# CLT for real $\beta$-ensembles at high temperature 

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# CLT for real $\beta$-Ensembles at High Temperature* 

Charlie Dworaczek Guera ${ }^{\dagger}$ Ronan Memin ${ }^{\ddagger}$


#### Abstract

We establish a central limit theorem for the fluctuations of the empirical measure in the betaensemble of dimension $N$ at a temperature proportional to $N$ and with convex, smooth potential. The space of test functions for which the CLT holds includes $\mathcal{C}^{1}$, vanishing functions at infinity. It is obtained by the inversion of an operator which is a pertubation of a Sturm-Liouville operator. The method that we use is based on a change of variables introduced in BFG15] and in Shc14].


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## 1 Introduction and main result

The beta-ensemble of dimension $N \geqslant 1$ with parameter $\beta>0$ and potential $V$ is the probability measure on $\mathbb{R}^{N}$ given by

$$
\begin{equation*}
d \mathbf{P}_{N}^{\beta, V}\left(x_{1}, \ldots, x_{N}\right)=\frac{1}{Z_{N}(V, \beta)} \prod_{i<j}\left|x_{i}-x_{j}\right|^{\beta} e^{-\sum_{i=1}^{N} V\left(x_{i}\right)} d x_{1} \ldots d x_{N} \tag{1}
\end{equation*}
$$

The potential $V$ has to be chosen so that the partition function

$$
Z_{N}(V, \beta)=\int_{\mathbb{R}^{N}} \prod_{i<j}\left|x_{i}-x_{j}\right|^{\beta} e^{-\sum_{i=1}^{N} V\left(x_{i}\right)} d x_{1} \ldots d x_{N}
$$

[^0]is finite. This is the case for example if for some $\beta^{\prime}>\max (1, \beta)$,
\[

$$
\begin{equation*}
\liminf _{|x| \rightarrow \infty} \frac{V(x)}{N \beta^{\prime} \ln |x|}>1 \tag{2}
\end{equation*}
$$

\]

see [AGZ10, equation (2.6.2)]. The parameter $\beta$, which is allowed to depend on $N$, is the so-called inverse temperature.

Under the special choice $V(x)=\frac{x^{2}}{2}$, the measure (1) can be seen as the joint law of the (unordered) eigenvalues of certain matrix models:

- For $\beta=1$ (resp. $\beta=2$ ), it is the law of the eigenvalues of the Gaussian Orthogonal Ensemble (resp. Gaussian Unitary Ensemble), AGZ10 [Theorem 2.5.2].
- For general $\beta>0$, potentially depending on $N$, it is the law of the spectrum of certain tri-diagonal random matrices as shown by Dumitriu and Edelman in DE02.

We consider here the high temperature regime where $\beta$ scales as $1 / N$, and write $\beta=\frac{2 P}{N}$ for some $P>0$. The corresponding measure is therefore

$$
\begin{equation*}
d \mathbb{P}_{N}^{V, P}\left(x_{1}, \ldots, x_{N}\right)=\frac{1}{Z_{N}^{V, P}} \prod_{i<j}\left|x_{i}-x_{j}\right|^{\frac{2 P}{N}} e^{-\sum_{i=1}^{N} V\left(x_{i}\right)} d x_{1} \ldots d x_{N} \tag{3}
\end{equation*}
$$

with partition function

$$
\begin{equation*}
Z_{N}^{V, P}=\int_{\mathbb{R}^{N}} \prod_{i<j}\left|x_{i}-x_{j}\right|^{\frac{2 P}{N}} e^{-\sum_{i=1}^{N} V\left(x_{i}\right)} d x_{1} \ldots d x_{N} \tag{4}
\end{equation*}
$$

It was shown in GZ19 that under $\mathbb{P}_{N}^{V, P}$, the sequence of empirical measures

$$
\hat{\mu}_{N}=\frac{1}{N} \sum_{i=1}^{N} \delta_{x_{i}}
$$

satisfies a large deviation principle at speed $N$ with strictly convex, good rate function. As a consequence, $\hat{\mu}_{N}$ converges almost surely in distribution towards a deterministic measure $\mu_{P}^{V}$ as $N$ goes to infinity, meaning that almost surely, for every bounded continuous $f: \mathbb{R} \rightarrow \mathbb{R}$,

$$
\int_{\mathbb{R}} f d \hat{\mu}_{N} \underset{N \rightarrow \infty}{\longrightarrow} \int_{\mathbb{R}} f d \mu_{P}^{V}
$$

The limiting measure $\mu_{P}^{V}$ can be seen to have a density $\rho_{P}^{V}$ which satisfies for almost every $x \in \mathbb{R}$

$$
\begin{equation*}
V(x)-2 P \int_{\mathbb{R}} \ln |x-y| \rho_{P}^{V}(y) d y+\ln \rho_{P}^{V}(x)=\lambda_{P}^{V} \tag{5}
\end{equation*}
$$

where $\lambda_{P}^{V}$ is constant (see GM22, Lemma 3.2] for example).
The $\beta$-ensemble in the regime $\beta N \underset{N \rightarrow \infty}{\longrightarrow} 2 P>0$ has drawn a lot of attention from the random matrix and statistical physics communities lately. The limiting density was first described in the case of the quadratic potential in ABG12, as a crossover between the Wigner semicircle law (fixed $\beta>0$ case) and the Gaussian density (case $\beta=0$ ). The fluctuations of the eigenvalues in the bulk and at the edge of a configuration were studied for example in [BGP15, NT18, NT20, Pak18, Lam21. These fluctuations were shown to be described by Poisson statistics in this regime. Recently, Spohn uncovered in Spo20] a link between the study of the Classical Toda chain and the $\beta$-ensemble in this regime, the result was later extended to more general potentials in GM22. It is explained in Spo21 [Section 6] how the central
limit theorem for the empirical measure in the high temperature $\beta$ ensemble is linked with the currents of the Toda chain.

The Central Limit Theorem for the fluctuations of the linear statistics of beta-ensembles was first established by Joh98 for $\beta=2$ polynomial potential, then generalized and further developed in the regime where $\beta$ is fixed in Shc13, BG13a, BG13b, BLS18. Also an optimal local law was found in this regime in BMP22. The CLT was obtained in the high-temperature regime $\beta N \rightarrow 2 P>0$ by Nakano and Trinh in [NT18, Theorem 4.9] for quadratic $V$, relying on the tridiagonal representation for the beta-ensemble with quadratic potential in [DE02. In [HL21], the authors prove the CLT in the case of the circular beta-ensemble at high temperature with general potential, using a normal approximation method involving the spectral analysis of an operator associated to the limiting covariance structure. Their method allowed them to derive a Berry-Esseen bound, i.e. a speed of convergence of the fluctuations towards a Gaussian variable.

In this paper, we adapt part of the arguments of HL21] to our setup. More precisely, we show that for a class of regular, convex potentials $V$ satisfying a growth condition of the type

$$
\lim _{|x| \rightarrow \infty} \frac{V^{\prime \prime}(x)}{V^{\prime}(x)^{2}}=0
$$

denoting $\nu_{N}=\hat{\mu}_{N}-\mu_{P}^{V}$ and considering test functions $f$ belonging to the range of a certain integrodifferential operator, the scaled fluctuations of $\hat{\mu}_{N}$, defined by

$$
\sqrt{N} \nu_{N}(f):=\sqrt{N}\left(\int_{\mathbb{R}} f d \mu_{N}-\int_{\mathbb{R}} f d \mu_{P}^{V}\right)
$$

converge in law towards centered Gaussian law with variance depending on $f$.
When considering the fixed temperature regime, i.e. $\beta$ fixed, one has to renormalize the $x_{i}$ 's by $\sqrt{N}$. It is shown in AGZ10] Theorem 2.6.1] that the measure

$$
\frac{1}{N} \sum_{i=1}^{N} \delta_{x_{i} / \sqrt{N}}
$$

satisfies a large deviation principle, and the limiting measure is characterized in AGZ10][Lemma 2.6.2] by an equation similar to (5). In fact, the term $\ln \rho_{P}^{V}$ in the left-hand side of (5) is the only difference in the equation characterizing the limiting measure in the fixed $\beta$ case. We point out the very similar characterization of the equilibrium measure corresponding to the minimization problem arising in [BGK16]. There again, the limiting measure is compactly supported. The term $\ln \rho_{P}^{V}$ is of prime importance because its presence implies that the support of $\rho_{P}^{V}$ is the whole real line. It leads to technicalities to deal with the behavior at infinity of most of the associated objects, namely dealing with weighted Lebesgue spaces $L^{2}\left(\rho_{P}^{V}\right)$ and the corresponding Sobolev spaces $H^{k}\left(\rho_{P}^{V}\right)$.

Our strategy is based on a change of variables in the partition function $Z_{N}^{V, P}(4)$, used for the betaensemble at fixed temperature introduced in BFG15 and Shc14, and used in Gui19 and in BGK16] to derive the loop equations and in BLS18 to derive a CLT in the $\beta$-ensemble with $\beta$ fixed. The outline of the argument goes as follows: Take $\phi: \mathbb{R} \rightarrow \mathbb{R}$ smooth, vanishing fast enough at infinity, and do the change of variables in $Z_{N}^{V, P}, x_{i}=y_{i}+\frac{t}{\sqrt{N}} \phi\left(y_{i}\right), 1 \leqslant i \leqslant N$, to get

$$
Z_{N}^{P, V}=\int_{\mathbb{R}^{N}} \prod_{i<j}\left|y_{i}-y_{j}+\frac{t}{\sqrt{N}}\left(\phi\left(y_{i}\right)-\phi\left(y_{j}\right)\right)\right|^{2 P / N} e^{-\sum_{i=1}^{N} V\left(y_{i}+\frac{t}{\sqrt{N}} \phi\left(y_{i}\right)\right)} \prod_{i=1}^{N}\left(1+\frac{t}{\sqrt{N}} \phi^{\prime}\left(y_{i}\right)\right) d^{N} \mathbf{y}
$$

Expanding the different terms in this integral, one gets

$$
Z_{N}^{P, V}=\int_{\mathbb{R}^{N}} \prod_{i<j}\left|y_{i}-y_{j}\right|^{\frac{2 P}{N}} e^{-\sum_{i=1}^{N} V\left(y_{i}\right)} e^{\frac{t}{\sqrt{N}}\left[\frac{2 P}{N} \sum_{i<j} \frac{\phi\left(y_{i}\right)-\phi\left(y_{j}\right)}{y_{i}-y_{j}}+\sum_{i=1}^{N}\left(\phi^{\prime}\left(y_{i}\right)-V^{\prime}\left(y_{i}\right) \phi\left(y_{i}\right)\right)\right]} e^{\frac{t^{2}}{2} \sigma_{N}^{2}(\phi)} d^{N} \mathbf{y}
$$

where the term $\sigma_{N}^{2}(\phi)$ converges towards a limiting variance $\sigma^{2}(\phi)$ depending on $\phi, P$ and $V$. After dividing both parts of the equation by $Z_{N}^{P, V}$, and because of equation (5) characterizing $\mu_{P}^{V}$, one can deduce from the last equation the convergence of the Laplace transform

$$
\begin{equation*}
\mathbb{E}\left[e^{t \sqrt{N}\left(\nu_{N}(\Xi \phi)+\text { error term }\right)}\right] \underset{N \rightarrow \infty}{\longrightarrow} e^{\frac{t^{2}}{2} \sigma^{2}(\phi)} \tag{6}
\end{equation*}
$$

where $\Xi$ is a linear operator acting on test functions and defined by

$$
\begin{equation*}
(\Xi \phi)(x)=2 P \int_{\mathbb{R}} \frac{\phi(x)-\phi(y)}{x-y} d \mu_{P}^{V}(y)+\phi^{\prime}(x)-V^{\prime}(x) \phi(x) \tag{7}
\end{equation*}
$$

Once the error term is taken care of, (6) shows the central limit theorem for test functions of the form $\Xi \phi$. Following [HL21], the operator $\mathcal{L}$ given by

$$
\begin{equation*}
\mathcal{L} \phi=\Xi \phi^{\prime} \tag{8}
\end{equation*}
$$

can be analyzed using Hilbert space techniques. In particular, the operator $\mathcal{L}$, seen as an unbounded operator of the Hilbert space

$$
\mathrm{H}=\left\{u \in L^{2}\left(\rho_{P}^{V}\right) \mid u^{\prime} \in L^{2}\left(\rho_{P}^{V}\right), \int_{\mathbb{R}} u \rho_{P}^{V} d x=0\right\}, \quad\langle u, v\rangle_{\mathrm{H}}=\left\langle u^{\prime}, v^{\prime}\right\rangle_{L^{2}\left(\rho_{P}^{V}\right)}
$$

can be decomposed as

$$
-\mathcal{L}=\mathcal{A}+2 P \mathcal{W}
$$

where $\mathcal{A}$ is a positive Sturm-Liouville operator and $\mathcal{W}$ is positive and self-adjoint. Such a writing allows us to show that $-\mathcal{L}$ is invertible, see Theorem 6.7.

We now state the assumptions we make on the potential $V$.
Assumptions 1.1. The potential $V$ satisfies:
i) $V \in \mathcal{C}^{3}(\mathbb{R})$, goes to infinity at infinity and is convex.
ii) For all polynomial $Q \in \mathbb{R}[X]$ and $\alpha>0, Q\left(V^{\prime}(x)\right) e^{-V(x)}=\underset{|x| \rightarrow \infty}{o}\left(x^{-\alpha}\right)$.
iii) $\left|V^{\prime}(x)\right| \underset{|x| \rightarrow+\infty}{\longrightarrow}+\infty$ and $x \mapsto \frac{1}{V^{\prime}(x)^{2}}$ is integrable at infinity. Furthermore, for any sequence $x_{N}$ such that $\left|x_{N}\right|$ goes to infinity, and for all real $a<b$, we have, as $N$ goes to infinity,

$$
\frac{1}{V^{\prime}\left(x_{N}\right)^{2}} \sup _{a \leqslant x \leqslant b}\left|V^{\prime \prime}\left(x_{N}+x\right)\right| \underset{N \rightarrow \infty}{\longrightarrow} 0
$$

iv) $\frac{V^{\prime \prime}(x)}{V^{\prime}(x)}=\underset{|x| \rightarrow \infty}{O}$ (1) and $\frac{V^{(3)}(x)}{V^{\prime}(x)}=\underset{|x| \rightarrow \infty}{O}$ (1).

The convexity assumption in $i$ ) will guarantee the Poincaré inequality, stated in Proposition 2.4 for the equilibrium measure $\mu_{P}^{V}$. Because $i$ ) implies that $V$ goes to infinity faster than linearly, we will see that it also ensures exponential decay at infinity of $\rho_{P}^{V}$. Recalling the sufficient condition for $\mathbb{P}_{N}^{V, P}$ of equation (2) to be defined, this first assumption implies that there exists $\alpha>0$ such that $\lim \inf |x| \rightarrow \infty \frac{V(x)}{|x|}>\alpha$. This guarantees in particular that the beta-ensemble (3) is well-defined for all $N \geqslant 1$ and $P \geqslant 0$.

The second assumption ensures that any power of $V^{\prime}$ and $V^{\prime \prime}$ is in $L^{2}\left(\rho_{P}^{V}\right)$ and $\rho_{P}^{V}$, which behaves like $e^{-V}$ up to a sub-exponential factor, belongs to the Sobolev space $H^{2}(\mathbb{R}) \subset \mathcal{C}^{1}(\mathbb{R})$. Indeed, for $k \leqslant 2$, using $i v$ ), $\rho_{P}^{V}{ }^{(k)}$ behaves at infinity like $\left(V^{\prime}\right)^{k} \rho_{P}^{V}$ as shown in 2.2 which is in $L^{2}(\mathbb{R})$ by assumption $\left.i i\right)$.

Assumption iii) will be used to localize the minimum/maximum point of a typical configuration $\left(x_{1}, \ldots, x_{N}\right)$ following the law $\mathbb{P}_{N}^{V, P}$ : this will be done in Corollary 4.2 which comes as a consequence of

Lam21 [Theorem 3.4]. More precisely, Corollary 4.2 establishes that for sequences $\left(\alpha_{N}^{+}\right)_{N},\left(\alpha_{N}^{-}\right)_{N}$ going to infinity, the random variables

$$
\alpha_{N}^{+}\left(\max _{1 \leqslant j \leqslant N} x_{j}-E_{N}^{+}\right) \quad \text { and } \quad \alpha_{N}^{-}\left(\max _{1 \leqslant j \leqslant N} x_{j}-E_{N}^{-}\right)
$$

converge in distribution. For large $N$, the scalars $E_{N}^{+}$and $E_{N}^{-}$can thus be seen as the edges of a typical configuration. Furthermore,

$$
\begin{equation*}
V\left(E_{N}^{ \pm}\right) \sim \ln N \tag{9}
\end{equation*}
$$

We refer to Section 4 for detailed statements. We will need another technical assumption to ensure that Taylor remainders arising in the proof of Theorem 5.2 are negligible. This final step in the proof of Theorem 1.3 consists in lifting the result of Proposition 5.1 from compactly supported functions to more general functions. We use Assumption $i v$ ) to ensure that $\mathcal{L}^{-1}$ is regular enough $i e$ that for sufficiently smooth functions $f,\left(\mathcal{L}^{-1} f\right)^{\prime} \in H^{2}(\mathbb{R})$.
Assumption 1.2. With the notations of Theorem 4.1, we have

$$
\sup _{d\left(x, I_{N}\right) \leqslant 1}\left|V^{(3)}(x)\right|=o\left(N^{1 / 2}\right)
$$

where $I_{N}=\left[E_{N}^{-}-2 ; E_{N}^{+}+2\right]$.
In view of the asymptotics (9), Assumption 1.2 is reasonable, at least in the cases where $V$ is of the form $|x|^{a}+R(x)$ for $a=2$ or $a>3$ (we choose $a$ so that $V$ is three times differentiable), and with $R \in \mathcal{C}^{3}(\mathbb{R})$, convex, small enough (see the comment following Theorem 1.3 or $V=$ cosh, for example. On the other hand a scaled potential like $e^{x^{2}}$ doesn't satisfy assumptions iii) ans iv)..

We are now able to state the main result, ie the central limit theorem for functions belonging to the image of the operator $\mathcal{L}$ introduced in (8).

