23), it is important to notice that the singularities of $\tilde{I}_{\tau}$ are far away from the support of $\chi_{L}^{\otimes N}$. Note that, even-though in [6] only $C^{2}$ regularity was on the right-hand side of $(2.23), C^{\infty}$ regularity can be proven with similar arguments. Indeed, the proof relies on Lemma 1.6, which works for derivatives of any order.

Using (2.23) and (2.21), we arrive at

$$
\begin{equation*}
\tilde{I}_{\tau} \chi_{L}^{\otimes N} \Psi=\frac{f_{(U, V)}}{L^{3}} \Psi+O_{\left(C^{\infty}, L^{2}\left(\mathbb{R}^{3 N}\right)\right)}\left(\frac{1}{L^{4}}\right) \tag{2.24}
\end{equation*}
$$

Since, $\left(H_{\infty}^{\perp}-E_{1}-E_{2}\right)^{-1}$ is bounded from $L^{2}\left(\mathbb{R}^{3 N}\right)$ to itself and independent of $\tau$, equations (2.6) and (2.24) imply that

$$
\begin{equation*}
\xi_{\tau} \Psi=O_{\left(C^{\infty}, L^{2}\left(\mathbb{R}^{3 N}\right)\right)}\left(\frac{1}{L^{3}}\right) \tag{2.25}
\end{equation*}
$$

Combining (2.20), (2.22), and (2.25), we obtain

$$
\begin{equation*}
W_{\tau} \Psi=\Psi+O_{\left(C^{\infty}, L^{2}\left(\mathbb{R}^{3 N}\right)\right)}\left(\frac{1}{L^{3}}\right), \quad\left\|W_{\tau} \Psi\right\|^{2}=1+O_{C^{\infty}}\left(\frac{1}{L^{6}}\right) \tag{2.26}
\end{equation*}
$$

We now calculate the numerator of the right-hand side of (2.16). We will start with a decomposition of $\left(H_{\infty}+\tilde{I}_{\tau}-E_{1}-E_{2}\right) W_{\tau} \Psi$. Using (2.7), we find

$$
\begin{equation*}
\left(H_{\infty}+\tilde{I}_{\tau}-E_{1}-E_{2}\right) W_{\tau} \Psi=\left(H_{\infty}+\tilde{I}_{\tau}-E_{1}-E_{2}\right) \chi_{L}^{\otimes N} \Psi-\left(H_{\infty}+\tilde{I}_{\tau}-E_{1}-E_{2}\right) \chi_{4 L / 3}^{\otimes N} \xi_{\tau} \Psi \tag{2.27}
\end{equation*}
$$

On the one hand, since $\left(H_{\infty}-E_{1}-E_{2}\right) \Psi=0$ and $\chi_{L}^{\otimes N}$ commutes with the potentials, we have

$$
\begin{equation*}
\left(H_{\infty}+\tilde{I}_{\tau}-E_{1}-E_{2}\right) \chi_{L}^{\otimes N} \Psi=\sum_{j=1}^{N}\left[T_{j}, \chi_{L}^{\otimes N}\right] \Psi+\tilde{I}_{\tau} \chi_{L}^{\otimes N} \Psi \tag{2.28}
\end{equation*}
$$

On the other hand,

$$
\begin{align*}
\left(H_{\infty}+\tilde{I}_{\tau}-E_{1}-E_{2}\right) \chi_{4 L / 3}^{\otimes N} \xi_{\tau} \Psi= & \sum_{j=1}^{N}\left[T_{j}, \chi_{4 L / 3}^{\otimes N}\right] \xi_{\tau} \Psi \\
& +\chi_{4 L / 3}^{\otimes N}\left(H_{\infty}-E_{1}-E_{2}\right) \xi_{\tau} \Psi+\tilde{I}_{\tau} \chi_{4 L / 3}^{\otimes N} \xi_{\tau} \Psi \tag{2.29}
\end{align*}
$$

Using (2.18) and that $H_{\infty}$ commutes with $\Pi^{\perp}$, we find that

$$
\begin{equation*}
\left(H_{\infty}-E_{1}-E_{2}\right) \xi_{\tau} \Psi=\left(H_{\infty}^{\perp}-E_{1}-E_{2}\right) \xi_{\tau} \Psi \tag{2.30}
\end{equation*}
$$

where recall that $H_{\infty}^{\perp}$ was defined in (1.41). From (2.6) and (2.30) we obtain that

$$
\begin{equation*}
\left(H_{\infty}-E_{1}-E_{2}\right) \xi_{\tau} \Psi=\Pi^{\perp} \chi_{L}^{\otimes N} \tilde{I}_{\tau} \Psi=\chi_{L}^{\otimes N} \tilde{I}_{\tau} \Psi-\Pi \chi_{L}^{\otimes N} \tilde{I}_{\tau} \Psi \tag{2.31}
\end{equation*}
$$

Since, due to (2.5), $\chi_{4 L / 3}$ is equal to 1 on the support of $\chi_{L}$, we have

$$
\begin{equation*}
\chi_{4 L / 3} \chi_{L}=\chi_{L} \tag{2.32}
\end{equation*}
$$

The equalities (2.31) and (2.32) imply

$$
\begin{equation*}
\chi_{4 L / 3}^{\otimes N}\left(H_{\infty}-E_{1}-E_{2}\right) \xi_{\tau} \Psi=\chi_{L}^{\otimes N} \tilde{I}_{\tau} \Psi-\chi_{4 L / 3}^{\otimes N} \Pi \chi_{L}^{\otimes N} \tilde{I}_{\tau} \Psi \tag{2.33}
\end{equation*}
$$

Inserting (2.33) in (2.29), we obtain

$$
\begin{align*}
\left(H_{\infty}+\tilde{I}_{\tau}-E_{1}-E_{2}\right) \chi_{4 L / 3}^{\otimes N} \xi_{\tau} \Psi=\sum_{j=1}^{N}\left[T_{j},\right. & \left.\chi_{4 L / 3}^{\otimes N}\right] \xi_{\tau} \Psi+\chi_{L}^{\otimes N} \tilde{I}_{\tau} \Psi \\
& -\chi_{4 L / 3}^{\otimes N} \Pi \chi_{L}^{\otimes N} \tilde{I}_{\tau} \Psi+\tilde{I}_{\tau} \chi_{4 L / 3}^{\otimes N} \xi_{\tau} \Psi \tag{2.34}
\end{align*}
$$

Inserting now (2.28) and (2.34) in (2.27), we arrive at

$$
\begin{equation*}
\left(H_{\infty}+\tilde{I}_{\tau}-E_{1}-E_{2}\right) W_{\tau} \Psi=\sum_{j=1}^{N}\left[T_{j}, \chi_{L}^{\otimes N}\right] \Psi-\sum_{j=1}^{N}\left[T_{j}, \chi_{4 L / 3}^{\otimes N}\right] \xi_{\tau} \Psi+\chi_{4 L / 3}^{\otimes N} \Pi \chi_{L}^{\otimes N} \tilde{I}_{\tau} \Psi-\tilde{I}_{\tau} \chi_{4 L / 3}^{\otimes N} \xi_{\tau} \Psi \tag{2.35}
\end{equation*}
$$

We remark that

$$
\left\|\sum_{j=1}^{N}\left[T_{j}, \chi_{L}^{\otimes N}\right] \Psi\right\|^{2}=\left\langle\sum_{j=1}^{N}\left[T_{j}, \chi_{L}^{\otimes N}\right] \Psi, \sum_{j=1}^{N}\left[T_{j}, \chi_{L}^{\otimes N}\right] \Psi\right\rangle
$$

But, according to Lemma A.2, the left-hand factor of the scalar product is in $L^{2}$ and its norm can be controlled uniformly in $L$. Therefore, due to (2.21), we can apply (A.9) and find that

$$
\begin{equation*}
\sum_{j=1}^{N}\left[T_{j}, \chi_{L}^{\otimes N}\right] \Psi=O_{\left(C^{\infty}, L^{2}\left(\mathbb{R}^{3 N}\right)\right)}\left(e^{-d L}\right) \tag{2.36}
\end{equation*}
$$

Here, the $C^{\infty}$ regularity is trivial since the left-hand side does not depend on $U$ and $V$. Using (2.35)-(2.36) we find

$$
\begin{align*}
&\left(H_{\infty}+\tilde{I}_{\tau}-E_{1}-E_{2}\right) W_{\tau} \Psi=-\sum_{j=1}^{N}\left[T_{j}, \chi_{4 L / 3}^{\otimes N}\right] \xi_{\tau} \Psi+\chi_{4 L / 3}^{\otimes N} \Pi \chi_{L}^{\otimes N} \tilde{I}_{\tau} \Psi \\
&-\tilde{I}_{\tau} \chi_{4 L / 3}^{\otimes N} \xi_{\tau} \Psi+O_{\left(C^{\infty}, L^{2}\left(\mathbb{R}^{3 N}\right)\right)}\left(e^{-d L}\right) \tag{2.37}
\end{align*}
$$

In order to estimate the numerator of the right-hand side of (2.16), we first take the scalar product of the terms of the right-hand side of (2.37) with $\chi_{4 L / 3}^{\otimes N} \xi_{\tau} \Psi$. We begin by observing that, thanks to Lemma A. 2 and (2.25),

$$
\begin{equation*}
\sum_{j=1}^{N}\left[T_{j}, \chi_{4 L / 3}^{\otimes N}\right] \xi_{\tau} \Psi=O_{\left(C^{\infty}, L^{2}\left(\mathbb{R}^{3 N}\right)\right)}\left(\frac{1}{L^{4}}\right) \tag{2.38}
\end{equation*}
$$

which, using (2.25) once more, gives

$$
\begin{equation*}
\left\langle\chi_{4 L / 3}^{\otimes N} \xi_{\tau} \Psi, \sum_{j=1}^{N}\left[T_{j}, \chi_{4 L / 3}^{\otimes N}\right] \xi_{\tau} \Psi\right\rangle=O_{C^{\infty}}\left(\frac{1}{L^{7}}\right) \tag{2.39}
\end{equation*}
$$

We find from the expression of $\tilde{I}_{\tau}$ in (2.4) that

$$
\tilde{I}_{\tau} \chi_{4 L / 3}^{\otimes N}=O_{\left(C^{\infty}, L^{\infty}\left(\mathbb{R}^{3 N}\right)\right)}\left(\frac{1}{L}\right)
$$

which combined with (2.25), gives that

$$
\begin{equation*}
\tilde{I}_{\tau} \chi_{4 L / 3}^{\otimes N} \xi_{\tau} \Psi=O_{\left(C^{\infty}, L^{2}\left(\mathbb{R}^{3 N}\right)\right)}\left(\frac{1}{L^{4}}\right) \tag{2.40}
\end{equation*}
$$

(2.25) and (2.40) imply

$$
\begin{equation*}
\left\langle\chi_{4 L / 3}^{\otimes N} \xi_{\tau} \Psi, \tilde{I}_{\tau} \chi_{4 L / 3}^{\otimes N} \xi_{\tau} \Psi\right\rangle=O_{C^{\infty}}\left(\frac{1}{L^{7}}\right) \tag{2.41}
\end{equation*}
$$

Moreover, by (2.18),

$$
\begin{equation*}
\left\langle\chi_{4 L / 3}^{\otimes N} \xi_{\tau} \Psi, \chi_{4 L / 3}^{\otimes N} \Pi \chi_{L}^{\otimes N} \tilde{I}_{\tau} \Psi\right\rangle=\left\langle\Pi^{\perp} \xi_{\tau} \Psi,\left[\left(\chi_{4 L / 3}^{\otimes N}\right)^{2}, \Pi\right] \chi_{L}^{\otimes N} \tilde{I}_{\tau} \Psi\right\rangle=O_{C^{\infty}}\left(e^{-d L}\right) \tag{2.42}
\end{equation*}
$$

since $\left[\left(\chi_{4 L / 3}^{\otimes N}\right)^{2}, \Pi\right]$ can be bounded as in (2.17). In the last step we have also used (2.24) and (2.25).
Collecting (2.39), (2.41) and (2.42), we find, with the help of (2.25) and (2.37), that

$$
\begin{equation*}
\left\langle\chi_{4 L / 3}^{\otimes N} \xi_{\tau} \Psi,\left(H_{\infty}+\tilde{I}_{\tau}-E_{1}-E_{2}\right) W_{\tau} \Psi\right\rangle=O_{C^{\infty}}\left(\frac{1}{L^{7}}\right) . \tag{2.43}
\end{equation*}
$$

Taking now a scalar product of the terms of the right-hand side of (2.37) with $\chi_{L}^{\otimes N} \Psi$, we first find

$$
\begin{equation*}
\left\langle\chi_{L}^{\otimes N} \Psi, \sum_{j=1}^{N}\left[T_{j}, \chi_{4 L / 3}^{\otimes N}\right] \xi_{\tau} \Psi\right\rangle=-\left\langle\sum_{j=1}^{N}\left[T_{j}, \chi_{4 L / 3}^{\otimes N}\right] \chi_{L}^{\otimes N} \Psi, \xi_{\tau} \Psi\right\rangle=O_{C \infty}\left(e^{-d L}\right) \tag{2.44}
\end{equation*}
$$

by arguing as in the proof of (2.36) and using (2.25). Then, we see that, by (2.6),

$$
\left\langle\chi_{L}^{\otimes N} \Psi, \tilde{I}_{\tau} \chi_{4 L / 3}^{\otimes N} \xi_{\tau} \Psi\right\rangle=\left\langle\tilde{I}_{\tau} \chi_{L}^{\otimes N} \Psi, \chi_{4 L / 3}^{\otimes N} \Pi^{\perp}\left(H_{\infty}^{\perp}-E_{1}-E_{2}\right)^{-1} \Pi^{\perp} \tilde{I}_{\tau} \chi_{L}^{\otimes N} \Psi\right\rangle
$$

As a consequence, using (2.32) and applying (2.24), we arrive at

$$
\left\langle\chi_{L}^{\otimes N} \Psi, \tilde{I}_{\tau} \chi_{4 L / 3}^{\otimes N} \xi_{\tau} \Psi\right\rangle=\frac{\left\langle f_{(U, V)} \Psi, \Pi^{\perp}\left(H_{\infty}^{\perp}-E_{1}-E_{2}\right)^{-1} \Pi^{\perp} f_{(U, V)} \Psi\right\rangle}{L^{6}}+O_{C^{\infty}}\left(\frac{1}{L^{7}}\right)
$$

which, in view of the definition of the van der Waals constant $C_{\mathrm{vdW}}$ given in (1.42), can be rewritten as

$$
\begin{equation*}
\left\langle\chi_{L}^{\otimes N} \Psi, \tilde{I}_{\tau} \chi_{4 L / 3}^{\otimes N} \xi_{\tau} \Psi\right\rangle=\frac{C_{\mathrm{vdW}}(\Psi, U, V)}{L^{6}}+O_{C^{\infty}}\left(\frac{1}{L^{7}}\right) . \tag{2.45}
\end{equation*}
$$

With (2.37), (2.44) and (2.45), we find that

$$
\left\langle\chi_{L}^{\otimes N} \Psi,\left(H_{\infty}+\tilde{I}_{\tau}-E_{1}-E_{2}\right) W_{\tau} \Psi\right\rangle=-\frac{C_{\mathrm{vdW}}(\Psi, U, V)}{L^{6}}+\left\langle\chi_{L}^{\otimes N} \Psi, \chi_{4 L / 3}^{\otimes N} \Pi_{L}^{\otimes N} \tilde{I}_{\tau} \Psi\right\rangle+O_{C^{\infty}}\left(\frac{1}{L^{7}}\right)
$$

which together with (2.43) and (2.7) gives

$$
\begin{equation*}
\left\langle W_{\tau} \Psi,\left(H_{\infty}+\tilde{I}_{\tau}-E_{1}-E_{2}\right) W_{\tau} \Psi\right\rangle=-\frac{C_{\mathrm{vdW}}(\Psi, U, V)}{L^{6}}+\left\langle\chi_{L}^{\otimes N} \Psi, \chi_{4 L / 3}^{\otimes N} \Pi \chi_{L}^{\otimes N} \tilde{I}_{\tau} \Psi\right\rangle+O_{C^{\infty}}\left(\frac{1}{L^{7}}\right) . \tag{2.46}
\end{equation*}
$$

Using (2.17), (2.32) and the fact that $\Pi \Psi=\Psi$,

$$
\begin{equation*}
\left\langle\chi_{L}^{\otimes N} \Psi, \chi_{4 L / 3}^{\otimes N} \Pi \chi_{L}^{\otimes N} \tilde{I}_{\tau} \Psi\right\rangle=\left\langle\chi_{L}^{\otimes N} \Psi, \tilde{I}_{\tau} \chi_{L}^{\otimes N} \Psi\right\rangle+O_{C^{\infty}}\left(e^{-d L}\right) \tag{2.47}
\end{equation*}
$$

(2.46) and (2.47) imply

$$
\begin{equation*}
\left\langle W_{\tau} \Psi,\left(H_{\infty}+\tilde{I}_{\tau}-E_{1}-E_{2}\right) W_{\tau} \Psi\right\rangle=\left\langle\chi_{L}^{\otimes N} \Psi, \tilde{I}_{\tau} \chi_{L}^{\otimes N} \Psi\right\rangle-\frac{C_{\mathrm{vdW}}(\Psi, U, V)}{L^{6}}+O_{C^{\infty}}\left(\frac{1}{L^{7}}\right) \tag{2.48}
\end{equation*}
$$

So far, we have worked with a general normalized state $\Psi \in \mathcal{G}_{1} \otimes \mathcal{G}_{2}$. Using Lemma 1.7 with the molecular charge densities defined by (1.34), we obtain in the special case $\Psi=\Psi_{1} \otimes \Psi_{2}$

$$
\begin{align*}
\left\langle\chi_{L}^{\otimes N} \Psi_{1} \otimes \Psi_{2}, \tilde{I}_{\tau} \chi_{L}^{\otimes N} \Psi_{1} \otimes \Psi_{2}\right\rangle & =\sum_{2 \leqslant m+n \leqslant 5} \frac{\mathcal{F}^{(n, m)}\left(U \rho_{\chi_{L}^{\otimes N_{1}} \Psi_{1}}, V \rho_{\chi_{L}^{\otimes N_{2}} \Psi_{2}}\right)}{L^{m+n+1}}+O_{C^{\infty}}\left(\frac{1}{L^{7}}\right) \\
& =\sum_{2 \leqslant m+n \leqslant 5} \frac{\mathcal{F}^{(n, m)}\left(U \rho_{\Psi_{1}}, V \rho_{\Psi_{2}}\right)}{L^{m+n+1}}+O_{C^{\infty}}\left(\frac{1}{L^{7}}\right), \tag{2.49}
\end{align*}
$$

where, in the last step we could remove the cut-off functions due to the exponential decay of the ground state. Indeed, (2.21) implies, in view of (1.35), that $\mathcal{F}^{(n, m)}\left(U \rho_{\chi_{L}^{\otimes N_{1} \Psi_{1}}}, V \rho_{\chi_{L}^{\otimes N_{2} \Psi_{2}}}\right)$ and $\mathcal{F}^{(n, m)}\left(U \rho_{\Psi_{1}}, V \rho_{\Psi_{2}}\right)$ are exponentially close to each other.

Inserting (2.49) in (2.48) gives, for $\Psi=\Psi_{1} \otimes \Psi_{2}$,

$$
\begin{align*}
\left\langle W_{\tau} \Psi_{1} \otimes \Psi_{2},\left(H_{\infty}+\tilde{I}_{\tau}-E_{1}-E_{2}\right) W_{\tau} \Psi_{1} \otimes \Psi_{2}\right\rangle= & \sum_{1 \leqslant m+n \leqslant 5} \frac{\mathcal{F}^{(n, m)}\left(U \rho_{\Psi_{1}}, V \rho_{\Psi_{2}}\right)}{L^{m+n+1}} \\
& -\frac{C_{\mathrm{vdW}}\left(\Psi_{1} \otimes \Psi_{2}, U, V\right)}{L^{6}}+O_{C^{\infty}}\left(\frac{1}{L^{7}}\right) . \tag{2.50}
\end{align*}
$$

Using (2.8), (2.16), and (2.50), we find, choosing $\Psi=\Psi_{1} \otimes \Psi_{2}$, that

$$
\begin{align*}
& \mathcal{E}_{\tau}-E_{1}-E_{2} \leqslant \frac{1}{\left\|W_{\tau} \Psi_{1} \otimes \Psi_{2}\right\|^{2}}\left(\sum_{1 \leqslant m+n \leqslant 5} \frac{\mathcal{F}^{(n, m)}\left(U \rho_{\Psi_{1}}, V \rho_{\Psi_{2}}\right)}{L^{m+n+1}}\right. \\
&\left.-\frac{C_{\mathrm{vdW}}\left(\Psi_{1} \otimes \Psi_{2}, U, V\right)}{L^{6}}+O_{C^{\infty}}\left(\frac{1}{L^{7}}\right)\right) . \tag{2.51}
\end{align*}
$$

Note that the norm of the trial function gives a negligible contribution because, from (2.26),

$$
\begin{equation*}
\frac{1}{\left\|W_{\tau} \Psi\right\|^{2}}=\frac{1}{1+O_{C^{\infty}}\left(\frac{1}{L^{6}}\right)}=1+O_{C^{\infty}}\left(\frac{1}{L^{6}}\right) . \tag{2.52}
\end{equation*}
$$

From (2.51) and (2.52) for $\Psi=\Psi_{1} \otimes \Psi_{2}$, we immediately arrive at (1.45).
The next step is to prove that, for all $\Psi \in \mathcal{G}_{1} \otimes \mathcal{G}_{2}, U$ and $V$, the van der Waals constant $C_{\mathrm{vdW}}(\Psi, U, V)$, defined in (1.42), is positive. The definition clearly implies that $C_{\mathrm{vdW}}(\Psi, U, V) \geqslant 0$ for all states $\Psi$ and for all $U$ and $V$. Moreover, it is zero if and only if $\Pi^{\perp} f_{(U, V)} \Psi=0$. This would mean that $f_{(U, V)} \Psi$ is a ground state of $\hat{H}_{1}+\hat{H}_{2}$. We will prove that it is not the case, with the help of the following lemma. It is analogous of [6, Lemma 2.4]. Nevertheless, since we cannot use the Leibniz rule in the (SR) case, we need a modification of the proof.
Lemma 2.1. Let $\Psi \in L^{2}\left(\mathbb{R}^{3 N_{k}}\right) \backslash\{0\}(k=1,2)$ be such that $H_{k} \Psi=E \Psi$ for some $E$, where the operators $H_{k}$ have been defined in (1.18). For a vector $\left(b_{1}, \ldots, b_{N_{k}}\right) \in \mathbb{R}^{3 N_{k}} \backslash\{0\}$, we define the operator $B$ by

$$
(B \Psi)(x):=\sum_{j=1}^{N_{k}} b_{j} \cdot x_{j} \Psi(x)
$$

Then, $H_{k} B \Psi \neq E B \Psi$.
Proof. In the non-relativistic case, the analogous result is proven in [6, Lemma 2.4]. Since it relies on the Leibniz rule, we have to rework the proof to obtain the same result in the semirelativistic case.