Theorem 1.3. Assume that $V$ satisfies Assumptions 1.1 and Assumption 1.2. Then for $\phi$ verifying the following conditions:

- $\phi \in \mathcal{C}^{1}(\mathbb{R})$
- $\phi(x)=O(1 / x)$ and $\phi^{\prime}(x)=O\left(1 / x^{2}\right)$ at infinity
- $\int_{\mathbb{R}} \phi(x) d \mu_{P}^{V}(x)=0$
we have the convergence in law

$$
\begin{equation*}
\sqrt{N} \nu_{N}(\phi) \rightarrow \mathcal{N}\left(0,\left(\sigma_{P}^{V}\right)^{2}(\phi)\right) \tag{10}
\end{equation*}
$$

where the limiting variance $\left(\sigma_{P}^{V}\right)^{2}(\phi)$ is given by

$$
\begin{align*}
& \left(\sigma_{P}^{V}\right)^{2}(\phi)=\int_{\mathbb{R}}\left(\left(\mathcal{L}^{-1} \phi\right)^{\prime \prime}(x)^{2}+V^{\prime \prime}(x)\left(\mathcal{L}^{-1} \phi\right)^{\prime}(x)^{2}\right) d \mu_{P}^{V}(x) \\
& \quad+P \iint_{\mathbb{R}^{2}}\left(\frac{\left(\mathcal{L}^{-1} \phi\right)^{\prime}(x)-\left(\mathcal{L}^{-1} \phi\right)^{\prime}(y)}{x-y}\right)^{2} d \mu_{P}^{V}(x) d \mu_{P}^{V}(y) \tag{11}
\end{align*}
$$

Remark 1.4. Since $\nu_{N}(\phi+c)=\nu_{N}(\phi)$ for all constant $c \in \mathbb{R}$, the assumption $\int_{\mathbb{R}} \phi(x) d \mu_{P}^{V}=0$ can be dropped by replacing $\phi$ by $\phi-\int_{\mathbb{R}} \phi(x) d \mu_{P}^{V}$ in the expression of the limiting variance.

As a tool to deal with the error term of equation (6), we establish a concentration inequality for the empirical measure. This inequality is stated in terms of the following distance over the set of probability distributions $\mathcal{P}(\mathbb{R})$.

For $\mu, \mu^{\prime} \in \mathcal{P}(\mathbb{R})$ we define the distance

$$
\begin{equation*}
d\left(\mu, \mu^{\prime}\right)=\sup _{\substack{\|f\|_{\text {Lip }} \leqslant 1 \\\|f\|_{1 / 2} \leqslant 1}}\left\{\left|\int f d \mu-\int f d \mu^{\prime}\right|\right\} \tag{12}
\end{equation*}
$$

where $\|f\|_{\text {Lip }}$ denotes the Lipschitz constant of $f$, and $\|f\|_{1 / 2}^{2}=\int_{\mathbb{R}}|t||\mathcal{F}[f](t)|^{2} d t$, where $\mathcal{F}$ denotes the Fourier transform on $L^{2}(\mathbb{R})$ which takes the following expression $\mathcal{F}[f](t)=\int_{\mathbb{R}} f(x) e^{-\mathrm{i} t x} d x$ for $f \in L^{1}(\mathbb{R}) \cap L^{2}(\mathbb{R})$.

We then have
Theorem 1.5. There exists $K \in \mathbb{R}$ (depending on $P$ and on $V$ ), such that for any $N \geqslant 1$ and $r>0$,

$$
\begin{equation*}
\mathbb{P}_{N}^{V, P}\left(d\left(\hat{\mu}_{N}, \mu_{P}^{V}\right)>r\right) \leqslant e^{-N r^{2} \frac{P \pi^{2}}{2}+5 P \ln N+K} \tag{13}
\end{equation*}
$$

This result is the analog of [HL21, Theorem 1.4].
The paper is organized as follows. In Section 2 we discuss the regularity of the equilibrium density $\rho_{P}^{V}$ under Assumption 1.1. In Section 3 we prove Theorem 1.5 Section 4 is dedicated to the localization of the edge of a typical configuration, mentioned in the discussion preceding the statement of Assumption 1.2. We next prove in Section 5 the convergence of the Laplace transform of $\sqrt{N} \nu_{N}(\mathcal{L} \phi)$ for general functions $\phi$ which establishes Theorem 1.3 for functions of the form $\mathcal{L} \phi$. Section 6 is dedicated to the diagonalization and inversion of $\mathcal{L}$ given by (8). In Section 7 . we show regularity properties of $\mathcal{L}^{-1}$ to establish Theorem 1.3 We detail in Appendix A elements of proof for the spectral theory of Schrödinger operators, used in Section 6

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## 2 Regularity of the equilibrium measure and Hilbert transform

In this section, we discuss the regularity properties of the equilibrium density $\rho_{P}^{V}$, namely its decay at infinity and its smoothness, and give formulas for its two first derivatives.

The Hilbert transform, whose definition we recall, plays a central role in the analysis of the equilibrium measure. It is first defined on the Schwartz class through $\forall \phi \in \mathcal{S}(\mathbb{R}), \forall x \in \mathbb{R}$,

$$
\begin{equation*}
\mathcal{H}[\phi](x):=f_{\mathbb{R}} \frac{\phi(t)}{t-x} d t=\lim _{\varepsilon \downarrow 0} \int_{|t-x|>\varepsilon} \frac{\phi(t)}{t-x} d t=\int_{0}^{+\infty} \frac{\phi(x+t)-\phi(x-t)}{t} d t \tag{14}
\end{equation*}
$$

where $f$ denotes the Cauchy principal value integral, and then extended to $L^{2}(\mathbb{R})$ thanks to property ii) of Lemma 2.1. $\|f\|_{L^{2}(d x)}=\frac{1}{\pi}\|\mathcal{H}[f]\|_{L^{2}(d x)}$. The last expression in 14 is a definition where the integral converges in the classical sense.

We also recall the definition of the logarithmic potential $U^{f}$ of a density of probability $f: \mathbb{R} \rightarrow \mathbb{R}$, given for $x \in \mathbb{R}$ by

$$
\begin{equation*}
U^{f}(x)=-\int_{\mathbb{R}} \ln |x-y| f(y) d y \tag{15}
\end{equation*}
$$

Because we assume $f \in L^{1}(\mathbb{R})$ to be nonnegative, $U^{f}$ takes values in $[-\infty,+\infty)$. If $f$ integrates the function $\ln$, i.e $\int_{\mathbb{R}} \ln |x| f(x) d x<+\infty$, then $U^{f}$ takes real values.

Additionally, one can check that the logarithmic potential and the Hilbert transform of $f$ are linked through the distributional identity $\left(U^{f}\right)^{\prime}=\mathcal{H}[f]$.

We recall in the next lemma some properties of the Hilbert transform that we will use in the rest of the paper.
Lemma 2.1 (Properties of the Hilbert transform).
i) Fourier transform: For all $\phi \in L^{2}(\mathbb{R}), \mathcal{F}[\mathcal{H}[\phi]](\omega)=\mathrm{i} \pi \operatorname{sgn}(\omega) \mathcal{F}[\phi](\omega)$ for all $\omega \in \mathbb{R}$.
ii) As a consequence, $\frac{1}{\pi} \mathcal{H}$ is an isometry of $L^{2}(\mathbb{R})$, and $\mathcal{H}$ satisfies on $L^{2}(\mathbb{R})$ the identity $\mathcal{H}^{2}=-\pi^{2} I$.
iii) Derivative: For any $f \in H^{1}(\mathbb{R}), \mathcal{H}[f]$ is also $H^{1}(\mathbb{R})$ and $\mathcal{H}[f]^{\prime}=\mathcal{H}\left[f^{\prime}\right]$.
iv) For all $p>1$, the Hilbert transform can be extended as a bounded operator $\mathcal{H}: L^{p}(\mathbb{R}) \rightarrow L^{p}(\mathbb{R})$.
v) Skew-self adjointness: For any $f, g \in L^{2}(\mathbb{R}),\langle\mathcal{H}[f], g\rangle_{L^{2}(\mathbb{R})}=-\langle f, \mathcal{H}[g]\rangle_{L^{2}(\mathbb{R})}$.

Proof. We refer to [Kin09] for the proofs of these properties.
As a consequence of [GZ19, $\hat{\mu}_{N}$ converges almost surely under $\mathbb{P}_{N}^{V, P}$ towards the unique minimizer of the energy-functional $\mathcal{E}_{P}^{V}$, defined for $\mu \in \mathcal{P}(\mathbb{R})$ by

$$
\mathcal{E}_{P}^{V}(\mu)=\left\{\begin{array}{l}
\int_{\mathbb{R}}\left[V+\ln \left(\frac{d \mu}{d x}\right)\right] d \mu-P \iint_{\mathbb{R}^{2}} \ln |x-y| d \mu(x) d \mu(y) \text { if } \mu \ll d x  \tag{16}\\
+\infty \text { otherwise }
\end{array}\right.
$$

(Here we wrote $\mu \ll d x$ for " $\mu$ is absolutely continuous with respect to Lebesgue measure")
Consequently, following [GM22, Lemma 3.2], the density $\rho_{P}^{V}$ of $\mu_{P}^{V}$ satisfies equation (5), which we rewrite here for convenience.

$$
\begin{equation*}
V(x)-2 P \int_{\mathbb{R}} \ln |x-y| \rho_{P}^{V}(y) d y+\ln \rho_{P}^{V}(x)=\lambda_{P}^{V} \tag{17}
\end{equation*}
$$

where $\lambda_{P}^{V}$ is a constant (depending on $V$ and $P$ ). Using this equation, we show in the next lemma that $\rho_{P}^{V}$ decays exponentially and is twice continuously differentiable.

We now drop the superscript of $\rho_{P}^{V}$ and $\mu_{P}^{V}$ and denote it $\rho_{P}$ and $\mu_{P}$ for convenience.
Lemma 2.2. Under Assumption 1.1.

- The support of $\mu_{P}$ is $\mathbb{R}$ and there exists a constant $C_{P}^{V}$ such that for all $x \in \mathbb{R}$,

$$
\rho_{P}(x) \leqslant C_{P}^{V}(1+|x|)^{2 P} e^{-V(x)}
$$

- The density $\rho_{P}$ is in $\mathcal{C}^{2}(\mathbb{R})$ and we have

$$
\begin{equation*}
\rho_{P}^{\prime}=-\left(V^{\prime}+2 P \mathcal{H}\left[\rho_{P}\right]\right) \rho_{P} \tag{18}
\end{equation*}
$$

and

$$
\begin{equation*}
\rho_{P}^{\prime \prime}=\left(-2 P \mathcal{H}\left[\rho_{P}\right]^{\prime}-V^{\prime \prime}+V^{\prime 2}+4 P^{2} \mathcal{H}\left[\rho_{P}\right]^{2}+4 P V^{\prime} \mathcal{H}\left[\rho_{P}\right]\right) \rho_{P} \tag{19}
\end{equation*}
$$

Proof. For the first point, GM22, Lemma 3.2] establishes that the support of $\mu_{P}$ is the whole real axis, and that under the first condition of 1.1 , we have the bound, valid for all $x \in \mathbb{R}$

$$
\begin{equation*}
\rho_{P}(x) \leqslant \frac{K_{P}^{V}}{(1+|x|)^{2}}, \tag{20}
\end{equation*}
$$

with $K_{P}^{V}$ a positive constant. Using (17) and the fact that

$$
\ln |x-y| \leqslant \ln (1+|x|)+\ln (1+|y|),
$$

we see that for all $x \in \mathbb{R}$,

$$
\begin{equation*}
\rho_{P}(x) \leqslant C_{P}^{V} \exp (-V(x)+2 P \ln (1+|x|)) \tag{21}
\end{equation*}
$$

with

$$
C_{P}^{V}=\exp \left(2 P \int_{\mathbb{R}} \ln (1+|y|) \rho_{P}(y) d y+\lambda_{P}^{V}\right)
$$

which is indeed finite by 20 .
For the second point, we use that $\left(U^{\rho_{P}}\right)^{\prime}=\mathcal{H}\left[\rho_{P}\right]$ weakly and equation (17) to conclude on the distributional identity

$$
\rho_{P}^{\prime}=\left(-V^{\prime}-2 P \mathcal{H}\left[\rho_{P}\right]\right) \rho_{P}
$$

By the second point of Assumption 1.1. $V^{\prime}(x) e^{-V(x)+2 P \ln (1+|x|)}=o\left(x^{-1}\right)$ as $|x| \rightarrow \infty$, thus by 21, $V^{\prime} \rho_{P} \in L^{2}(\mathbb{R})$. Also since $\rho_{P}$ is $L^{2}(\mathbb{R})$ and bounded, we deduce, by using that $\mathcal{H}\left[L^{2}(\mathbb{R})\right]=L^{2}(\mathbb{R})$, that $\mathcal{H}\left[\rho_{P}\right] \rho_{P} \in L^{2}(\mathbb{R})$. Adding up these terms we get $\rho_{P} \in H^{1}(\mathbb{R})$. Because $\mathcal{H}\left[\rho_{P}\right]^{\prime}=\mathcal{H}\left[\rho_{P}^{\prime}\right]$ in a weak sense by Lemma 2.1 $\mathcal{H}\left[\rho_{P}\right] \in H^{1}(\mathbb{R})$. By the classical fact that $H^{1}(\mathbb{R})$ is contained in the set of $1 / 2$-Hölder functions $\mathcal{C}^{1 / 2}(\mathbb{R})$, we have $\mathcal{H}\left[\rho_{P}\right] \in \mathcal{C}^{1 / 2}(\mathbb{R})$ and so $U^{\rho_{P}} \in \mathcal{C}^{1,1 / 2}(\mathbb{R})$, the set of functions in $\mathcal{C}^{1}(\mathbb{R})$ with derivative of class $1 / 2$-Hölder.
Using the fact that $V$ is continuously differentiable, the previous equation for the weak derivative of $\rho_{P}$ then ensures that $\rho_{P} \in \mathcal{C}^{1}(\mathbb{R})$ and equation 18 holds in the strong sense.

Differentiating (in a weak sense) equation (18) we obtain

$$
\rho_{P}^{\prime \prime}=\left(-2 P \mathcal{H}\left[\rho_{P}\right]^{\prime}-V^{\prime \prime}+V^{\prime 2}+4 P^{2} \mathcal{H}\left[\rho_{P}\right]^{2}+4 P V^{\prime} \mathcal{H}\left[\rho_{P}\right]\right) \rho_{P}
$$

The three first terms belong to $L^{2}(\mathbb{R})$ for the same reasons as before. By equation $21, \rho_{P} \in L^{4}(\mathbb{R})$ and by lemma 1.1, so is $\mathcal{H}\left[\rho_{P}\right]$, thus using the boundedness of $\rho_{P}$ we see that $\rho_{P} \mathcal{H}\left[\rho_{P}\right]^{2} \in L^{2}(\mathbb{R})$. For the last term, we use that $V^{\prime} \rho_{P}$ and $\mathcal{H}\left[\rho_{P}\right]$ belong to $L^{4}(\mathbb{R})$ to ensure by Cauchy-Schwarz inequality that it is in $L^{2}(\mathbb{R})$. Finally, we can conclude that $\rho_{P} \in H^{2}(\mathbb{R})$ and so that $\mathcal{H}\left[\rho_{P}\right] \in H^{2}(\mathbb{R})$ with $\mathcal{H}\left[\rho_{P}\right]^{\prime \prime}=\mathcal{H}\left[\rho_{P}^{\prime \prime}\right]$ (in a weak sense). As before, we conclude that $\rho_{P} \in \mathcal{C}^{2}(\mathbb{R})$ and that equation 19 holds in a strong sense.

We next show that the Hilbert transform of $\rho_{P}$ is continuous and decays at infinity.
Lemma 2.3. Let $u \in L^{2}(\mathbb{R})$ such that $\int_{\mathbb{R}} u(t) d t$ exists and $f: t \mapsto t u(t) \in H^{1}(\mathbb{R})$ then

$$
\mathcal{H}[u](x) \underset{|x| \rightarrow \infty}{\sim} \frac{-\int_{\mathbb{R}} u(t) d t}{x} .
$$

Moreover if $\int_{\mathbb{R}} u(t) d t=0, \int_{\mathbb{R}} f(t) d t$ exists and $g: t \mapsto t^{2} u(t) \in H^{1}(\mathbb{R})$, then

$$
\mathcal{H}[u](x) \underset{|x| \rightarrow \infty}{\sim} \frac{-\int_{\mathbb{R}} t u(t) d t}{x^{2}} .
$$

As a consequence, we obtain that $\mathcal{H}\left[\rho_{P}\right](x) \underset{|x| \rightarrow \infty}{\sim} x^{-1}$ and the logarithmic potential $U^{\rho_{P}}$ is Lipschitz bounded, with bounded derivative $\mathcal{H}\left[\rho_{P}\right]$.

Proof. Let $u \in L^{2}(\mathbb{R})$, such that $\int_{\mathbb{R}} u(t) d t$ exists and $f: t \mapsto t u(t) \in H^{1}(\mathbb{R})$. Then

$$
x \mathcal{H}[u](x)+\int_{\mathbb{R}} u(t) d t=\int_{\mathbb{R}}\left[\frac{x u(x+t)-x u(x-t)}{2 t}+\frac{u(x+t)}{2}+\frac{u(x-t)}{2}\right] d t=\mathcal{H}[f](x) .
$$

Since $f \in H^{1}(\mathbb{R})$, so is $\mathcal{H}[f]$, proving that it goes to zero at infinity. Hence

$$
\mathcal{H}[u](x) \underset{|x| \rightarrow \infty}{\sim} \frac{-\int_{\mathbb{R}} u(t) d t}{x}
$$

Moreover if $\int_{\mathbb{R}} u(t) d t=0, \int_{\mathbb{R}} f(t) d t$ exists and $g: t \mapsto t^{2} u(t) \in H^{1}(\mathbb{R})$, then by the same argument:

$$
x^{2} \mathcal{H}[u](x)=x \mathcal{H}[f](x)=\mathcal{H}[g](x)-\int_{\mathbb{R}} f(t) d t
$$

where $g(t)=t^{2} u(t)$. We deduce that $\mathcal{H}[u](x) \underset{|x| \rightarrow \infty}{\sim} \frac{-\int_{\mathbb{R}} t u(t) d t}{x^{2}}$ since $\mathcal{H}[g]$ goes to zero at infinity.
We conclude this section by stating the Poincaré inequality for the measure $\mu_{P}$ under the assumption that $V$ is convex.

Proposition 2.4. The measure $\mu_{P}$ satisfies the following Poincaré inequality: There exists a constant $C$ such that for all $f \in \mathcal{C}_{c}^{\infty}(\mathbb{R})$,

$$
\begin{equation*}
\operatorname{Var}_{\mu_{P}}(f) \leqslant C \int_{\mathbb{R}}\left|f^{\prime}\right|^{2} d \mu_{P} \tag{22}
\end{equation*}
$$

This fact comes as a direct consequence of [BBCG08] [Corollary 1.9], which states that if the probability measure $\mu$ has a log-concave density on $\mathbb{R}$, then it satisfies 22 for $f$ smooth enough (actually the result is true in $\mathbb{R}^{d}$, replacing $f^{\prime}$ by $\left.\nabla f\right)$. Indeed, by convexity of $V$ and concavity of $x \mapsto \int_{\mathbb{R}} \ln |x-y| \rho_{P}(y) d y$, equation 17 implies that $\ln \rho_{P}$ is concave.