Assume that $H_{k} B \Psi=E B \Psi$. Then, $\Psi$ satisfies the equations

$$
H_{k} \Psi=\sum_{j=1}^{N_{k}} T_{j} \Psi+I_{k} \Psi=E \Psi
$$

and

$$
H_{k} B \Psi=\sum_{j=1}^{N_{k}} T_{j} B \Psi+I_{k} B \Psi=E B \Psi
$$

where $I_{k}:=H_{k}-\sum_{j=1}^{N_{k}} T_{j}$ contains all the interaction terms. Subtracting the last two equations, we find that

$$
\begin{equation*}
\sum_{j=1}^{N_{k}}\left[T_{j}, B\right] \Psi=0 \tag{2.53}
\end{equation*}
$$

Let us denote, for each $j, b_{j}=\left(b_{j}^{1}, b_{j}^{2}, b_{j}^{3}\right)$ and $x_{j}=\left(x_{j}^{1}, x_{j}^{2}, x_{j}^{3}\right)$. We have then that $\left[T_{j}, B\right]=\sum_{k=1}^{3} b_{j}^{k}\left[T_{j}, x_{j}^{k}\right]$. Let us compute the commutator $\left[T_{j}, x_{j}^{k}\right]$. Using the definition of $T_{j}$ through the Fourier transformation in (1.3), we find that that

$$
\left[T_{j}, x_{j}^{k}\right] \psi=-\mathrm{i}\left(\sqrt{-\Delta_{j}+1}\right)^{-1} \partial_{x_{j}^{k}}
$$

so $(2.53)$ is equivalent to

$$
\sum_{j=1}^{N_{k}}-\mathrm{i}\left(\sqrt{-\Delta_{j}+1}\right)^{-1} b_{j} \cdot \nabla_{j} \Psi=0
$$

Applying the Fourier transformation (1.4), we find

$$
\sum_{j=1}^{N_{k}} \frac{b_{j} \cdot p_{j}}{\sqrt{1+\left|p_{j}\right|^{2}}} \mathcal{F} \Psi(p)=0
$$

Assuming without loss of generality that $b_{1} \neq 0$, we find that $\partial_{p_{1}} \sum_{j=1}^{N_{k}} \frac{b_{j} \cdot p_{j}}{\sqrt{1+\left|p_{j}\right|^{2}}} \neq 0$. As a consequence, by the implicit function theorem, $\sum_{j=1}^{N_{k}} \frac{b_{j} \cdot p_{j}}{\sqrt{1+\left|p_{j}\right|^{2}}}$ is, when restricted on compact sets, zero only in a
finite union of graphs. Therefore, it is nonzero almost everywhere. Consequently, $\mathcal{F} \Psi(p)=0$ almost everywhere, and thus $\Psi=0$, which contradicts the assumption $\Psi \neq 0$.

We are now ready to prove the positivity of the van der Waals constant, following the proof of $[6$, Proposition 2.3]. For the sake of completeness, we give it here. Due to the definition of $f_{(U, V)}$ in (1.39), we have that

$$
f_{(U, V)} \Psi=\sum_{n, m \geqslant 0} a_{n m} \Phi_{1, n} \otimes \Phi_{2, m}
$$

where $\left(a_{n m}\right)$ are some real coefficients, $\Phi_{j, 0} \in \mathcal{G}_{j}$, and, for all $n>0, \Phi_{j, n} \in x_{i}^{k} \mathcal{G}_{j}$, for some $i \in\{1, \ldots, N\}$ and $k \in\{1,2,3\}$. Assuming that $f_{(U, V)} \Psi$ is a ground state of $H_{\infty}$, we find that

$$
\sum_{n, m \geqslant 1} a_{n m} \Phi_{1, n}^{\perp} \otimes \Phi_{2, m}^{\perp}=0
$$

with $\Phi_{k, n}^{\perp}=\Pi_{k}^{\perp} \Phi_{k, n}$ where $\Pi_{k}$ is the orthogonal projection on $\mathcal{G}_{k}$. As a consequence, the functions $\left(\Phi_{1, n}^{\perp} \otimes \Phi_{2, m}^{\perp}\right)_{n, m \geqslant 1}$ are linearly dependent and we have that, for dimensional reasons, at least one of the sets $\left(\Phi_{1, n}^{\perp}\right)_{n \geqslant 1}$ or $\left(\Phi_{2, m}^{\perp}\right)_{m \geqslant 1}$ is linearly dependent. If we assume, for example, that there exists a finite family of complex numbers $\left(b_{n}\right)$ such that

$$
\sum_{n \geqslant 1} b_{n} \Phi_{1, n}^{\perp}=0
$$

then $\sum_{n \geqslant 1} b_{n} \Phi_{1, n}$ is an eigenvector of $H_{1}$ associated with the eigenvalue $E_{1}$, which contradicts Lemma 2.1.
2.2. Lower bound: general case. In this section, we prove the following proposition, for which we do not assume Condition (Poles). We will use it as an intermediate result to prove Theorems 1.10 (b) and 1.11 in the next subsection.

Proposition 2.2. Assume Condition (Neutr). Then

$$
\begin{equation*}
\mathcal{E}_{\tau}=E_{\infty}+\min _{\substack{\Psi \in \mathcal{G}_{1} \otimes \mathcal{G}_{2} \\\|\Psi\|=1}}\left(\left\langle\chi_{L}^{\otimes N} \Psi, \tilde{I}_{\tau} \chi_{L}^{\otimes N} \Psi\right\rangle-\frac{C_{\mathrm{vdW}}(\Psi, U, V)}{L^{6}}\right)+O_{L^{\infty}}\left(\frac{1}{L^{7}}\right) \tag{2.54}
\end{equation*}
$$

Recall that, if (Neutr) holds, then $E_{\infty}=E_{1}+E_{2}$. We will use this equality without mentioning it, when we refer to equations of the previous subsection. Note that, in the (NR) case, the proposition easily follows from the arguments of [6], even though it was not stated in this form. As a consequence, we consider here only the (SR) case.

In order to prove a lower bound on the ground state energy, we want to make use of the Feshbach map method, which we recall here: let $H$ be a self-adjoint operator on a Hilbert space $\mathscr{H}$ with domain $\mathcal{D}(H)$, $P$ an orthogonal projection on $\mathscr{H}$ such that $\operatorname{Ran} P \subset \mathcal{D}(H)$ and $H^{\perp}:=P^{\perp} H P^{\perp}$, where $P^{\perp}:=1-P$. The Feshbach map is an operator on $\operatorname{Ran} P$ defined in the following way: for any $E \in \mathbb{R}$ such that $\left(H^{\perp}-E\right)$ is invertible,

$$
F_{P}(E):=\left.\left(P H P-P H P^{\perp}\left(H^{\perp}-E\right)^{-1} P^{\perp} H P\right)\right|_{\operatorname{Ran} P .}
$$

The following well-known theorem, see e.g. [12, Section IV], shows how the Feshbach map is a useful tool to estimate the ground state energy of the system that we consider.

Theorem 2.3. Suppose $H, \mathscr{D}(H), \mathscr{H}, P$ to be as above and that, for some $E \in \mathbb{R}$, there exists $C>0$ such that

$$
H^{\perp}-E \geqslant C>0
$$

Then

$$
E \text { is an eigenvalue of } H \quad \Longleftrightarrow \quad E \text { is an eigenvalue of } F_{P}(E)
$$

Moreover, if, for some $\phi \in \operatorname{Ran} P \backslash\{0\}$, we have $F_{P}(E) \phi=E \phi$, then

$$
\phi-\left(H^{\perp}-E\right)^{-1} P^{\perp} H \phi,
$$

is an eigenfunction of $H$ to the eigenvalue $E$.

Furthermore, it turns out that, if $E$ is the ground state energy of $H$, then $E$ is the ground state energy of $F_{P}(E)$ and therefore there exists $\psi_{0} \in \operatorname{Ran} P$ with $\left\|\psi_{0}\right\|=1$ such that

$$
\begin{equation*}
E=\left\langle\psi_{0}, F_{P}(E) \psi_{0}\right\rangle=\min _{\substack{\psi \in \operatorname{Ran} P \\\|\psi\|=1}}\left\langle\psi, F_{P}(E) \psi\right\rangle, \tag{2.55}
\end{equation*}
$$

see e.g. [5, Lemma 5.6].
In order to apply the Feshbach method to our case, we choose $\mathscr{H}=Q_{N} \mathcal{H}, H=\hat{H}_{\tau}$, recall (1.12), (1.26), and $P=P_{\tau}$ the orthogonal projection with

$$
\begin{equation*}
\operatorname{Ran} P_{\tau}:=\left\{Q_{N} K_{\tau} W_{\tau} \Psi \mid \Psi \in \mathcal{G}_{1} \otimes \mathcal{G}_{2}\right\} \tag{2.56}
\end{equation*}
$$

where $Q_{N}, K_{\tau}, W_{\tau}, \mathcal{G}_{j}$ were respectively defined in (1.10), (2.1), (2.7) and (1.40).
The next lemma is the analogous of [8, Proposition 4.1] in the semirelativistic case and shows that we are in the assumptions of Theorem 2.3 thanks to Condition (Neutr).

Lemma 2.4. Let $P_{\tau}$ be the orthogonal projection defined by (2.56) and assume Condition (Neutr). Then there exists $C>0$ such that, for $L$ large enough,

$$
\begin{equation*}
P_{\tau}^{\perp} H_{\tau} P_{\tau}^{\perp}-\mathcal{E}_{\tau} \geqslant C . \tag{2.57}
\end{equation*}
$$

Proof. If we follow the proof of [8, Proposition 4.1], we get that there exists $C, C^{\prime}>0$ such that for all $L$ large enough

$$
\Pi_{\tau}^{\perp} H_{\tau} \Pi_{\tau}^{\perp}-E_{\infty} \geqslant\left(C-\frac{C^{\prime}}{L}\right)
$$

on $Q_{N} \mathcal{H}$, where recall that $E_{\infty}$ is given by (1.44). Here $\Pi_{\tau}$ is the orthogonal projection on

$$
\begin{equation*}
\operatorname{Ran} \Pi_{\tau}:=\left\{Q_{N} K_{\tau} \chi_{L}^{\otimes N} \Psi, \Psi \in \mathcal{G}_{1} \otimes \mathcal{G}_{2}\right\} \tag{2.58}
\end{equation*}
$$

This together with (1.44) and (1.45) implies that

$$
\begin{equation*}
\Pi_{\tau}^{\perp} H_{\tau} \Pi_{\tau}^{\perp}-\mathcal{E}_{\tau} \geqslant\left(C-\frac{C^{\prime}}{L}\right) \tag{2.59}
\end{equation*}
$$

on $Q_{N} \mathcal{H}$. Even if the result of [8] is about non-relativistic Hamiltonians, the proof can easily be adapted to the (SR) case. Indeed, the only part of the proof where the fact that the kinetic energy operator is a Laplacian matters is the application of the IMS localization formula, which is replaced in our case by equation (A.8) in Appendix A. In fact, [14] provided in Section 3 finer estimates that work for a broader class of localization functions in the (SR) case.

To extract (2.57) from (2.59), we need to bound

$$
\begin{equation*}
H_{\tau}\left(\Pi_{\tau}^{\perp}-P_{\tau}^{\perp}\right)=H_{\tau}\left(P_{\tau}-\Pi_{\tau}\right)=-\left(H_{\tau}-E_{\infty}\right) \Pi_{\tau}+\left(H_{\tau}-E_{\infty}\right) P_{\tau}-E_{\infty}\left(\Pi_{\tau}-P_{\tau}\right) \tag{2.60}
\end{equation*}
$$

Let us begin by bounding $\left(H_{\tau}-E_{\infty}\right) \Pi_{\tau}$. Due to (2.58), we have

$$
\begin{equation*}
\left\|\left(H_{\tau}-E_{\infty}\right) \Pi_{\tau}\right\|=\max _{\Psi \in \mathcal{G}_{1} \otimes \mathcal{G}_{2},\|\Psi\|=1} \frac{\left\|\left(H_{\tau}-E_{\infty}\right) Q_{N} K_{\tau} \chi_{L}^{\otimes N} \Psi\right\|}{\left\|Q_{N} K_{\tau} \chi_{L}^{\otimes N} \Psi\right\|} \tag{2.61}
\end{equation*}
$$

Let us consider some $\Psi \in \mathcal{G}_{1} \otimes \mathcal{G}_{2}$ such that $\|\Psi\|=1$. We prove, as we have done in (2.14), that

$$
\begin{equation*}
\left\|Q_{N} K_{\tau} \chi_{L}^{\otimes N} \Psi\right\|^{2}=\frac{1}{\binom{N}{N_{1}}}\left\|K_{\tau} \chi_{L}^{\otimes N} \Psi\right\|^{2}=\frac{1}{\binom{N}{N_{1}}}\left\|\chi_{L}^{\otimes N} \Psi\right\|^{2}=\frac{1}{\binom{N}{N_{1}}}+O_{C^{\infty}}\left(e^{-d L}\right) \tag{2.62}
\end{equation*}
$$

where in the second step we used that $K_{\tau}$ is unitary and in the last step we used (2.22).
We next bound $\left(H_{\tau}-E_{\infty}\right) Q_{N} K_{\tau} \chi_{L}^{\otimes N} \Psi$. From (2.2), (2.3), (1.23) and the fact that $Q_{N}$ commutes with $H_{\tau}$, we have

$$
\begin{equation*}
\left(H_{\tau}-E_{\infty}\right) Q_{N} K_{\tau}=Q_{N} K_{\tau}\left(H_{\infty}+\tilde{I}_{\tau}-E_{\infty}\right) \tag{2.63}
\end{equation*}
$$

As $Q_{N}$ is an orthogonal projection and $K_{\tau}$ is unitary, (2.61) and (2.63) imply

$$
\begin{equation*}
\left\|\left(H_{\tau}-E_{\infty}\right) \Pi_{\tau}\right\| \leqslant \max _{\Psi \in \mathcal{G}_{1} \otimes \mathcal{G}_{2},\|\Psi\|=1} \frac{1}{\left\|Q_{N} K_{\tau} \chi_{L}^{\otimes N} \Psi\right\|}\left\|\left(H_{\infty}+\tilde{I}_{\tau}-E_{\infty}\right) \chi_{L}^{\otimes N} \Psi\right\| \tag{2.64}
\end{equation*}
$$

But we know from (1.44), (2.28), (2.36) and (2.24) that

$$
\begin{equation*}
\left(H_{\infty}-E_{\infty}+\tilde{I}_{\tau}\right) \chi_{L}^{\otimes N} \Psi=O_{\left(C^{\infty}, L^{2}\left(\mathbb{R}^{3 N}\right)\right)}\left(\frac{1}{L^{3}}\right) \tag{2.65}
\end{equation*}
$$

By (2.62), (2.64) and (2.65), we find

$$
\begin{equation*}
\left\|\left(H_{\tau}-E_{\infty}\right) \Pi_{\tau}\right\|=O_{L^{\infty}}\left(\frac{1}{L^{3}}\right) \tag{2.66}
\end{equation*}
$$

Note that, because of the max in (2.64), we do not have $C^{\infty}$ estimate on the right hand side of (2.66).
Let us now bound $\left(H_{\tau}-E_{\infty}\right) P_{\tau}$. Using (2.56) and arguing as in the proof of (2.64), we find that

$$
\begin{equation*}
\left\|\left(H_{\tau}-E_{\infty}\right) P_{\tau}\right\| \leqslant \max _{\Psi \in \mathcal{G}_{1} \otimes \mathcal{G}_{2},\|\Psi\|=1} \frac{\left\|\left(H_{\infty}+\tilde{I}_{\tau}-E_{\infty}\right) W_{\tau} \Psi\right\|}{\left\|Q_{N} K_{\tau} W_{\tau} \Psi\right\|} . \tag{2.67}
\end{equation*}
$$

Using (2.14), (2.26) and that $K_{\tau}$ is unitary, we arrive at

$$
\begin{equation*}
\left\|Q_{N} K_{\tau} W_{\tau} \Psi\right\|^{2}=\frac{1}{\binom{N}{N_{1}}}+O_{C \infty}\left(\frac{1}{L^{6}}\right) \tag{2.68}
\end{equation*}
$$

uniformly on $\Psi \in \mathcal{G}_{1} \otimes \mathcal{G}_{2}$ with $\|\Psi\|=1$.
We have, from (2.37), (2.38) and (2.40), that

$$
\begin{equation*}
\left(H_{\infty}+\tilde{I}_{\tau}-E_{\infty}\right) W_{\tau} \Psi=\chi_{4 L / 3}^{\otimes N} \Pi \chi_{L}^{\otimes N} \tilde{I}_{\tau} \Psi+O_{\left(C^{\infty}, L^{2}\left(\mathbb{R}^{3 N}\right)\right)}\left(\frac{1}{L^{4}}\right) \tag{2.69}
\end{equation*}
$$

and therefore, by (2.24),

$$
\begin{equation*}
\left(H_{\infty}+\tilde{I}_{\tau}-E_{\infty}\right) W_{\tau} \Psi=O_{\left(C^{\infty}, L^{2}\left(\mathbb{R}^{3 N}\right)\right)}\left(\frac{1}{L^{3}}\right) \tag{2.70}
\end{equation*}
$$

Inserting (2.68) and (2.70) in (2.67), we obtain

$$
\begin{equation*}
\left\|\left(H_{\tau}-E_{\infty}\right) P_{\tau}\right\|=O_{L^{\infty}}\left(\frac{1}{L^{3}}\right) \tag{2.71}
\end{equation*}
$$

Finally, let us bound $\left\|P_{\tau}-\Pi_{\tau}\right\|$. To this end, we remind the following well-known linear algebra lemma:

Lemma 2.5. Let $V$ be finite-dimensional subspace of a Hilbert space with inner product $\langle\cdot \mid \cdot\rangle$ and $\left(v_{i}\right)_{i=1, \ldots, n}$ a basis of $V$. Let us denote by $G$ the Gram matrix, defined by $G_{i j}=\left\langle v_{i} \mid v_{j}\right\rangle$. Then $P_{V}$, the orthogonal projection on $V$, satisfies

$$
P_{V}=\sum_{i=1}^{n}\left(G^{-1}\right)_{i j}\left|v_{i}\right\rangle\left\langle v_{j}\right| .
$$

Let now $\Psi_{k}, k=1, \ldots, \operatorname{dim} \mathcal{G}_{1} \otimes \mathcal{G}_{2}$, be an orthonormal basis of $\mathcal{G}_{1} \otimes \mathcal{G}_{2}$. Arguing similarly as in the proof of (2.68), we can find

$$
\begin{equation*}
\left\langle Q_{N} K_{\tau} W_{\tau} \Psi_{k}, Q_{N} K_{\tau} W_{\tau} \Psi_{l}\right\rangle=\frac{\delta_{k l}}{\binom{N}{N_{1}}}+O_{C^{\infty}}\left(\frac{1}{L^{6}}\right) \tag{2.72}
\end{equation*}
$$

A similar argument gives that

$$
\begin{equation*}
\left\langle Q_{N} K_{\tau} \chi_{L}^{\otimes N} \Psi_{k}, Q_{N} K_{\tau} \chi_{L}^{\otimes N} \Psi_{l}\right\rangle=\frac{\delta_{k l}}{\binom{N}{N_{1}}}+O_{C^{\infty}}\left(e^{-d L}\right) \tag{2.73}
\end{equation*}
$$

Note that the different order of the error terms in (2.72) and (2.73) originates from the difference of the order in the error terms of the right-hand sides of (2.22) and (2.26). As a consequence, for $L$ large enough, the functions $\left(Q_{N} K_{\tau} \chi_{L}^{\otimes N} \Psi_{k}\right)$ (resp. $\left.\left(Q_{N} K_{\tau} W_{\tau} \Psi_{k}\right)\right), k=1, \ldots, \operatorname{dim} \mathcal{G}_{1} \otimes \mathcal{G}_{2}$, are linearly independent. Thus, due to (2.58), (2.56), they form a basis of $\operatorname{Ran} \Pi_{\tau}$ (resp. Ran $P_{\tau}$ ).

Therefore, we can apply Lemma 2.5 to find

$$
\begin{equation*}
\Pi_{\tau}=\sum_{k, l=1}^{\operatorname{dim} \mathcal{G}_{1} \otimes \mathcal{G}_{2}}\left|Q_{N} K_{\tau} \chi_{L}^{\otimes N} \Psi_{k}\right\rangle g_{\tau}^{k l}\left\langle Q_{N} K_{\tau} \chi_{L}^{\otimes N} \Psi_{l}\right|, \tag{2.74}
\end{equation*}
$$

and

$$
\begin{equation*}
P_{\tau}=\sum_{k, l=1}^{\operatorname{dim} \mathcal{G}_{1} \otimes \mathcal{G}_{2}}\left|Q_{N} K_{\tau} W_{\tau} \Psi_{k}\right\rangle G_{\tau}^{k l}\left\langle Q_{N} K_{\tau} W_{\tau} \Psi_{l}\right| \tag{2.75}
\end{equation*}
$$

where $g_{\tau}^{k l}, G_{\tau}^{k l}$ are the elements of the inverses of the respective Gram matrices. From (2.72) and (2.73) we find

$$
\begin{equation*}
g_{\tau}^{k l}=\binom{N}{N_{1}} \delta_{k l}+O_{C^{\infty}}\left(e^{-c L}\right), \quad G_{\tau}^{k l}=\binom{N}{N_{1}} \delta_{k l}+O_{C^{\infty}}\left(\frac{1}{L^{6}}\right) \tag{2.76}
\end{equation*}
$$

By (2.74), (2.75) and (2.7) we find

$$
\begin{aligned}
P_{\tau}-\Pi_{\tau}= & -\sum_{k, l=1}^{\operatorname{dim} \mathcal{G}_{1} \otimes \mathcal{G}_{2}}\left|Q_{N} K_{\tau} \chi_{4 L / 3}^{\otimes N} \xi_{\tau} \Psi_{k}\right\rangle G_{\tau}^{k l}\left\langle Q_{N} K_{\tau} W_{\tau} \Psi_{l}\right| \\
& +\sum_{k, l=1}^{\operatorname{dim} \mathcal{G}_{1} \otimes \mathcal{G}_{2}}\left|Q_{N} K_{\tau} \chi_{L}^{\otimes N} \Psi_{k}\right\rangle\left(G_{\tau}^{k l}-g_{\tau}^{k l}\right)\left\langle Q_{N} K_{\tau} W_{\tau} \Psi_{l}\right| \\
& -\sum_{k, l=1}^{\operatorname{dim} \mathcal{G}_{1} \otimes \mathcal{G}_{2}}\left|Q_{N} K_{\tau} \chi_{L}^{\otimes N} \Psi_{k}\right\rangle g_{\tau}^{k l}\left\langle Q_{N} K_{\tau} \chi_{4 L / 3}^{\otimes N} \xi_{\tau} \Psi_{l}\right|,
\end{aligned}
$$

which together with (2.25) and (2.76) implies

$$
\begin{equation*}
\left\|P_{\tau}-\Pi_{\tau}\right\|=O_{L^{\infty}}\left(\frac{1}{L^{3}}\right) \tag{2.77}
\end{equation*}
$$

Inserting (2.66), (2.71), and (2.77) into (2.60), we find that

$$
\left\|H_{\tau}\left(\Pi_{\tau}^{\perp}-P_{\tau}^{\perp}\right)\right\|=O_{L^{\infty}}\left(\frac{1}{L^{3}}\right)
$$

which, together with (2.59), proves that (2.57) is satisfied for $L$ large enough. This concludes the proof of Lemma 2.4.