Remark 2.5. We will apply later inequality $(22)$ to more general functions than $\mathcal{C}_{c}^{\infty}(\mathbb{R})$, namely functions of the weighted Sobolev space $H^{1}\left(\rho_{P}\right)$, defined in Section $\sigma$; which can be seen as the completion of $\mathcal{C}_{c}^{\infty}(\mathbb{R})$ with respect to the norm $\|u\|_{L^{2}\left(\rho_{P}\right)}+\left\|u^{\prime}\right\|_{L^{2}\left(\rho_{P}\right)}$.

## 3 Concentration inequality, proof of Theorem 1.5

We prove in this section the concentration Theorem 1.5. Its proof is a direct adaptation of Theorem 1.4 of [HL21], which shows the analogous estimate in the circular setup. It is inspired by MMS14 and based on a comparison between a configuration $\mathbf{x}_{N}=\left(x_{1}, \ldots, x_{N}\right)$ sampled with $\mathbb{P}_{N}^{V, P}$ and a regularized version $\mathbf{y}_{N}=\left(y_{1}, \ldots, y_{N}\right)$, which we describe here.

Definition 3.1. Let $\mathbf{x}_{N}=\left(x_{1}, \ldots, x_{N}\right) \in \mathbb{R}^{N}$, and suppose (up to reordering) that $x_{1} \leqslant x_{2} \ldots \leqslant x_{N}$. We define $\mathbf{y}_{N} \in \mathbb{R}^{N}$ by:
$y_{1}:=x_{1}$, and for $0 \leqslant k \leqslant N-1, y_{k+1}:=y_{k}+\max \left\{x_{k+1}-x_{k}, N^{-3}\right\}$.
Note that the configuration $\mathbf{y}_{N}$ given by the previous definition satisfies $y_{k+1}-y_{k} \geqslant N^{-3}$, and $\mathbf{y}_{N}$ is close to $\mathbf{x}_{N}$ in the sense that

$$
\begin{equation*}
\sum_{k=1}^{N}\left|x_{k}-y_{k}\right| \leqslant \frac{1}{2 N} \tag{23}
\end{equation*}
$$

Indeed, by construction we have $\left|x_{k}-y_{k}\right|=y_{k}-x_{k} \leqslant(k-1) N^{-3}$, and we get the bound by summing these inequalities.

The key point of the proof of Theorem 1.5 is comparing the empirical measure $\hat{\mu}_{N}=\frac{1}{N} \sum_{i=1}^{N} \delta_{x_{i}}$, where $\mathbf{x}_{N}$ follows $\mathbb{P}_{N}^{V, P}$, to the regularized measure

$$
\begin{equation*}
\widetilde{\mu}_{N}:=\lambda_{N^{-5}} * \frac{1}{N} \sum_{i=1}^{N} \delta_{y_{i}} \tag{24}
\end{equation*}
$$

$i e$ the convolution of $\lambda_{N^{-5}}$ and the empirical measure, where $\lambda_{N^{-5}}$ is the uniform measure on $\left[0, N^{-5}\right]$. The interest of introducing the measure $\widetilde{\mu}_{N}$ is that it is close to $\hat{\mu}_{N}$, while having a finite energy $\mathcal{E}_{P}^{V}\left(\widetilde{\mu}_{N}\right)$, given by 16 . Finally, notice that the empirical measure doesn't change when reordering $x_{1}, \ldots, x_{N}$, and thus we do not lose in generality for our purposes in assuming that $x_{1} \leqslant \ldots \leqslant x_{N}$ in definition 3.1

We now introduce a distance on $\mathcal{P}(\mathbb{R})$ which is well-suited to our context.
Definition 3.2. For $\mu, \mu^{\prime} \in \mathcal{P}(\mathbb{R})$ we define the distance (possibly infinite) $D\left(\mu, \mu^{\prime}\right)$ by

$$
\begin{align*}
D\left(\mu, \mu^{\prime}\right) & :=\left(-\int \ln |x-y| d\left(\mu-\mu^{\prime}\right)(x) d\left(\mu-\mu^{\prime}\right)(y)\right)^{1 / 2}  \tag{25}\\
& =\left(\int_{0}^{+\infty} \frac{1}{t}\left|\mathcal{F}\left[\mu-\mu^{\prime}\right]\right|^{2} d t\right)^{1 / 2} \tag{26}
\end{align*}
$$

where the Fourier transform of a signed measure $\nu$ is defined by $\mathcal{F}[\nu](x):=\int e^{-i t x} d\left(\mu-\mu^{\prime}\right)(x)$
Let $f: \mathbb{R} \rightarrow \mathbb{R}$ with finite $1 / 2$ norm $\|f\|_{1 / 2}:=\left(\int_{\mathbb{R}}|t||\mathcal{F}[f](t)|^{2} d t\right)^{1 / 2}$. By Plancherel theorem and Hölder inequality, for any $\mu, \mu^{\prime} \in \mathcal{P}(\mathbb{R})$, setting $\nu=\mu-\mu^{\prime}$,

$$
\left|\int_{\mathbb{R}} f d \mu-\int_{\mathbb{R}} f d \mu^{\prime}\right|^{2}=\left.\left.\left|\frac{1}{2 \pi} \int_{\mathbb{R}}\right| t\right|^{1 / 2} \mathcal{F}[f](t) \frac{\overline{\mathcal{F}[\nu](t)}}{|t|^{1 / 2}} d t\right|^{2} \leqslant \frac{1}{2 \pi^{2}}\|f\|_{1 / 2}^{2} D^{2}\left(\mu, \mu^{\prime}\right)
$$

Therefore the metric $d$ defined in $\sqrt{12}$ is dominated by $D$ :

$$
\begin{equation*}
d\left(\mu, \mu^{\prime}\right) \leqslant \frac{1}{\sqrt{2} \pi} D\left(\mu, \mu^{\prime}\right) \tag{27}
\end{equation*}
$$

The following lemma shows how the distance $D$ is related to the energy-functional $\mathcal{E}_{P}^{V}$ defined in 16, we will write $\mathcal{E}_{P}$ for simplicity.
Lemma 3.3. We have for any absolutely continuous $\mu \in \mathcal{P}(\mathbb{R})$ with finite energy $\mathcal{E}_{P}^{V}(\mu)$,

$$
\begin{equation*}
\mathcal{E}_{P}(\mu)-\mathcal{E}_{P}\left(\mu_{P}\right)=P D^{2}\left(\mu, \mu_{P}\right)+\int \ln \left(\frac{d \mu}{d \mu_{P}}\right) d \mu \tag{28}
\end{equation*}
$$

Proof of Lemma 3.3. Subtracting $\mathcal{E}_{P}(\mu)-\mathcal{E}_{P}\left(\mu_{P}\right)$ we find

$$
\begin{array}{r}
\mathcal{E}_{P}(\mu)-\mathcal{E}_{P}\left(\mu_{P}\right)=\int V d\left(\mu-\mu_{P}\right)+\int \ln \frac{d \mu}{d x} d \mu-\int \ln \rho_{P} d \mu_{P}-P \iint \ln |x-y| d \mu(x) d \mu(y) \\
+P \iint \ln |x-y| d \mu_{P}(x) d \mu_{P}(y) \tag{29}
\end{array}
$$

Now, if $\nu$ is a signed measure of mass zero, integrating 17 we get

$$
\int V(x) d \nu(x)-2 P \iint \ln |x-y| d \nu(x) d \mu_{P}(y)+\int \ln \left(\rho_{P}\right)(x) d \nu(x)=0
$$

We take $\nu=\mu-\mu_{P}$, and get

$$
\begin{aligned}
\int V(x) d\left(\mu-\mu_{P}\right)(x)=2 P \iint \ln |x-y| d \mu(x) d \mu_{P}(y)-2 P & \iint \ln |x-y| d \mu_{P}(x) d \mu_{P}(y) \\
& -\int \ln \left(\rho_{P}\right)(x) d \mu(x)+\int \ln \left(\rho_{P}\right)(x) d \mu_{P}(x) .
\end{aligned}
$$

Plugging this last identity in 29), we find

$$
\mathcal{E}_{P}(\mu)-\mathcal{E}_{P}\left(\mu_{P}\right)=-P \iint \ln |x-y| d \nu(x) d \nu(y)+\int \ln \left(\frac{d \mu}{d \mu_{P}}\right)(x) d \mu(x)
$$

which establishes the result.
Proof of Theorem 1.5. We first give a lower bound for the partition function $Z_{N}^{V, P}$ (4) of $\mathbb{P}_{N}^{V, P}$. We rewrite it as

$$
Z_{N}^{V, P}=\int_{\mathbb{R}^{N}} \exp \left(\frac{2 P}{N} \sum_{i<j} \ln \left|x_{i}-x_{j}\right|-\sum_{i=1}^{N}\left[V\left(x_{i}\right)+\ln \rho_{P}\left(x_{i}\right)\right]\right) d \rho_{P}\left(x_{1}\right) \ldots d \rho_{P}\left(x_{N}\right),
$$

and apply Jensen inequality to obtain:

$$
\begin{aligned}
\ln Z_{N}^{V, P} & \geqslant \int_{\mathbb{R}^{N}}\left(\frac{2 P}{N} \sum_{i<j} \ln \left|x_{i}-x_{j}\right|-\sum_{i=1}^{N}\left[V\left(x_{i}\right)+\ln \rho_{P}\left(x_{i}\right)\right]\right) d \rho_{P}\left(x_{1}\right) \ldots d \rho_{P}\left(x_{N}\right) \\
& \geqslant P(N-1) \iint \ln |x-y| d \rho_{P}(x) d \rho_{P}(y)-N \int_{\mathbb{R}}\left[V+\ln \rho_{P}\right] d \rho_{P} \\
& \geqslant-N \mathcal{E}_{P}^{V}\left[\mu_{P}\right]-P \iint \ln |x-y| d \rho_{P}(x) d \rho_{P}(y) .
\end{aligned}
$$

Using this estimate and the fact that for $1 \leqslant i, j \leqslant N$ we have $\left|x_{i}-x_{j}\right| \leqslant\left|y_{i}-y_{j}\right|$, with $\mathbf{y}_{N}=\left(y_{1}, \ldots, y_{N}\right)$ of definition 3.1 we deduce the bound on the density of probability

$$
\begin{equation*}
\frac{d \mathbb{P}_{N}^{V, P}}{d \mathbf{x}}\left(x_{1}, \ldots, x_{N}\right) \leqslant e^{N \mathcal{E}_{P}\left(\mu_{P}\right)+P} \iint \ln |x-y| d \mu_{P}(x) d \mu_{P}(y)+\frac{P}{N} \sum_{i \neq j} \ln \left|y_{i}-y_{j}\right|-\sum_{i=1}^{N} V\left(x_{i}\right) . \tag{30}
\end{equation*}
$$

Recalling (24), we now show the following estimate

$$
\begin{equation*}
\sum_{i \neq j} \ln \left|y_{i}-y_{j}\right| \leqslant 2 N+N^{2} \iint \ln |x-y| d \widetilde{\mu}_{N}(x) d \widetilde{\mu}_{N}(y)+3 N \ln N+\frac{3}{2} N . \tag{31}
\end{equation*}
$$

Let $i \neq j$ and $u, v \in\left[0, N^{-5}\right]$. Since for $x \neq 0$ and $|h| \leqslant \frac{|x|}{2}$, we have $|\ln | x+h|-\ln | x| | \leqslant \frac{2|h|}{|x|}$, we deduce

$$
|\ln | y_{i}-y_{j}+u-v|-\ln | y_{i}-y_{j}| | \leqslant \frac{2|u-v|}{\left|y_{i}-y_{j}\right|} \leqslant \frac{2 N^{-5}}{N^{-3}}=\frac{2}{N^{2}} .
$$

Thus, summing over $i \neq j$ and integrating with respect to $u$ and $v$, we get

$$
\begin{aligned}
\sum_{i \neq j} \ln \left|y_{i}-y_{j}\right| & \leqslant 2+\sum_{i \neq j} \iint \ln \left|y_{i}-y_{j}+u-v\right| d \lambda_{N^{-5}}(u) d \lambda_{N^{-5}}(v) \\
& =2+N^{2} \iint \ln |x-y| d \widetilde{\mu}_{N}(x) d \widetilde{\mu}_{N}(y)-N \iint \ln |u-v| d \lambda_{N^{-5}}(u) d \lambda_{N^{-5}}(v)
\end{aligned}
$$

The last integral is equal to $-\frac{3}{2}-5 \ln N$, so we deduce (31). We now combine (30) and (31). Recall (16) and set

$$
c_{N}=P\left(\iint \ln |x-y| d \mu_{P}(x) d \mu_{P}(y)+3 / 2+2 / N\right)
$$

We get

$$
\begin{aligned}
\frac{d \mathbb{P}_{N}^{V, P}}{d \mathbf{x}}\left(x_{1}, \ldots, x_{N}\right) & \leqslant e^{c_{N}+5 P \ln N} \exp \left[N\left\{\mathcal{E}_{P}\left(\mu_{P}\right)-\mathcal{E}_{P}\left(\widetilde{\mu}_{N}\right)+\int\left(V+\ln \frac{d \widetilde{\mu}_{N}}{d x}\right) d \widetilde{\mu}_{N}\right\}-\sum_{i=1}^{N} V\left(x_{i}\right)\right] \\
& =e^{c_{N}+5 P \ln N} \exp \left[-N P D^{2}\left(\widetilde{\mu}_{N}, \mu_{P}\right)+N \int\left(V+\ln \rho_{P}\right) d \widetilde{\mu}_{N}-\sum_{i=1}^{N} V\left(x_{i}\right)\right]
\end{aligned}
$$

where we used equation (28) in the last equality. Using again equation (17) we then see that the density $\frac{d \mathbb{P}_{N}^{V, P}}{d \mathbf{x}}\left(x_{1}, \ldots, x_{N}\right)$ is bounded by

$$
e^{c_{N}+5 P \ln N} \exp \left[-N P D^{2}\left(\widetilde{\mu}_{N}, \mu_{P}\right)+2 P N \iint \ln |x-y| d\left(\widetilde{\mu}_{N}-\hat{\mu}_{N}\right)(x) d \mu_{P}(y)\right] \prod_{i=1}^{N} \rho_{P}\left(x_{i}\right)
$$

Recalling (15), $\iint \ln |x-y| d\left(\widetilde{\mu}_{N}-\hat{\mu}_{N}\right)(x) d \mu_{P}(y)=-\int U^{\rho_{P}} d\left(\widetilde{\mu}_{N}-\hat{\mu}_{N}\right)$. As a consequence of the bound on the density $\frac{d \mathbb{P}_{N}^{V, P}}{d \mathbf{x}}\left(x_{1}, \ldots, x_{N}\right)$ we established, we have for all $r>0$

$$
\begin{equation*}
\mathbb{P}_{N}^{V, P}\left(D^{2}\left(\widetilde{\mu}_{N}, \mu_{P}\right)>r\right) \leqslant e^{-N P r+c_{N}+5 P \ln N} \int_{\mathbb{R}^{N}} \exp \left\{-2 P N \int U^{\rho_{P}} d\left(\widetilde{\mu}_{N}-\hat{\mu}_{N}\right)\right\} \prod_{i=1}^{N} \rho_{P}\left(x_{i}\right) d x_{i} \tag{32}
\end{equation*}
$$

Next, we show that $-N \int U^{\rho_{P}} d\left(\widetilde{\mu}_{N}-\hat{\mu}_{N}\right)$ is bounded.
By Lemma 2.3. $U^{\rho_{P}}$ is differentiable with bounded derivative $\mathcal{H}\left[\rho_{P}\right]$ on $\mathbb{R}$. As a consequence,

$$
\begin{aligned}
\left|N \int U^{\rho_{P}} d\left(\widetilde{\mu}_{N}-\hat{\mu}_{N}\right)\right| & \leqslant \sum_{i=1}^{N} \int\left|U^{\rho_{P}}\left(y_{i}+u\right)-U^{\rho_{P}}\left(x_{i}\right)\right| d \lambda_{N^{-5}}(u) \\
& \leqslant\left\|\mathcal{H}\left[\rho_{P}\right]\right\|_{\infty}\left(\sum_{i=1}^{N}\left|y_{i}-x_{i}\right|+N \int u d \lambda_{N^{-5}}(u)\right) \\
& \leqslant\left\|\mathcal{H}\left[\rho_{P}\right]\right\|_{\infty}\left(\frac{1}{2 N}+N^{-4} / 2\right)
\end{aligned}
$$

where we used (23) in the last inequality. Therefore, we deduce from 32)

$$
\begin{equation*}
\mathbb{P}_{N}^{V, P}\left(D^{2}\left(\widetilde{\mu}_{N}, \mu_{P}\right)>r\right) \leqslant e^{-N P r+c_{N}+5 P \ln N+2 P\left\|\mathcal{H}\left[\rho_{P}\right]\right\|_{\infty}}=e^{-N P r+5 P \ln N+K_{N}} \tag{33}
\end{equation*}
$$

with $K_{N}:=c_{N}+2 P\left\|\mathcal{H}\left[\rho_{P}\right]\right\|_{\infty}$. Since $c_{N}$ is bounded, so is $K_{N}$.
Finally, let $f$ be a Lipschitz bounded function with $\|f\|_{\text {Lip }} \leqslant 1$, then, we have (as we did for $U^{\rho_{P}}$ )

$$
\left|\int f d \hat{\mu}_{N}-\int f d \widetilde{\mu}_{N}\right| \leqslant N^{-2}
$$

Thus,

$$
d\left(\hat{\mu}_{N}, \mu_{P}\right) \leqslant d\left(\hat{\mu}_{N}, \widetilde{\mu}_{N}\right)+d\left(\widetilde{\mu}_{N}, \mu_{P}\right) \leqslant N^{-2}+\frac{1}{\sqrt{2} \pi} D\left(\widetilde{\mu}_{N}, \mu_{P}\right)
$$

and for any $N$ such that $r-N^{-2} \geqslant r / 2$ (in particular $r-N^{-2}>0$ ) we get

$$
\mathbb{P}_{N}^{V, P}\left(d\left(\hat{\mu}_{N}, \mu_{P}\right)>r\right) \leqslant \mathbb{P}_{N}^{V, P}\left(\frac{1}{2 \pi^{2}} D^{2}\left(\widetilde{\mu}_{N}, \mu_{P}\right)>\left(r-N^{-2}\right)^{2}\right) \leqslant \mathbb{P}_{N}^{V, P}\left(\frac{1}{2 \pi^{2}} D^{2}\left(\widetilde{\mu}_{N}, \mu_{P}\right)>r^{2} / 4\right)
$$

and the last term is bounded by $e^{-N r^{2} \frac{P \pi^{2}}{2}+5 P \ln N+K}$ for some $K$ large enough, which concludes the proof.