Due to Lemma 2.4, we can use Theorem 2.3, (2.56) and (2.55) to obtain that

$$
\begin{equation*}
\mathcal{E}_{\tau}=\min _{\substack{\Psi \in \mathcal{G}_{1} \otimes \mathcal{G}_{2} \\\|\Psi\|=1}} \frac{\left\langle Q_{N} K_{\tau} W_{\tau} \Psi, F_{P_{\tau}}\left(\mathcal{E}_{\tau}\right) Q_{N} K_{\tau} W_{\tau} \Psi\right\rangle}{\left\|Q_{N} K_{\tau} W_{\tau} \Psi\right\|^{2}} \tag{2.78}
\end{equation*}
$$

where the Feshbach map $F_{P_{\tau}}$ is given in our setting by

$$
\begin{equation*}
F_{P_{\tau}}(E):=P_{\tau} H_{\tau} P_{\tau}-P_{\tau} H_{\tau} P_{\tau}^{\perp}\left(H_{\tau}^{\perp}-E\right)^{-1} P_{\tau}^{\perp} H_{\tau} P_{\tau} . \tag{2.79}
\end{equation*}
$$

Here $H_{\tau}^{\perp}=P_{\tau}^{\perp} H_{\tau} P_{\tau}^{\perp}$. The first term of the right-hand side of (2.79) is not a function of $E$ and the second one is a non-linear function of $E$. In order to find a lower bound for $\mathcal{E}_{\tau}$, we will bound from below the quadratic form in the right-hand side of (2.78). We start by bounding from below

$$
\begin{equation*}
\frac{\left\langle Q_{N} K_{\tau} W_{\tau} \Psi, P_{\tau} H_{\tau} P_{\tau} Q_{N} K_{\tau} W_{\tau} \Psi\right\rangle}{\left\|Q_{N} K_{\tau} W_{\tau} \Psi\right\|^{2}} \tag{2.80}
\end{equation*}
$$

We use (2.16) and (2.56) which give that for all $\Psi \in \mathcal{G}_{1} \otimes \mathcal{G}_{2}$,

$$
\begin{equation*}
\frac{\left\langle Q_{N} K_{\tau} W_{\tau} \Psi, P_{\tau} H_{\tau} P_{\tau} Q_{N} K_{\tau} W_{\tau} \Psi\right\rangle}{\left\|Q_{N} K_{\tau} W_{\tau} \Psi\right\|^{2}}=\frac{\left\langle W_{\tau} \Psi,\left(H_{\infty}+\tilde{I}_{\tau}\right) W_{\tau} \Psi\right\rangle}{\left\|W_{\tau} \Psi\right\|^{2}} \tag{2.81}
\end{equation*}
$$

Inserting (2.48) and (2.52) in the right-hand side of (2.81), we find that

$$
\begin{equation*}
\frac{\left\langle Q_{N} K_{\tau} W_{\tau} \Psi, P_{\tau} H_{\tau} P_{\tau} Q_{N} K_{\tau} W_{\tau} \Psi\right\rangle}{\left\|Q_{N} K_{\tau} W_{\tau} \Psi\right\|^{2}}=E_{\infty}+\left\langle\chi_{L}^{\otimes N} \Psi, \tilde{I}_{\tau} \chi_{L}^{\otimes N} \Psi\right\rangle-\frac{C_{\mathrm{vdW}}(\Psi, U, V)}{L^{6}}+O_{C^{\infty}}\left(\frac{1}{L^{7}}\right) . \tag{2.82}
\end{equation*}
$$

We continue by bounding from above

$$
\begin{equation*}
\frac{\left\langle Q_{N} K_{\tau} W_{\tau} \Psi, P_{\tau} H_{\tau} P_{\tau}^{\perp}\left(H_{\tau}^{\perp}-\mathcal{E}_{\tau}\right)^{-1} P_{\tau}^{\perp} H_{\tau} P_{\tau} Q_{N} K_{\tau} W_{\tau} \Psi\right\rangle}{\left\|Q_{N} K_{\tau} W_{\tau} \Psi\right\|^{2}} \tag{2.83}
\end{equation*}
$$

Let $Q_{N} K_{\tau} W_{\tau} \Psi \in \operatorname{Ran} P_{\tau}$, for $\Psi \in \mathcal{G}_{1} \otimes \mathcal{G}_{2}$. We observe that, thanks to (2.56), (2.57) and the equality $P_{\tau}^{\perp} H_{\tau} P_{\tau}=P_{\tau}^{\perp}\left(H_{\tau}-E_{\infty}\right) P_{\tau}$, there exists $C>0$ such that, if $L$ is large enough,

$$
\begin{equation*}
\frac{\left\langle Q_{N} K_{\tau} W_{\tau} \Psi, P_{\tau} H_{\tau} P_{\tau}^{\perp}\left(H_{\tau}^{\perp}-\mathcal{E}_{\tau}\right)^{-1} P_{\tau}^{\perp} H_{\tau} P_{\tau} Q_{N} K_{\tau} W_{\tau} \Psi\right\rangle}{\left\|Q_{N} K_{\tau} W_{\tau} \Psi\right\|^{2}} \leqslant C^{-1} \frac{\left\|P_{\tau}^{\perp}\left(H_{\tau}-E_{\infty}\right) Q_{N} K_{\tau} W_{\tau} \Psi\right\|^{2}}{\left\|Q_{N} K_{\tau} W_{\tau} \Psi\right\|^{2}} \tag{2.84}
\end{equation*}
$$

Using (2.63), we find

$$
P_{\tau}^{\perp}\left(H_{\tau}-E_{\infty}\right) Q_{N} K_{\tau} W_{\tau} \Psi=P_{\tau}^{\perp} Q_{N} K_{\tau}\left(H_{\infty}+\tilde{I}_{\tau}-E_{\infty}\right) W_{\tau} \Psi
$$

which, together with (2.69), gives

$$
\begin{equation*}
P_{\tau}^{\perp}\left(H_{\tau}-E_{\infty}\right) Q_{N} K_{\tau} W_{\tau} \Psi=P_{\tau}^{\perp} Q_{N} K_{\tau} \chi_{4 L / 3}^{\otimes N} \Pi \chi_{L}^{\otimes N} \tilde{I}_{\tau} \Psi+O_{\left(L^{\infty}, L^{2}\left(\mathbb{R}^{3 N}\right)\right)}\left(\frac{1}{L^{4}}\right) \tag{2.85}
\end{equation*}
$$

Let $\Psi=\Psi_{1}$ be a normalized element of $\mathcal{G}_{1} \otimes \mathcal{G}_{2}$. In order to bound $P_{\tau}^{\perp} Q_{N} K_{\tau} \chi_{4 L / 3}^{\otimes N} \Pi$, we will prove that

$$
P_{\tau}^{\perp} Q_{N} K_{\tau} \chi_{4 L / 3}^{\otimes N} \Psi_{1}=O_{\left(L^{\infty}, L^{2}\left(\mathbb{R}^{3 N}\right)\right)}\left(\frac{1}{L^{3}}\right),
$$

uniformly on $\Psi_{1} \in \mathcal{G}_{1} \otimes \mathcal{G}_{2}$ with $\left\|\Psi_{1}\right\|=1$. We begin by extending $\Psi_{1}$ to an orthonormal basis $\Psi_{1}, \ldots, \Psi_{n}$ of $\mathcal{G}_{1} \otimes \mathcal{G}_{2}$. Using (2.75), we see that

$$
\begin{equation*}
P_{\tau} Q_{N} K_{\tau} \Psi_{1}=\sum_{k, l=1}^{n}\left|Q_{N} K_{\tau} W_{\tau} \Psi_{k}\right\rangle G_{\tau}^{k l}\left\langle Q_{N} K_{\tau} W_{\tau} \Psi_{l}, Q_{N} K_{\tau} \chi_{4 L / 3}^{\otimes N} \Psi_{1}\right\rangle \tag{2.86}
\end{equation*}
$$

Arguing as in the proof of (2.68), we find

$$
\begin{equation*}
\left\langle Q_{N} K_{\tau} W_{\tau} \Psi_{l}, Q_{N} K_{\tau} \chi_{4 L / 3}^{\otimes N} \Psi_{1}\right\rangle=\frac{\left\langle W_{\tau} \Psi_{l}, \chi_{4 L / 3}^{\otimes N} \Psi_{1}\right\rangle}{\binom{N}{N_{1}}} \tag{2.87}
\end{equation*}
$$

But we have, due to (2.21),

$$
\begin{equation*}
\left\langle\chi_{L}^{\otimes N} \Psi_{l}, \chi_{4 L / 3}^{\otimes N} \Psi_{1}\right\rangle=\delta_{l 1}+O_{C^{\infty}}\left(e^{-d L}\right) \tag{2.88}
\end{equation*}
$$

and, due to (2.18), (2.17) and (2.25),

$$
\begin{equation*}
\left\langle\chi_{4 L / 3}^{\otimes N} \xi_{\tau} \Psi_{l}, \chi_{4 L / 3}^{\otimes N} \Psi_{1}\right\rangle=\left\langle\left[\left(\chi_{4 L / 3}^{\otimes N}\right)^{2}, \Pi^{\perp}\right] \xi_{\tau} \Psi_{l}, \Psi_{1}\right\rangle=O_{C^{\infty}}\left(e^{-d L}\right), \tag{2.89}
\end{equation*}
$$

since $\Pi^{\perp} \Psi_{1}=0$. Using (2.87) (2.88), (2.89) and (2.7), we find that

$$
\begin{equation*}
\left\langle Q_{N} K_{\tau} W_{\tau} \Psi_{l}, Q_{N} K_{\tau} \chi_{4 L / 3}^{\otimes N} \Psi_{1}\right\rangle=\frac{\delta_{1 l}}{\binom{N}{N_{1}}}+O_{C \infty}\left(e^{-d L}\right) \tag{2.90}
\end{equation*}
$$

With the help of (2.76), (2.86) and (2.90), we see that

$$
\begin{equation*}
P_{\tau} Q_{N} K_{\tau} \Psi_{1}=Q_{N} K_{\tau} W_{\tau} \Psi_{1}+O_{\left(L^{\infty}, L^{2}\left(\mathbb{R}^{3 N}\right)\right)}\left(\frac{1}{L^{6}}\right) \tag{2.91}
\end{equation*}
$$

We obtain from (2.91) and (2.7) that
$P_{\tau}^{\perp} Q_{N} K_{\tau} \Psi_{1}=Q_{N} K_{\tau}\left(1-\chi_{L}^{\otimes N}\right) \Psi_{1}+Q_{N} K_{\tau} \chi_{4 L / 3}^{\otimes N} \xi_{\tau} \Psi_{1}+O_{\left(L^{\infty}, L^{2}\left(\mathbb{R}^{3 N}\right)\right)}\left(\frac{1}{L^{6}}\right)=O_{\left(L^{\infty}, L^{2}\left(\mathbb{R}^{3 N}\right)\right)}\left(\frac{1}{L^{3}}\right)$,
where in the last equality we used (2.25) and (2.22). As a consequence,

$$
\begin{equation*}
\left\|P_{\tau}^{\perp} Q_{N} K_{\tau} \chi_{4 L / 3}^{\otimes N} \Pi\right\|=O_{L^{\infty}}\left(\frac{1}{L^{3}}\right) . \tag{2.92}
\end{equation*}
$$

Inserting (2.92) and (2.23) in (2.85), we obtain

$$
\begin{equation*}
\left\|P_{\tau}^{\perp}\left(H_{\tau}-E_{\infty}\right) Q_{N} K_{\tau} W_{\tau} \Psi\right\|=O_{L^{\infty}}\left(\frac{1}{L^{4}}\right) \tag{2.93}
\end{equation*}
$$

uniformly for all normalized $\Psi \in \mathcal{G}_{1} \otimes \mathcal{G}_{2}$.
Using (2.93), (2.68) and (2.84), we arrive at

$$
\begin{equation*}
\frac{\left\langle Q_{N} K_{\tau} W_{\tau} \Psi, P_{\tau} H_{\tau} P_{\tau}^{\perp}\left(H_{\tau}^{\perp}-\mathcal{E}_{\tau}\right)^{-1} P_{\tau}^{\perp} H_{\tau} P_{\tau} Q_{N} K_{\tau} W_{\tau} \Psi\right\rangle}{\left\|Q_{N} K_{\tau} W_{\tau} \Psi\right\|^{2}}=O_{L^{\infty}}\left(\frac{1}{L^{8}}\right) . \tag{2.94}
\end{equation*}
$$

Using (2.82), (2.94), (2.79) and (2.78), we obtain (2.54). This concludes the proof of Proposition 2.2.
2.3. Proof of Theorems 1.10 (b) and 1.11. In this section, we prove Theorems 1.10 (b) and 1.11 using Proposition 2.2. As we explained below the statement of Proposition 2.2, the latter holds in both the (SR) and the (NR) cases. Thus, we consider here both cases simultaneously. Note that, in this step, the analysis of [6] has to be changed, due to the weaker assumption (Poles) for the ground state eigenspaces.

In view of Proposition 2.2, our goal is to estimate the quadratic form

$$
\begin{equation*}
Q_{\tau}(\Psi):=\left\langle\chi_{L}^{\otimes N} \Psi, \tilde{I}_{\tau} \chi_{L}^{\otimes N} \Psi\right\rangle \tag{2.95}
\end{equation*}
$$

on the unit ball of $\mathcal{G}_{1} \otimes \mathcal{G}_{2}$, under one of the two versions of Condition (Poles). Let us first assume Condition (Poles-all). Together with (2.49), it directly implies that, for all normalized $\Psi_{1} \in \mathcal{G}_{1}$ and $\Psi_{2} \in \mathcal{G}_{2}$,

$$
\begin{equation*}
Q_{\tau}\left(\Psi_{1} \otimes \Psi_{2}\right)=\sum_{1 \leqslant m+n \leqslant 5} \frac{\mathcal{F}^{(n, m)}(U, V)}{L^{m+n+1}}+O_{C^{\infty}}\left(\frac{1}{L^{7}}\right), \tag{2.96}
\end{equation*}
$$

where the notation $\mathcal{F}^{(n, m)}(U, V)$ was introduced right after Equation (1.43). Let us now consider the general case, i.e. where $\Psi$ is not a factorized state. For such a $\Psi$, since $\mathcal{G}_{j}$ are finite dimensional, we know that there exist orthonormal vectors $\left\{\Phi_{1}^{(k)}\right\}_{k=1}^{\operatorname{dim} \mathcal{G}_{1}} \subseteq \mathcal{G}_{1}$ and $\left\{\Phi_{2}^{(l)}\right\}_{l=1}^{\operatorname{dim} \mathcal{G}_{2}} \subseteq \mathcal{G}_{2}$ such that

$$
\begin{equation*}
\Psi=\sum_{k=1}^{\operatorname{dim} \mathcal{G}_{1}} \sum_{l=1}^{\operatorname{dim} \mathcal{G}_{2}} a_{k, l} \Phi_{1}^{(k)} \otimes \Phi_{2}^{(l)}, \quad\left\{a_{k, l}\right\}_{(k, l)} \subseteq \mathbb{C}, \sum_{k, l}\left|a_{k, l}\right|^{2}=1 \tag{2.97}
\end{equation*}
$$

As a consequence, we can write that

$$
\begin{align*}
Q_{\tau}(\Psi)= & \sum_{p, q} \sum_{l, k} \overline{a_{p, q}} a_{l, k}\left\langle\chi_{L}^{\otimes N} \Phi_{1}^{(p)} \otimes \Phi_{2}^{(q)}, \tilde{I}_{\tau} \chi_{L}^{\otimes N} \Phi_{1}^{(l)} \otimes \Phi_{2}^{(k)}\right\rangle \\
= & \sum_{p, q}\left|a_{p, q}\right|^{2} Q_{\tau}\left(\Phi_{1}^{(p)} \otimes \Phi_{2}^{(q)}\right) \\
& +\sum_{(p, q) \neq(k, l)} \overline{a_{p, q}} a_{l, k}\left\langle\chi_{L}^{\otimes N} \Phi_{1}^{(p)} \otimes \Phi_{2}^{(q)}, \tilde{I}_{\tau} \chi_{L}^{\otimes N} \Phi_{1}^{(k)} \otimes \Phi_{2}^{(l)}\right\rangle . \tag{2.98}
\end{align*}
$$

For the diagonal part of (2.98), we find, with the help of (2.96) and (2.97), that

$$
\begin{equation*}
\sum_{p, q}\left|a_{p, q}\right|^{2} Q_{\tau}\left(\Phi_{1}^{(p)} \otimes \Phi_{2}^{(q)}\right)=\sum_{1 \leqslant n+m \leqslant 5} \frac{\mathcal{F}^{(n, m)}(U, V)}{L^{m+n+1}}+O_{C^{\infty}}\left(\frac{1}{L^{7}}\right) \tag{2.99}
\end{equation*}
$$

We now prove that, for all $(p, q) \neq(k, l)$,

$$
\begin{equation*}
\left\langle\chi_{L}^{\otimes N} \Phi_{1}^{(p)} \otimes \Phi_{2}^{(q)}, \tilde{I}_{\tau} \chi_{L}^{\otimes N} \Phi_{1}^{(k)} \otimes \Phi_{2}^{(l)}\right\rangle=O_{C^{\infty}}\left(\frac{1}{L^{7}}\right) \tag{2.100}
\end{equation*}
$$

By the complex polarization identity, we find that

$$
\begin{align*}
& 4\left\langle\chi_{L}^{\otimes N} \Phi_{1}^{(p)} \otimes \Phi_{2}^{(q)}, \tilde{I}_{\tau} \chi_{L}^{\otimes N} \Phi_{1}^{(k)} \otimes \Phi_{2}^{(l)}\right\rangle=\left\langle\chi_{L}^{\otimes N}\left(\Phi_{1}^{(p)}+\Phi_{1}^{(k)}\right) \otimes \Phi_{2}^{(q)}, \tilde{I}_{\tau} \chi_{L}^{\otimes N}\left(\Phi_{1}^{(p)}+\Phi_{1}^{(k)}\right) \otimes \Phi_{2}^{(l)}\right\rangle \\
&-\left\langle\chi_{L}^{\otimes N}\left(\Phi_{1}^{(p)}-\Phi_{1}^{(k)}\right) \otimes \Phi_{2}^{(q)}, \tilde{I}_{\tau} \chi_{L}^{\otimes N}\left(\Phi_{1}^{(p)}-\Phi_{1}^{(k)}\right) \otimes \Phi_{2}^{(l)}\right\rangle \\
&+\mathrm{i}\left\langle\chi_{L}^{\otimes N}\left(\mathrm{i} \Phi_{1}^{(p)}+\Phi_{1}^{(k)}\right) \otimes \Phi_{2}^{(q)}, \tilde{I}_{\tau} \chi_{L}^{\otimes N}\left(\mathrm{i} \Phi_{1}^{(p)}+\Phi_{1}^{(k)}\right) \otimes \Phi_{2}^{(l)}\right\rangle \\
&-\mathrm{i}\left\langle\chi_{L}^{\otimes N}\left(-\mathrm{i} \Phi_{1}^{(p)}+\Phi_{1}^{(k)}\right) \otimes \Phi_{2}^{(q)}, \tilde{I}_{\tau} \chi_{L}^{\otimes N}\left(-\mathrm{i} \Phi_{1}^{(p)}+\Phi_{1}^{(k)}\right) \otimes \Phi_{2}^{(l)}\right\rangle . \tag{2.101}
\end{align*}
$$

Let us first consider the case where $q=l$, and thus $p \neq k$. We can rewrite the previous equation in the following way:

$$
\begin{aligned}
4\left\langle\chi_{L}^{\otimes N} \Phi_{1}^{(p)} \otimes \Phi_{2}^{(q)}, \tilde{I}_{\tau} \chi_{L}^{\otimes N} \Phi_{1}^{(k)} \otimes \Phi_{2}^{(q)}\right\rangle= & Q_{\tau}\left(\left(\Phi_{1}^{(p)}+\Phi_{1}^{(k)}\right) \otimes \Phi_{2}^{(q)}\right)-Q_{\tau}\left(\left(\Phi_{1}^{(p)}-\Phi_{1}^{(k)}\right) \otimes \Phi_{2}^{(q)}\right) \\
& +\mathrm{i} Q_{\tau}\left(\left(\mathrm{i} \Phi_{1}^{(p)}+\Phi_{1}^{(k)}\right) \otimes \Phi_{2}^{(q)}\right)-\mathrm{i} Q_{\tau}\left(\left(-\mathrm{i} \Phi_{1}^{(p)}+\Phi_{1}^{(k)}\right) \otimes \Phi_{2}^{(q)}\right)
\end{aligned}
$$

But, since the four vectors $\Phi_{1}^{(p)}+\Phi_{1}^{(k)}, \Phi_{1}^{(p)}-\Phi_{1}^{(k)}, \Phi_{1}^{(p)}+\Phi_{1}^{(k)}$, and $-\mathrm{i} \Phi_{1}^{(p)}+\Phi_{1}^{(k)}$ are in $\mathcal{G}_{1}$ and have the same norm $\sqrt{2}$, we have, by $(2.96)$, that the four values $Q_{\tau}\left(\left(\Phi_{1}^{(p)}+\Phi_{1}^{(k)}\right) \otimes \Phi_{2}^{(q)}\right), Q_{\tau}\left(\left(\Phi_{1}^{(p)}-\Phi_{1}^{(k)}\right) \otimes \Phi_{2}^{(q)}\right)$, $Q_{\tau}\left(\left(\mathrm{i} \Phi_{1}^{(p)}+\Phi_{1}^{(k)}\right) \otimes \Phi_{2}^{(q)}\right)$, and $Q_{\tau}\left(\left(-\mathrm{i} \Phi_{1}^{(p)}+\Phi_{1}^{(k)}\right) \otimes \Phi_{2}^{(q)}\right)$ are all equal up to an error of order $O_{C^{\infty}}\left(1 / L^{7}\right)$. As a consequence,

$$
\begin{equation*}
\left\langle\chi_{L}^{\otimes N} \Phi_{1}^{(p)} \otimes \Phi_{2}^{(q)}, \tilde{I}_{\tau} \chi_{L}^{\otimes N} \Phi_{1}^{(k)} \otimes \Phi_{2}^{(q)}\right\rangle=O_{C^{\infty}}\left(\frac{1}{L^{7}}\right) \tag{2.102}
\end{equation*}
$$

Let us now consider the case $q \neq l$ and $p \neq k$. In view of (2.101), we will prove that

$$
\left\langle\chi_{L}^{\otimes N}\left(\Phi_{1}^{(p)}+\Phi_{1}^{(k)}\right) \otimes \Phi_{2}^{(q)}, \tilde{I}_{\tau} \chi_{L}^{\otimes N}\left(\Phi_{1}^{(p)}+\Phi_{1}^{(k)}\right) \otimes \Phi_{2}^{(l)}\right\rangle=O_{C^{\infty}}\left(\frac{1}{L^{7}}\right)
$$

All the other terms of the right-hand side of (2.101) can be similarly bounded, since $-\Phi_{1}^{(k)}$ and $\pm \mathrm{i} \Phi_{1}^{(k)}$ are normalized ground states of $\hat{H}_{1}$ too. We apply again a polarization identity, which gives

$$
\begin{aligned}
4\left\langle\chi _ { L } ^ { \otimes N } \left(\Phi_{1}^{(p)}\right.\right. & \left.\left.+\Phi_{1}^{(k)}\right) \otimes \Phi_{2}^{(q)}, \tilde{I}_{\tau} \chi_{L}^{\otimes N}\left(\Phi_{1}^{(p)}+\Phi_{1}^{(k)}\right) \otimes \Phi_{2}^{(l)}\right\rangle \\
= & Q_{\tau}\left(\left(\Phi_{1}^{(p)}+\Phi_{1}^{(k)}\right) \otimes\left(\Phi_{2}^{(q)}+\Phi_{2}^{(l)}\right)\right)-Q_{\tau}\left(\left(\Phi_{1}^{(p)}+\Phi_{1}^{(k)}\right) \otimes\left(\Phi_{2}^{(q)}-\Phi_{2}^{(l)}\right)\right) \\
& +\mathrm{i} Q_{\tau}\left(\left(\Phi_{1}^{(p)}+\Phi_{1}^{(k)}\right) \otimes\left(\mathrm{i} \Phi_{2}^{(q)}+\Phi_{2}^{(l)}\right)\right)-\mathrm{i} Q_{\tau}\left(\left(\Phi_{1}^{(p)}+\Phi_{1}^{(k)}\right) \otimes\left(-\mathrm{i} \Phi_{2}^{(q)}+\Phi_{2}^{(l)}\right)\right)
\end{aligned}
$$

With the same argument as for (2.102), we find that the 4 quadratic forms in the right-hand side of the previous equation are equal up to $O_{C}^{\infty}\left(1 / L^{7}\right)$ and therefore

$$
\left\langle\chi_{L}^{\otimes N} \Phi_{1}^{(p)} \otimes \Phi_{2}^{(q)}, \tilde{I}_{\tau} \chi_{L}^{\otimes N} \Phi_{1}^{(k)} \otimes \Phi_{2}^{(l)}\right\rangle=O_{C^{\infty}}\left(\frac{1}{L^{7}}\right)
$$

The remaining case $p=k, q \neq l$ can be handled just in the same way as the case $p \neq k, q=l$. This concludes the proof of (2.100).

From (2.95), (2.98), (2.99), and (2.100), we find that

$$
\begin{equation*}
\left\langle\chi_{L}^{\otimes N} \Psi, \tilde{I}_{\tau} \chi_{L}^{\otimes N} \Psi\right\rangle=\sum_{1 \leqslant m+n \leqslant 5} \frac{\mathcal{F}^{(n, m)}(U, V)}{L^{m+n+1}}+O_{C^{\infty}}\left(\frac{1}{L^{7}}\right) . \tag{2.103}
\end{equation*}
$$

Inserting (2.103) in (2.54), we find

$$
\mathcal{E}_{\tau}=E_{\infty}+\sum_{1 \leqslant m+n \leqslant 5} \frac{\mathcal{F}^{(n, m)}(U, V)}{L^{m+n+1}}-\frac{C_{\mathrm{vdW}}(U, V)}{L^{6}}+O_{L^{\infty}}\left(\frac{1}{L^{7}}\right),
$$

where we recall that $C_{\mathrm{vdW}}(U, V)$ was defined in (1.43). This concludes the proof of point (b) of Theorem 1.10.

We now prove Theorem 1.11. In the case where $n_{1}+n_{2}<5$ and we assume only Condition (Poleslead), the proof of (1.47) is almost the same as the one of (1.46). The only change is that (2.96), (2.99) and therefore (2.103) are no longer valid. We have the following modifications: (2.96) becomes

$$
Q_{\tau}\left(\Psi_{1} \otimes \Psi_{2}\right)=\frac{\mathcal{F}^{\left(n_{1}, n_{2}\right)}(U, V)}{L^{n_{1}+n_{2}+1}}+O_{C^{\infty}}\left(\frac{1}{L^{n_{1}+n_{2}+2}}\right)
$$

and thus (2.99) becomes

$$
\sum_{p, q}\left|a_{p, q}\right|^{2} Q_{\tau}\left(\Phi_{1}^{(p)} \otimes \Phi_{2}^{(q)}\right)=\frac{\mathcal{F}^{\left(n_{1}, n_{2}\right)}(U, V)}{L^{n_{1}+n_{2}+1}}+O_{C^{\infty}}\left(\frac{1}{L^{n_{1}+n_{2}+2}}\right)
$$

Similarly, (2.100) becomes

$$
\left\langle\chi_{L}^{\otimes N} \Phi_{1}^{(p)} \otimes \Phi_{2}^{(q)}, \tilde{I}_{\tau} \chi_{L}^{\otimes N} \Phi_{1}^{(k)} \otimes \Phi_{2}^{(l)}\right\rangle=O_{C^{\infty}}\left(\frac{1}{L^{n_{1}+n_{2}+2}}\right)
$$

for all $(p, q) \neq(k, l)$. As a consequence, instead of (2.103) we have

$$
\begin{equation*}
\left\langle\chi_{L}^{\otimes N} \Psi, \tilde{I}_{\tau} \chi_{L}^{\otimes N} \Psi\right\rangle=\frac{\mathcal{F}^{\left(n_{1}, n_{2}\right)}(U, V)}{L^{n_{1}+n_{2}+1}}+O_{C \infty}\left(\frac{1}{L^{n_{1}+n_{2}+2}}\right) \tag{2.104}
\end{equation*}
$$

Inserting (2.104) in (2.54), we arrive at (1.47).

## 3. Proof of Proposition 1.19: A smoothed version of the ground state energy and EXPANSION OF ITS DERIVATIVES

In order to prove Proposition 1.19, we would like to use that, for each $L$ large enough, any local pseudo-minimum of the function $(U, V) \mapsto \mathcal{E}_{L}(U, V)$, recall (1.22) and (1.54), should be a stationary point of it as well, but, in principle, $\mathcal{E}_{L}$ could be non-differentiable, due to the infimum in its definition.

We treat both the (SR) and (NR) cases. We could not adapt the analysis of [6] in Section 3 so that it works for both cases as, in the (SR) case, major additional difficulties occur. The first one is, like before, that the kinetic energy is a nonlocal operator, but this is not all: the singularities of the Coulomb potentials in $H_{\tau}$ move when $\tau$ changes and, therefore, the second derivative of the interaction terms with respect to rotations gives a term that is not locally $L^{1}$. In [6], this problem was compensated using that the ground states of $H_{\tau}$ are in $H^{2}$ in the (NR) case. This is not at all true in the (SR) case. In fact, in this case, it is not even clear whether the ground states are in $H^{1}$, since, as we explained below

Equation (1.8), the domain of $H_{\tau}$ may contain functions that are not in $H^{1}$. This is another reason why the strategy of [6] is not applicable in our setting.

To overcome these difficulties, we will use ideas of the proof of Theorem 1 in [37]. Hunziker proved in the (NR) case that there is a family of unitary transformations $U(\Xi)$ such that $U^{-1}(\Xi) H_{N}(Y+\Xi, \mathcal{Z}) U(\Xi)$ depends smoothly on $\Xi$, when $\Xi$ is in a neighborhood of 0 . The unitary transformations were chosen so that the motion of the singularities of the Coulomb potentials was reversed. Here, we will adapt this method so that it firstly provides quantitative information for the derivatives and secondly becomes applicable in the (SR) case as well.

In fact, in the (NR) case, our strategy provides a simpler proof than the one of Section 3 of [6] and makes it possible to provide information not only for the first two derivatives, but for higher derivatives as well, see Proposition 3.1 below. Note that, in our proof, we treat both cases at the same time and we specify when we consider only one of them.

We assume that Condition (Neutr) is satisfied. Let $\tau_{0}:=(L, \bar{U}, \bar{V})$ be such that the function $\mathcal{E}_{L}$, defined in (1.54), has a local pseudo-minimum in $(\bar{U}, \bar{V})$ - recall Definition 1.17. By Theorem 2.3 and (2.56), there exists some $\Psi_{\tau_{0}} \in \mathcal{G}_{1} \otimes \mathcal{G}_{2}$ with $\left\|\Psi_{\tau_{0}}\right\|=1$ such that the function

$$
\begin{equation*}
\Phi_{\tau_{0}}:=Q_{N} K_{\tau_{0}} W_{\tau_{0}} \Psi_{\tau_{0}}-R_{\tau_{0}} \tag{3.1}
\end{equation*}
$$

with

$$
\begin{equation*}
R_{\tau_{0}}:=\left(H_{\tau_{0}}^{\perp}-\mathcal{E}_{\tau_{0}}\right)^{-1} P_{\tau_{0}}^{\perp}\left(H_{\tau_{0}}-E_{\infty}\right) Q_{N} K_{\tau_{0}} W_{\tau_{0}} \Psi_{\tau_{0}} \tag{3.2}
\end{equation*}
$$

is a ground state of $\hat{H}_{\tau_{0}}$, recall (1.26). In other words,

$$
\mathcal{E}_{\tau_{0}}=\inf \sigma\left(\hat{H}_{\tau_{0}}\right)=\frac{\left\langle\Phi_{\tau_{0}}, H_{\tau_{0}} \Phi_{\tau_{0}}\right\rangle}{\left\|\Phi_{\tau_{0}}\right\|^{2}}
$$

Note that in (3.2) we could subtract $E_{\infty}$ since $P_{\tau_{0}}^{\perp} Q_{N} K_{\tau_{0}} W_{\tau_{0}} \Psi_{\tau_{0}}=0$ by (2.56). A key observation here is that, if for some open neighborhood $\Omega \subset \mathrm{SO}(3) \times \operatorname{SO}(3),\left(\Phi_{\tau}\right)_{\tau=(L, U, V),(U, V) \in \Omega}$ is a family of functions so that $\Phi_{\tau_{0}}$ is given by (3.1), then $(\bar{U}, \bar{V})$ is a local pseudo-minimum not only of $\mathcal{E}_{\tau}$ but also of the new function

$$
\begin{equation*}
\tilde{\mathcal{E}}_{\tau}:=\frac{\left\langle\Phi_{\tau}, H_{\tau} \Phi_{\tau}\right\rangle}{\left\|\Phi_{\tau}\right\|^{2}} \tag{3.3}
\end{equation*}
$$

The observation simply follows from the fact that $\tilde{\mathcal{E}}_{\tau} \geqslant \mathcal{E}_{\tau}$, by the definition of $\mathcal{E}_{\tau}$ in (1.22), and the equality $\mathcal{E}_{\tau_{0}}=\tilde{\mathcal{E}}_{\tau_{0}}$.

We want to choose the family $\Phi_{\tau}$ appropriately so that $\tilde{\mathcal{E}}_{\tau}$ has the same leading order as $\mathcal{E}_{\tau}$ when $L$ is large and, at the same time, is smooth. This is an idea that already appeared in [6], but here we choose a simpler family. Moreover, the rest of the proof completely changes.

We consider the family of vector fields $m_{\tau}: \mathbb{R}^{3} \rightarrow \mathbb{R}^{3}$ with

$$
\begin{equation*}
m_{\tau}(x):=U \bar{U}^{-1} x_{\frac{8 L}{3}}(x)+\left[L e_{1}+V \bar{V}^{-1}\left(x-L e_{1}\right)\right] \chi_{\frac{8 L}{3}}\left(x-L e_{1}\right)+x\left(1-\chi_{\frac{8 L}{3}}(x)-\chi_{\frac{8 L}{3}}\left(x-L e_{1}\right)\right), \tag{3.4}
\end{equation*}
$$

where $\chi_{\frac{8 L}{3}}$ was defined in (2.5). Note that, by (2.5),

$$
\begin{equation*}
\chi_{\frac{8 L}{3}} \chi_{2 L}=\chi_{2 L}, \quad \chi_{2 L} \chi_{\frac{4 L}{3}}=\chi_{\frac{4 L}{3}} \tag{3.5}
\end{equation*}
$$

and

$$
\begin{equation*}
\operatorname{supp} \chi_{\frac{8 L}{3}} \cap\left(\operatorname{supp} \chi_{\frac{8 L}{3}}-L e_{1}\right)=\emptyset \tag{3.6}
\end{equation*}
$$

The idea of the definition (3.4) is that, if $x$ is close to the $j$-th molecule ( $j=1$ or 2 ), then it is rotated together with the nuclei around the center of the $j$-th molecule. Indeed, using (3.4)-(3.6), we find that

$$
\begin{equation*}
m_{\tau}(x)=U \bar{U}^{-1} x \text { on } \operatorname{supp} \chi_{2 L}, m_{\tau}(x)=L e_{1}+V \bar{V}^{-1}\left(x-L e_{1}\right) \text { on }\left(\operatorname{supp} \chi_{2 L}-L e_{1}\right) \tag{3.7}
\end{equation*}
$$

compare with (1.21). If $x$ is far from both molecules, namely outside of the supports of the functions $\chi_{\frac{8 L}{3}}$ and $\chi_{\frac{8 L}{3}}\left(.-L e_{1}\right)$, then $m_{\tau}(x)=x$ by (3.4).

A simple computation gives that

$$
\begin{equation*}
m_{\tau}(x)=x+v_{\tau}(x) \tag{3.8}
\end{equation*}
$$

where

$$
\begin{equation*}
v_{\tau}(x):=(U-\bar{U}) \bar{U}^{-1} x \chi_{\frac{8 L}{3}}(x)+(V-\bar{V}) \bar{V}^{-1}\left(x-L e_{1}\right) \chi_{\frac{8 L}{3}}\left(x-L e_{1}\right) . \tag{3.9}
\end{equation*}
$$

We have that, for all $(U, V) \in \mathrm{SO}(3) \times \mathrm{SO}(3)$ and $L>1$,

$$
\begin{equation*}
\left\|v_{\tau}^{\prime}\right\|_{L^{\infty}} \leqslant C \sup _{y \in \mathbb{R}^{3},\|y\|=1}\left(\|(U-\bar{U}) y\|_{\mathbb{R}^{3}}+\|(V-\bar{V}) y\|_{\mathbb{R}^{3}}\right) \tag{3.10}
\end{equation*}
$$

where

$$
C:=\left\|\left(x \chi_{\frac{8 L}{3}}(x)\right)^{\prime}\right\|_{L^{\infty}}
$$

is uniformly bounded with respect to $L$. In particular, there exists some open neighborhood $\Omega$ of $(\bar{U}, \bar{V})$ such that

$$
\begin{equation*}
(U, V) \in \Omega \Longrightarrow v_{\tau}(\cdot) \text { is a contraction uniformly in } L>1 . \tag{3.11}
\end{equation*}
$$

It follows that, if $(U, V) \in \Omega$, then $m_{\tau}(\cdot)$ is a bijection. Indeed, for any $y \in \mathbb{R}^{3}$, the vector field $x \mapsto y-v_{\tau}(x)$ is a contraction and therefore has a unique fixed point. As a consequence, $m_{\tau}$ has an inverse map. Since $m_{\tau} \in C^{\infty}$ and $m_{\tau}^{\prime}$ is invertible, it follows, by the implicit function theorem, that $m_{\tau}$ is a $C^{\infty}$ diffeomorphism.

Following ideas of [37], we consider the family of unitary transformations $u_{\tau}: L^{2}\left(\mathbb{R}^{3 N}\right) \rightarrow L^{2}\left(\mathbb{R}^{3 N}\right)$ defined for $(U, V) \in \Omega$, through

$$
\begin{equation*}
\left(u_{\tau} \Psi\right)\left(x_{1}, \ldots, x_{N}\right):=\left(\prod_{j=1}^{N} J_{m_{\tau}}\left(x_{j}\right)\right) \Psi\left(m_{\tau}\left(x_{1}\right), \ldots, m_{\tau}\left(x_{N}\right)\right), \tag{3.12}
\end{equation*}
$$

where $J_{f}(x)$ denotes the determinant of the Jacobian matrix of $f$ at the point $x$. We now define

$$
\begin{equation*}
\Phi_{\tau}:=Q_{N} K_{\tau} W_{\tau} \Psi_{\tau_{0}}-u_{\tau}^{-1} R_{\tau_{0}} \tag{3.13}
\end{equation*}
$$

where recall that $R_{\tau_{0}}$ was defined in (3.2). In [37], a more general family of unitary transformations was considered by choosing the family $m_{\tau}$ in a more general way. Here, we choose $m_{\tau}$ more specifically such that the electrons close to a molecule rotate together with its nuclei. This observation will enable us to obtain quantitative information for the derivatives with the help, e.g., of equation (3.22) below.

We remark that, because of (2.93) and the boundedness of the resolvent $\left(H_{\tau_{0}}^{\perp}-\mathcal{E}_{\tau_{0}}\right)^{-1}$ from $L^{2}\left(\mathbb{R}^{3 N}\right)$ to the form domain of $H_{\tau_{0}}$, which is $H^{1 / 2}\left(\mathbb{R}^{3 N}\right)$ in the ( $\mathbf{S R}$ ) case and $H^{1}\left(\mathbb{R}^{3 N}\right)$ in the (NR) case,

$$
R_{\tau_{0}}=O_{\left(C^{\infty}, H^{q}\right)}\left(\frac{1}{L^{4}}\right), \quad q= \begin{cases}1, & \text { in the (NR) case }  \tag{3.14}\\ \frac{1}{2}, & \text { in the (SR) case }\end{cases}
$$

and thus $\Phi_{\tau}$ is in the form domain of $H_{\tau}$. The $C^{\infty}$ estimate follows here for free, as $R_{\tau_{0}}$ does not depend on $\tau$. Note that, in the semirelativistic case, the uniformity with respect to $\tau_{0}$ for large $L$ follows from (1.8).

We choose to insert the family defined by (3.13) in (3.3). We find that

$$
\begin{equation*}
\tilde{\mathcal{E}}_{\tau}-E_{\infty}=\frac{\left\langle\Phi_{\tau},\left(H_{\tau}-E_{\infty}\right) \Phi_{\tau}\right\rangle}{\left\|\Phi_{\tau}\right\|^{2}} . \tag{3.15}
\end{equation*}
$$

We emphasize the fact that, on the right-hand side of (3.13), $\Psi_{\tau_{0}}$ and $R_{\tau_{0}}$ remain fixed when $\tau$ varies. This will help us to prove that $\tilde{\mathcal{E}}_{\tau}$ is a smooth function of $\tau$ and we will moreover provide quantitative information for its derivatives. We prove the following result:

Proposition 3.1. Let us assume that Conditions (Poles-lead) and (Neutr) are satisfied. We assume that $n_{1}+n_{2}<5$, where recall that $n_{1}$ and $n_{2}$ are the indices of the first non-vanishing multipole moments of the molecules. For $\tilde{\mathcal{E}}_{\tau}$ defined by (3.3) and (3.13), we have

$$
\begin{equation*}
\tilde{\mathcal{E}}_{\tau}=E_{\infty}+\frac{\mathcal{F}^{\left(n_{1}, n_{2}\right)}(U, V)}{L^{n_{1}+n_{2}+1}}+O_{C^{\infty}}\left(\frac{1}{L^{n_{1}+n_{2}+2}}\right) \tag{3.16}
\end{equation*}
$$

in the sense of (1.30)-(1.31).
Proof of Proposition 3.1. In view of (3.15), we begin by estimating $\left\|\Phi_{\tau}\right\|^{2}$. Using (3.13) together with (3.14) and the fact that $u_{\tau}$ commutes with $Q_{N}$, we find

$$
\begin{equation*}
\left\|\Phi_{\tau}\right\|^{2}=\left\|Q_{N} K_{\tau} W_{\tau} \Psi_{\tau_{0}}\right\|^{2}-2 \operatorname{Re}\left\langle R_{\tau_{0}}, Q_{N} u_{\tau} K_{\tau} W_{\tau} \Psi_{\tau_{0}}\right\rangle+O_{C^{\infty}}\left(\frac{1}{L^{8}}\right) \tag{3.17}
\end{equation*}
$$

From (2.1), (3.4) and (3.12), it follows that

$$
\begin{equation*}
\left(u_{\tau} K_{\tau} \Psi\right)\left(x_{1}, \ldots, x_{N}\right)=\left(\prod_{j=1}^{N} J_{m_{\tau}}\left(x_{j}\right)\right)\left(K_{\tau_{0}} \Psi\right)\left(l_{\tau}\left(x_{1}\right), \ldots, l_{\tau}\left(x_{N_{1}}\right), n_{\tau}\left(x_{N_{1}+1}\right), \ldots, n_{\tau}\left(x_{N}\right)\right) \tag{3.18}
\end{equation*}
$$

where

$$
\begin{equation*}
l_{\tau}(x):=\bar{U} U^{-1} m_{\tau}(x), \quad n_{\tau}(x):=\bar{V} V^{-1}\left(m_{\tau}(x)-L e_{1}\right)+L e_{1} . \tag{3.19}
\end{equation*}
$$

Using (3.7) and (3.19), it follows that

$$
\begin{equation*}
l_{\tau}(x)=x, \quad \forall x \in \operatorname{supp} \chi_{2 L}, \quad n_{\tau}(x)=x, \quad \forall x \in\left(\operatorname{supp} \chi_{2 L}-L e_{1}\right) \tag{3.20}
\end{equation*}
$$

Moreover, (3.7) implies the equality

$$
\begin{equation*}
J_{m_{\tau}}(x)=1, \quad \forall x \in \operatorname{supp} \chi_{2 L} \cup\left(\operatorname{supp} \chi_{2 L}-L e_{1}\right) \tag{3.21}
\end{equation*}
$$

Combining (3.18), (3.20) and (3.21), we arrive at

$$
\begin{equation*}
u_{\tau} K_{\tau} \Phi=K_{\tau_{0}} \Phi, \text { when } \operatorname{supp} \Phi \subset \operatorname{supp} \chi_{2 L}^{\otimes N} \tag{3.22}
\end{equation*}
$$

In particular, since by $(2.7) W_{\tau} \Psi_{\tau_{0}}$ is supported on $\left(\operatorname{supp} \chi_{\frac{4 L}{3}}\right)^{\otimes N} \subset \operatorname{supp} \chi_{2 L}^{\otimes N}$, (3.22) implies that

$$
\begin{equation*}
u_{\tau} K_{\tau} W_{\tau} \Psi_{\tau_{0}}=K_{\tau_{0}} W_{\tau} \Psi_{\tau_{0}} \tag{3.23}
\end{equation*}
$$

Using (2.22), (2.26), (3.14) and (3.23), we find

$$
\begin{equation*}
\left\langle R_{\tau_{0}}, Q_{N} u_{\tau} K_{\tau} W_{\tau} \Psi_{\tau_{0}}\right\rangle=\left\langle R_{\tau_{0}}, Q_{N} K_{\tau_{0}} \chi_{L}^{\otimes N} \Psi_{\tau_{0}}\right\rangle+O_{C^{\infty}}\left(\frac{1}{L^{7}}\right) \tag{3.24}
\end{equation*}
$$

Since by (3.2) $R_{\tau_{0}}=P_{\tau_{0}}^{\perp} R_{\tau_{0}}$, using (2.77) and (3.14), we see that

$$
\begin{equation*}
\left\langle R_{\tau_{0}}, Q_{N} K_{\tau_{0}} \Psi_{\tau_{0}}\right\rangle=\left\langle R_{\tau_{0}}, \Pi_{\tau_{0}}^{\perp} Q_{N} K_{\tau_{0}} \chi_{L}^{\otimes N} \Psi_{\tau_{0}}\right\rangle+O_{C^{\infty}}\left(\frac{1}{L^{7}}\right)=O_{C^{\infty}}\left(\frac{1}{L^{7}}\right) \tag{3.25}
\end{equation*}
$$

where in the last step we used that $\Pi_{\tau_{0}}^{\perp} Q_{N} K_{\tau_{0}} \chi_{L}^{\otimes N} \Psi_{\tau_{0}}=0$, which is immediate consequence of (2.58). The $C^{\infty}$ estimate follows for free, since the terms appearing in (3.25) are $\tau$ independent. From (3.17), (3.24) and (3.25), we find that

$$
\begin{equation*}
\left\|\Phi_{\tau}\right\|^{2}=\left\|Q_{N} K_{\tau} W_{\tau} \Psi_{\tau_{0}}\right\|^{2}+O_{C^{\infty}}\left(\frac{1}{L^{7}}\right) \tag{3.26}
\end{equation*}
$$

We now estimate $\left\langle\Phi_{\tau},\left(H_{\tau}-E_{\infty}\right) \Phi_{\tau}\right\rangle$. By (3.13), we can decompose

$$
\begin{align*}
\left\langle\Phi_{\tau},\left(H_{\tau}-E_{\infty}\right) \Phi_{\tau}\right\rangle= & \left\langle Q_{N} K_{\tau} W_{\tau} \Psi_{\tau_{0}},\left(H_{\tau}-E_{\infty}\right) Q_{N} K_{\tau} W_{\tau} \Psi_{\tau_{0}}\right\rangle  \tag{3.27}\\
& -2 \operatorname{Re}\left\langle u_{\tau}^{-1} R_{\tau_{0}},\left(H_{\tau}-E_{\infty}\right) Q_{N} K_{\tau} W_{\tau} \Psi_{\tau_{0}}\right\rangle  \tag{3.28}\\
& +\left\langle R_{\tau_{0}}, u_{\tau}\left(H_{\tau}-E_{\infty}\right) u_{\tau}^{-1} R_{\tau_{0}}\right\rangle . \tag{3.29}
\end{align*}
$$

Equations (2.82), (2.104) and (3.26) imply that

$$
\begin{align*}
\frac{\left\langle Q_{N} K_{\tau} W_{\tau} \Psi_{\tau_{0}},\left(H_{\tau}-E_{\infty}\right) Q_{N} K_{\tau} W_{\tau} \Psi_{\tau_{0}}\right\rangle}{\left\|\Phi_{\tau}\right\|^{2}} & =\frac{\left\langle Q_{N} K_{\tau} W_{\tau} \Psi_{\tau_{0}},\left(H_{\tau}-E_{\infty}\right) Q_{N} K_{\tau} W_{\tau} \Psi_{\tau_{0}}\right\rangle}{\left\|Q_{N} K_{\tau} W_{\tau} \Psi_{\tau_{0}}\right\|^{2}+O_{C^{\infty}}\left(\frac{1}{L^{7}}\right)} \\
& =\frac{\mathcal{F}^{\left(n_{1}, n_{2}\right)}(U, V)}{L^{n_{1}+n_{2}+1}}+O_{C^{\infty}}\left(\frac{1}{L^{n_{1}+n_{2}+2}}\right) \tag{3.30}
\end{align*}
$$