As a consequence of Theorem 1.5 we are able to control the quantities

$$
\begin{equation*}
\zeta_{N}(\phi):=\iint_{\mathbb{R}^{2}} \frac{\phi(x)-\phi(y)}{x-y} d\left(\hat{\mu}_{N}-\mu_{P}\right)(x) d\left(\hat{\mu}_{N}-\mu_{P}\right)(y) \tag{34}
\end{equation*}
$$

for a certain class of test functions $\phi$.
Corollary 3.4. There exists $C, K>0$ such that for all $\phi \in \mathcal{C}^{2}(\mathbb{R}) \cap H^{2}(\mathbb{R})$ with bounded second derivative, we have for $\varepsilon>0$ and $N$ large enough,

$$
\mathbb{P}_{N}^{V, P}\left(\sqrt{N}\left|\zeta_{N}(\phi)\right| \leqslant N^{-1 / 2+\varepsilon}\right) \geqslant 1-\exp \left\{-\frac{P N^{\varepsilon}}{2 C N_{2}(\phi)}+5 P \ln N+K\right\}
$$

with $N_{2}(\phi)=\left\|\phi^{\prime}\right\|_{L^{2}(d x)}+\left\|\phi^{\prime \prime}\right\|_{L^{2}(d x)}$.
Proof. We follow the proof given in [Gui19][Cor. 4.16] and adapt it to our setting. Let us denote by $\widetilde{\zeta_{N}}(\phi)$ the quantity

$$
\iint_{\mathbb{R}^{2}} \frac{\phi(x)-\phi(y)}{x-y} d\left(\widetilde{\mu}_{N}-\mu_{P}\right)(x) d\left(\widetilde{\mu}_{N}-\mu_{P}\right)(y) .
$$

We have the almost sure inequality, by a Taylor estimate

$$
\begin{equation*}
\left|\zeta_{N}(\phi)-\widetilde{\zeta_{N}}(\phi)\right| \leqslant 2 N^{-2}\left\|\phi^{\prime \prime}\right\|_{\infty} \tag{35}
\end{equation*}
$$

Thus, for any $\delta>0$,

$$
\begin{aligned}
\mathbb{P}_{N}^{V, P}\left(\left|\zeta_{N}(\phi)\right|>\delta\right) & \leqslant \mathbb{P}_{N}^{V, P}\left(\left|\zeta_{N}(\phi)-\widetilde{\zeta_{N}}(\phi)\right|>\delta / 2\right)+\mathbb{P}_{N}^{V, P}\left(\left|\widetilde{\zeta_{N}}(\phi)\right|>\delta / 2\right) \\
& \leqslant \mathbb{P}_{N}^{V, P}\left(2 N^{-2}\left\|\phi^{\prime \prime}\right\|_{\infty}>\delta / 2\right)+\mathbb{P}_{N}^{V, P}\left(\left|\widetilde{\zeta_{N}}(\phi)\right|>\delta / 2\right)
\end{aligned}
$$

where the first term of the right-hand side is either 0 or 1 . With $\delta=N^{-1+\varepsilon}, \varepsilon>0$, it is zero for $N$ large enough. For such a choice of $\delta$, and for $N$ large enough,

$$
\mathbb{P}_{N}^{V, P}\left(\left|\zeta_{N}(\phi)\right|>N^{-1+\varepsilon}\right) \leqslant \mathbb{P}_{N}^{V, P}\left(\left|\widetilde{\zeta_{N}}(\phi)\right|>\frac{1}{2} N^{-1+\varepsilon}\right)
$$

We next show that, for some $C>0$ independent of $\phi$, we have

$$
\begin{equation*}
\left|\widetilde{\zeta_{N}}(\phi)\right| \leqslant C D^{2}\left(\widetilde{\mu}_{N}, \mu_{P}\right) N_{2}(\phi) . \tag{36}
\end{equation*}
$$

We begin by showing this inequality for $\psi \in \mathcal{S}(\mathbb{R})$. By using the inverse Fourier transform we have

$$
\begin{aligned}
& \widetilde{\zeta}_{N}(\psi)=\frac{1}{2 \pi} \iint \frac{\int d t \mathcal{F}[\psi](t) e^{-i t x}-\int d t \mathcal{F}[\psi](t) e^{-i t y}}{x-y} d\left(\widetilde{\mu}_{N}-\mu_{P}\right)(x) d\left(\widetilde{\mu}_{N}-\mu_{P}\right)(y) \\
& \quad=\frac{-1}{2 \pi} \int d t \mathrm{i} t \mathcal{F}[\psi](t) \iint e^{-i t y} \frac{e^{-i t(x-y)}-1}{-i t(x-y)} d\left(\widetilde{\mu}_{N}-\mu_{P}\right)(x) d\left(\widetilde{\mu}_{N}-\mu_{P}\right)(y) \\
& \quad=\frac{-1}{2 \pi} \int d t \mathrm{i} t \mathcal{F}[\psi](t) \iint e^{-\mathrm{i} t y} \int_{0}^{1} d \alpha e^{-\mathrm{i} \alpha t(x-y)} d\left(\widetilde{\mu}_{N}-\mu_{P}\right)(x) d\left(\widetilde{\mu}_{N}-\mu_{P}\right)(y) \\
& \quad=\frac{-1}{2 \pi} \int d t \mathrm{i} t \mathcal{F}[\psi](t) \int_{0}^{1} d \alpha \int e^{-\mathrm{i} \alpha t x} d\left(\widetilde{\mu}_{N}-\mu_{P}\right)(x) \int e^{-\mathrm{i}(1-\alpha) t y} d\left(\widetilde{\mu}_{N}-\mu_{P}\right)(y)
\end{aligned}
$$

We then apply in order the triangular inequality, Cauchy-Schwarz inequality, a change of variable and the fact that $\left|\mathcal{F}\left[\widetilde{\mu}_{N}-\mu_{P}\right]\right|^{2}$ is an even function.

$$
\begin{aligned}
\left|\widetilde{\zeta}_{N}(\psi)\right| & \leqslant \frac{1}{2 \pi} \int_{\mathbb{R}} d t|t \mathcal{F}[\psi](t)| \int_{0}^{1} d \alpha\left|\mathcal{F}\left[\widetilde{\mu}_{N}-\mu_{P}\right](\alpha t)\right| \cdot\left|\mathcal{F}\left[\widetilde{\mu}_{N}-\mu_{P}\right]((1-\alpha) t)\right| \\
& \leqslant \frac{1}{2 \pi} \int_{\mathbb{R}} d t|t \mathcal{F}[\psi](t)|\left(\int_{0}^{1} d \alpha\left|\mathcal{F}\left[\widetilde{\mu}_{N}-\mu_{P}\right](\alpha t)\right|^{2}\right)^{\frac{1}{2}}\left(\int_{0}^{1} d \alpha\left|\mathcal{F}\left[\widetilde{\mu}_{N}-\mu_{P}\right]((1-\alpha) t)\right|^{2}\right)^{\frac{1}{2}} \\
& \leqslant \frac{1}{2 \pi} \int_{\mathbb{R}} d t|t \mathcal{F}[\psi](t)| \int_{0}^{1} d \alpha\left|\mathcal{F}\left[\widetilde{\mu}_{N}-\mu_{P}\right](\alpha t)\right|^{2} \\
& \leqslant \frac{1}{2 \pi} \int_{0}^{+\infty} d t|t \mathcal{F}[\psi](t)| \int_{0}^{1} \frac{t d \alpha}{t \alpha}\left|\mathcal{F}\left[\widetilde{\mu}_{N}-\mu_{P}\right](\alpha t)\right|^{2}+\frac{1}{2 \pi} \int_{-\infty}^{0} d t|t \mathcal{F}[\phi](t)| \int_{0}^{1} \frac{-t d \alpha}{-t \alpha}\left|\mathcal{F}\left[\widetilde{\mu}_{N}-\mu_{P}\right](\alpha t)\right|^{2} \\
& \leqslant \frac{1}{2 \pi} \int_{\mathbb{R}} d t|t \mathcal{F}[\psi](t)| D^{2}\left(\widetilde{\mu}_{N}, \mu_{P}\right) \\
& \leqslant \frac{1}{2 \pi}\left(\int_{\mathbb{R}} d t|t \mathcal{F}[\psi](t)|^{2}\left(1+t^{2}\right)\right)^{\frac{1}{2}}\left(\int_{\mathbb{R}} \frac{d t}{1+t^{2}}\right)^{\frac{1}{2}} D^{2}\left(\widetilde{\mu}_{N}, \mu_{P}\right) \\
& \leqslant \frac{1}{2 \sqrt{\pi}} D^{2}\left(\widetilde{\mu}_{N}, \mu_{P}\right) N_{2}(\psi)
\end{aligned}
$$

By density of $\mathcal{S}(\mathbb{R})$ in $L^{2}(\mathbb{R})$, and since $\widetilde{\zeta}_{N}:\left(H^{2}(\mathbb{R}), N_{2}\right) \rightarrow \mathbb{R}$ is continuous, the inequality still holds for $\phi$. Thus, using equation (33),

$$
\mathbb{P}_{N}^{V_{N} P}\left(\left|\widetilde{\zeta_{N}}(\phi)\right|>\frac{1}{2} N^{-1+\varepsilon}\right) \leqslant \mathbb{P}_{N}^{V, P}\left(D^{2}\left(\widetilde{\mu}_{N}, \mu_{P}\right)>\frac{N^{-1+\varepsilon}}{2 C N_{2}(\phi)}\right) \leqslant \exp \left\{-P \frac{N^{\varepsilon}}{2 C N_{2}(\phi)}+5 P \ln N+K\right\},
$$

which concludes the proof.

## 4 Localization of the edge of a configuration

In Lam21][Theorem 1.8, Theorem 3.4], Lambert was able to control the edge (i.e the minimum and the maximum) of a typical configuration $\left(x_{1}, \ldots, x_{N}\right)$ distributed according to $\mathbb{P}_{N}^{V_{N}, P}$, by showing that the random measure

$$
\Xi_{N}:=\sum_{j=1}^{N} \delta_{\varphi_{N}^{-1}\left(x_{j}\right)}
$$

converges in distribution towards a Poisson point process for a function $\varphi_{N}$ which takes the form

$$
\varphi_{N}(x):=E_{N}+\alpha_{N}^{-1} x .
$$

Before being more precise on the construction of $\left(E_{N}\right)_{N}$ and $\left(\alpha_{N}\right)_{N}$, we explain, following Lam21, how one can use this convergence to localize the edge of a typical configuration $\left(x_{1}, \ldots, x_{N}\right)$. Let us assume for a moment that $\Xi_{N}$ converges towards a Poisson point process with intensity $\theta(x)=e^{-x}$, with $E_{N} \rightarrow+\infty$. In particular, the random variable

$$
\Xi_{N}(t,+\infty)
$$

converges in distribution towards a Poisson random variable with mean $\int_{t}^{+\infty} e^{-x} d x$. Combined with the equalities

$$
\begin{aligned}
\mathbb{P}_{N}^{V, P}\left(\Xi_{N}(t,+\infty)=0\right) & =\mathbb{P}_{N}^{V, P}\left(\forall 1 \leqslant j \leqslant N, \varphi_{N}^{-1}\left(x_{j}\right)=\alpha_{N}\left(x_{j}-E_{N}\right) \leqslant t\right) \\
& =\mathbb{P}_{N}^{V, P}\left(\alpha_{N}\left(\max _{1 \leqslant j \leqslant N} x_{j}-E_{N}\right) \leqslant t\right),
\end{aligned}
$$

we deduce that for all $t \in \mathbb{R}$

$$
\mathbb{P}_{N}^{V, P}\left(\alpha_{N}\left(\max _{1 \leqslant j \leqslant N} x_{j}-E_{N}\right) \leqslant t\right) \underset{N \rightarrow \infty}{\longrightarrow} \exp \left(-e^{-t}\right)
$$

Therefore, the random variable

$$
\alpha_{N}\left(\max _{1 \leqslant j \leqslant N} x_{j}-E_{N}\right)
$$

converges in distribution to the Gumbel law, showing that the maximum of a configuration is of order $E_{N}$. Furthermore, as will be clear from the construction of $\alpha_{N}$ and $E_{N}, \alpha_{N}$ is positive, and goes to infinity as $N$ goes to infinity.

Replacing in the previous analysis $\theta(x)=e^{x}$ and $E_{N} \rightarrow-\infty$, we would have deduced in the same fashion that

$$
\alpha_{N}\left(\min _{1 \leqslant j \leqslant N} x_{j}-E_{N}\right)
$$

converges in law.
With the above notations, we can apply [Lam21] [Theorem 3.4] to our context.
Theorem 4.1. Let $v= \pm$. There exists $\left(E_{N}^{v}\right)_{N},\left(\alpha_{N}^{v}\right)_{N}$ sequences of real numbers with $\left|E_{N}^{v}\right| \rightarrow+\infty$, $\alpha_{N}^{v}>0$ for large enough $N$, satisfying $V^{\prime}\left(E_{N}^{v}\right)=\alpha_{N}^{v} v$, such that:
a) $\frac{N e^{-V\left(E_{N}^{v}\right)+2 P \ln \left|E_{N}^{v}\right|+\lambda_{V}^{P}}}{\alpha_{N}^{v}} \underset{N \rightarrow \infty}{\longrightarrow} 1$ (recall $\lambda_{P}^{V}$ is defined through equation (5)),
b) $\frac{\ln \left(\alpha_{N}^{v}\right)}{N} \underset{N \rightarrow \infty}{\longrightarrow} 0$ and $\alpha_{N}^{v}\left|E_{N}^{v}\right| \underset{N \rightarrow \infty}{\longrightarrow}+\infty$,
c) For all compact $K \subset \mathbb{R}$,

$$
\left(\alpha_{N}^{v}\right)^{-2} \sup _{x \in K}\left|V^{\prime \prime}\left(\varphi_{N}(x)\right)\right| \underset{N \rightarrow \infty}{\longrightarrow} 0
$$

As a consequence, the random measure $\Xi_{N}$ converges in distribution as $N \rightarrow \infty$ to a Poisson point process with intensity $\theta(x)=e^{-v x}$.

Proof. We prove it in the case $v=+$, the case where $v=-$ being similar. We show that there exists a sequence $\left(E_{N}^{+}\right)_{N}$ going to $+\infty$ satisfying $f\left(E_{N}^{+}\right)=-\ln N$, where we defined the function $f$ by

$$
f(x)=-V(x)+2 P \ln |x|+\lambda_{P}^{V}-\ln \left|V^{\prime}(x)\right|
$$

(we will then have $\alpha_{N}^{+}=V^{\prime}\left(E_{N}^{+}\right)>0$. In the case $v=-1$ we would have looked for a sequence $E_{N}^{-}$going to $-\infty$ and $\left.\alpha_{N}^{-}=-V^{\prime}\left(E_{N}^{-}\right)\right)$
As a consequence of Assumptions 1.1, ii), one shows that $\ln \left|V^{\prime}\right|$ is negligible with respect to $V$ at infinity. Therefore, because $\frac{\ln |x|}{V(x)} \underset{|x| \rightarrow \infty}{\longrightarrow} 0$,

$$
f(x)=-V(x)+o(V(x))
$$

at infinity (in particular, $f(x) \underset{x \rightarrow+\infty}{\longrightarrow}-\infty$ ). We deduce that for $0<\varepsilon<1$ fixed, there exists $A>0$ such that for all $x>A$,

$$
\begin{equation*}
-(1+\varepsilon) V(x)<f(x)<-(1-\varepsilon) V(x) \tag{37}
\end{equation*}
$$

and because $f(x) \underset{x \rightarrow+\infty}{\longrightarrow}-\infty$ there exists $\left(E_{N}^{+}\right)_{N}$ going to infinity such that for all $N \geqslant 1, f\left(E_{N}^{+}\right)=-\ln N$. Setting $x=E_{N}^{+}$in (37), we obtain that $-V\left(E_{N}^{+}\right) \sim f\left(E_{N}^{+}\right)=-\ln N$. By convexity of $V$ and the fact that it goes to infinity at infinity, $V$ is increasing on some $[M,+\infty[$, where $M \geqslant 0$. Thus $-(1 \pm \varepsilon) V(x)=-\ln N$
iff $x=V^{-1}\left(\frac{\ln N}{1 \pm \varepsilon}\right)$, where $V^{-1}$ denotes $\left(V_{\mid[M,+\infty[ }\right)^{-1}$. We conclude by 37] that such an $E_{N}^{+}$must satisfy

$$
\begin{equation*}
V^{-1}\left(\frac{\ln N}{1+\varepsilon}\right) \leqslant E_{N}^{+} \leqslant V^{-1}\left(\frac{\ln N}{1-\varepsilon}\right) \tag{38}
\end{equation*}
$$

By convexity of $V$ and the fact that it goes to infinity at infinity, $\left(\alpha_{N}^{+}\right)_{N}$ is non-decreasing and goes to infinity. It is thus positive for $N$ large enough, ensuring that $\alpha_{N}\left|E_{N}^{+}\right| \underset{N \rightarrow \infty}{\longrightarrow}+\infty$. Property c) of the theorem follows from Assumptions 1.1. point ii). It remains to show that $\frac{\ln \left(\alpha_{N}^{+}\right)}{N}=\frac{\ln \left|V^{\prime}\left(E_{N}^{+}\right)\right|}{N} \underset{N \rightarrow \infty}{\longrightarrow} 0$. By construction, we have

$$
\frac{\ln \left|V^{\prime}\left(E_{N}^{+}\right)\right|}{N}=\frac{\ln \left(N e^{-V\left(E_{N}^{+}\right)+2 P \ln N+\lambda_{P}}\right)}{N}=-\frac{V\left(E_{N}^{+}\right)}{N}+o(1)
$$

Using that $V\left(E_{N}^{+}\right) \sim \ln N$, we can conclude that $\ln \left|V^{\prime}\left(E_{N}^{+}\right)\right|=o(N)$ which concludes the proof.
By the discussion preceding Theorem 4.1. we deduce
Corollary 4.2 (Edge of a configuration). Let $E_{N}^{ \pm}, \alpha_{N}^{ \pm}:=\left|V^{\prime}\left(E_{N}^{ \pm}\right)\right|$be the sequences of Theorem 4.1 associated with $v= \pm 1$. Then, both random variables

$$
\alpha_{N}^{+}\left(\max _{1 \leqslant j \leqslant N} x_{j}-E_{N}^{+}\right)
$$

and

$$
\alpha_{N}^{-}\left(\min _{1 \leqslant j \leqslant N} x_{j}-E_{N}^{-}\right)
$$

converge to a Gumbel law, whose distribution function is given for $t \geqslant 0$ by $\mathcal{G}([0, t])=\exp \left(e^{-t}\right)$. Furthermore, $V\left(E_{N}^{ \pm}\right) \sim \ln N$ and $\alpha_{N}^{ \pm} \underset{N \rightarrow \infty}{\longrightarrow} \pm \infty$.

Remark 4.3. Note that Lam21][Theorem 3.4] applies for $V$ of class $\mathcal{C}^{2}$ outside of a compact set, allowing to take $V(x)=|x|^{a}$ for $a \geqslant 1$. Furthermore, if $V(x)=|x|^{a}+R(x)$ for $a \geqslant 1, R \in \mathcal{C}^{2}(\mathbb{R})$ and convex and where $R(x), R^{\prime}(x)$ and $R^{\prime \prime}(x)$ are negligible respectively with respect to $x^{a}$, $x^{a-1}$ and $x^{a-2}$, we find $E_{N}^{ \pm} \sim \pm(\ln N)^{1 / a}$.

If $V(x)=\cosh (x)$, we find $E_{N}^{+} \sim-E_{N}^{-} \sim \arg \cosh (\ln N) \sim \ln \ln N$.
The next lemma will be convenient in the proof of Theorem 5.2 when dealing with error terms.
Lemma 4.4. With the notations of Corollary 4.2, we have

$$
\mu_{P}\left(\left[E_{N}^{-}, E_{N}^{+}\right]^{c}\right)=o\left(N^{-1 / 2}\right)
$$

Proof. Let $0<\delta<1$, to be specified later. We have

$$
\int_{E_{N}^{+}}^{+\infty} \rho_{P} d x=\int_{E_{N}^{+}}^{+\infty}\left(\rho_{P}\right)^{\delta}\left(\rho_{P}\right)^{1-\delta} d x \leqslant \int_{\mathbb{R}}\left(\rho_{P}\right)^{\delta} d x \sup _{\left[E_{N}^{+},+\infty[ \right.}\left(\rho_{P}\right)^{1-\delta}
$$

By the first inequality of Lemma 2.2 the integral is finite. Also from the same inequality, we have for some constant $C^{\prime}$ and $x$ big enough $\rho_{P}(x) \leqslant C^{\prime} e^{-\frac{3}{4} V(x)}$. Because $V$ is increasing in a neighborhood of $+\infty$, we get for $N$ large enough

$$
\sup _{\left[E_{N}^{+},+\infty[ \right.}\left(\rho_{P}\right)^{1-\delta} \leqslant C^{\prime 1-\delta} e^{-(1-\delta) \frac{3}{4} V\left(E_{N}^{+}\right)} .
$$

Taking $\delta>0$ such that $\frac{1}{2}-(1-\delta) \frac{3}{4}=:-\gamma<0$ and using that $V\left(E_{N}^{+}\right)=\ln N+o(\ln N)$ (established in the proof of Theorem 4.1,

$$
\sqrt{N} \int_{E_{N}^{+}}^{+\infty} \rho_{P} d x \leqslant K e^{-\gamma \ln N+(1-\delta) \frac{3}{4} o(\ln N)}
$$

and the right-hand side goes to zero as $N$ goes to infinity. We deal with the integral $\int_{-\infty}^{E_{N}^{-}} \rho_{P} d x$ in the same way.

Remark 4.5. We could improve the proof to show that $\mu_{P}\left(\left[E_{N}^{-}, E_{N}^{+}\right]^{c}\right) \sim \frac{1}{N}$ but showing that it is o( $\left.N^{\frac{1}{2}}\right)$ is sufficient for what we need and requires less carefulness.

## 5 Laplace transform for smooth test functions, proof of Theorem 1.3

Section 3 allows us to justify in Proposition 5.1 the heuristics we gave in equation (6) for $\phi$ having compact support. We will then extend in Theorem 5.2 this result to a more general set of functions, by an approximation by compactly supported functions, using Corollary 4.2.

Proposition 5.1. For $\phi \in \mathcal{C}^{1}(\mathbb{R}, \mathbb{R})$ with compact support, we have for any real $t$, as $N$ goes to infinity,

$$
\begin{equation*}
\mathbb{E}_{N}^{V, P}\left[e^{t \sqrt{N} \nu_{N}(\Xi \phi)}\right] \rightarrow \exp \left\{\frac{t^{2}}{2} q_{P}(\phi)\right\} \tag{39}
\end{equation*}
$$

where $\Xi \phi$ is given by equation (7), and $q_{P}(\phi)$ is given by

$$
\begin{equation*}
q_{P}(\phi):=\int_{\mathbb{R}}\left(\phi^{\prime}(x)^{2}+V^{\prime \prime}(x) \phi(x)^{2}\right) d \mu_{P}(x)+P \iint_{\mathbb{R}^{2}}\left(\frac{\phi(x)-\phi(y)}{x-y}\right)^{2} d \mu_{P}(x) d \mu_{P}(y) \tag{40}
\end{equation*}
$$

Proof. Let $\phi \in \mathcal{C}_{c}^{1}(\mathbb{R}, \mathbb{R})$, and let $t \in \mathbb{R}$. We perform in equation (4) the change of variables $x_{i}=y_{i}+\frac{t}{\sqrt{N}} \phi\left(y_{i}\right), 1 \leqslant i \leqslant N$, which is a diffeomorphism for $N$ big enough. We thus have

$$
\begin{equation*}
Z_{N}^{V, P}=\int \prod_{1 \leqslant i<j \leqslant N}\left|y_{i}-y_{j}+\frac{t}{\sqrt{N}}\left(\phi\left(y_{i}\right)-\phi\left(y_{j}\right)\right)\right|^{2 P / N} . e^{-\sum_{i=1}^{N} V\left(y_{i}+\frac{t}{\sqrt{N}} \phi\left(y_{i}\right)\right)} \cdot \prod_{i=1}^{N}\left(1+\frac{t}{\sqrt{N}} \phi^{\prime}\left(y_{i}\right)\right) d^{N} \mathbf{y} \tag{41}
\end{equation*}
$$

and we develop separately the different terms of this integral. The first term can be written as:

$$
\prod_{i<j}\left|y_{i}-y_{j}\right|^{2 P / N} \prod_{i<j}\left|1+\frac{t}{\sqrt{N}} \frac{\phi\left(y_{i}\right)-\phi\left(y_{j}\right)}{y_{i}-y_{j}}\right|^{2 P / N}
$$

The second product above, setting $\Delta \phi_{i, j}:=\frac{\phi\left(y_{i}\right)-\phi\left(y_{j}\right)}{y_{i}-y_{j}}$ and using Taylor-Lagrange theorem, equals

$$
\exp \left(\frac{2 P}{N} \sum_{i<j} \ln \left|1+\frac{t}{\sqrt{N}} \frac{\phi\left(y_{i}\right)-\phi\left(y_{j}\right)}{y_{i}-y_{j}}\right|\right)=\exp \left(\frac{2 P}{N} \sum_{i<j}\left(\frac{t}{\sqrt{N}} \Delta \phi_{i, j}-\frac{t^{2}}{2 N}\left(\Delta \phi_{i, j}\right)^{2}+R_{N, 1}(i, j)\right)\right)
$$

where we noticed that $1+\frac{t}{\sqrt{N}} \Delta \phi_{i, j} \geqslant 1-\frac{t}{\sqrt{N}}\left\|\phi^{\prime}\right\|_{\infty}>0$ if $N$ is big enough, and where

$$
\left|R_{N, 1}(i, j)\right| \leqslant \frac{|t|^{3}}{3 N^{3 / 2}}\left\|\phi^{\prime}\right\|_{\infty}^{3}
$$

Again by Taylor-Lagrange theorem, the second term in 41) equals

$$
\exp \left(-\sum_{i=1}^{N}\left(V\left(y_{i}\right)+\frac{t}{\sqrt{N}} V^{\prime}\left(y_{i}\right) \phi\left(y_{i}\right)+\frac{t^{2}}{2 N} V^{\prime \prime}\left(y_{i}\right) \phi\left(y_{i}\right)^{2}+R_{N, 2}(i)\right)\right)
$$

where $R_{N, 2}(i)=\frac{t^{3}}{6 N^{3 / 2}} V^{(3)}\left(y_{i}+\frac{t \theta_{i}}{\sqrt{N}} \phi\left(y_{i}\right)\right) \phi\left(y_{i}\right)^{3}$ for some $\theta_{i} \in[0,1]$, thus for $N$ large enough

$$
\left|R_{N, 2}(i)\right| \leqslant \frac{|t|^{3}}{6 N^{3 / 2}}\|\phi\|_{\infty}^{3} \sup _{d(x, \text { supp } \phi) \leqslant 1}\left|V^{(3)}(x)\right|
$$

The last term reads

$$
\prod_{i=1}^{N}\left(1+\frac{t}{\sqrt{N}} \phi^{\prime}\left(y_{i}\right)\right)=\exp \left(\sum_{i=1}^{N}\left(\frac{t}{\sqrt{N}} \phi^{\prime}\left(y_{i}\right)-\frac{t^{2}}{2 N} \phi^{\prime}\left(y_{i}\right)^{2}+R_{N, 3}(i)\right)\right)
$$

with $\left|R_{N, 3}(i)\right| \leqslant \frac{t^{3}}{3 N^{3 / 2}}\left\|\phi^{\prime}\right\|_{\infty}^{3}$. Dividing both sides of equation 41) by $Z_{N}^{V, P}$ we get

$$
\begin{aligned}
\mathbb{E}_{N}^{V, P}[\exp \{t & \left.\sqrt{N}\left(P \iint_{\mathbb{R}^{2}} \frac{\phi(x)-\phi(y)}{x-y} d \hat{\mu}_{N}(x) d \hat{\mu}_{N}(y)+\int_{\mathbb{R}}\left(\phi^{\prime}-V^{\prime} \phi\right) d \hat{\mu}_{N}\right)\right\} \times \exp \left\{K_{N}(t, \phi)\right\} \\
& \left.\times \exp \left\{\frac{t^{2}}{2}\left(-P \iint_{\mathbb{R}^{2}}\left(\frac{\phi(x)-\phi(y)}{x-y}\right)^{2} d \hat{\mu}_{N}(x) d \hat{\mu}_{N}(y)-\int_{\mathbb{R}}\left(V^{\prime \prime} \phi^{2}+\phi^{\prime 2}\right) d \hat{\mu}_{N}\right)\right\}\right]=1
\end{aligned}
$$

with $\left|K_{N}(t, \phi)\right| \leqslant \frac{c(t, \phi)}{\sqrt{N}}$ where $c(t, \phi) \geqslant 0$ is independent of $N$. This bound shows that taking the limit $N \rightarrow \infty$ we can get rid of $K_{N}$ :

$$
\begin{aligned}
\lim _{N \rightarrow \infty} \mathbb{E}_{N}^{V, P}[ & \exp \left\{t \sqrt{N}\left(P \iint_{\mathbb{R}^{2}} \frac{\phi(x)-\phi(y)}{x-y} d \hat{\mu}_{N}(x) d \hat{\mu}_{N}(y)+\int_{\mathbb{R}}\left(\phi^{\prime}-V^{\prime} \phi\right) d \hat{\mu}_{N}\right)\right\} \\
& \left.\times \exp \left\{\frac{t^{2}}{2}\left(-P \iint_{\mathbb{R}^{2}}\left(\frac{\phi(x)-\phi(y)}{x-y}\right)^{2} d \hat{\mu}_{N}(x) d \hat{\mu}_{N}(y)-\int_{\mathbb{R}}\left(V^{\prime \prime} \phi^{2}+\phi^{2}\right) d \hat{\mu}_{N}\right)\right\}\right]=1
\end{aligned}
$$

Using Fubini's theorem (the function $(x, y) \mapsto \frac{\phi(x)-\phi(y)}{x-y}$ being bounded continuous on $\mathbb{R}^{2}$ ), the first line in the expectation value can be rewritten as $e^{t \sqrt{N} \Lambda_{N}}$ with

$$
\begin{equation*}
\Lambda_{N}:=2 P \iint_{\mathbb{R}^{2}} \frac{\phi(x)-\phi(y)}{x-y} d \mu_{P}(x) d\left(\hat{\mu}_{N}-\mu_{P}\right)(y)+\int_{\mathbb{R}}\left(\phi^{\prime}-V^{\prime} \phi\right) d\left(\hat{\mu}_{N}-\mu_{P}\right)+P \zeta_{N}(\phi) \tag{42}
\end{equation*}
$$

where we used equation (5) and $\zeta_{N}(\phi)$ is given by (34). Let $F: \mathcal{P}(\mathbb{R}) \rightarrow \mathbb{R}$ be defined by

$$
\begin{equation*}
F(\mu)=-P \iint_{\mathbb{R}^{2}}\left(\frac{\phi(x)-\phi(y)}{x-y}\right)^{2} d \mu(x) d \mu(y)-\int_{\mathbb{R}}\left(V^{\prime \prime} \phi^{2}+\phi^{\prime 2}\right) d \mu \tag{43}
\end{equation*}
$$

It is continuous for the topology of weak convergence since all the functions in the integrals are bounded continuous. So far we have established that

$$
\lim _{N \rightarrow \infty} \mathbb{E}_{N}^{V, P}\left[e^{t \sqrt{N} \Lambda_{N}+\frac{t^{2}}{2} F\left(\hat{\mu}_{N}\right)}\right]=1
$$

with $\Lambda_{N}$ given by 42 . We now replace in the latter equation the term $F\left(\hat{\mu}_{N}\right)$ by its limiting expression, $F\left(\mu_{P}\right)$. Fix a metric that is compatible with the weak convergence of probability measures on $\mathbb{R}$. For example,

$$
\begin{equation*}
d_{\mathrm{Lip}}(\mu, \nu)=\sup \left|\int f d \mu-\int f d \nu\right| \tag{44}
\end{equation*}
$$

where the supremum runs over $f: \mathbb{R} \rightarrow \mathbb{R}$ bounded and Lipschitz with $\|f\|_{\infty} \leqslant 1$ and Lipschitz constant $|f|_{\text {Lip }} \leqslant 1$. By the large deviations principle for $\left(\hat{\mu}_{N}\right)_{N}$ under the probability (3) established by [GZ19, Theorem 1.1], for all $\delta>0$ the event $\left\{d_{\text {Lip }}\left(\hat{\mu}_{N}, \mu_{P}\right)>\delta\right\}$ has (for $N$ big enough) probability smaller than $e^{-N c_{\delta}}$ where $c_{\delta}>0$. Hence,

$$
\lim _{N \rightarrow \infty} \mathbb{E}_{N}^{V, P}\left[e^{t \sqrt{N} \Lambda_{N}+\frac{t^{2}}{2} F\left(\hat{\mu}_{N}\right)}\right]=\lim _{N \rightarrow \infty} \mathbb{E}_{N}^{V, P}\left[\mathbf{1}_{\left\{d_{\text {Lip }}\left(\hat{\mu}_{N}, \mu_{P}\right) \leqslant \delta\right\}} e^{t \sqrt{N} \Lambda_{N}+\frac{t^{2}}{2} F\left(\hat{\mu}_{N}\right)}\right]
$$

By continuity of $F$ there is some $\varepsilon(\delta)$ which goes to 0 as $\delta \rightarrow 0$ such that, for $d_{\text {Lip }}\left(\nu, \mu_{P}\right) \leqslant \delta$, we have $\left|F(\nu)-F\left(\mu_{P}\right)\right| \leqslant \varepsilon(\delta)$. Taking the (decreasing) limit as $\delta$ goes to zero we deduce

$$
\lim _{N \rightarrow \infty} \mathbb{E}_{N}^{V, P}\left[e^{t \sqrt{N} \Lambda_{N}+\frac{t^{2}}{2} F\left(\hat{\mu}_{N}\right)}\right]=\lim _{\delta \rightarrow 0} \lim _{N \rightarrow \infty} \mathbb{E}_{N}^{V, P}\left[\mathbf{1}_{\left\{d_{\text {Lip }}\left(\hat{\mu}_{N}, \mu_{P}\right) \leqslant \delta\right\}} e^{t \sqrt{N} \Lambda_{N}}\right] e^{\frac{t^{2}}{2} F\left(\mu_{P}\right)}
$$

But the same large deviations argument shows that

$$
\lim _{\delta \rightarrow 0} \lim _{N \rightarrow \infty} \mathbb{E}_{N}^{V, P}\left[\mathbf{1}_{\left\{d_{\text {Lip }}\left(\hat{\mu}_{N}, \mu_{P}\right) \leqslant \delta\right\}} e^{t \sqrt{N} \Lambda_{N}}\right]=\lim _{N \rightarrow \infty} \mathbb{E}_{N}^{V, P}\left[e^{t \sqrt{N} \Lambda_{N}}\right]
$$

Thus, we have shown that

$$
\begin{equation*}
\lim _{N \rightarrow \infty} \mathbb{E}_{N}^{V, P}\left[e^{t \sqrt{N}\left(2 P \iint_{\mathbb{R}^{2}} \frac{\phi(x)-\phi(y)}{x-y} d \mu_{P}(x) d\left(\hat{\mu}_{N}-\mu_{P}\right)(y)+\int_{\mathbb{R}}\left(\phi^{\prime}-V^{\prime} \phi\right) d\left(\hat{\mu}_{N}-\mu_{P}\right)+P \zeta_{N}(\phi)\right)}\right]=e^{-\frac{t^{2}}{2} F\left(\mu_{P}\right)}, \tag{45}
\end{equation*}
$$

Which establishes that $\sqrt{N} \Lambda_{N}=\sqrt{N}\left(\nu_{N}(\Xi \phi)+P \zeta_{N}(\phi)\right)$ converges in law towards a centered Gaussian random variable with announced variance. We finally get rid of the remaining term $\zeta_{N}(\phi)$, using Corollary 3.4 taking $\varepsilon=1 / 4$ for example, we see in particular that $\sqrt{N} \zeta_{N}(\phi)$ converges in probability towards zero. The conclusion follows from Slutsky's lemma.

We now extend the result of Proposition 5.1 to a more general set of functions. With the notations of Proposition 5.1. we have

Theorem 5.2. Let $\phi \in H^{2}(\mathbb{R}) \cap \mathcal{C}^{2}(\mathbb{R})$ such that $\phi^{\prime \prime}$ is bounded. Additionally, suppose that $V^{(3)} \phi^{2}$, $V^{\prime \prime} \phi \phi^{\prime}, V^{\prime \prime} \phi^{2}$ and $V^{\prime} \phi$ are bounded. Then, recalling 40 we have the convergence in distribution as $N$ goes to infinity

$$
\sqrt{N} \nu_{N}(\Xi \phi) \rightarrow \mathcal{N}\left(0, q_{P}(\phi)\right)
$$

Proof. For $N \geqslant 1$, let $E_{N}^{-}, E_{N}^{+}$be given by Corollary 4.2 . Let $\chi_{N}: \mathbb{R} \rightarrow[0,1]$ be $\mathcal{C}^{2}$ with compact support such that

$$
\chi_{N}(x)=1 \text { for } x \in\left[E_{N}^{-}-1, E_{N}^{+}+1\right] \text { and } \chi_{N}(x)=0 \text { for } x \in\left[E_{N}^{-}-2, E_{N}^{+}+2\right]^{c}
$$

and such that, denoting $\phi_{N}=\phi \chi_{N}, \sup _{N}\left\|\phi_{N}^{\prime}\right\|_{\infty}+\left\|\phi_{N}^{\prime}\right\|_{L^{2}(\mathbb{R})}, \sup _{N}\left\|\phi_{N}^{\prime \prime}\right\|_{\infty}+\left\|\phi_{N}^{\prime \prime}\right\|_{L^{2}(\mathbb{R})}<+\infty$ (we assumed $\phi \in H^{2}(\mathbb{R})$, in particular $\phi^{\prime}$ is bounded and such a $\chi_{N}$ exists). The point of cutting $\phi$ outside the set $\left[E_{N}^{-}-1, E_{N}^{+}+1\right]$ is that with high probability, the empirical measure $\hat{\mu}_{N}$ doesn't see the difference between $\phi$ and $\phi_{N}$.