Note that this is the only point where we use Condition (Poles) in the proof of Proposition 3.1.
We will now estimate the terms (3.28) and (3.29). We begin with (3.28). By (2.63) and the fact that $u_{\tau}$ is unitary and commutes with $Q_{N}$, we have

$$
\begin{align*}
\left\langle u_{\tau}^{-1} R_{\tau_{0}},\left(H_{\tau}-E_{\infty}\right) Q_{N} K_{\tau} W_{\tau} \Psi_{\tau_{0}}\right\rangle & =\left\langle R_{\tau_{0}}, Q_{N} u_{\tau} K_{\tau}\left(H_{\infty}+\widetilde{I}_{\tau}-E_{\infty}\right) W_{\tau} \Psi_{\tau_{0}}\right\rangle \\
& =\left\langle R_{\tau_{0}}, u_{\tau} K_{\tau} \chi_{2 L}^{\otimes N}\left(H_{\infty}+\widetilde{I}_{\tau}-E_{\infty}\right) W_{\tau} \Psi_{\tau_{0}}\right\rangle  \tag{3.31}\\
& +\left\langle R_{\tau_{0}}, u_{\tau} K_{\tau}\left(1-\chi_{2 L}^{\otimes N}\right)\left(H_{\infty}+\widetilde{I}_{\tau}-E_{\infty}\right) W_{\tau} \Psi_{\tau_{0}}\right\rangle \tag{3.32}
\end{align*}
$$

where in the last step we used the equality $Q_{N} R_{\tau_{0}}=R_{\tau_{0}}$ and the decomposition $1=\chi_{2 L}^{\otimes N}+\left(1-\chi_{2 L}^{\otimes N}\right)$. On the one hand, using (3.22), we find

$$
\begin{equation*}
\left\langle R_{\tau_{0}}, u_{\tau} K_{\tau} \chi_{2 L}^{\otimes N}\left(H_{\infty}+\widetilde{I}_{\tau}-E_{\infty}\right) W_{\tau} \Psi_{\tau_{0}}\right\rangle=\left\langle R_{\tau_{0}}, K_{\tau_{0}} \chi_{2 L}^{\otimes N}\left(H_{\infty}+\widetilde{I}_{\tau}-E_{\infty}\right) W_{\tau} \Psi_{\tau_{0}}\right\rangle \tag{3.33}
\end{equation*}
$$

The only part which depends on $\tau$ in the right-hand side of (3.33) is $\left(H_{\infty}+\widetilde{I}_{\tau}-E_{\infty}\right) W_{\tau} \Psi_{\tau_{0}}$. Therefore, inserting (2.70) and (3.14) in (3.33), we find

$$
\begin{equation*}
\left\langle R_{\tau_{0}}, K_{\tau_{0}} \chi_{2 L}^{\otimes N}\left(H_{\infty}+\widetilde{I}_{\tau}-E_{\infty}\right) W_{\tau} \Psi_{\tau_{0}}\right\rangle=O_{C^{\infty}}\left(\frac{1}{L^{7}}\right) \tag{3.34}
\end{equation*}
$$

Let us now bound (3.32). Here, we have to distinguish between the two types of kinetic energy. In the (NR) case, since $W_{\tau} \Psi_{\tau_{0}}$ and $\left(1-\chi_{2 L}^{\otimes N}\right)$ have disjoint supports and the Hamiltonian is a local operator, we immediately see that this term is 0 . For the (SR) case, we need to work more. The only nonlocal terms in the Hamiltonian are the kinetic energy operators. Furthermore, the functions $W_{\tau} \Psi_{\tau_{0}}$
and $\left(1-\chi_{2 L}^{\otimes N}\right) K_{\tau}^{-1} u_{\tau}^{-1} R_{\tau_{0}}$ are antisymmetric with respect to permutations of variables leaving invariant the set $\left\{1, \ldots, N_{1}\right\}$. Therefore, we can write

$$
\begin{align*}
\left\langle R_{\tau_{0}}, u_{\tau} K_{\tau}\left(1-\chi_{2 L}^{\otimes N}\right)\left(H_{\infty}+\widetilde{I}_{\tau}-E_{\infty}\right) W_{\tau} \Psi_{\tau_{0}}\right\rangle= & N_{1}\left\langle R_{\tau_{0}}, u_{\tau} K_{\tau}\left(1-\chi_{2 L}^{\otimes N}\right) T_{1} W_{\tau} \Psi_{\tau_{0}}\right\rangle  \tag{3.35}\\
& +N_{2}\left\langle R_{\tau_{0}}, u_{\tau} K_{\tau}\left(1-\chi_{2 L}^{\otimes N}\right) T_{N} W_{\tau} \Psi_{\tau_{0}}\right\rangle \tag{3.36}
\end{align*}
$$

where recall that $T_{i}$ is the semirelativistic kinetic energy operator applied to the $i$-th electron.
We begin by taking some fixed $\Phi \in L^{2}\left(\mathbb{R}^{3 N}\right)$. By disjointness of supports of $\chi_{4 L / 3}^{\otimes N}$ and of $1-\chi_{2 L}^{\otimes N}$, we see with Lemma A. 1 that $\left(1-\chi_{2 L}^{\otimes N}\right) T_{1} \chi_{4 L / 3}^{\otimes N} \Phi$ is in $L^{2}\left(\mathbb{R}^{3 N}\right)$. Moreover, (A.1) implies that, for almost all $\left(x_{1}, X\right) \in \mathbb{R}^{3} \times \mathbb{R}^{3(N-1)}$ where $X:=\left(x_{2}, \ldots, x_{N}\right)$, we have

$$
\begin{equation*}
\left(\left(1-\chi_{2 L}^{\otimes N}\right) T_{1} \chi_{4 L / 3}^{\otimes N} \Phi\right)\left(x_{1}, X\right)=-\frac{1-\chi_{2 L}\left(x_{1}\right)}{2 \pi^{2}} \int_{\mathbb{R}^{3}} \frac{K_{2}\left(\left|x_{1}-y\right|\right)}{\left|x_{1}-y\right|^{2}}\left(\chi_{4 L / 3}^{\otimes N} \Phi\right)(y, X) \mathrm{d} y \tag{3.37}
\end{equation*}
$$

where $K_{2}$ is the modified Bessel function of the second kind. Note that we could replace $1-\chi_{2 L}^{\otimes N}\left(x_{1}, X\right)$ with $1-\chi_{2 L}\left(x_{1}\right)$ in the right hand side of (3.37) due to (3.5). Equation (3.37) together with (3.18), (3.20) and (3.21) implies
$u_{\tau} K_{\tau}\left(\left(1-\chi_{2 L}^{\otimes N}\right) T_{1} \chi_{4 L / 3}^{\otimes N} \Phi\right)\left(x_{1}, X\right)=-\frac{1-\chi_{2 L}\left(x_{1}\right)}{2 \pi^{2}} J_{m_{\tau}}\left(x_{1}\right) \int_{\mathbb{R}^{3}} \frac{K_{2}\left(\left|l_{\tau}\left(x_{1}\right)-y\right|\right)}{\left|l_{\tau}\left(x_{1}\right)-y\right|^{2}}\left[K_{\tau_{0}}\left(\chi_{4 L / 3}^{\otimes N} \Phi\right)\right](y, X) \mathrm{d} y$.

Let us now consider the integral in the right-hand side of (3.38). Note that, by (2.5), for $\left|x_{1}\right| \leqslant L / 5$, we have $1-\chi_{2 L}\left(x_{1}\right)=0$ and thus the right-hand side of (3.38) is zero. Therefore, we consider only the case $\left|x_{1}\right| \geqslant L / 5$. This implies that

$$
\begin{equation*}
\left|l_{\tau}\left(x_{1}\right)\right| \geqslant \frac{L}{5} \tag{3.39}
\end{equation*}
$$

because $l_{\tau}$ is bijective and, due to $(3.20), l_{\tau}(B(0, L / 5))=B(0, L / 5)$.
On the other hand, $\left[K_{\tau_{0}}\left(\chi_{4 L / 3}^{\otimes N} \Phi\right)\right](y, X)$ is not 0 only if $y$ is in the support of $\chi_{4 L / 3}$, i.e., when $|y| \leqslant L / 6$. Therefore by (3.39), it follows that $\left|l_{\tau}\left(x_{1}\right)-y\right| \geqslant \frac{L}{30}$ in all cases where the right-hand side of (3.38) is not 0 .

In order to prove regularity, let us choose arbitrary $A, B \in \Gamma$, as defined in (1.29), and introduce

$$
\begin{equation*}
\tau_{t}:=\left(L, e^{A t} \bar{U}, e^{B t} \bar{V}\right), \quad t \in \mathbb{R} \tag{3.40}
\end{equation*}
$$

Recall that $\Omega$ is a neighborhood of $(\bar{U}, \bar{V})$ such that (3.11) holds. There exists $\eta>0$ such that, for all $t \in(-\eta, \eta)$ and all $A, B \in \Gamma,\left(e^{A t} \bar{U}, e^{B t} \bar{V}\right) \in \Omega$, so we can apply what we have proven until now. Using [1, Formula 9.6.29] and [1, Formula 9.7.2], we find that, for all $n \in \mathbb{N}_{0}:=\mathbb{N} \cup\{0\}$, there exists $C_{n}>0$ such that for all $z \geqslant 1$

$$
\left|K_{2}^{(n)}(z)\right| \leqslant C_{n} e^{-z / 2}
$$

Moreover, by (3.4) and (3.19), for all $n \in \mathbb{N}_{0}$, there exist $D_{n}>0$, independent of $A$ and $B$, such that for all $t \in(-\eta, \eta)$,

$$
\left|\frac{\mathrm{d}^{n}}{\mathrm{~d} t^{n}} l_{\tau_{t}}(x)\right| \leqslant D_{n}(1+|x|), \quad \forall x \in \mathbb{R}^{3} .
$$

As a consequence, for all $x_{1}$ and $y$ such that $\left|x_{1}-y\right| \geqslant L / 30$, the function

$$
\begin{equation*}
g_{\left(x_{1}, y\right)}(t):=\frac{K_{2}\left(\left|l_{\tau_{t}}\left(x_{1}\right)-y\right|\right)}{\left|l_{\tau_{t}}\left(x_{1}\right)-y\right|^{2}}, \tag{3.41}
\end{equation*}
$$

is infinitely differentiable. Moreover, there exists a $d>0$ such that, for all $n \in \mathbb{N}_{0}$, there exists $M_{n}>0$, such that, for all $L$ large enough, for all $\left|x_{1}\right| \geqslant L / 5$ and $|y| \leqslant L / 6$, all $A, B \in \Gamma$ and all $t \in(-\eta, \eta)$,

$$
\begin{equation*}
\left|g_{\left(x_{1}, y\right)}^{(n)}(t)\right| \leqslant M_{n} e^{-\left|l_{\tau_{t}}\left(x_{1}\right)-y\right| / 2} \leqslant M_{n} e^{-d\left|x_{1}-y\right|} \tag{3.42}
\end{equation*}
$$

where, in the last step, we used the fact that, if $\left|x_{1}\right| \geqslant L / 5$ and $|y| \leqslant L / 6$, then $\left|l_{\tau_{t}}\left(x_{1}\right)-y\right| \geqslant \frac{1}{2}\left|x_{1}-y\right|$, provided that $t$ is in a small neighborhood of zero, which is independent of $A$ and $B$. In fact, we have $\left|l_{\tau_{t}}\left(x_{1}\right)-y\right| \geqslant \frac{1}{11}\left|x_{1}-y\right|$ uniformly for all times. Moreover, from (3.4), it follows that

$$
\begin{equation*}
J_{\tau}=O_{\left(C^{\infty}, L^{\infty}\right)}(1) \tag{3.43}
\end{equation*}
$$

Equation (3.41) together with (3.42) enables the application of the dominated convergence theorem in (3.38) and together with (3.43) we can extract that

$$
\begin{equation*}
u_{\tau} K_{\tau}\left(1-\chi_{2 L}^{\otimes N}\right) T_{1} \chi_{4 L / 3}^{\otimes N}=O_{\left(C^{\infty}, \mathcal{B}\left(L^{2}\right)\right)}\left(e^{-d L}\right) \tag{3.44}
\end{equation*}
$$

because of Young's inequality and the fact that $\left\|e^{-d|\cdot|} \chi_{|x| \geqslant \frac{L}{30}}\right\|_{L^{1}}$ is exponentially decaying in $L$.
From (3.44), (2.26) and the Leibniz rule, we find that,

$$
\begin{equation*}
\left\langle R_{\tau_{0}}, u_{\tau} K_{\tau}\left(1-\chi_{2 L}^{\otimes N}\right) T_{1} W_{\tau} \Psi_{\tau_{0}}\right\rangle=O_{C^{\infty}}\left(e^{-d L}\right) \tag{3.45}
\end{equation*}
$$

This bounds the derivatives of the term (3.35). It similarly follows that

$$
\begin{equation*}
\left\langle R_{\tau_{0}}, u_{\tau} K_{\tau}\left(1-\chi_{2 L}^{\otimes N}\right) T_{N} W_{\tau} \Psi_{\tau_{0}}\right\rangle=O_{C \infty}\left(e^{-d L}\right) \tag{3.46}
\end{equation*}
$$

Inserting (3.45) and (3.46) in (3.35)-(3.36), we find that

$$
\begin{equation*}
\left\langle R_{\tau_{0}}, u_{\tau} K_{\tau}\left(1-\chi_{2 L}^{\otimes N}\right)\left(H_{\infty}+\widetilde{I}_{\tau}-E_{\infty}\right) W_{\tau} \Psi_{\tau_{0}}\right\rangle=O_{C \infty}\left(e^{-d L}\right) \tag{3.47}
\end{equation*}
$$

Inserting (3.34) and (3.47) into (3.31)-(3.32), we find that

$$
\begin{equation*}
\left\langle u_{\tau}^{-1} R_{\tau_{0}},\left(H_{\tau}-E_{\infty}\right) Q_{N} K_{\tau} W_{\tau} \Psi_{\tau_{0}}\right\rangle=O_{C^{\infty}}\left(\frac{1}{L^{7}}\right) \tag{3.48}
\end{equation*}
$$

for both (NR) and the (SR) cases, because in the (NR) case the left-hand sides of (3.45) and (3.46) are zero, as we discussed above.

Let us now bound (3.29). Here, we treat the (NR) and (SR) cases simultaneously. We directly see with (1.5) that

$$
\begin{equation*}
u_{\tau}\left(H_{\tau}-E_{\infty}\right) u_{\tau}^{-1}=\sum_{i=1}^{N} u_{\tau} T_{i} u_{\tau}^{-1}+\alpha u_{\tau} I_{N}(Y(\tau), \mathcal{Z}) u_{\tau}^{-1}-E_{\infty} \tag{3.49}
\end{equation*}
$$

Let us begin with the potentials, which are the same in the (NR) and in the (SR) cases. We start with the electron-electron interactions. From (3.12) and (3.8), we see that, for all $x_{i}$ and $x_{j}, i, j=1, \ldots, N$, $i \neq j$ and for all $\tau$ close enough to $\tau_{0}$,

$$
u_{\tau} \frac{1}{\left|x_{i}-x_{j}\right|} u_{\tau}^{-1}=\frac{1}{\left|x_{i}-x_{j}+v_{\tau}\left(x_{i}\right)-v_{\tau}\left(x_{j}\right)\right|}=(1+\omega)^{-1 / 2} \frac{1}{\left|x_{i}-x_{j}\right|}
$$

where

$$
\omega:=\frac{1}{\left|x_{i}-x_{j}\right|^{2}}\left(2\left(x_{i}-x_{j}\right) \cdot\left(v_{\tau}\left(x_{i}\right)-v_{\tau}\left(x_{j}\right)\right)+\left|v_{\tau}\left(x_{i}\right)-v_{\tau}\left(x_{j}\right)\right|^{2}\right)
$$

Since $\omega$ is a polynomial function of values of $v_{\tau}$, it depends smoothly on $U$ and $V$. Moreover, due to (3.10), there exists an open neighborhood $\Omega^{\prime}$ of $(\bar{U}, \bar{V})$ such that, for all $x, y \in \mathbb{R}^{3}$ and all $(U, V) \in \Omega^{\prime}$, $\left|v_{\tau}(x)-v_{\tau}(y)\right|<|x-y| / 3$, which implies that $|\omega|<7 / 9$. Consequently, $(1+\omega)^{-1 / 2}$ is a smooth function of $U$ and $V$ on $\Omega^{\prime}$. Note that every derivative of $v_{\tau}$ with respect to $x$ is bounded uniformly in $L>1$. As a consequence, for all $n \in \mathbb{N}_{0}$, there exists $C_{n}>0$ such that, for all $t$ close enough to 0 and all $A, B \in \Gamma$,

$$
\begin{equation*}
\left|\frac{\mathrm{d}^{n}}{\mathrm{~d} t^{n}} u_{\tau_{t}} \frac{1}{\left|x_{i}-x_{j}\right|} u_{\tau_{t}}^{-1}\right| \leqslant \frac{C_{n}}{\left|x_{i}-x_{j}\right|}, \quad \forall i, j \in\{1, \ldots, N\}, \text { with } i \neq j \tag{3.50}
\end{equation*}
$$

where recall that $\tau_{t}$ was defined in (3.40). To deal with the electron-nucleus attractions, we observe, with the help of (3.7), (1.21) and recalling that in the current setting $U=e^{t A} \bar{U}, V=e^{t B} \bar{V}$, that

$$
\begin{equation*}
Y(\tau)=\left(m_{\tau_{t}}\left(y_{1}\right), \ldots, m_{\tau_{t}}\left(y_{M}\right)\right) \tag{3.51}
\end{equation*}
$$

Moreover,

$$
u_{\tau_{t}} \frac{1}{\left|x_{i}-m_{\tau_{t}}\left(y_{j}\right)\right|} u_{\tau_{t}}^{-1}=\frac{1}{\left|m_{\tau_{t}}\left(x_{i}\right)-m_{\tau_{t}}\left(y_{j}\right)\right|},
$$

so we can argue as in the proof of (3.50) to conclude that, for all $n \in \mathbb{N}_{0}$, there exists $D_{n}>0$ such that, for all $t$ close enough to 0 ,

$$
\begin{equation*}
\left|\frac{\mathrm{d}^{n}}{\mathrm{~d} t^{n}} u_{\tau_{t}} \frac{1}{\left|x_{i}-m_{\tau}\left(y_{j}\right)\right|} u_{\tau_{t}}^{-1}\right| \leqslant \frac{D_{n}}{\left|x_{i}-y_{j, 0}\right|}, \quad \forall i \in\{1, \ldots, N\}, \forall j \in\{1, \ldots, M\} \tag{3.52}
\end{equation*}
$$

where $y_{j, 0}$ is the position of the $j$-th nucleus in the configuration $Y\left(\tau_{0}\right)$. We now consider the interactions between the nuclei. Since we can write that

$$
\left\langle R_{\tau_{0}}, u_{\tau_{t}} \frac{1}{\left|m_{\tau_{t}}\left(y_{k}\right)-m_{\tau_{t}}\left(y_{l}\right)\right|} u_{\tau_{t}}^{-1} R_{\tau_{0}}\right\rangle=\frac{1}{\left|m_{\tau_{t}}\left(y_{k}\right)-m_{\tau_{t}}\left(y_{l}\right)\right|}\left\langle R_{\tau_{0}}, R_{\tau_{0}}\right\rangle,
$$

we obtain that

$$
\begin{equation*}
\left\langle R_{\tau_{0}}, u_{\tau} \frac{1}{\left|m_{\tau}\left(y_{k}\right)-m_{\tau}\left(y_{l}\right)\right|} u_{\tau}^{-1} R_{\tau_{0}}\right\rangle=O_{C^{\infty}}\left(\frac{1}{L^{8}}\right) . \tag{3.53}
\end{equation*}
$$

Using (1.6), (3.50), (3.51), (3.52) and (3.53), we can derive that

$$
\begin{equation*}
\left\langle R_{\tau_{0}}, u_{\tau} I_{N}(Y(\tau), \mathcal{Z}) u_{\tau}^{-1} R_{\tau_{0}}\right\rangle \leqslant O_{C^{\infty}}\left(\left\|R_{\tau_{0}}\right\|_{H^{\frac{1}{2}}}^{2}\right) \tag{3.54}
\end{equation*}
$$

where in the last step we also used Kato's inequality $\frac{1}{|x|} \leqslant \frac{\pi}{2} \sqrt{1-\Delta}$, see [34, Theorem 2.1]. Using (3.54) and (3.14), we find that

$$
\begin{equation*}
\left\langle R_{\tau_{0}}, u_{\tau} I_{N}(Y(\tau), \mathcal{Z}) u_{\tau}^{-1} R_{\tau_{0}}\right\rangle=O_{C^{\infty}}\left(\frac{1}{L^{8}}\right) \tag{3.55}
\end{equation*}
$$

for both the (NR) and (SR) cases.
Finally, we bound the first term of the right-hand side of (3.49). Let us take $t \in \mathbb{R}$ and $A, B \in \Gamma$. For $t$ small enough so that $\left(e^{t A}, e^{t B}\right) \in \Omega$, we define the following operator on $L^{2}\left(\mathbb{R}^{3}\right)$ :

$$
D(t):=1-u_{\tau_{t}} \Delta u_{\tau_{t}}^{-1}=u_{\tau_{t}}(1-\Delta) u_{\tau_{t}}^{-1}
$$

where recall that $\tau_{t}$ was defined in (3.40). Since, for all $t, u_{\tau_{t}}$ is unitary and $-\Delta \geqslant 0$, we have that $D(t) \geqslant 1$ in the sense of quadratic forms. We find, following [37, Theorem 1], that there exist functions $\left(a_{t}^{i j}\right),\left(b_{t}^{k}\right), c_{t}$, which are all $C_{c}^{\infty}$ with compact support uniformly in $t$ and depend smoothly on $t$, such that

$$
\begin{equation*}
D(t)=-\Delta+\sum_{i, j} a_{t}^{i j}(x) \frac{\partial}{\partial x_{i}} \frac{\partial}{\partial x_{j}}+\sum_{k} b_{t}^{k}(x) \frac{\partial}{\partial x_{k}}+c_{t}(x) . \tag{3.56}
\end{equation*}
$$

Moreover, when $t$ goes to 0 , all the functions $a_{t}^{i j}, b_{t}^{k}, c_{t}$ converge to 0 in $C^{n}$ norm for all $n \in \mathbb{N}_{0}$, uniformly in $A, B \in \Gamma$. This together with (3.14) implies that, in the (NR) case, each of the terms $\left\langle R_{\tau_{0}}, u_{\tau} T_{i} u_{\tau}^{-1} R_{\tau_{0}}\right\rangle$ depends smoothly on $t$. Moreover, for all $n \geqslant 1$,

$$
\frac{\mathrm{d}^{n}}{\mathrm{~d} t^{n}}\left\langle R_{\tau_{0}}, u_{\tau_{t}} T_{i} u_{\tau_{t}}^{-1} R_{\tau_{0}}\right\rangle=\sum_{i, j}\left\langle R_{\tau_{0}}, \frac{\mathrm{~d}^{n} a_{t}^{i j}}{\mathrm{~d} t^{n}} \frac{\partial}{\partial x_{i}} \frac{\partial}{\partial x_{j}} R_{\tau_{0}}\right\rangle+\sum_{k}\left\langle R_{\tau_{0}}, \frac{\mathrm{~d}^{n} b_{t}^{k}}{\mathrm{~d} t^{n}} \frac{\partial}{\partial x_{k}} R_{\tau_{0}}\right\rangle+\left\langle R_{\tau_{0}}, \frac{\mathrm{~d}^{n} c_{t}}{\mathrm{~d} t^{n}} R_{\tau_{0}}\right\rangle
$$

As a consequence, in the (NR) case, we immediately obtain that

$$
\begin{equation*}
\left\langle R_{\tau_{0}}, u_{\tau} T_{i} u_{\tau}^{-1} R_{\tau_{0}}\right\rangle=O_{C^{\infty}}\left(\frac{1}{L^{8}}\right), \tag{3.57}
\end{equation*}
$$

due to (3.14).
For the (SR) case, we use the well known identity $\sqrt{u}=\frac{1}{\pi} \int_{0}^{\infty} \frac{u}{u+w} \frac{\mathrm{~d} w}{\sqrt{w}}$ for all $u>0$. It follows, by the spectral theorem, that for all $t$

$$
\begin{equation*}
u_{\tau_{t}} \sqrt{1-\Delta} u_{\tau_{t}}^{-1}=\sqrt{D(t)}=\frac{1}{\pi} \int_{0}^{\infty} D(t)(D(t)+w)^{-1} \frac{\mathrm{~d} w}{\sqrt{w}} . \tag{3.58}
\end{equation*}
$$