The support of $\phi_{N}$ is then contained in $\left[E_{N}^{-}-2, E_{N}^{+}+2\right]$, and we now argue that the proof of Proposition 5.1 can be adapted so that

$$
\begin{equation*}
\sqrt{N} \nu_{N}\left(\Xi \phi_{N}\right) \rightarrow \mathcal{N}\left(0, q_{P}(\phi)\right) \tag{46}
\end{equation*}
$$

Similarly as in Proposition 5.1 we perform in $Z_{N}^{V, P}$ the change of variables $x_{i}=y_{i}+\frac{t}{\sqrt{N}} \phi_{N}\left(y_{i}\right)$, $1 \leqslant i \leqslant N$, which is the same as before, but with $\phi$ replaced by $\phi_{N}$. First, with $I_{N}:=\left[E_{N}^{-}-2, E_{N}^{+}+2\right]$, the error term

$$
K_{N}\left(t, \phi_{N}\right) \leqslant 2 \frac{t^{3}}{3 N^{1 / 2}}\left\|\phi_{N}^{\prime}\right\|_{\infty}^{3}+\frac{t^{3}}{6 N^{1 / 2}}\left\|\phi_{N}\right\|_{\infty} \sup _{d\left(x, I_{N}\right) \leqslant 1}\left|V^{(3)}(x)\right|
$$

of the proof of Proposition 5.1 is still going to zero, because of our choice of $\chi_{N}$ and Assumption 1.2 As previously, we then have

$$
\begin{equation*}
\lim _{N \rightarrow \infty} \mathbb{E}_{N}^{V, P}\left[e^{t \sqrt{N} \Lambda_{N}\left(\phi_{N}\right)+\frac{t^{2}}{2} F_{N}\left(\hat{\mu}_{N}\right)}\right]=1 \tag{47}
\end{equation*}
$$

with

$$
\Lambda_{N}\left(\phi_{N}\right):=2 P \iint_{\mathbb{R}^{2}} \frac{\phi_{N}(x)-\phi_{N}(y)}{x-y} d \mu_{P}(x) d\left(\hat{\mu}_{N}-\mu_{P}\right)(y)+\int_{\mathbb{R}}\left(\phi_{N}^{\prime}-V^{\prime} \phi_{N}\right) d\left(\hat{\mu}_{N}-\mu_{P}\right)+P \zeta_{N}\left(\phi_{N}\right)
$$

where $\zeta_{N}$ is given by (34), and

$$
F_{N}\left(\hat{\mu}_{N}\right)=-P \iint_{\mathbb{R}^{2}}\left(\frac{\phi_{N}(x)-\phi_{N}(y)}{x-y}\right)^{2} d \hat{\mu}_{N}(x) d \hat{\mu}_{N}(y)-\int_{\mathbb{R}}\left(V^{\prime \prime} \phi_{N}^{2}+\phi_{N}^{\prime 2}\right) d \hat{\mu}_{N}
$$

Taking again the distance $d_{\text {Lip }}$ defined in (44), one can check that for $\mu, \nu$ probability measures over $\mathbb{R}$,

$$
\left|F_{N}(\mu)-F_{N}(\nu)\right| \leqslant C_{N} d_{\operatorname{Lip}}(\mu, \nu)
$$

where $C_{N}$ is a term depending on the norms $\left\|\phi_{N}^{\prime}\right\|_{\infty},\left\|\phi_{N}^{\prime \prime}\right\|_{\infty},\left\|V^{\prime \prime} \phi_{N}^{2}\right\|_{\infty}$ and $\left\|\left(V^{\prime \prime} \phi_{N}^{2}\right)^{\prime}\right\|_{\infty}$. The choice of $\chi_{N}$ and the fact that $\phi$ is chosen so that $V^{(3)} \phi^{2}$ and $V^{\prime \prime} \phi \phi^{\prime}$ are bounded guarantee that $\left\|\left(V^{\prime \prime} \phi_{N}^{2}\right)^{\prime}\right\|_{\infty}$ is bounded in $N$. The other norms are easily bounded by hypothesis. Therefore $C_{N}$ can be seen to be uniformly bounded in $N$, and we find some $C \geqslant 0$ independent of $N$ such that

$$
\left|F_{N}(\mu)-F_{N}(\nu)\right| \leqslant C d_{\operatorname{Lip}}(\mu, \nu)
$$

As in proposition 5.1 we use the large deviation principle for $\left(\hat{\mu}_{N}\right)$ to deduce

$$
\lim _{N \rightarrow+\infty} \mathbb{E}_{N}^{V, P}\left[e^{t \sqrt{N} \Lambda_{N}\left(\phi_{N}\right)+\frac{t^{2}}{2} F_{N}\left(\hat{\mu}_{N}\right)}\right]=\lim _{N \rightarrow+\infty} \mathbb{E}_{N}^{V, P}\left[e^{t \sqrt{N} \Lambda_{N}\left(\phi_{N}\right)}\right] e^{\frac{t^{2}}{2} F_{N}\left(\mu_{P}\right)}
$$

By dominated convergence, $F_{N}\left(\mu_{P}\right)$ converges to $F\left(\mu_{P}\right)$, the function $F$ being given by 43). This shows the convergence as $N$ goes to infinity

$$
\lim _{N \rightarrow+\infty} \mathbb{E}_{N}^{V, P}\left[e^{t \sqrt{N} \Lambda_{N}\left(\phi_{N}\right)}\right]=e^{-\frac{t^{2}}{2} F\left(\mu_{P}\right)}
$$

and $\sqrt{N}\left(\nu_{N}\left(\Xi \phi_{N}\right)+P \zeta_{N}\left(\phi_{N}\right)\right)$ converges towards a centered Gaussian variable with variance $-F\left(\mu_{P}\right)=$ $q_{P}(\phi)$. Because $\sup _{N}\left\|\phi_{N}^{\prime}\right\|_{L^{2}(d x)}+\left\|\phi_{N}^{\prime \prime}\right\|_{L^{2}(d x)}$ is finite, we can apply again Corollary 3.4 to deduce the convergence in law 46).
We now have the ingredients to conclude, by showing that the characteristic function

$$
\mathbb{E}_{N}^{V, P}\left[e^{\mathrm{i} t \sqrt{N} \nu_{N}(\Xi \phi)}\right]=\mathbb{E}_{N}^{V, P}\left[e^{\mathrm{i} t \sqrt{N} \int \Xi \phi d \hat{\mu}_{N}}\right] e^{-\mathrm{i} t \sqrt{N} \int \Xi \phi d \mu_{P}}
$$

converges to the characteristic of a Gaussian variable with appropriate variance. By Corollary 4.2 the probability under $\mathbb{P}_{N}^{V, P}$ of the event $\mathcal{E}_{N}=\left\{x_{1}, \ldots, x_{N} \in\left[E_{N}^{-}-1, E_{N}^{+}+1\right]\right\}$ converges to 1 . Along with the convergence (46), we deduce

$$
e^{-\frac{t^{2}}{2} q_{P}(\phi)}=\lim _{N} \mathbb{E}_{N}^{V, P}\left[e^{\mathrm{i} t \sqrt{N} \int \Xi \phi_{N} d \hat{\mu}_{N}}\right] e^{-\mathrm{i} t \sqrt{N} \int \Xi \phi_{N} d \mu_{P}}=\lim _{N} \mathbb{E}_{N}^{V, P}\left[\mathbf{1}_{\mathcal{E}_{N}} e^{\mathrm{i} t \sqrt{N} \int \Xi \phi_{N} d \hat{\mu}_{N}}\right] e^{-\mathrm{i} t \sqrt{N} \int \Xi \phi_{N} d \mu_{P}}
$$

Where we used

$$
\left|\mathbb{E}_{N}^{V, P}\left[\mathbf{1}_{\mathcal{E}_{N}^{c}} e^{\mathrm{i} t \sqrt{N} \int \Xi \phi_{N} d \hat{\mu}_{N}}\right] e^{-\mathrm{i} t \sqrt{N} \int \Xi \phi_{N} d \mu_{P}}\right| \leqslant \mathbb{P}_{N}^{V, P}\left(\mathcal{E}_{N}^{c}\right) \xrightarrow[N \rightarrow+\infty]{ } 0
$$

Using that $\phi_{N}=\phi$ on $J_{N}=\left[E_{N}^{-}-1, E_{N}^{+}+1\right]$,

$$
\begin{aligned}
\int \Xi \phi_{N} d \mu_{P} & =2 P \iint^{\frac{\phi_{N}(x)-\phi_{N}(y)}{x-y} d \mu_{P}(x) d \mu_{P}(y)+\int\left(\phi_{N}^{\prime}-V^{\prime} \phi_{N}\right) d \mu_{P}} \\
& =2 P \iint_{J_{N}^{2}} \frac{\phi(x)-\phi(y)}{x-y} d \mu_{P}(x) d \mu_{P}(y)+2 P \iint_{\left(J_{N}^{2}\right)^{c}} \frac{\phi_{N}(x)-\phi_{N}(y)}{x-y} d \mu_{P}(x) d \mu_{P}(y) \\
& +\int_{J_{N}}\left(\phi^{\prime}-V^{\prime} \phi\right) d \mu_{P}+\int_{J_{N}^{c}}\left(\phi \chi_{N}^{\prime}+\phi^{\prime} \chi_{N}-V^{\prime} \phi \chi_{N}\right) d \mu_{P}
\end{aligned}
$$

By boundedness of $\left(\left\|\phi_{N}^{\prime}\right\|_{\infty}\right)_{N}$, the second term is bounded by

$$
C_{P} \iint_{\left(J_{N}^{2}\right)^{c}} d \mu_{P} d \mu_{P} \leqslant 2 C_{P} \mu_{P}\left(J_{N}^{c}\right)=o\left(N^{-1 / 2}\right)
$$

where we used the union bound and Lemma 4.4. By the same estimate and the fact that $\chi_{N}$ can be chosen so that $\left(\left\|\chi_{N}^{\prime}\right\|_{\infty}\right)_{N}$ is bounded, and because $\phi^{\prime}, V^{\prime} \phi$ are bounded, the last term is also $o\left(N^{-1 / 2}\right)$. By the previous arguments, we also conclude that

$$
2 P \iint_{\left(J_{N}^{2}\right)^{c}} \frac{\phi(x)-\phi(y)}{x-y} d \mu_{P}(x) d \mu_{P}(y)+\int_{J_{N}^{c}}\left(\phi^{\prime}-V^{\prime} \phi\right) d \mu_{P}=o\left(N^{-1 / 2}\right),
$$

thus

$$
\int \Xi \phi_{N} d \mu_{P}=\int \Xi \phi d \mu_{P}+o\left(N^{-1 / 2}\right)
$$

and so far we have

$$
e^{-\frac{t^{2}}{2} q_{P}(\phi)}=\lim _{N} \mathbb{E}_{N}^{V, P}\left[\mathbf{1}_{\mathcal{E}_{N}} e^{\mathrm{i} t \sqrt{N} \int \Xi \phi_{N} d \hat{\mu}_{N}}\right] e^{-\mathrm{i} t \sqrt{N} \int \Xi \phi d \mu_{P}}
$$

Finally, on $\mathcal{E}_{N}$, using $\phi_{N}=\phi$ and that $\hat{\mu}_{N}$ is supported in $J_{N}$,

$$
\begin{aligned}
\int \Xi \phi_{N} d \hat{\mu}_{N} & =2 P \iint_{J_{N}^{2}} \frac{\phi(x)-\phi(y)}{x-y} d \mu_{P}(x) d \hat{\mu}_{N}(y)+2 P \iint_{\left(J_{N}^{2}\right)^{c}} \frac{\phi_{N}(x)-\phi_{N}(y)}{x-y} d \mu_{P}(x) d \hat{\mu}_{N}(y)+\int_{J_{N}}\left(\phi^{\prime}-V^{\prime} \phi\right) d \hat{\mu}_{N} \\
& =2 P \iint^{\phi(x)-\phi(y)} \\
x-y & \mu_{P}(x) d \hat{\mu}_{N}(y)+\int\left(\phi^{\prime}-V^{\prime} \phi\right) d \hat{\mu}_{N}+o\left(N^{-1 / 2}\right)
\end{aligned}
$$

Where in the second line we used, using Lemma 4.4 again, that

$$
\iint_{\left(J_{N}^{2}\right)^{c}} \frac{\phi_{N}(x)-\phi_{N}(y)}{x-y} d \mu_{P}(x) d \hat{\mu}_{N}(y)=\iint_{J_{N} \times J_{N}^{c}} \frac{\phi_{N}(x)-\phi_{N}(y)}{x-y} d \mu_{P}(x) d \hat{\mu}_{N}(y)=o\left(N^{-1 / 2}\right),
$$

and the same estimate holds for $\phi_{N}$ replaced by $\phi$. Therefore,

$$
e^{-\frac{t^{2}}{2} q_{P}}(\phi)=\lim _{N} \mathbb{E}_{N}^{V, P}\left[\mathbf{1}_{\mathcal{E}_{N}} e^{\mathrm{i} t \sqrt{N} \int \Xi \phi d \hat{\mu}_{N}}\right] e^{-\mathrm{i} t \sqrt{N} \int \Xi \phi d \mu_{P}}
$$

This establishes that

$$
\lim _{N} \mathbb{E}_{N}^{V, P}\left[e^{\mathrm{i} t \sqrt{N} \int \Xi \phi d \hat{\nu}_{N}}\right]=e^{-\frac{t^{2}}{2} q_{P}(\phi)},
$$

which concludes the proof.
Remark 5.3. Taking $\phi$ such that $\phi^{\prime}$ satisfies the conditions of Theorem 5.2, we then have

$$
\begin{equation*}
\mathbb{E}_{N}^{V, P}\left[e^{t \sqrt{N} \nu_{N}(\mathcal{L} \phi)}\right] \underset{N \rightarrow \infty}{\longrightarrow} \exp \left\{\frac{t^{2}}{2} q_{P}\left(\phi^{\prime}\right)\right\} \tag{48}
\end{equation*}
$$

where the operator $\mathcal{L}$ is defined as $\mathcal{L} \phi:=\Xi \phi^{\prime}$, ie

$$
\begin{equation*}
\mathcal{L} \phi=2 P \int_{\mathbb{R}} \frac{\phi^{\prime}(x)-\phi^{\prime}(y)}{x-y} d \mu_{P}(y)+\phi^{\prime \prime}(x)-V^{\prime}(x) \phi^{\prime}(x) \tag{49}
\end{equation*}
$$

Note that $q_{P}^{V}\left(\phi^{\prime}\right)=\left(\sigma_{P}^{V}\right)^{2}(\mathcal{L} \phi)$ where $\sigma_{P}^{V}$ is defined in 11. By Theorem 7.1. the class of functions in $\mathcal{L}^{-1}(\mathcal{T})$ where

$$
\mathcal{T}:=\left\{f \in \mathcal{C}^{1}(\mathbb{R}), f=O(1 / x), f^{\prime}=O\left(1 / x^{2}\right), \int_{\mathbb{R}} f \rho_{P}=0\right\}
$$

satisfies 48. This proves Theorem 1.3 .

## 6 Inversion of $\mathcal{L}$

This section is dedicated to the definition of $\mathcal{L}$ given by 8 and its domain and then we focus on its inversion. We rely heavily on results of Appendix A the diagonalization of the operator $\mathcal{A}$ by the use of the theory of Schrödinger operators.

Let $P>0$ be fixed. We introduce the operators $\mathcal{A}$ and $\mathcal{W}$, acting on sufficiently smooth functions of $L^{2}\left(\rho_{P}\right)$, by

$$
\begin{equation*}
\mathcal{A} \phi=-\frac{\left(\phi^{\prime} \rho_{P}\right)^{\prime}}{\rho_{P}}=-\left(\phi^{\prime \prime}+\frac{\rho_{P}^{\prime}}{\rho_{P}} \phi^{\prime}\right) \quad \text { and } \quad \mathcal{W} \phi=-\mathcal{H}\left[\phi^{\prime} \rho_{P}\right] \tag{50}
\end{equation*}
$$

We first show the following decomposition of $\mathcal{L}$.
Lemma 6.1. For $\phi$ twice differentiable we have the following pointwise identity

$$
\begin{equation*}
-\mathcal{L} \phi=\mathcal{A} \phi+2 P \mathcal{W} \phi \tag{51}
\end{equation*}
$$

Proof. We write for $x \in \mathbb{R}$

$$
\begin{equation*}
2 P \int_{\mathbb{R}} \frac{\phi^{\prime}(x)-\phi^{\prime}(y)}{x-y} \rho_{P}(y) d y=-2 P \phi^{\prime}(x) \mathcal{H}\left[\rho_{P}\right](x)+2 P \mathcal{H}\left[\phi^{\prime} \rho_{P}\right](x) \tag{52}
\end{equation*}
$$

Then,

$$
\mathcal{L} \phi=\phi^{\prime \prime}-V^{\prime} \phi^{\prime}-2 P \phi^{\prime} \mathcal{H}\left[\rho_{P}\right]+2 P \mathcal{H}\left[\phi^{\prime} \rho_{P}\right] .
$$

By (18) we have $-V^{\prime}-2 P \mathcal{H}\left[\rho_{P}\right]=\frac{\rho_{P}^{\prime}}{\rho_{P}}$, which concludes the proof.
In order to state the next theorem, whose proof we detail in the Appendix, we introduce the following Sobolev-type spaces. Let

$$
H_{V^{\prime}}^{2}(\mathbb{R}):=\left\{u \in H^{2}(\mathbb{R}), u V^{\prime} \in L^{2}(\mathbb{R})\right\}
$$

We now define

$$
H_{V^{\prime}}^{2}\left(\rho_{P}\right):=\rho_{P}^{-1 / 2} H_{V^{\prime}}^{2}(\mathbb{R})
$$

and its homogeneous counterpart

$$
H_{V^{\prime}, 0}^{2}\left(\rho_{P}\right):=\left\{u \in H_{V^{\prime}}^{2}\left(\rho_{P}\right), \int_{\mathbb{R}} u \rho_{P} d x=0\right\}
$$

Finally, we let $L_{0}^{2}\left(\rho_{P}\right)$ be the subset of $L^{2}\left(\rho_{P}\right)$ of zero mean functions with respect to $\rho_{P}$.
We detail the proof of the following theorem in Appendix A which is based on Schrödinger operators theory.

Theorem 6.2 (Diagonalization of $\mathcal{A}$ in $L_{0}^{2}\left(\rho_{P}\right)$ ). There exists a sequence $0<\lambda_{1}<\lambda_{2}<\ldots$ going to infinity, and a complete orthonormal set $\left(\phi_{n}\right)_{n \geqslant 1}$ of $L_{0}^{2}\left(\rho_{P}\right)$ of associated eigenfunctions for $\mathcal{A}$, meaning that

- $\operatorname{span}\left\{\phi_{n}, n \geqslant 1\right\}$ is dense in $L_{0}^{2}\left(\rho_{P}\right)$,
- For all $i, j,\left\langle\phi_{i}, \phi_{j}\right\rangle_{L^{2}\left(\rho_{P}\right)}=\delta_{i, j}$,
- For all $n \geqslant 1, \mathcal{A} \phi_{n}=\lambda_{n} \phi_{n}$.

Furthermore, each $\phi_{n}$ is in $H_{V^{\prime}, 0}^{2}\left(\rho_{P}\right), \mathcal{A}: H_{V^{\prime}, 0}^{2}\left(\rho_{P}\right) \rightarrow L_{0}^{2}\left(\rho_{P}\right)$ is bijective, and we have the writing, for $u \in L_{0}^{2}\left(\rho_{P}\right)$

$$
\mathcal{A}^{-1} u=\sum_{n \geqslant 1} \lambda_{n}^{-1}\left\langle u, \phi_{n}\right\rangle_{L^{2}\left(\rho_{P}\right)} \phi_{n}
$$

We see the operators $\mathcal{A}$ and $\mathcal{W}$ as unbounded operators on the space

$$
\mathrm{H}=\left\{u \in H^{1}\left(\rho_{P}\right) \mid \int_{\mathbb{R}} u \rho_{P} d x=0\right\}
$$

endowed with the inner product $\langle u, v\rangle_{\mathrm{H}}=\left\langle u^{\prime}, v^{\prime}\right\rangle_{L^{2}\left(\rho_{P}\right)}$. This defines an inner product on H and makes it a complete space: it can be seen that $H^{1}\left(\rho_{P}\right)$ is the completion of $\mathcal{C}_{c}^{\infty}(\mathbb{R})$ with respect to the inner product $\langle u, v\rangle_{L^{2}\left(\rho_{P}\right)}+\left\langle u^{\prime}, v^{\prime}\right\rangle_{L^{2}\left(\rho_{P}\right)}$. The space H is then the kernel of the bounded (with respect to $\|\cdot\|_{\mathrm{H}}$ ) linear form, $\langle\widetilde{1}, \cdot\rangle_{L^{2}\left(\rho_{P}\right)}$ on $H^{1}\left(\rho_{P}\right)$, and both inner products are equivalent on H because of the Poincaré inequality, Proposition 2.4 The use of H is motivated by the fact that both $\mathcal{A}$ and $\mathcal{W}$ are self-adjoint positive on this space as we show in Lemma 6.4.

In the next proposition, we deduce from Theorem 6.2 the diagonalization of $\mathcal{A}$ in H .
Proposition 6.3 (Diagonalization of $\mathcal{A}$ in H ). With the same eigenvalues $0<\lambda_{1}<\lambda_{2}<\ldots$ as in Theorem 6.2. there exists a complete orthonormal set $\left(\psi_{n}\right)_{n \geqslant 1}$ of H formed by eigenfunctions of $\mathcal{A}$.
Proof. With $\left(\phi_{n}\right)_{n \geqslant 1}$ of Theorem 6.2 .

$$
\begin{aligned}
\delta_{i, j}=\left\langle\phi_{i}, \phi_{j}\right\rangle_{L^{2}\left(\rho_{P}\right)} & =\frac{1}{\lambda_{j}}\left\langle\phi_{i}, \mathcal{A} \phi_{j}\right\rangle_{L^{2}\left(\rho_{P}\right)} \\
& =\frac{1}{\lambda_{j}}\left\langle\phi_{i}^{\prime}, \phi_{j}^{\prime}\right\rangle_{L^{2}\left(\rho_{P}\right)} \\
& =\frac{1}{\lambda_{j}}\left\langle\phi_{i}, \phi_{j}\right\rangle_{\mathrm{H}} .
\end{aligned}
$$

With $\psi_{n}=\frac{1}{\sqrt{\lambda_{n}}} \phi_{n},\left(\psi_{n}\right)_{n \geqslant 1}$ is then orthonormal with respect to the inner product of H . To show that $\operatorname{span}\left\{\psi_{n}, n \geqslant 1\right\}$ is dense in H , let $u \in \mathrm{H}$ be such that for all $j \geqslant 1,\left\langle u, \phi_{j}\right\rangle_{\mathrm{H}}=0$. In the last series of equalities, replace $\phi_{i}$ by $u$ : we see that $u$ is orthogonal to each $\phi_{j}$ in $L^{2}\left(\rho_{P}\right)$, thus $u$ is a constant as shown in the proof of Lemma A.10 and because $u \in \mathrm{H}$ it has zero mean against $\rho_{P}$. This shows that $u=0$.

We set for what follows $\mathcal{D}(\mathcal{A})=\left\{u \in H_{V^{\prime}, 0}^{2}\left(\rho_{P}\right) \mid \mathcal{A} u \in \mathrm{H}\right\}$ and $\mathcal{D}(\mathcal{W})=\{u \in \mathrm{H} \mid \mathcal{W} u \in \mathrm{H}\}$.
Lemma 6.4. The following properties hold:

- The operator $\mathcal{W}: \mathcal{D}(\mathcal{W}) \rightarrow \mathrm{H}$ is positive: for all $\phi \in \mathcal{D}(\mathcal{W})$,

$$
\langle\mathcal{W} \phi, \phi\rangle_{\mathrm{H}}=\frac{1}{2}\left\|\phi^{\prime} \rho_{P}\right\|_{1 / 2}^{2} \geqslant 0
$$

with equality only for $\phi=0$, where the $1 / 2$-norm of $u$ is given by

$$
\|u\|_{1 / 2}^{2}=\int_{\mathbb{R}}|x| \cdot|\mathcal{F}[u](x)|^{2} d x
$$

- Both $\mathcal{A}$ and $\mathcal{W}$ are self-adjoint for the inner product of H .

Proof. To prove the first point, let $\phi \in \mathcal{D}(\mathcal{W})$. Then,

$$
\begin{aligned}
2 \pi\langle\mathcal{W} \phi, \phi\rangle_{\mathrm{H}}=-2 \pi\left\langle\mathcal{H}\left[\phi^{\prime} \rho_{P}\right]^{\prime}, \phi^{\prime} \rho_{P}\right\rangle_{L^{2}(d x)}=- & \left\langle i x \mathcal{F}\left[\mathcal{H}\left[\phi^{\prime} \rho_{P}\right]\right], \mathcal{F}\left[\phi^{\prime} \rho_{P}\right]\right\rangle_{L^{2}(d x)} \\
& =\pi\langle | x\left|\mathcal{F}\left[\phi^{\prime} \rho_{P}\right], \mathcal{F}\left[\phi^{\prime} \rho_{P}\right]\right\rangle_{L^{2}(d x)}=\pi\left\|\phi^{\prime} \rho_{P}\right\|_{1 / 2}^{2} \geqslant 0
\end{aligned}
$$

and because $\phi$ is in H , this last quantity is zero if and only if $\phi$ vanishes.
For the second point, let $u, v \in \mathcal{D}(\mathcal{W})$. Using Plancherel's isometry and $i$ ) of Lemma 2.1.

$$
\langle\mathcal{W} u, v\rangle_{\mathrm{H}}=\left\langle(\mathcal{W} u)^{\prime}, v^{\prime} \rho_{P}\right\rangle_{L^{2}(d x)}=\frac{1}{2}\langle | x\left|\mathcal{F}\left[u^{\prime} \rho_{P}\right], \mathcal{F}\left[v^{\prime} \rho_{P}\right]\right\rangle_{L^{2}(d x)}
$$

and this last expression is symmetric in $(u, v)$. The proof of the self-adjointness of $\mathcal{A}$ follows from integration by partss.

Definition 6.5 (Quadratic form associated to $-\mathcal{L}$ ). We define for all $u, v \in \mathrm{H} \cap \mathcal{C}_{c}^{\infty}(\mathbb{R})$ the quadratic form associated to $-\mathcal{L}$ by

$$
q_{-\mathcal{L}}(u, v)=\langle\mathcal{A} u, \mathcal{A} v\rangle_{L^{2}\left(\rho_{P}\right)}+2 P\left\langle\mathcal{F}\left[u^{\prime} \rho_{P}\right], \mathcal{F}\left[v^{\prime} \rho_{P}\right]\right\rangle_{L^{2}(|x| d x)}
$$

Note that for all $u, v \in \mathrm{H} \cap \mathcal{C}_{c}^{\infty}(\mathbb{R}), q_{-\mathcal{L}}(u, v)=\langle-\mathcal{L} u, v\rangle_{\mathrm{H}}$ and that whenever $u \in \mathcal{D}(\mathcal{A}) \cap \mathcal{D}(\mathcal{W})$,

$$
\begin{equation*}
q_{-\mathcal{L}}(u, u)=\langle\mathcal{A} u, u\rangle_{\mathrm{H}}+2 P\langle\mathcal{W} u, u\rangle_{\mathrm{H}} \geqslant \lambda_{1}(\mathcal{A})\|u\|_{\mathrm{H}}^{2} \tag{53}
\end{equation*}
$$

by Proposition 6.3 and Lemma 6.4 After extending the $q_{-\mathcal{L}}$ to its form domain $Q(\mathcal{L})$ which is equal to $\left\{u \in \mathrm{H}, \mathcal{A} u \in L^{2}\left(\rho_{P}\right), \mathcal{F}\left[u^{\prime} \rho_{P}\right] \in L^{2}(|x| d x)\right\}=H_{V^{\prime}, 0}^{2}\left(\rho_{P}\right)$. The equality comes from the fact that $\mathcal{A}^{-1}\left(L_{0}^{2}\left(\rho_{P}\right)\right)=H_{V^{\prime}, 0}^{2}\left(\rho_{P}\right)$, that $\mathrm{H} \subset H_{V^{\prime}, 0}^{2}\left(\rho_{P}\right)$ and that $\mathcal{F}\left[u^{\prime} \rho_{P}\right] \in L^{2}\left(x^{2} d x\right)$ whenever $u \in H_{V^{\prime}, 0}^{2}\left(\rho_{P}\right)$, indeed $u^{\prime} \rho_{P} \in H^{1}(\mathbb{R})$ because $\left(u^{\prime} \rho_{P}\right)^{\prime}=-\rho_{P} \mathcal{A} u \in L^{2}(\mathbb{R})$. We now define $\mathcal{D}(\mathcal{L})$ the domain of definition of $-\mathcal{L}$ by:

$$
\mathcal{D}(\mathcal{L}):=\left\{u \in Q(\mathcal{L}), v \mapsto q_{-\mathcal{L}}(u, v) \text { can be extended to a continuous linear form on } \mathrm{H}\right\}
$$

Proposition 6.6. $\mathcal{D}(\mathcal{L})=\mathcal{D}(\mathcal{A}) \cap \mathcal{D}(\mathcal{W})$.
Proof. Let $u \in \mathcal{D}(\mathcal{L})$, by Riesz's theorem there exists $f_{u} \in \mathrm{H}$, such that $q_{-\mathcal{L}}(u, v)=\left\langle f_{u}, v\right\rangle_{\mathrm{H}}$ for all $v \in \mathrm{H}$, we set $-\mathcal{L} u:=f_{u}$, it is called the Friedrichs extension of $-\mathcal{L}$. Then for all $v \in \mathrm{H} \cap \mathcal{C}_{c}^{\infty}(\mathbb{R})$, by integration by part we get:

$$
\langle-\mathcal{L} u, v\rangle_{\mathrm{H}}=q_{-\mathcal{L}}(u, v)=\langle u, \mathcal{A} v\rangle_{\mathrm{H}}+2 P\langle u, \mathcal{W} v\rangle_{\mathrm{H}}
$$

hence we deduce the distributionnal identity $-\mathcal{L} u=\mathcal{A} u+2 P \mathcal{W} u$. Since $u \in H_{V^{\prime}, 0}^{2}\left(\rho_{P}\right), \mathcal{W} u \in H^{1}\left(\rho_{P}\right)$ implying that $\mathcal{A} u \in \mathrm{H}$ and then that $\mathcal{W} u \in \mathrm{H}$.

We are now ready to state the main theorem of this section, that is the inversion of $\mathcal{L}$ on $\mathcal{D}(\mathcal{L})$.
Theorem 6.7 (Inversion of $\mathcal{L}) .-\mathcal{L}: \mathcal{D}(\mathcal{L}) \longrightarrow \mathrm{H}$ is bijective. Furthermore, $(-\mathcal{L})^{-1}$ is continuous from $\left(\mathrm{H},\|.\|_{\mathrm{H}}\right)$ to $\left(\mathcal{D}(\mathcal{L}), q_{-\mathcal{L}}\right)$.

Proof. Let $f \in \mathrm{H}$, since $\langle f, .\rangle_{\mathrm{H}}$ is a linear form on $Q(\mathcal{L})=H_{V^{\prime}, 0}^{2}\left(\rho_{P}\right)$ which is, by 53 , continuous with respect to $q_{-\mathcal{L}}$, one can applies Riesz's theorem so there exists a unique $u_{f} \in H_{V^{\prime}, 0}^{2}(\mathbb{R})$, such that for all $v \in \mathrm{H},\langle f, v\rangle_{\mathrm{H}}=q_{-\mathcal{L}}\left(u_{f}, v\right)$. Since, $u_{f}$ is clearly in $\mathcal{D}(\mathcal{L})$ by definition of the Friedrichs extension of $-\mathcal{L}$, we have $-\mathcal{L} u=f$.

Remark 6.8. We can diagonalize $\mathcal{L}$ by the same argument we used in Appendix $A$ to diagonalize $\mathcal{A}$ in $L_{0}^{2}\left(\rho_{P}\right)$.

We now state a result that could allow one to extract more regularity for $\mathcal{L}^{-1}$ by the help of an explicit form that uses Fredholm determinant theory for Hilbert-Schmidt operators, the reader can refer to GGK12.

Definition 6.9 (Fredholm determinant). Let $\mathcal{U}$ be a self-adjoint Hilbert-Schmidt operator, we denote the Fredholm determinant by $\operatorname{det}(I+\mathcal{U})$.

Theorem 6.10 (Determinant formula for $\left.\mathcal{L}^{-1}\right)$. For all $u \in \mathrm{H}$, such that $x \mapsto \frac{1}{\rho_{P}(x)} \int_{x}^{+\infty} u(t) \rho_{P}(t) d t$ is integrable at $+\infty$, we have:

$$
\begin{equation*}
\mathcal{L}^{-1} u=\mathcal{A}^{-1} u-\rho_{P}^{-1 / 2} \mathcal{R}\left[\rho_{P}^{1 / 2} \mathcal{A}^{-1} u\right] \tag{54}
\end{equation*}
$$

where $\mathcal{R}$ is the kernel operator defined for all $v \in L^{2}(\mathbb{R})$ by:

$$
\mathcal{R}[v](x)=\int_{\mathbb{R}} R(x, y) v(y) d y
$$

where

$$
R(x, y)=\frac{1}{\operatorname{det}(I+\mathcal{K})} \sum_{n \geqslant 0} \frac{1}{n!} \int_{\mathbb{R}^{n}} \operatorname{det}_{n+1}\left[\begin{array}{cc}
K(x, y) & K\left(x, \lambda_{b}\right) \\
K\left(\lambda_{a}, y\right) & K\left(\lambda_{a}, \lambda_{b}\right)
\end{array}\right]_{a, b=1 \ldots n} d \lambda_{1} \ldots d \lambda_{n}
$$

where $\mathcal{K}$ is the kernel operator defined for all $w \in L^{2}\left(\rho_{P}\right)$ by:

$$
\begin{equation*}
\mathcal{K}[v](x)=\int_{\mathbb{R}} K(x, y) w(y) d y \tag{55}
\end{equation*}
$$

with

$$
\begin{equation*}
K(x, y)=-2 P \rho_{P}(x) \rho_{P}(y) \ln \left|1-\frac{y}{x}\right| . \tag{56}
\end{equation*}
$$

Proof. Let $f \in \mathrm{H}$, there exists a unique $u \in \mathcal{D}(\mathcal{A})$ such that $\mathcal{A} u=f$. Since $\left(u^{\prime} \rho_{P}\right)^{\prime}=\rho_{P} \mathcal{A} u \in L^{2}(\mathbb{R})$, hence $u^{\prime} \rho_{P} \in H^{1}(\mathbb{R})$ so $u^{\prime}(x) \rho_{P}(x) \underset{|x| \rightarrow+\infty}{\longrightarrow} 0-\frac{\left(u^{\prime} \rho_{P}\right)^{\prime}}{\rho_{P}}=f$ hence

$$
\begin{equation*}
\left(\mathcal{A}^{-1} f\right)^{\prime}(x) \rho_{P}(x)=u^{\prime}(x) \rho_{P}(x)=\int_{x}^{+\infty} f(t) \rho_{P}(t) d t \tag{57}
\end{equation*}
$$

Using the fact that $\int_{\mathbb{R}} u(x) \rho_{P}(x) d x=0$, integrating again we get:

$$
u(x)=-\int_{x}^{+\infty} \frac{d s}{\rho_{P}(s)} \int_{s}^{+\infty} f(t) \rho_{P}(t) d t+C
$$

where $C=\int_{\mathbb{R}} \rho_{P}(x) d x \int_{x}^{+\infty} \frac{d s}{\rho_{P}(s)} \int_{s}^{+\infty} f(t) \rho_{P}(t) d t$. Now let $g \in \mathrm{H}$, there exists a unique $v \in \mathcal{D}(\mathcal{L})$, such that $-\mathcal{L} v=\mathcal{A} v+2 P \mathcal{W} v=g$ and then $v+2 P \mathcal{W} \mathcal{A}^{-1} v=\mathcal{A}^{-1} g$. using (57), we get:

$$
\mathcal{W} \mathcal{A}^{-1} v(x)=f_{\mathbb{R}} \frac{d s}{s-x} \int_{s}^{+\infty} d t v(t) \rho_{P}(t)
$$