Therefore, with the help of the product rule, we find

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{~d} t} \sqrt{D(t)}=\frac{1}{\pi} \int_{0}^{\infty} D^{\prime}(t)(D(t)+w)^{-1} \frac{\mathrm{~d} w}{\sqrt{w}}-\frac{1}{\pi} \int_{0}^{\infty} D(t)(D(t)+w)^{-1} D^{\prime}(t)(D(t)+w)^{-1} \frac{\mathrm{~d} w}{\sqrt{w}} \tag{3.59}
\end{equation*}
$$

in the sense of quadratic forms on $H^{\frac{1}{2}}\left(\mathbb{R}^{3}\right)$. A rigorous justification can be achieved if we first consider quadratic forms on $H^{2}\left(\mathbb{R}^{3}\right)$ and apply the dominated convergence theorem. We omit the details. Observe now that, due to (3.58),

$$
\begin{equation*}
\frac{1}{\pi} \int_{0}^{\infty} D^{\prime}(t)(D(t)+w)^{-1} \frac{\mathrm{~d} w}{\sqrt{w}}=D^{\prime}(t) D^{-1}(t) \sqrt{D(t)} \tag{3.60}
\end{equation*}
$$

Using the decomposition $D(t)=(D(t)+w)-w$ and (3.60), we find that

$$
\begin{align*}
\frac{1}{\pi} \int_{0}^{\infty} D(t)(D(t)+w)^{-1} & D^{\prime}(t)(D(t)+w)^{-1} \frac{\mathrm{~d} w}{\sqrt{w}} \\
& =D^{\prime}(t) D^{-1}(t) \sqrt{D(t)}-\frac{1}{\pi} \int_{0}^{\infty}(D(t)+w)^{-1} D^{\prime}(t)(D(t)+w)^{-1} \sqrt{w} \mathrm{~d} w \tag{3.61}
\end{align*}
$$

Combining (3.59)-(3.61), we obtain that

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{~d} t} \sqrt{D(t)}=\frac{1}{\pi} \int_{0}^{\infty}(D(t)+w)^{-1} D^{\prime}(t)(D(t)+w)^{-1} \sqrt{w} \mathrm{~d} w \tag{3.62}
\end{equation*}
$$

Using the commutator formula $\left[A^{-1}, C\right]=A^{-1}[C, A] A^{-1}$, we find

$$
(D(t)+w)^{-1} D^{\prime}(t)=D^{\prime}(t)(D(t)+w)^{-1}+(D(t)+w)^{-1}\left[D^{\prime}(t), D(t)\right](D(t)+w)^{-1}
$$

and, using the same commutator formula once more, we arrive at

$$
\begin{align*}
(D(t)+w)^{-1} D^{\prime}(t)=D^{\prime}(t)(D(t)+w)^{-1}+ & {\left[D^{\prime}(t), D(t)\right](D(t)+w)^{-2} } \\
& +(D(t)+w)^{-1}\left[\left[D^{\prime}(t), D(t)\right], D(t)\right](D(t)+w)^{-2} \tag{3.63}
\end{align*}
$$

Inserting the equalities $\int_{0}^{\infty}(D(t)+w)^{-2} \sqrt{w} \mathrm{~d} w=\frac{\pi}{2}(D(t))^{-\frac{1}{2}}$ and $\int_{0}^{\infty}(D(t)+w)^{-3} \sqrt{w} \mathrm{~d} w=\frac{3 \pi}{8}(D(t))^{-\frac{3}{2}}$ together with (3.63) in (3.62), we find that

$$
\frac{\mathrm{d}}{\mathrm{~d} t} \sqrt{D(t)}=F(t) \sqrt{D(t)}+C(t)
$$

where

$$
F(t):=\frac{1}{2} D^{\prime}(t) D^{-1}(t)+\frac{3}{8}\left[D^{\prime}(t), D(t)\right] D(t)^{-2}
$$

and

$$
C(t):=\frac{1}{\pi} \int_{0}^{\infty}(D(t)+w)^{-1}\left[\left[D^{\prime}(t), D(t)\right], D(t)\right](D(t)+w)^{-3} \sqrt{w} \mathrm{~d} w
$$

But for all $s \in \mathbb{R}$ the operator $D^{\prime}(t)$ is bounded from $H^{s+2}$ to $H^{s}$ due to (3.56), as it is a second order differential operator with $C_{c}^{\infty}$ coefficients. On the other hand, for $t$ in a neighborhood of 0 independent of the matrices $A$ and $B$, the operator $D(t)$ is elliptic, since it is a small perturbation of the Laplacian. Thus, by elliptic regularity, $D^{-1}(t)$ is bounded from $H^{s}$ to $H^{s+2}$ for all $s \in \mathbb{R}$. Consequently, the operator $F(t)$ is bounded from $H^{1 / 2}\left(\mathbb{R}^{3}\right)$ to itself for $t$ small enough.

Concerning $C(t)$, one can see that the second commutator $\left[\left[D^{\prime}(t), D(t)\right], D(t)\right]$ is a 4 -th order differential operator and it is bounded from $H^{s}\left(\mathbb{R}^{3}\right)$ to $H^{s-4}\left(\mathbb{R}^{3}\right)$, for all $s \in \mathbb{R}$. Consequently, the operator $D(t)^{-1}\left[\left[D^{\prime}(t), D(t)\right], D(t)\right] D(t)^{-1}$ is bounded from $H^{1 / 2}\left(\mathbb{R}^{3}\right)$ to itself. Moreover, $(D(t)+w)^{-1} D(t) \leqslant 1$ for all $w \geqslant 0$. There exists, therefore, a $\Theta>0$ such that, for all $\psi \in L^{2}\left(\mathbb{R}^{3}\right)$,

$$
\|C(t) \psi\| \leqslant \Theta\left\|\int_{0}^{\infty}(D(t)+w)^{-2} \sqrt{w} \mathrm{~d} w \psi\right\|
$$

Since $D(t) \geqslant 1$ for all $t$ and $\int_{0}^{\infty}(1+w)^{-2} \sqrt{w} \mathrm{~d} w<\infty$, the operator $C(t)$ is bounded.
Let $F_{1}(t):=F(t)$ and $C_{1}(t):=C(t)$. In order to bound the $n$-th derivative of $\sqrt{D(t)}$ for all $n \in \mathbb{N}$, we define inductively $F_{n+1}(t):=F_{n}^{\prime}(t)+F_{n}(t) F(t)$ and $C_{n+1}(t):=F_{n}(t) C(t)+C_{n}^{\prime}(t)$. Then, for each $n$, $F_{n}$ and $C_{n}$ are smooth families of operators. We get, by induction, that:

$$
\frac{\mathrm{d}^{n}}{\mathrm{~d} t^{n}} \sqrt{D(t)}=F_{n}(t) \sqrt{D(t)}+C_{n}(t)
$$

Moreover, for each $t$ close enough to $0, F_{n}(t)$ is bounded from $H^{1 / 2}\left(\mathbb{R}^{3}\right)$ to itself and $C_{n}(t)$ is bounded from $L^{2}\left(\mathbb{R}^{3}\right)$ to itself, as one can show by induction too, since $D^{(n)}(t)$ is bounded from $H^{s+2}\left(\mathbb{R}^{3}\right)$ to $H^{s}\left(\mathbb{R}^{3}\right)$ for all $s \in \mathbb{R}$ and $n \in \mathbb{N}$. This implies that there exists a constant $\Lambda_{n}$ such that

$$
\begin{equation*}
\left|\frac{\mathrm{d}^{n}}{\mathrm{~d} t^{n}}\left\langle R_{\tau_{0}}, \sqrt{D(t)} R_{\tau_{0}}\right\rangle\right| \leqslant \Lambda_{n}\left\|R_{\tau_{0}}\right\|_{H^{1 / 2}}^{2}=O\left(\frac{1}{L^{8}}\right) \tag{3.64}
\end{equation*}
$$

where we used again (3.14). Inserting (3.55) and (3.57) (in the (NR) case) or (3.64) (in the (SR) case) in (3.49), we find that

$$
\begin{equation*}
\left\langle R_{\tau_{0}}, u_{\tau}\left(H_{\tau}-E_{\infty}\right) u_{\tau}^{-1} R_{\tau_{0}}\right\rangle=O_{C^{\infty}}\left(\frac{1}{L^{8}}\right) \tag{3.65}
\end{equation*}
$$

Finally, inserting (3.30), (3.48) and (3.65) in (3.27)-(3.29), we find that

$$
\begin{equation*}
\frac{\left\langle\Phi_{\tau},\left(H_{\tau}-E_{\infty}\right) \Phi_{\tau}\right\rangle}{\left\|\Phi_{\tau}\right\|^{2}}=\frac{\mathcal{F}^{\left(n_{1}, n_{2}\right)}(U, V)}{L^{n_{1}+n_{2}+1}}+O_{C^{\infty}}\left(\frac{1}{L^{n_{1}+n_{2}+2}}\right)+\frac{1}{\left\|\Phi_{\tau}\right\|^{2}} O_{C^{\infty}}\left(\frac{1}{L^{7}}\right) \tag{3.66}
\end{equation*}
$$

which together with $(2.72),(3.26)$ and the hypothesis $n_{1}+n_{2}<5$ leads to (3.16). This concludes the proof of Proposition 3.1.

To conclude the proof of Proposition 1.19, we use the fact that $\tilde{\mathcal{E}}_{\tau}$ has a local pseudo-minimum at $(\bar{U}, \bar{V})$. Thus, if we take any matrices $A$ and $B$ as in (1.53), and define $\tau_{t}:=\left(L, e^{t A} \bar{U}, e^{t B} \bar{V}\right)$, we find that

$$
\left.\frac{\mathrm{d} \tilde{\mathcal{E}}_{\tau_{t}}}{\mathrm{~d} t}\right|_{t=0}=0,\left.\quad \frac{\mathrm{~d}^{2} \tilde{\mathcal{E}}_{\tau_{t}}}{\mathrm{~d} t^{2}}\right|_{t=0} \geqslant 0
$$

This implies, together with Proposition 3.1, that, for any $\delta>0$, we can find an $L_{m}>0$ such that, for all $L \geqslant L_{m}$,

$$
\left|\frac{\mathrm{d} \mathcal{F}^{\left(n_{1}, n_{2}\right)}\left(e^{t A} \bar{U}, e^{t B} \bar{V}\right)}{\mathrm{d} t}\right|_{t=0}\left|\leqslant \delta, \quad \frac{\mathrm{~d}^{2} \mathcal{F}^{\left(n_{1}, n_{2}\right)}\left(e^{t A} \bar{U}, e^{t B} \bar{V}\right)}{\mathrm{d} t^{2}}\right|_{t=0} \geqslant-\delta
$$

for all $A, B \in \Gamma$. This concludes the proof of Proposition 1.19.

## Appendix A. Some properties of the semirelativistic kinetic energy operator

In this appendix, we give some properties of the semirelativistic kinetic energy operator $T:=\sqrt{1-\Delta}-$ 1 on $L^{2}\left(\mathbb{R}^{3}\right)$, as defined in (1.3). Since it is a pseudodifferential and not a differential operator, it is not as easy to deal with as the Laplace operator. First, it is not local. The following lemma enables us to control this non-locality.
Lemma A.1. There exists $C>0$ such that, for $R$ large enough, the following holds: for all $f \in L^{2}\left(\mathbb{R}^{3}\right)$, the linear form $T f \in H^{-1}\left(\mathbb{R}^{3}\right)$ can be extended to a bounded linear form on

$$
\left\{g \in L^{2}\left(\mathbb{R}^{3}\right), \operatorname{dist}(\operatorname{supp} f, \operatorname{supp} g) \geqslant R\right\}
$$

with norm less than $C e^{-R}\|f\|_{L^{2}}$. Moreover, for all $f, g \in L^{2}\left(\mathbb{R}^{3}\right)$ such that $\operatorname{dist}(\operatorname{supp} f, \operatorname{supp} g) \geqslant R$,

$$
\begin{equation*}
\langle f, T g\rangle=-\frac{1}{2 \pi^{2}} \int_{\mathbb{R}^{3}} \int_{\mathbb{R}^{3}} \frac{K_{2}(|x-y|)}{|x-y|^{2}} \overline{f(x)} g(y) \mathrm{d} x \mathrm{~d} y \tag{A.1}
\end{equation*}
$$

where $K_{2}$ is the modified Bessel function of the second kind.
Proof. We know from [46, Section 7.11] that, for all $f, g \in L^{2}\left(\mathbb{R}^{3}\right)$ and all $t>0$,

$$
\left\langle f, e^{-t T} g\right\rangle=\frac{e^{t}}{2 \pi^{2}} \int_{\mathbb{R}^{3}} \int_{\mathbb{R}^{3}} \frac{t}{t^{2}+|x-y|^{2}} K_{2}\left(\sqrt{|x-y|^{2}+t^{2}}\right) \overline{f(x)} g(y) \mathrm{d} x \mathrm{~d} y
$$

where the factor $e^{t}$ originated from the 1 that was subtracted from the operator $\sqrt{1-\Delta}$.
If $f$ and $g$ have disjoint supports, we can write that

$$
\begin{align*}
\frac{1}{t} & \left(\langle f, g\rangle-\left\langle f, e^{-t T} g\right\rangle\right)=-\frac{1}{t}\left\langle f, e^{-t T} g\right\rangle \\
& =\frac{-e^{t}}{2 \pi^{2}} \int_{\operatorname{supp} g} \int_{\operatorname{supp} f} \frac{1}{t^{2}+|x-y|^{2}} K_{2}\left(\sqrt{|x-y|^{2}+t^{2}}\right) \overline{f(x)} g(y) \mathrm{d} x \mathrm{~d} y \tag{A.2}
\end{align*}
$$

If we assume that $\operatorname{dist}(\operatorname{supp} f, \operatorname{supp} g) \geqslant R$, then, if $x \in \operatorname{supp} f$ and $y \in \operatorname{supp} g, 1 /\left(t^{2}+|x-y|^{2}\right) \leqslant$ $1 /\left(t^{2}+R^{2}\right)$. Moreover, as we can see for example in the result 9.7.2 of [1], $K_{2}(x) \sim \sqrt{\frac{\pi}{2 x}} e^{-x}$ for large $x>0$. We have thus that there exists a constant $C>0$ such that, if $R$ is large enough,

$$
\begin{equation*}
\left|\frac{1}{t}\left\langle f, e^{-t T} g\right\rangle\right| \leqslant C e^{t} \int_{\operatorname{supp} g} \int_{\operatorname{supp} f} \frac{1}{t^{2}+R^{2}} \sqrt{\frac{\pi}{2 \sqrt{|x-y|^{2}+t^{2}}}} e^{-\sqrt{t^{2}+|x-y|^{2}}}|f(x)||g(y)| \mathrm{d} x \mathrm{~d} y \tag{A.3}
\end{equation*}
$$

We can thus use the dominated convergence theorem to take the limit when $t \rightarrow 0$ : we find from (A.2) that, for all such $f$ and $g,\langle f, T g\rangle$ is well-defined and

$$
\langle f, T g\rangle=-\frac{1}{2 \pi^{2}} \int_{\mathbb{R}^{3}} \int_{\mathbb{R}^{3}} \frac{K_{2}(|x-y|)}{|x-y|^{2}} \overline{f(x)} g(y) \mathrm{d} x \mathrm{~d} y .
$$

We get from (A.3) that

$$
\begin{aligned}
|\langle f, T g\rangle| & \leqslant C \int_{\operatorname{supp} g} \int_{\operatorname{supp} f} \frac{1}{R^{2}} \frac{1}{\sqrt{|x-y|}} e^{-|x-y|}|f(x)||g(y)| \mathrm{d} x \mathrm{~d} y \\
& \leqslant \frac{C}{R^{5 / 2}} \int_{\operatorname{supp} g} \int_{\operatorname{supp} f} e^{-|x-y|}|f(x) \| g(y)| \mathrm{d} x \mathrm{~d} y \\
& \leqslant \frac{C}{R^{5 / 2}}\|f\|_{L^{2}}\|g\|_{L^{2}} \int_{|z| \geqslant R} e^{-|z|} \mathrm{d} z .
\end{aligned}
$$

Calculating the integral, we arrive at

$$
|\langle f, T g\rangle| \leqslant C e^{-R}\|f\|_{L^{2}}\|g\|_{L^{2}}
$$

as desired.

Another difficulty with the semirelativistic kinetic energy operator is that there is no Leibniz rule. As a consequence, commutators between this operator and a multiplication operator cannot be explicitly computed. Nevertheless, we are able to bound some commutators with the following lemma.

Lemma A.2. Let $\zeta \in C_{0}^{\infty}\left(\mathbb{R}^{3}\right)$, with supp $\zeta \subset B_{1 / 4}(0)$ and denote $\zeta_{R}(x):=\zeta\left(\frac{x}{R}\right)$. Then, there exists $C>0$ such that we have for all $R>1$, in operator norm,

$$
\left\|\left[T, \zeta_{R}\right]\right\|_{\mathcal{B}\left(L^{2}\right)} \leqslant \frac{C}{R}
$$

Proof. We first consider $\psi \in H^{1}\left(\mathbb{R}^{3}\right)$. In view of the definition of $T$ in (1.3), we use the Fourier transformation as defined in (1.4). Let us denote, for all $p \in \mathbb{R}^{3}, \widetilde{T}(p):=\sqrt{|p|^{2}+1}$. We define

$$
\Xi_{R}:=\frac{1}{(2 \pi)^{3 / 2}}\left(\widetilde{T} \mathcal{F}\left(\zeta_{R} * \psi\right)-\mathcal{F}\left(\zeta_{R}\right) *(\widetilde{T} \mathcal{F}(\psi))\right)
$$

Then $\Xi_{R}=\mathcal{F}\left(\left[T, \zeta_{R}\right] \psi\right)$ so, by Plancherel identity, $\left\|\left[T, \zeta_{R}\right] \psi\right\|_{L^{2}\left(\mathbb{R}^{3}\right)}=\left\|\Xi_{R}\right\|_{L^{2}\left(\mathbb{R}^{3}\right)}$.
For almost every $p$, we have

$$
\begin{aligned}
(2 \pi)^{3 / 2} \Xi_{R}(p) & =\widetilde{T}(p)\left(\mathcal{F}\left(\zeta_{R}\right) * \mathcal{F}(\psi)\right)(p)-\mathcal{F}\left(\zeta_{R}\right) *(\widetilde{T} \mathcal{F}(\psi))(p) \\
& =\widetilde{T}(p) \int_{\mathbb{R}^{3}} \mathcal{F}\left(\zeta_{R}\right)(q) \mathcal{F}(\psi)(p-q) \mathrm{d} q-\int_{\mathbb{R}^{3}} \mathcal{F}\left(\zeta_{R}\right)(q) \widetilde{T}(p-q) \mathcal{F}(\psi)(p-q) \mathrm{d} q \\
& =\int_{\mathbb{R}^{3}} \mathcal{F}\left(\zeta_{R}\right)(q) \mathcal{F}(\psi)(p-q)(\widetilde{T}(p)-\widetilde{T}(p-q)) \mathrm{d} q
\end{aligned}
$$

Therefore, since the function $\widetilde{T}$ is Lipschitz with Lipschitz constant 1,

$$
(2 \pi)^{3 / 2}\left|\Xi_{R}(p)\right| \leqslant \int_{\mathbb{R}^{3}}\left|\mathcal{F}\left(\zeta_{R}\right)(q)\|\mathcal{F}(\psi)(p-q)\| \widetilde{T}(p)-\widetilde{T}(p-q)\right| \mathrm{d} q \leqslant \int_{\mathbb{R}^{3}}\left|q\left\|\mathcal{F}\left(\zeta_{R}\right)(q)\right\| \mathcal{F}(\psi)(p-q)\right| \mathrm{d} q
$$

so we have

$$
(2 \pi)^{3 / 2}\left\|\Xi_{R}\right\|_{L^{2}\left(\mathbb{R}^{3}\right)} \leqslant\| \| \cdot \mid \mathcal{F}\left(\zeta_{R}\right)\left\|_{L^{1}\left(\mathbb{R}^{3}\right)}\right\| \mathcal{F}(\psi) \|_{L^{2}\left(\mathbb{R}^{3}\right)}
$$

Since $\|\mathcal{F}(\psi)\|_{L^{2}\left(\mathbb{R}^{3}\right)}=\|\psi\|_{L^{2}\left(\mathbb{R}^{3}\right)}$,

$$
\left\|\left[T, \zeta_{R}\right] \psi\right\|_{L^{2}\left(\mathbb{R}^{3}\right)} \leqslant \frac{1}{(2 \pi)^{3 / 2}}\left\||\cdot| \mathcal{F}\left(\zeta_{R}\right)\right\|_{L^{1}\left(\mathbb{R}^{3}\right)}\|\psi\|_{L^{2}\left(\mathbb{R}^{3}\right)}
$$

By density of $H^{1}$ in $L^{2}$, we can extend this inequality to all $\psi \in L^{2}\left(\mathbb{R}^{3}\right)$, which implies that

$$
\begin{equation*}
\left\|\left[T, \zeta_{R}\right]\right\|_{\mathcal{B}\left(L^{2}\right)} \leqslant \frac{1}{(2 \pi)^{3 / 2}}\left\||\cdot| \mathcal{F}\left(\zeta_{R}\right)\right\|_{L^{1}\left(\mathbb{R}^{3}\right)} \tag{A.4}
\end{equation*}
$$

Let us estimate $\left\||\cdot| \mathcal{F}\left(\zeta_{R}\right)\right\|_{L^{1}\left(\mathbb{R}^{3}\right)}$. We have that

$$
\int_{\mathbb{R}^{3}}\left|q \| \mathcal{F}\left(\zeta_{R}\right)(q)\right| \mathrm{d} q=\frac{1}{(2 \pi)^{3 / 2}} \int_{\mathbb{R}^{3}}|q|\left|\int_{\mathbb{R}^{3}} e^{-\mathrm{i} q \cdot x} \zeta\left(\frac{x}{R}\right) \mathrm{d} x\right| \mathrm{d} q=\frac{1}{(2 \pi)^{3 / 2}} \int_{\mathbb{R}^{3}}\left|\frac{p}{R}\right|\left|\int_{\mathbb{R}^{3}} e^{-\mathrm{i} p \cdot y} \zeta(y) \mathrm{d} y\right| \mathrm{d} p,
$$

where we have performed the changes of variable $y=\frac{x}{R}$ and $p=R q$. We find that the norm is equal to $R^{-1}(2 \pi)^{-3 / 2} \int_{\mathbb{R}^{3}}|p|\left|\int_{\mathbb{R}^{3}} e^{-\mathrm{i} p \cdot y} \zeta(y) \mathrm{d} y\right| \mathrm{d} p$. Since $\zeta$ is in the Schwartz space, we know that this integral is finite. As a consequence,

$$
\begin{equation*}
\int_{\mathbb{R}^{3}}\left|q \| \mathcal{F}\left(\zeta_{R}\right)(q)\right| \mathrm{d} q \leqslant \frac{C}{R} \tag{A.5}
\end{equation*}
$$

which together with (A.4) concludes the proof.
An important consequence of this result is that we can get a result analogous to what is called IMS localization formula in the non-relativistic case, see e.g. [25, Theorem 3.2].

Recall that $\chi_{1}$ was defined in (2.5). Let us define $v_{1}:=\chi_{1}, v_{2}:=v_{1}\left(\cdot-e_{1}\right)$ and $v_{0}:=1-v_{1}-v_{2}$. We define too $V:=\sqrt{v_{0}^{2}+v_{1}^{2}+v_{2}^{2}}$ : by construction, $V \geqslant \sqrt{3} / 2$. For $R>0$ and $i=0,1,2$, we define the functions

$$
\begin{equation*}
J_{i, R}(x):=\frac{v_{i}(x / R)}{V(x / R)} \tag{A.6}
\end{equation*}
$$

We have that

$$
\sum_{i=0}^{2} J_{i, R}(x)^{2}=1
$$

for all $x$ and $R$.

We define the one-particle localization error for $\psi \in H^{1 / 2}\left(\mathbb{R}^{3}\right)$ by

$$
\operatorname{Err}[\psi]:=\langle\psi, T \psi\rangle-\sum_{i=0}^{2}\left\langle J_{i, R} \psi, T J_{i, R} \psi\right\rangle .
$$

Lemma A.3. There exists $C>0$ such that, for all $R>0$ and $\psi \in H^{1 / 2}\left(\mathbb{R}^{3}\right)$,

$$
|\operatorname{Err}[\psi]| \leqslant \frac{C}{R}\|\psi\|_{L^{2}\left(\mathbb{R}^{3}\right)}^{2}
$$

Proof. We remark that $\langle\psi, T \psi\rangle=\sum_{i=0}^{2}\left\langle\psi, T J_{i, R}^{2} \psi\right\rangle$ so

$$
|\operatorname{Err}[\psi]| \leqslant \sum_{i=0}^{2}\left|\left\langle\psi,\left[T, J_{i, R}\right] J_{i, R} \psi\right\rangle\right| \leqslant \sum_{i=0}^{2}\left\|\left[T, J_{i, R}\right]\right\|_{\mathcal{B}\left(L^{2}\right)}\|\psi\|_{L^{2}\left(\mathbb{R}^{3}\right)}^{2}
$$

since $\left\|J_{i, R}\right\|_{L^{\infty}} \leqslant 1$ for all $i$ and $R$. For $i=1$ or 2 , the function $J_{i, R}$ satisfies the hypothesis of Lemma A. 2 so $\left\|\left[T, J_{i, R}\right]\right\|_{\mathcal{B}\left(L^{2}\right)} \leqslant C / R$. Concerning the term $J_{0, R}$, we use the fact that $\left[T, J_{0, R}\right]=-\left[T, 1-J_{0, R}\right]$. But $1-J_{0, R}$ is compactly supported. Thus, $1-J_{0, R}$ satisfies the hypotheses of Lemma A.2, and we can conclude with the same argument.

This result can be extended to the case of $N$ electrons. Let $\gamma=\left(\gamma_{1}, \ldots, \gamma_{N}\right) \in\{0,1,2\}^{N}$.
We define

$$
\begin{equation*}
J_{\gamma, R}\left(x_{1}, \ldots, x_{N}\right):=\prod_{i=1}^{N} J_{\gamma_{i}, R}\left(x_{i}\right), \tag{A.7}
\end{equation*}
$$

where $J_{\gamma_{i}, R}$ was defined in (A.6).
We define the $N$-particle localization error for $\psi \in H^{1 / 2}\left(\mathbb{R}^{3 N}\right)$ by

$$
\operatorname{Err}[\psi]:=\left\langle\psi, \sum_{i_{1}}^{N} T_{i} \psi\right\rangle-\sum_{\gamma \in\{0,1,2\}^{N}}\left\langle J_{\gamma, R} \psi, \sum_{i_{1}}^{N} T_{i} J_{\gamma, R} \psi\right\rangle .
$$

Proposition A.4. There exists $C>0$ such that, for all $R>0$ and $\psi \in H^{1 / 2}\left(\mathbb{R}^{3 N}\right)$,

$$
\begin{equation*}
|\operatorname{Err}[\psi]| \leqslant \frac{C}{R}\|\psi\|_{L^{2}\left(\mathbb{R}^{3 N}\right)}^{2} \tag{A.8}
\end{equation*}
$$

Proof. The proof is similar to the one-particle case.
Remark A.5. As explained in [14, Theorem 3.1], it is even possible to have a control of the localization error by $C / R^{2}$. Nevertheless, we do not need it in our work and it requires a longer proof.

The next lemma states that, even if the commutator is not local, we are able to control this nonlocality in the sense that, if we apply it to a function which has some decay at infinity, then we get something which is decaying, at least in a quadratic form sense.

Lemma A.6. Let us consider, for $R>0$, the function $\chi_{R}$ defined in (2.5). Then there exists $C>0$ such that, for all functions $F: \mathbb{R}^{+} \rightarrow[1, \infty)$, all $R$ large enough and all $\xi, \eta \in L^{2}\left(\mathbb{R}^{3}\right)$ such that $F(|\cdot|) \eta \in L^{2}\left(\mathbb{R}^{3}\right)$,

$$
\left|\left\langle\eta,\left[T, \chi_{R}\right] \xi\right\rangle\right| \leqslant C \max \left\{\frac{1}{R} \sup _{r \in[R / 20,+\infty)} \frac{1}{F(r)}, e^{-3 R / 80}\right\}\|F(|\cdot|) \eta\|_{L^{2}\left(\mathbb{R}^{3}\right)}\|\xi\|_{L^{2}\left(\mathbb{R}^{3}\right)}
$$

Remark A.7. Applying Lemma (A.6) for $F(x)=e^{a x}, a>0$ we find

$$
\begin{equation*}
\left|\left\langle\eta,\left[T, \chi_{R}\right] \xi\right\rangle\right| \leqslant C e^{-\min (a / 20,3 / 80) R}\left\|e^{a|\cdot|} \eta\right\|_{L^{2}\left(\mathbb{R}^{3}\right)}\|\xi\|_{L^{2}\left(\mathbb{R}^{3}\right)} \tag{A.9}
\end{equation*}
$$

If $F(x)=1+x^{n}$ for some $n \in \mathbb{N}$, we have

$$
\begin{equation*}
\left|\left\langle\eta,\left[T, \chi_{R}\right] \xi\right\rangle\right| \leqslant C \frac{1}{R^{n+1}}\left\|\left(1+|\cdot|^{n}\right) \eta\right\|_{L^{2}\left(\mathbb{R}^{3}\right)}\|\xi\|_{L^{2}\left(\mathbb{R}^{3}\right)} \tag{A.10}
\end{equation*}
$$

Proof. We introduce the function $\chi_{R / 2}^{c}:=1-\chi_{R / 2}$ :

$$
\begin{equation*}
\left\langle\eta,\left[T, \chi_{R}\right] \xi\right\rangle=\left\langle\chi_{R / 2} \eta,\left[T, \chi_{R}\right] \xi\right\rangle+\left\langle\chi_{R / 2}^{c} \eta,\left[T, \chi_{R}\right] \xi\right\rangle \tag{A.11}
\end{equation*}
$$

On the one hand,

$$
\begin{align*}
\left|\left\langle\chi_{R / 2}^{c} \eta,\left[T, \chi_{R}\right] \xi\right\rangle\right| & \leqslant\|F(|\cdot|) \eta\|_{L^{2}\left(\mathbb{R}^{3}\right)}\left\|\frac{1}{F(|\cdot|)} \chi_{R / 2}^{c}\right\|_{L^{\infty}\left(\mathbb{R}^{3}\right)}\left\|\left[T, \chi_{R}\right] \xi\right\|_{L^{2}\left(\mathbb{R}^{3}\right)} \\
& \leqslant \frac{C}{R}\left\|\frac{1}{F(|\cdot|)} \chi_{R / 2}^{c}\right\|_{L^{\infty}\left(\mathbb{R}^{3}\right)}\|F(|\cdot|) \eta\|_{L^{2}\left(\mathbb{R}^{3}\right)}\|\xi\|_{L^{2}\left(\mathbb{R}^{3}\right)} \\
& \leqslant \frac{C}{R} \sup _{[R / 20,+\infty)} \frac{1}{F(r)}\|F(|\cdot|) \eta\|_{L^{2}\left(\mathbb{R}^{3}\right)}\|\xi\|_{L^{2}\left(\mathbb{R}^{3}\right)} \tag{A.12}
\end{align*}
$$

where we used Lemma A. 2 to go from the first to the second line.
On the other hand,

$$
\left\langle\chi_{R / 2} \eta,\left[T, \chi_{R}\right] \xi\right\rangle=-\left\langle\chi_{R / 2} \eta,\left[T, \chi_{R}^{c}\right] \xi\right\rangle=-\left\langle\chi_{R / 2} \eta, T \chi_{R}^{c} \xi\right\rangle
$$

since $\chi_{R}^{c} \chi_{R / 2}=0$. But $\chi_{R / 2} \eta$ and $\chi_{R}^{c} \xi$ are both $L^{2}$ functions and the distance between their support is $3 R / 80$. Hence, we get from Lemma A. 1 that, for $R$ large enough:

$$
\begin{equation*}
\left|\left\langle\chi_{R / 2} \eta, T \chi_{R}^{c} \xi\right\rangle\right| \leqslant C e^{-3 R / 80}\|\eta\|_{L^{2}\left(\mathbb{R}^{3}\right)}\|\xi\|_{L^{2}\left(\mathbb{R}^{3}\right)} \tag{A.13}
\end{equation*}
$$

Inserting (A.12) and (A.13) into (A.11) concludes the proof.

## Appendix B. Proof of Theorems 1.1, 1.2 and 1.4

B.1. Proof of Theorem 1.1. The proof is split into 2 parts. In the first one, we prove that

$$
\left[\inf \sigma\left(\hat{H}_{N-1}(Y, \mathcal{Z})\right), \infty\right) \subset \sigma_{\mathrm{ess}}\left(\hat{H}_{N}(Y, \mathcal{Z})\right)
$$

and in the second one we prove the reverse inclusion.
"Easy part": We prove here

$$
\begin{equation*}
\left[\inf \sigma\left(\hat{H}_{N-1}(Y, \mathcal{Z})\right), \infty\right) \subset \sigma\left(\hat{H}_{N}(Y, \mathcal{Z})\right) \tag{B.1}
\end{equation*}
$$

Since the left-hand side is an interval, it will be in the essential spectrum.
We denote $E_{N-1}:=\inf \sigma\left(\hat{H}_{N-1}(Y, \mathcal{Z})\right)$. Let us pick $y \in\left[E_{N-1}, \infty\right)$. We want to prove that $y \in$ $\sigma\left(\hat{H}_{N}(Y, \mathcal{Z})\right)$. By assumption on $y$, there exists $a \geqslant 0$ such that $y=E_{N-1}+a$. Let $\epsilon>0$. Then, since $E_{N-1} \in \sigma\left(\hat{H}_{N-1}(Y, \mathcal{Z})\right)$, there exists by Weyl's criterion (cf. [59, Theorem VII.12]) a $\phi_{N-1, \epsilon} \in$ $\mathcal{D}\left(\hat{H}_{N-1}(Y, \mathcal{Z})\right)$, normalized, such that

$$
\begin{equation*}
\left\|\left(H_{N-1}(Y, \mathcal{Z})-E_{N-1}\right) \phi_{N-1, \epsilon}\right\| \leqslant \frac{\epsilon}{3 \sqrt{N}} \tag{B.2}
\end{equation*}
$$

But compactly supported functions are dense in $\mathcal{D}\left(\hat{H}_{N-1}(Y, \mathcal{Z})\right)$ in the graph norm. Indeed, if we consider a function $\zeta \in \mathcal{D}\left(\hat{H}_{N-1}(Y, \mathcal{Z})\right)$ and, for $R>0$, the function $\chi_{R}$ defined in (2.5), we have that

$$
\left\|\hat{H}_{N-1}(Y, \mathcal{Z})\left(\chi_{R}^{\otimes N-1} \zeta-\zeta\right)\right\| \leqslant\left\|\left(1-\chi_{R}^{\otimes N-1}\right) \hat{H}_{N-1}(Y, \mathcal{Z}) \zeta\right\|+\left\|\sum_{i=1}^{N}\left[T_{i}, \chi_{R}^{\otimes N-1}\right] \zeta\right\|
$$

Since $\zeta \in \mathcal{D}\left(\hat{H}_{N-1}(Y, \mathcal{Z})\right)$, we have $\hat{H}_{N-1}(Y, \mathcal{Z}) \zeta \in L^{2}\left(\mathbb{R}^{3(N-1)}\right)$ and therefore

$$
\left\|\left(1-\chi_{R}^{\otimes N-1}\right) \hat{H}_{N-1}(Y, \mathcal{Z}) \zeta\right\| \rightarrow 0, \quad \text { when } R \rightarrow \infty
$$

On the other hand, by Lemma A.2, $\left\|\sum_{i=1}^{N}\left[T_{i}, \chi_{R}^{\otimes N-1}\right] \zeta\right\| \leqslant C / R\|\zeta\|$, which goes to 0 as well when $R$ goes to $\infty$. As a consequence of this density argument, we can assume, without loss of generality, that $\phi_{N-1, \epsilon}$ is compactly supported. Moreover, since $\sigma(T)=[0,+\infty)$ and therefore $a \in \sigma(T)$, there exists $f_{\epsilon} \in C_{c}^{\infty}\left(\mathbb{R}^{3}\right)$, normalized, such that

$$
\begin{equation*}
\left\|(T-a) f_{\epsilon}\right\| \leqslant \frac{\epsilon}{3 \sqrt{N}} \tag{B.3}
\end{equation*}
$$

For $h \in \mathbb{R}^{3}$, we consider the function defined on $\mathbb{R}^{3}$ by $f_{\epsilon, h}(x):=f_{\epsilon}(x-h)$. We further consider the function

$$
\begin{equation*}
g_{\epsilon, h}:=\frac{Q_{N}\left(\phi_{N-1, \epsilon} \otimes f_{\epsilon, h}\right)}{\left\|Q_{N}\left(\phi_{N-1, \epsilon} \otimes f_{\epsilon, h}\right)\right\|} . \tag{B.4}
\end{equation*}
$$

We will prove that, if we choose $|h|$ large enough, $g_{\epsilon, h}$ is in the domain of $\hat{H}_{N}(Y, \mathcal{Z})$ and

$$
\left\|\left(H_{N}(Y, \mathcal{Z})-y\right) g_{\epsilon, h}\right\|<\epsilon
$$

Let us first compute the norm which appears in the denominator of (B.4). We observe that, if $|h|$ is large enough, then the electronic densities (defined similarly to (1.33)) of $\phi_{N-1, \epsilon}$ and $f_{\epsilon, h}$ have disjoint supports. We can therefore argue similarly to the proof of (2.14), with $N_{1}=N-1$, to find

$$
\left\|Q_{N}\left(\phi_{N-1, \epsilon} \otimes f_{\epsilon, h}\right)\right\|=\frac{1}{\sqrt{N}}
$$

As a consequence we have that

$$
\begin{equation*}
\left\|\left(H_{N}(Y, \mathcal{Z})-y\right) g_{\epsilon, h}\right\| \leqslant \sqrt{N}\left\|\left(H_{N}(Y, \mathcal{Z})-y\right)\left(\phi_{N-1, \epsilon} \otimes f_{\epsilon, h}\right)\right\| \tag{B.5}
\end{equation*}
$$

where in the last step we used that $Q_{N}$ commutes with $H_{N}(Y, \mathcal{Z})$ and that it is an orthogonal projection. But we have that

$$
\left(H_{N}(Y, \mathcal{Z})-y\right)=\left(H_{N-1}(Y, \mathcal{Z})-E_{N-1}\right)+\left(T_{N}-a\right)+\sum_{i=1}^{N-1} \frac{\alpha}{\left|x_{i}-x_{N}\right|}-\sum_{j=1}^{M} \frac{Z_{j} \alpha}{\left|y_{j}-x_{N}\right|}
$$

so we get from (B.2), (B.3) and (B.5) that

$$
\left\|\left(H_{N}(Y, \mathcal{Z})-y\right) g_{\epsilon, h}\right\| \leqslant \frac{2 \epsilon}{3}+\sqrt{N}\left\|\left(\sum_{i=1}^{N-1} \frac{\alpha}{\left|x_{i}-x_{N}\right|}-\sum_{j=1}^{M} \frac{Z_{j} \alpha}{\left|y_{j}-x_{N}\right|}\right)\left(\phi_{N-1, \epsilon} \otimes f_{\epsilon, h}\right)\right\|
$$

We can choose $|h|$ large enough so that, for all $x_{N} \in \operatorname{supp} f_{\epsilon, h}$ and $\left(x_{1}, \ldots, x_{N-1}\right) \in \operatorname{supp} \phi_{N-1, \epsilon}$, $\sum_{i=1}^{N-1} \frac{\alpha}{\left|x_{i}-x_{N}\right|}<\epsilon /(6 \sqrt{N})$ and $\sum_{j=1}^{M} \frac{Z_{j} \alpha}{\left|y_{j}-x_{N}\right|}<\epsilon /(6 \sqrt{N})$. In such a case, we have that $g_{\epsilon, h}$ is in the domain of $\hat{H}_{N}(Y, \mathcal{Z})$ and

$$
\left\|\left(\hat{H}_{N}(Y, \mathcal{Z})-y\right) g_{\epsilon, h}\right\|<\epsilon
$$

concluding the proof of (B.1) since we can do it for arbitrarily small $\epsilon$.
"Hard part": Let us prove that, for all $N, Y, \mathcal{Z}$,

$$
\sigma_{\mathrm{ess}}\left(\hat{H}_{N}(Y, \mathcal{Z})\right) \subset\left[E_{N-1}, \infty\right)
$$

Let $\psi_{n} \in \operatorname{Ran} Q_{N}$ be a Weyl sequence for $\hat{H}_{N}(Y, \mathcal{Z})$ : for some $\lambda \in \sigma_{\text {ess }}\left(\hat{H}_{N}(Y, \mathcal{Z})\right),\left\|\left(\hat{H}_{N}(Y, \mathcal{Z})-\lambda\right) \psi_{n}\right\| \rightarrow$ 0 , with $\left\|\psi_{n}\right\|=1$ for all $n$ and $\psi_{n} \rightharpoonup 0$. We will show that

$$
\begin{equation*}
\lim _{n \rightarrow \infty}\left\langle\psi_{n}, \hat{H}_{N}(Y, \mathcal{Z}) \psi_{n}\right\rangle \geqslant E_{N-1} \tag{B.6}
\end{equation*}
$$

To this purpose, we use a method similar to the one given in Proposition A.4. The difference is that we consider only 2 functions instead of 3 . We define again $v_{1}:=\chi_{1}$ and we choose $v_{0}:=1-v_{1}$. We define then, for $\beta \in\{0,1\}^{N}$, $J_{\beta, R}$ similarly as in (A.7). We have the same bound on the error as in Proposition A.4, implying that there exist $C>0$ such that, for all $n \in \mathbb{N}$ and $R>1$,

$$
\begin{equation*}
\left\langle\psi_{n}, H_{N}(Y, \mathcal{Z}) \psi_{n}\right\rangle \geqslant \sum_{\beta \subset\{1, \ldots, N\}}\left\langle J_{\beta, R} \psi_{n}, \hat{H}_{N}(Y, \mathcal{Z}) J_{\beta, R} \psi_{n}\right\rangle-\frac{C}{R} \tag{B.7}
\end{equation*}
$$

Let us estimate the sum in the right-hand side of (B.7). First, we can observe that, for $B=\{1, \ldots, N\}$, the function $J_{B, R} \psi_{n} \in \operatorname{Ran} Q_{N}$ and thus

$$
\begin{equation*}
\left\langle J_{B, R} \psi_{n}, H_{N}(Y, \mathcal{Z}) J_{B, R} \psi_{n}\right\rangle \geqslant E_{N}\left\|J_{B, R} \psi_{n}\right\|^{2} \tag{B.8}
\end{equation*}
$$

where we recall that $E_{N}=\inf \sigma\left(\hat{H}_{N}(Y, \mathcal{Z})\right)$.
Now, for $\operatorname{Card}(\beta)=K \leqslant N-1$,

$$
\begin{align*}
\left\langle J_{\beta, R} \psi_{n}, H_{N}(Y, \mathcal{Z}) J_{\beta, R} \psi_{n}\right\rangle=\left\langle J_{\beta, R} \psi_{n}, H_{\beta}(Y, \mathcal{Z}) J_{\beta, R}\right. & \left.\psi_{n}\right\rangle \\
& +\left\langle J_{\beta, R} \psi_{n},\left(H_{N}(Y, \mathcal{Z})-H_{\beta}(Y, \mathcal{Z})\right) J_{\beta, R} \psi_{n}\right\rangle \tag{B.9}
\end{align*}
$$

where $H_{\beta}(Y, \mathcal{Z})$ is the Hamiltonian describing the electrons in the subset $\beta$, interacting with all nuclei, namely

$$
H_{\beta}(Y, \mathcal{Z}):=\sum_{i \in \beta}\left(T_{i}-\sum_{j=1}^{M} \frac{Z_{j} \alpha}{\left|x_{i}-y_{j}\right|}\right)+\sum_{\substack{i, j \in \beta \\ i<j}} \frac{\alpha}{\left|x_{i}-x_{j}\right|}+\sum_{1 \leqslant k<l \leqslant M} \frac{Z_{k} Z_{l} \alpha}{\left|y_{k}-y_{l}\right|}
$$

On the one hand,

$$
H_{N}(Y, \mathcal{Z})-H_{\beta}(Y, \mathcal{Z}) \geqslant-\sum_{i \notin \beta} \sum_{j=1}^{M} \frac{Z_{j} \alpha}{\left|x_{i}-y_{j}\right|}
$$

But, for $R$ large enough, if $x_{i}$ is in the support of $v_{0}$, there exists a constant $C^{\prime}$ such that, for all $j$, $\left|x_{i}-y_{j}\right| \geqslant C^{\prime} R$. As a consequence, there exists $C>0$ such that

$$
\begin{equation*}
\left\langle J_{\beta, R} \psi_{n},\left(H_{N}(Y, \mathcal{Z})-H_{\beta}(Y, \mathcal{Z})\right) J_{\beta, R} \psi_{n}\right\rangle \geqslant-\frac{C}{R} \tag{B.10}
\end{equation*}
$$

On the other hand, $J_{\beta, R} \psi_{n}$ is antisymmetric with respect to the variables in $\beta$. As a consequence,

$$
\begin{equation*}
\left\langle J_{\beta, R} \psi_{n}, H_{\beta}(Y, \mathcal{Z}) J_{\beta, R} \psi_{n}\right\rangle \geqslant E_{K}\left\|J_{\beta, R} \psi_{n}\right\|^{2} \geqslant E_{N-1}\left\|J_{\beta, R} \psi_{n}\right\|^{2} \tag{B.11}
\end{equation*}
$$

since, by the first part of the theorem, the sequence $\left(E_{k}\right)$ is nonincreasing. Inserting (B.10) and (B.11) into (B.9) and combining it with (B.7) and (B.8), we obtain that

$$
\left\langle\psi_{n}, H_{N}(Y, \mathcal{Z}) \psi_{n}\right\rangle \geqslant E_{N}\left\|J_{B, R} \psi_{n}\right\|^{2}+E_{N-1} \sum_{\beta \subsetneq\{1, \ldots, N\}}\left\|J_{\beta, R} \psi_{n}\right\|^{2}-\frac{C}{R}
$$

and, since $\sum_{\beta \subsetneq\{1, \ldots, N\}} J_{\beta, R}^{2}=1-J_{B, R}^{2}$ and $\left\|\psi_{n}\right\|=1$, we find that

$$
\left\langle\psi_{n}, H_{N}(Y, \mathcal{Z}) \psi_{n}\right\rangle \geqslant E_{N-1}+\left(E_{N}-E_{N-1}\right)\left\|J_{B, R} \psi_{n}\right\|^{2}-\frac{C}{R}
$$

We will now prove that, for all $R>0$, we have that

$$
\lim _{n \rightarrow \infty}\left\|J_{B, R} \psi_{n}\right\|=0
$$

and this will conclude the proof of (B.6) and therefore of Theorem 1.1. This comes from the fact that, for $c>0$ large, we have the following decomposition:

$$
\begin{gathered}
J_{B, R} \psi_{n}=\underbrace{J_{B, R}\left(\sum_{i=1}^{N} T_{i}+c\right)^{-\frac{1}{2}}}_{\text {compact operator }} \times \\
\underbrace{\left(\sum_{i=1}^{N} T_{i}+c\right)^{\frac{1}{2}}\left(H_{N}(Y, \mathcal{Z})+c\right)^{-\frac{1}{2}}}_{\text {bounded operator }} \underbrace{\left(H_{N}(Y, \mathcal{Z})+c\right)^{-\frac{1}{2}}}_{\text {bounded }} \underbrace{\left(H_{N}(Y, \mathcal{Z})+c\right) \psi_{n}}_{\rightarrow 0} .
\end{gathered}
$$

The weak convergence $\left(H_{N}(Y, \mathcal{Z})+c\right) \psi_{n} \rightharpoonup 0$ follows from the decomposition

$$
H_{N}(Y, \mathcal{Z})+c=\left(H_{N}(Y, \mathcal{Z})-\lambda\right)+(\lambda+c)
$$

We know from the definition of $\left(\psi_{n}\right)$ that $\left(H_{N}(Y, \mathcal{Z})-\lambda\right) \psi_{n} \rightarrow 0$ and that $(\lambda+c) \psi_{n} \rightharpoonup 0$. The operator $\left(H_{N}(Y, \mathcal{Z})+c\right)^{-\frac{1}{2}}$ is bounded since $H_{N}(Y, \mathcal{Z})+c$ is strictly positive for $c$ large enough. The boundedness of $\left(\sum_{i=1}^{N} T_{i}+c\right)^{\frac{1}{2}}\left(H_{N}(Y, \mathcal{Z})+c\right)^{-\frac{1}{2}}$ follows from (1.8). Finally, the operator $J_{B, R}\left(\sum_{i=1}^{N} T_{i}+c\right)^{-\frac{1}{2}}$ is compact, since it is the norm limit of the sequence $J_{B, R}\left(\sum_{i=1}^{N} T_{i}+c\right)^{-\frac{1}{2}} 1_{\sum_{i=1}^{N} T_{i}+c<n}$ of Hilbert-Schmidt operators. We thus obtain that $\left\|J_{B, R} \psi_{n}\right\| \rightarrow 0$, concluding the proof of the theorem.
B.2. Proof of Theorem 1.2. We prove in the following exponential decay of approximate eigenfunctions for the Hamiltonian $\hat{H}_{N}$.

The proof closely follows the proof of Theorem 2.1 in [14].
Let $\left(\gamma_{s}\right)$ be as in the statement of Theorem 1.2. For $\nu \geqslant 0$ and $\epsilon \geqslant 0$, we set for $r \in \mathbb{R}^{+}$

$$
G_{\nu, \epsilon}(r):=\frac{\nu r}{1+\epsilon r}
$$

We then define on $\mathbb{R}^{3 N}$ the function $F_{\nu, \epsilon}$ by

$$
F_{\nu, \epsilon}\left(x_{1}, \ldots, x_{N}\right):=\sum_{i=1}^{N} G_{\nu, \epsilon}\left(\left|x_{i}\right|\right)
$$

Moreover, for the function $\chi_{R}$ defined in (2.5), we introduce

$$
\xi_{\varepsilon}:=\left(1-\chi_{R}^{\otimes N}\right) e^{F_{\nu, \varepsilon}}
$$

By the definition of $\tilde{\Sigma}$ in (1.15), there exists a function $\theta(R)$ with $\lim _{R \rightarrow+\infty} \theta(R)=0$ such that for all $s \in \mathcal{I}$

$$
(\tilde{\Sigma}-\mu-\theta(R))\left\|\xi_{\varepsilon} \gamma_{s}\right\|^{2} \leqslant\left\langle\xi_{\varepsilon} \gamma_{s},\left(H_{N}(Y, \mathcal{Z})-\mu\right) \xi_{\varepsilon} \gamma_{s}\right\rangle=\operatorname{Re}\left\langle\gamma_{s}, \xi_{\varepsilon}^{2} \Gamma_{s}\right\rangle+\operatorname{Re}\left\langle\xi_{\varepsilon} \gamma_{s}, \sum_{i=1}^{N}\left[T_{i}, \xi_{\varepsilon}\right] \gamma_{s}\right\rangle
$$

where in the last step we used (1.16) and that $\xi_{\epsilon}$ commutes with potentials. Using the assumption that, for all $s \in \mathcal{I},\left\|e^{a|\cdot|} \Gamma_{s}\right\|_{L^{2}\left(\mathbb{R}^{3 N}\right)} \leqslant C_{1}$, we can bound, for all $R, \epsilon, \nu<a$ and $s$, the first term on the right-hand side:

$$
\begin{equation*}
(\tilde{\Sigma}-\mu-\theta(R))\left\|\xi_{\varepsilon} \gamma_{s}\right\|^{2} \leqslant C_{1}\left\|\gamma_{s}\right\|+\operatorname{Re}\left\langle\xi_{\varepsilon} \gamma_{s}, \sum_{i=1}^{N}\left[T_{i}, \xi_{\varepsilon}\right] \gamma_{s}\right\rangle \tag{B.12}
\end{equation*}
$$

By [14, Lemma C.10, Remark C.7], we know that for all $\nu \in(0,1)$, there exist $C_{\nu}>0$ such that, for all $R$,

$$
\begin{equation*}
\left|\operatorname{Re}\left\langle\xi_{\varepsilon} \gamma_{s}, \sum_{i=1}^{N}\left[T_{i}, \xi_{\varepsilon}\right] \gamma_{s}\right\rangle\right| \leqslant N C_{\nu}\left(\frac{L_{\chi}}{R}+\nu\right)^{2}\left\|e^{F_{\nu, \varepsilon}} \gamma_{s}\right\|^{2} \tag{B.13}
\end{equation*}
$$

where $L_{\chi}$ is the Lipschitz constant of $\chi$, recall (2.5). Moreover, $C_{\nu}$ goes to a finite value when $\nu$ goes to 0 . Let us introduce the abbreviation $\delta:=N C_{\nu}\left(\frac{L_{\chi}}{R}+\nu\right)^{2}$. Plugging (B.13) into (B.12) yields

$$
\begin{aligned}
(\tilde{\Sigma}-\mu-\theta(R))\left\|\xi_{\varepsilon} \gamma_{s}\right\|^{2} & \leqslant C_{1}\left\|\gamma_{s}\right\|+\delta\left\|e^{F_{\nu, \varepsilon}} \gamma_{s}\right\|^{2} \\
& \leqslant C_{1}\left\|\gamma_{s}\right\|+2 \delta\left\|\xi_{\varepsilon} \gamma_{s}\right\|^{2}+2 \delta\left\|\chi_{R}^{\otimes N} e^{F_{\nu, \varepsilon}} \gamma_{s}\right\|^{2} \\
& \leqslant C_{1}\left\|\gamma_{s}\right\|+2 \delta\left\|\xi_{\varepsilon} \gamma_{s}\right\|^{2}+2 \delta e^{N \nu R / 4}\left\|\gamma_{s}\right\|^{2}
\end{aligned}
$$

where in the last line we used that, for all $x \in \mathbb{R}^{3 N}$,

$$
\begin{equation*}
\chi_{R}^{\otimes N}(x) e^{F_{\nu, \varepsilon}(x)} \leqslant e^{N \nu R / 8} \tag{B.14}
\end{equation*}
$$

We have therefore obtained

$$
\begin{equation*}
(\tilde{\Sigma}-\mu-\theta(R)-2 \delta)\left\|\xi_{\varepsilon} \gamma_{s}\right\|^{2} \leqslant C_{1}\left\|\gamma_{s}\right\|+2 \delta e^{N \nu R / 4}\left\|\gamma_{s}\right\|^{2} \tag{B.15}
\end{equation*}
$$

But

$$
\lim _{R \rightarrow \infty} \theta(R)+2 \delta=2 N C_{\nu} \nu^{2}
$$

so, for all $\nu$ such that $2 N C_{\nu} \nu^{2}<\tilde{\Sigma}-\mu$, we can find a radius $R$ such that

$$
g:=(\tilde{\Sigma}-\mu-\theta(R)-2 \delta)>0
$$

and together with (B.15) we obtain:

$$
g\left\|\xi_{\varepsilon} \gamma_{s}\right\|^{2}<C_{1}\left\|\gamma_{s}\right\|+2 \delta e^{N \nu R / 4}\left\|\gamma_{s}\right\|^{2} \leqslant C_{R}
$$

where in the last step we used the uniform boundedness of $\left\|\gamma_{s}\right\|$, see (1.16). The constant $C_{R}>0$ does not depend on $\varepsilon$. Since $\xi_{\varepsilon}$ converges monotonically to $\left(1-\chi_{R}^{\otimes N}\right) e^{\nu|x|}$ for $\varepsilon \rightarrow 0$, the monotone convergence theorem for the previous inequality gives that for all $s$ and $\nu$, we have $\left\|\left(1-\chi_{R}^{\otimes N}\right) e^{\nu|x|} \gamma_{s}\right\|^{2} \leqslant g^{-1} C_{R}$, which together with (B.14) concludes the proof of Theorem 1.2.
B.3. Proof of Theorem 1.4. We begin by showing that there exists at least one eigenvalue below $\sigma_{\text {ess }}\left(\hat{H}_{N}(Y, \mathcal{Z})\right)$. As e.g. in [14, Appendix B], we prove it by induction on the number of electrons $N$.

Let us first consider the case $N=1$ : here, we do not need to take symmetry into account. Let $f \in C_{0}^{\infty}\left(\mathbb{R}^{3}\right)$ be such that $\|f\|_{L^{2}}=1$ and $\operatorname{supp} f \subset\left\{x \in \mathbb{R}^{3}: 1<|x|<2\right\}$. For $R>0$, we define

$$
f_{R}(x):=\frac{1}{R^{\frac{3}{2}}} f\left(\frac{x}{R}\right)
$$

Then, for all $R>0,\left\|f_{R}\right\|=1$ and

$$
\begin{aligned}
\left\langle f_{R}, H_{1}(Y, \mathcal{Z}) f_{R}\right\rangle & =\left\langle f_{R},\left(T_{1}-\sum_{j=1}^{M} \frac{Z_{j} \alpha}{\left|x_{1}-y_{j}\right|}\right) f_{R}\right\rangle+\sum_{1 \leqslant i<j \leqslant M} \frac{Z_{i} Z_{j} \alpha}{\left|y_{i}-y_{j}\right|} \\
& \leqslant\left\langle f_{R},\left(-\Delta_{1}-\sum_{j=1}^{M} \frac{Z_{j} \alpha}{\left|x-y_{j}\right|}\right) f_{R}\right\rangle+\sum_{1 \leqslant i<j \leqslant M} \frac{Z_{i} Z_{j} \alpha}{\left|y_{i}-y_{j}\right|} \\
& \leqslant \frac{1}{R^{2}}\left\langle f,-\Delta_{1} f\right\rangle-\sum_{j=1}^{M}\left\langle f, \frac{Z_{j} \alpha}{\left|R x_{1}-y_{j}\right|} f\right\rangle+\sum_{1 \leqslant i<j \leqslant M} \frac{Z_{i} Z_{j} \alpha}{\left|y_{i}-y_{j}\right|}
\end{aligned}
$$

where we have used that $T_{1} \leqslant-\Delta_{1}$ and a change of variables in the second and third line, respectively. With a Taylor expansion, we see that the second term decays as $-\sum_{j=1}^{M} \alpha Z_{j} / R$ when $R$ becomes large. This proves that, for $R$ large enough,

$$
\left\langle f_{R}, H_{1}(Y, \mathcal{Z}) f_{R}\right\rangle<\sum_{1 \leqslant i<j \leqslant M} \frac{Z_{i} Z_{j} \alpha}{\left|y_{i}-y_{j}\right|}=\inf \sigma_{\text {ess }}\left(H_{1}(Y, \mathcal{Z})\right)
$$

Then, $\hat{H}_{1}(Y, \mathcal{Z})$ has at least one eigenvalue below the infimum of its essential spectrum.
We consider now the induction step. Let us assume, for some $N \in \mathbb{N}$ with $N<|Z|+1$,

$$
E_{N-1}:=\inf \sigma\left(\hat{H}_{N-1}(Y, \mathcal{Z})\right)<\inf \sigma_{\mathrm{ess}}\left(\hat{H}_{N-1}(Y, \mathcal{Z})\right)
$$

which implies that there exists a ground state $\Psi_{N-1} \in \mathcal{D}\left(\hat{H}_{N-1}(Y, \mathcal{Z})\right)$ with $\left\|\psi_{N-1}\right\|=1$ such that $H_{N-1}(Y, \mathcal{Z}) \Psi_{N-1}=E_{N-1} \Psi_{N-1}$. By the assumption of Theorem 1.4, we can choose and fix $\delta>0$ with

$$
\begin{equation*}
\delta<\frac{|\mathcal{Z}|+1-N}{2(N-1+|\mathcal{Z}|)} \tag{B.16}
\end{equation*}
$$

For $R$ large enough, $\chi_{\delta R}^{\otimes N-1} \Psi_{N-1}$ is not identically 0 , where $\chi_{\delta R}$ was defined in (2.5). Thus, we can define the sequence of trial functions:

$$
\begin{equation*}
\Phi_{N, R}:=\frac{Q_{N}\left(\left(\chi_{\delta R}^{\otimes N-1} \Psi_{N-1}\right) \otimes f_{R}\right)}{\left\|Q_{N}\left(\left(\chi_{\delta R}^{\otimes N-1} \Psi_{N-1}\right) \otimes f_{R}\right)\right\|} \tag{B.17}
\end{equation*}
$$

For all $R$, we have $\Phi_{N, R} \in H^{1 / 2}\left(\mathbb{R}^{3 N}\right) \cap \operatorname{Ran} Q_{N}$.
Arguing as in the proof of (2.15), we find that

$$
\begin{equation*}
\left\langle\Phi_{N, R},\left(\hat{H}_{N}-E_{N-1}\right) \Phi_{N, R}\right\rangle=\frac{\left\langle\left(\chi_{\delta R}^{\otimes N-1} \Psi_{N-1}\right) \otimes f_{R},\left(H_{N}-E_{N-1}\right)\left(\chi_{\delta R}^{\otimes N-1} \Psi_{N-1}\right) \otimes f_{R}\right\rangle}{\left\|\left(\chi_{\delta R}^{\otimes N-1} \Psi_{N-1}\right) \otimes f_{R}\right\|^{2}} \tag{B.18}
\end{equation*}
$$

Let us first estimate the denominator of the right-hand side of (B.18). Since, by Corollary 1.3, $\Psi_{N-1}$ is exponentially decaying and we have chosen $\Psi_{N-1}$ and $f_{R}$ with norm 1 , we find that $\|\left(\chi_{\delta R}^{\otimes N-1} \Psi_{N-1}\right) \otimes$ $f_{R} \|=1+O\left(e^{-c R}\right)$.

To estimate the numerator of the right-hand side of (B.18), let us split

$$
\begin{equation*}
H_{N}-E_{N-1}=H_{N}-H_{N-1}+H_{N-1}-E_{N-1} \tag{B.19}
\end{equation*}
$$

Using that $\left(H_{N-1}-E_{N-1}\right)$ does not act on the coordinate $x_{N}$ and that $\left(H_{N-1}-E_{N-1}\right) \Psi_{N-1}=0$, we find

$$
\left\langle\left(\chi_{\delta R}^{\otimes N-1} \Psi_{N-1} \otimes f_{R}\right),\left(H_{N-1}-E_{N-1}\right)\left(\chi_{\delta R}^{\otimes N-1} \Psi_{N-1} \otimes f_{R}\right)\right\rangle=\left\langle\chi_{\delta R}^{\otimes N-1} \Psi_{N-1}, \sum_{j=1}^{N-1}\left[T_{j}, \chi_{\delta R}^{\otimes N-1}\right] \Psi_{N-1}\right\rangle
$$

As a consequence, since, by Corollary $1.3, \Psi_{N-1}$ is exponentially decaying, we see, using (A.9), that

$$
\begin{equation*}
\left\langle\left(\chi_{\delta R}^{\otimes N-1} \Psi_{N-1} \otimes f_{R}\right),\left(H_{N-1}-E_{N-1}\right)\left(\chi_{\delta R}^{\otimes N-1} \Psi_{N-1} \otimes f_{R}\right)\right\rangle=O\left(e^{-d R}\right) \tag{B.20}
\end{equation*}
$$

Let us now bound

$$
\begin{equation*}
\left\langle\left(\chi_{\delta R}^{\otimes N-1} \Psi_{N-1} \otimes f_{R}\right),\left(H_{N}-H_{N-1}\right)\left(\chi_{\delta R}^{\otimes N-1} \Psi_{N-1} \otimes f_{R}\right)\right\rangle \tag{B.21}
\end{equation*}
$$

Since $T \leqslant \frac{-\Delta}{2}$, we can do it in the following way:

$$
\begin{align*}
& \left\langle\left(\chi_{\delta R}^{\otimes N-1} \Psi_{N-1} \otimes f_{R}\right),\left(H_{N}-H_{N-1}\right)\left(\chi_{\delta R}^{\otimes N-1} \Psi_{N-1} \otimes f_{R}\right)\right\rangle \\
& \quad \leqslant\left\langle\chi_{\delta R}^{\otimes N-1} \Psi_{N-1} \otimes f_{R},\left(-\frac{\Delta_{N}}{2}-\sum_{j=1}^{M} \frac{Z_{j} \alpha}{\left|x_{N}-y_{j}\right|}+\sum_{j=1}^{N-1} \frac{\alpha}{\left|x_{N}-x_{j}\right|}\right)\left(\chi_{\delta R}^{\otimes N-1} \Psi_{N-1} \otimes f_{R}\right)\right\rangle . \tag{B.22}
\end{align*}
$$

Assume that $R$ is large enough such that for all $j,\left|y_{j}\right| \leqslant \delta R$. For $x_{N}$ in the support of $f_{R}$, we have that $\left|x_{N}\right|=(1+\theta) R$, for some $\theta \in[0,1]$. Therefore,

$$
\begin{equation*}
\frac{1}{\left|x_{N}-y_{j}\right|} \geqslant \frac{1}{\left|x_{N}\right|+\left|y_{j}\right|} \geqslant \frac{1}{(1+\theta+\delta) R} . \tag{B.23}
\end{equation*}
$$

Similarly, for $\left(x_{1}, \ldots, x_{N-1}\right)$ in the support of $\chi_{\delta R}^{\otimes N-1} \Psi_{N-1}$, we have $\left|x_{j}\right|<\delta R$ for all $j=1, \ldots, N-1$. Therefore,

$$
\begin{equation*}
\frac{1}{\left|x_{N}-x_{j}\right|} \leqslant \frac{1}{\left|x_{N}\right|-\left|x_{j}\right|} \leqslant \frac{1}{(1+\theta-\delta) R} \tag{B.24}
\end{equation*}
$$

Using (B.23) and (B.24) and recalling (B.17) we find that, on $\operatorname{supp} \Phi_{N, R}$,

$$
-\sum_{j=1}^{M} \frac{Z_{j}}{\left|x_{N}-y_{j}\right|}+\sum_{j=1}^{N-1} \frac{1}{\left|x_{N}-x_{j}\right|} \leqslant \frac{(1+\theta)(N-1-|\mathcal{Z}|)+\delta(N-1+|\mathcal{Z}|)}{(1+\theta+\delta)(1+\theta-\delta) R}
$$

As a consequence, using (B.16), we find, on $\operatorname{supp} \Phi_{N, R}$,

$$
\begin{equation*}
-\sum_{j=1}^{N} \frac{Z_{j}}{\left|x_{N}-y_{j}\right|}+\sum_{j=1}^{N-1} \frac{1}{\left|x_{N}-x_{j}\right|} \leqslant \frac{\left(\frac{1}{2}+\theta\right)(N-1-|\mathcal{Z}|)}{(1+\theta+\delta)(1+\theta-\delta) R} \leqslant \frac{N-1-|\mathcal{Z}|}{8 R} \tag{B.25}
\end{equation*}
$$

where in the last step we used that $\frac{1}{2}+\theta \geqslant \frac{1}{2}$ and that $(1+\theta+\delta)(1+\theta-\delta) \leqslant 4$ for all $\theta \in[0,1]$ as well as the assumed inequality $N-1-|\mathcal{Z}|<0$.

Inserting first $\left\langle f_{R},\left(-\Delta_{x}\right) f_{R}\right\rangle=\frac{1}{R^{2}}\left\langle f_{1},\left(-\Delta_{x}\right) f_{1}\right\rangle$ and (B.25) into (B.22) and then combining it with (B.20) and (B.19), we find that

$$
\begin{equation*}
\left\langle\left(\chi_{\delta R}^{\otimes N-1} \Psi_{N-1} \otimes f_{R}\right),\left(H_{N}-E_{N-1}\right)\left(\chi_{\delta R}^{\otimes N-1} \Psi_{N-1} \otimes f_{R}\right)\right\rangle \leqslant-\frac{\alpha(|\mathcal{Z}|+1-N)}{8 R}+\frac{D_{1}}{R^{2}} \tag{B.26}
\end{equation*}
$$

for some constant $D_{1}>0$.
We have thus that, for $R$ large enough,

$$
\begin{equation*}
\left\langle\Phi_{N, R},\left(\hat{H}_{N}-E_{N-1}\right) \Phi_{N, R}\right\rangle<0 \tag{B.27}
\end{equation*}
$$

which implies that $\hat{H}_{N}$ has an eigenvalue below its essential spectrum since $E_{N-1}=\inf \sigma_{\text {ess }}\left(\hat{H}_{N}\right)$.
Let us now prove that there are infinitely many eigenvalues below $E_{N-1}$. Let $n_{0}$ be an integer large enough such that, for $\Phi_{N, 2^{n_{0}}}$ defined in (B.17),

$$
\begin{equation*}
\left\langle\Phi_{N, 2^{n_{0}}},\left(\hat{H}_{N}-E_{N-1}\right) \Phi_{N, 2^{n_{0}}}\right\rangle \leqslant-\frac{C}{2^{n_{0}}}+\frac{D_{1}}{4^{n_{0}}} \tag{B.28}
\end{equation*}
$$

for some $C>0$ and $D_{1}$ as in (B.26).
Let $n \geqslant n_{0}$ : we will prove that $\hat{H}_{N}-E_{N-1}$ has at least $n$ negative eigenvalues. As far as we know, this part of the proof is new, even for atoms, since [14] does not consider this question. Let us consider a normed linear combination of $\Phi_{N, 2^{k}}$ 's with $k$ between $2 n$ and $3 n$ :

$$
\begin{equation*}
\Xi:=\sum_{i=1}^{n} c_{i} \Phi_{N, 2^{2 n+i}}, \quad \text { with } \quad \sum_{i=1}^{n}\left|c_{i}\right|^{2}=1 \tag{B.29}
\end{equation*}
$$

Then,

$$
\begin{align*}
& \left\langle\Xi,\left(\hat{H}_{N}-E_{N-1}\right) \Xi\right\rangle=\sum_{i=1}^{n}\left|c_{i}\right|^{2}\left\langle\Phi_{N, 2^{2 n+i}},\left(\hat{H}_{N}-E_{N-1}\right) \Phi_{N, 2^{2 n+i}}\right\rangle \\
&  \tag{B.30}\\
& \quad+\sum_{i \neq j}\left\langle\Phi_{N, 2^{2 n+i}},\left(\hat{H}_{N}-E_{N-1}\right) \Phi_{N, 2^{2 n+j}}\right\rangle
\end{align*}
$$

Since (B.28) remains satisfied if we replace $2^{n_{0}}$ by $k$ for any $k \geqslant 2^{n_{0}}$, we get that

$$
\begin{equation*}
\sum_{i=1}^{n}\left|c_{i}\right|^{2}\left\langle\Phi_{N, 2^{2 n+i}},\left(\hat{H}_{N}-E_{N-1}\right) \Phi_{N, 2^{2 n+i}}\right\rangle \leqslant-\sum_{i=1}^{n}\left|c_{i}\right|^{2} \frac{C}{2^{2 n+i}}+\frac{D_{1}}{4^{2 n}} \leqslant-\frac{C}{8^{n}}+\frac{D_{1}}{16^{n}} \tag{B.31}
\end{equation*}
$$

For the remaining terms, for support reasons, the only parts which are not 0 are of the form

$$
\left\langle\Phi_{N, 2^{2 n+i}}, T_{k} \Phi_{N, 2^{2 n+j}}\right\rangle
$$

Assume without loss of generality that $i<j$. Since the supports of $\Phi_{N, 2^{2 n+i}}$ and $\Phi_{N, 2^{2 n+j}}$ have distance of order $2^{n+j}$, we get from Lemma A. 1 that

$$
\begin{equation*}
\left\langle\Phi_{N, 2^{2 n+i}}, T_{k} \Phi_{N, 2^{2 n+j}}\right\rangle \leqslant C^{\prime} e^{-2^{n+j}} \leqslant C e^{-2^{n}} \tag{B.32}
\end{equation*}
$$

Inserting (B.31) and (B.32) into (B.30), we get that

$$
\left\langle\Xi,\left(\hat{H}_{N}-E_{N-1}\right) \Xi\right\rangle<0
$$

for all such $\Xi$ in the form (B.29), given that $n$ is large enough.
The linear space spanned by the $\phi_{N, 2^{k}}$ for $k$ between $2 n$ and $3 n$ is thus an $n+1$-dimensional space on which the quadratic form associated with $\hat{H}_{N}-E_{N-1}$ is negative. Therefore, $\hat{H}_{N}$ has at least $n+1$ eigenvalues below its essential spectrum. Since this is true for arbitrary large $n$, there are infinitely many eigenvalues below the bottom of the essential spectrum.

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[^0]:    ${ }^{1}$ Two directions are needed since the path could be already tangent to one direction of descent, as it is for a mountain pass.