By Sokhotski-Plejmel formula, we have:

$$
f_{\mathbb{R}} \frac{d s}{s-x} \int_{s}^{+\infty} d t v(t) \rho_{P}(t)=\lim _{M \rightarrow+\infty} \lim _{\varepsilon \rightarrow 0} \int_{-M}^{M} \frac{d s}{2}\left\{\frac{1}{x-s+\mathrm{i} \varepsilon}+\frac{1}{x-s-\mathrm{i} \varepsilon}\right\} \int_{s}^{+\infty} d t v(t) \rho_{P}(t)
$$

We then proceed to an integration by part:

$$
\begin{aligned}
f_{\mathbb{R}} \frac{d s}{s-x} \int_{s}^{+\infty} d t v(t) \rho_{P}(t)= & \lim _{M \rightarrow+\infty} \lim _{\varepsilon \rightarrow 0}\left[-\frac{\ln \left((x-s)^{2}+\varepsilon^{2}\right)}{2} \int_{s}^{+\infty} d t v(t) \rho_{P}(t)\right]_{-M}^{M} \\
& -\int_{\mathbb{R}} d s \ln |x-s| v(s) \rho_{P}(s)
\end{aligned}
$$

To conclude that $\mathcal{W} \mathcal{A}^{-1} v(x)=-\int_{\mathbb{R}} d s \ln |x-s| v(s) \rho_{P}(s)$, we just need to show that

$$
\ln (x) \int_{x}^{+\infty} d t v(t) \rho_{P}(t) \underset{|x| \rightarrow \infty}{\longrightarrow} 0
$$

which can be seen by Cauchy-Schwarz inequality:

$$
\left|\ln (x) \int_{x}^{+\infty} d t v(t) \rho_{P}(t)\right| \leqslant|\ln (x)|\|v\|_{L^{2}\left(\rho_{P}\right)} \cdot \rho_{P}(x)^{1 / 4}\left(\int_{\mathbb{R}} \rho_{P}(t)^{1 / 2} d t\right)^{1 / 2} .
$$

In this inequality, we used that $\rho_{P}$ is decreasing in a neighborhood of $+\infty$, hence

$$
\ln |x| \int_{s}^{+\infty} d t v(t) \rho_{P}(t) \underset{x \rightarrow+\infty}{\longrightarrow} 0
$$

the exact same argument allows us to conclude when $x$ goes to $-\infty$. Using the fact that $\int_{\mathbb{R}} v(t) \rho_{P}(t) d t=0$, we obtain the following equality:

$$
v-2 P \int_{\mathbb{R}} d s \ln |x-s| v(s) \rho_{P}(s)=\mathcal{A}^{-1} g:=h
$$

Now setting $\tilde{v}(t)=\rho_{P}^{1 / 2}(t) v(t)$ and $\tilde{h}=\rho_{P}^{1 / 2}(t) h(t)$, we obtain $\tilde{v}+\mathcal{K}[\tilde{v}]=\tilde{h}$ where $\mathcal{K}$ is defined in (55). Since its kernel (defined in (56)) belongs to $L^{2}\left(\mathbb{R}^{2}\right), \mathcal{K}$ is Hilbert-Schmidt. Hence by Fredholm determinant theory:

$$
\tilde{v}=\tilde{h}-\mathcal{R}[\tilde{h}]
$$

or $\mathcal{L}^{-1} g=\mathcal{A}^{-1} g-\rho_{P}^{-1 / 2} \mathcal{R}\left[\rho_{P}^{1 / 2} \mathcal{A}^{-1} g\right]$ as expected.

## 7 Regularity of the inverse of $\mathcal{L}$ and completion of the proof of Theorem 1.3

Since we have proven the central limit theorem for functions of the type $\mathcal{L} \phi$ with $\phi$ regular enough and satisfying vanishing asymptotic conditions at infinity, we exhibit a class of functions $f$ for which $\mathcal{L}^{-1} f$ is regular enough to satisfy conditions of Theorem 5.2. We define $\mathcal{T}$ the subset of H defined by

$$
\mathcal{T}:=\left\{f \in \mathcal{C}^{1}(\mathbb{R}), f(x)=\underset{|x| \rightarrow \infty}{O}\left(x^{-1}\right), f^{\prime}(x)=\underset{|x| \rightarrow \infty}{O}\left(x^{-2}\right), \int_{\mathbb{R}} f \rho_{P}=0\right\}
$$

Theorem 7.1. For all $f \in \mathcal{T}$, there exists a unique $u \in \mathcal{C}^{3}(\mathbb{R})$ such that $u^{\prime} \in H^{2}(\mathbb{R})$ with $u^{(3)}$ bounded wich verifies:

- $u^{\prime}(x)=\underset{|x| \rightarrow \infty}{O}\left(\frac{1}{x V^{\prime}(x)}\right)$
- $u^{\prime \prime}(x)=\underset{|x| \rightarrow \infty}{O}\left(\frac{1}{x V^{\prime}(x)}\right)$
- $u^{(3)}(x)=\underset{|x| \rightarrow \infty}{O}\left(\frac{1}{x}\right)$
such that $f=\mathcal{L} u$.
Proof. Let $f \in \mathcal{T} \subset H$, then since $-\mathcal{L}$ is bijective from $\mathcal{D}(\mathcal{L}) \rightarrow \mathrm{H}$, there exists a unique $u \in \mathcal{D}(\mathcal{L})$ such that $-\mathcal{L} u=f i e$ :

$$
\begin{equation*}
-u^{\prime \prime}-\frac{\rho_{P}^{\prime}}{\rho_{P}} u^{\prime}-2 P \mathcal{H}\left[u^{\prime} \rho_{P}\right]=f \tag{58}
\end{equation*}
$$

Hence we have

$$
\begin{equation*}
-\left(u^{\prime} \rho_{P}\right)^{\prime}=\rho_{P}\left(f+2 P \mathcal{H}\left[u^{\prime} \rho_{P}\right]\right) . \tag{59}
\end{equation*}
$$

Since $u \in \mathcal{D}(\mathcal{L}) \subset\left\{u \in H_{V^{\prime}, 0}^{2}(\mathbb{R}), \mathcal{A} u \in \mathcal{H}\right\}$, the functions $u^{\prime} \rho_{P}$ and its derivatives $\left(u^{\prime} \rho_{P}\right)^{\prime}=-\rho_{P} \mathcal{A} u$ and $\left(u^{\prime} \rho_{P}\right)^{\prime \prime}=-\frac{\rho_{P}^{\prime}}{\rho_{P}} \rho_{P}^{1 / 2} \cdot\left(\rho_{P}^{1 / 2} \mathcal{A} u\right)-\rho_{P}(\mathcal{A} u)^{\prime}$ are in $L^{2}(d x)$. In particular $u^{\prime} \rho_{P}$ goes to zero at infinity, and $\mathcal{H}\left[u^{\prime} \rho_{P}\right] \in H^{2}(\mathbb{R}) \subset C^{1}(\mathbb{R})$. So we can integrate 59$]$ on $[x,+\infty[$, since by Lemma 2.3 the right-hand side behaves like a $\underset{|x| \rightarrow \infty}{O}\left(\frac{\rho_{P}(x)}{x}\right)$, to get the following expression

$$
\begin{equation*}
u^{\prime}(x) \rho_{P}(x)=\int_{x}^{+\infty} \frac{\rho_{P}(t)}{\rho_{P}^{\prime}(t)}\left(f+2 P \mathcal{H}\left[u^{\prime} \rho_{P}\right]\right) \cdot \rho_{P}^{\prime}(t) d t \tag{60}
\end{equation*}
$$

From this expression, we can see that $u^{\prime} \in \mathcal{C}^{2}(\mathbb{R})$ so we just have to check the integrability condition at infinity and the fact that $u^{(3)}$ is bounded. After proceeding to an integration by parts, which is permitted by the previous argument, we obtain:

$$
\begin{equation*}
u^{\prime}(x)=-\frac{\rho_{P}(x)}{\rho_{P}^{\prime}(x)}\left(f(x)+2 P \mathcal{H}\left[u^{\prime} \rho_{P}\right](x)\right)-\frac{1}{\rho_{P}(x)} \int_{x}^{+\infty}\left(\frac{\rho_{P}(t)}{\rho_{P}^{\prime}(t)}\left(f+2 P \mathcal{H}\left[u^{\prime} \rho_{P}\right]\right)\right)^{\prime} \rho_{P}(t) d t \tag{61}
\end{equation*}
$$

and we define $R_{1}(x):=\frac{1}{\rho_{P}(x)} \int_{x}^{+\infty}\left(\frac{\rho_{P}(t)}{\rho_{P}^{\prime}(t)}\left(f+2 P \mathcal{H}\left[u^{\prime} \rho_{P}\right]\right)\right)^{\prime} \rho_{P}(t) d t$, we will need to show that $R_{1}$ is a remainder of order $O\left(\frac{1}{x V^{\prime}(x)^{2}}\right)$ at infinity. In this case we will have $u^{\prime}(x)=O\left(\frac{1}{x V^{\prime}(x)}\right)$ which will be useful for the following. If we reinject (61) in (58), we find:

$$
\begin{equation*}
u^{\prime \prime}=-\left(f+2 P \mathcal{H}\left[u^{\prime} \rho_{P}\right]\right)-\frac{\rho_{P}^{\prime}}{\rho_{P}}\left(-\frac{\rho_{P}}{\rho_{P}^{\prime}}\left(f+2 P \mathcal{H}\left[u^{\prime} \rho_{P}\right]\right)-R_{1}\right)=\frac{\rho_{P}^{\prime}}{\rho_{P}} R_{1} \tag{62}
\end{equation*}
$$

Hence

$$
u^{\prime \prime}(x)=\frac{\rho_{P}^{\prime}}{\rho_{P}^{2}}(x) \int_{x}^{+\infty} \rho_{P}(t) d t\{\underbrace{\left(\frac{\rho_{P}}{\rho_{P}^{\prime}}\right)^{\prime}(t)}_{\substack{o \\ t \rightarrow+\infty}} \underbrace{0}_{\substack{V^{\prime \prime}(t) \\ V_{t \rightarrow+\infty}}}\left(\frac{1}{t}\right)=\underbrace{\left[f+2 P \mathcal{H}\left[u^{\prime} \rho_{P}\right]\right](t)}_{\substack{o \\ t \rightarrow+\infty}}+\underbrace{\frac{\rho_{P}}{\rho_{P}^{\prime}}(t)}_{\left.=\underset{t \rightarrow+\infty}{V^{\prime}(t)}\right)} \underbrace{\left[f^{\prime}-2 P \mathcal{H}\left[\rho_{P} \mathcal{A} u\right]\right](t)}_{\left.t^{2}\right)}\} .
$$

The fact that $\mathcal{H}\left[\rho_{P} \mathcal{A} u\right](t)=\underset{t \rightarrow+\infty}{O}\left(t^{-2}\right)$ comes again from lemma 2.3 Finally we have that,

$$
\begin{aligned}
u^{(3)}(x)=\left(\frac{\rho_{P}^{\prime}}{\rho_{P}^{2}}\right)^{\prime}(x) \rho_{P}(x) R_{1}(x)- & \left(\frac{\rho_{P}^{\prime}}{\rho_{P}^{2}}\right)(x)\left(\frac{\rho_{P}}{\rho_{P}^{\prime}}\left(f+2 P \mathcal{H}\left[u^{\prime} \rho_{P}\right]\right)\right)^{\prime}(x) \rho_{P}(x) \\
& =\underbrace{\left.\left(\frac{\rho_{P}^{\prime \prime}}{\rho_{P}}-2 \frac{\rho_{P}^{\prime 2}}{\rho_{P}^{2}}\right)(x)^{2}\right)}_{\substack{o \\
\rho_{x \rightarrow+\infty}}}(x) \\
\rho_{1}(x) & -\underbrace{\left(\frac{\rho_{P}^{\prime}}{\rho_{P}}\right)(x)\left(\frac{\rho_{P}}{\rho_{P}^{\prime}}\left(f+2 P \mathcal{H}\left[u^{\prime} \rho_{P}\right]\right)\right)^{\prime}(x)}_{=\underset{x \rightarrow+\infty}{ }\left(\frac{V^{\prime \prime}(x)}{x V^{\prime}(x)}+x^{-2}\right)}
\end{aligned}
$$

The second term is $\underset{x \rightarrow+\infty}{O}\left(\frac{1}{x}\right)$ by the assumption that $\frac{V^{\prime \prime}}{V^{\prime}}(x)=\underset{|x| \rightarrow \infty}{O}$ (1). Hence, we just have to check that $R_{1}(x)=\underset{x \rightarrow+\infty}{O}\left(\frac{1}{x V^{\prime}(x)^{2}}\right)$ to establish that $u^{\prime}, u^{\prime \prime}, u^{(3)}$ are in $L^{2}(\mathbb{R})$. By a comparison argument, we control $R_{1}$ by controlling

$$
I_{1}(x):=\int_{x}^{+\infty} \frac{\rho_{P}(t)}{t V^{\prime}(t)} d t
$$

By integration by parts:

$$
\begin{align*}
I_{1}(x):= & -\underbrace{\frac{\rho_{P}(x)}{x V^{\prime}(x)} \frac{\rho_{P}}{\rho_{P}^{\prime}}(x)}-\int_{x}^{+\infty} \rho_{P}(t)\left(\frac{1}{t V^{\prime}} \frac{\rho_{P}}{\rho_{P}^{\prime}}\right)^{\prime}(t) d t \\
& =\underset{x \rightarrow+\infty}{O}\left(\frac{\rho_{P}(x)}{x V^{\prime}(x)^{2}}\right) \\
& \left.=\frac{\rho_{P}(x)}{x V^{\prime}(x)^{2}}\right)-\int_{x}^{+\infty} \rho_{P}(t) d t\left\{-\frac{1}{t^{2} V^{\prime}(t)} \frac{\rho_{P}}{\rho_{P}^{\prime}}(t)-\frac{V^{\prime \prime}(t)}{t V^{\prime}(t)^{2}} \frac{\rho_{P}}{\rho_{P}^{\prime}}+\frac{1}{t V^{\prime}(t)}\left(\frac{\rho_{P}}{\rho_{P}^{\prime}}\right)^{\prime}(t)\right\} \tag{63}
\end{align*}
$$

By the same argument as before, the last integral is of the form $\int_{x}^{+\infty} \underset{t \rightarrow+\infty}{O}\left(\frac{\rho_{P}(t)}{t V^{\prime}(t)^{2}}\right) d t$ so if $I_{2}(x):=\int_{x}^{+\infty} \frac{\rho_{P}(t)}{t V^{\prime}(t)^{2}} d t=\underset{x \rightarrow+\infty}{O}\left(\frac{\rho_{P}(x)}{x V^{\prime}(x)^{2}}\right)$ then so is $I_{1}$. By integration by parts, we obtain:

$$
\begin{aligned}
& I_{2}(x)=\rho_{P}(x) \frac{1}{x V^{\prime}(x)^{2}} \frac{\rho_{P}}{\rho_{P}^{\prime}}(x)-\int_{x}^{+\infty} \rho_{P}(t) d t\left\{\left(\frac{\rho_{P}}{\rho_{P}^{\prime}}\right)^{\prime}(t) \frac{1}{t V^{\prime}(t)^{2}}-\frac{\rho_{P}}{\rho_{P}^{\prime}}(t)\left(\frac{1}{t^{2} V^{\prime}(t)^{2}}+\frac{2 V^{\prime \prime}(t)}{t V^{\prime}(t)^{3}}\right)\right\} \\
&=\underset{x \rightarrow+\infty}{O}\left(\frac{\rho_{P}(x)}{x V^{\prime}(x)^{2}}\right)+\int_{x}^{+\infty} \underset{t \rightarrow+\infty}{O}\left(\frac{\rho_{P}(t)}{t V^{\prime}(t)^{3}}\right) d t
\end{aligned}
$$

The last integral is a $\underset{x \rightarrow+\infty}{O}\left(\frac{\rho_{P}(x)}{x V^{\prime}(x)^{2}}\right)$ because, again, by integration by part:

$$
\int_{x}^{+\infty} \frac{\rho_{P}(t)}{t V^{\prime}(t)^{3}} d t=\rho_{P}(x) \frac{1}{x V^{\prime}(x)^{3}} \frac{\rho_{P}}{\rho_{P}^{\prime}}(x)-\int_{x}^{+\infty} \underset{t \rightarrow+\infty}{O}\left(\frac{\rho_{P}(t)}{t V^{\prime}(t)^{4}}\right)
$$

and finally

$$
\int_{x}^{+\infty} \frac{\rho_{P}(t)}{t V^{\prime}(t)^{4}} d t \leqslant \frac{\rho_{P}(x)}{x V^{\prime}(x)^{2}} \int_{x}^{+\infty} \frac{d t}{V^{\prime}(t)^{2}}=\underset{x \rightarrow+\infty}{O}\left(\frac{\rho_{P}(x)}{x V^{\prime}(x)^{2}}\right)
$$

In the final step, we used the fact that $x \mapsto \frac{\rho_{P}(x)}{x V^{\prime}(x)}$ is decreasing in a neighborhood of $+\infty$ (which can be checked by differentiating) and that $x \mapsto \frac{1}{V^{\prime}(x)^{2}}$ is integrable at $\infty$ by assumption iii). Hence $R_{1}(x)=\underset{x \rightarrow+\infty}{O}\left(\frac{1}{x V^{\prime}(x)^{2}}\right)$ (the exact same result can be shown at $-\infty$ ), which leads to the fact

$$
\begin{equation*}
u^{\prime}(x)=\underset{|x| \rightarrow+\infty}{O}\left(\frac{1}{x V^{\prime}(x)}\right), \quad u^{\prime \prime}(x)=\underset{|x| \rightarrow+\infty}{O}\left(\frac{1}{x V^{\prime}(x)}\right) \quad \text { and } \quad u^{(3)}(x)=\underset{|x| \rightarrow+\infty}{O}\left(\frac{1}{x}\right) \tag{64}
\end{equation*}
$$

which establishes that these functions are in $L^{2}$ in a neighborhood of $\infty$. Since we already showed that $u \in \mathcal{C}^{3}(\mathbb{R}) \subset H_{\text {loc }}^{3}(\mathbb{R})$, it establishes that $u \in H^{3}(\mathbb{R}) \cap \mathcal{C}^{3}(\mathbb{R})$ with $u^{(3)}$ bounded. To complete the proof we just have to show that $\left(u^{\prime}\right)^{2} V^{(3)}, u^{\prime} u^{\prime \prime} V^{\prime \prime},\left(u^{\prime}\right)^{2} V^{\prime \prime}$ and $u^{\prime} V^{\prime}$ are bounded which is easily checked by (64) and Assumption 1.1 iv).

## A Appendix: proof of Theorem 6.2

In order to analyze $\mathcal{A}$, we let, for $u \in L^{2}(d x)$,

$$
\mathcal{S} u:=\rho_{P}^{1 / 2} \mathcal{A} \rho_{P}^{-1 / 2} u
$$

Note that $u \in\left(L^{2}(d x),\|\cdot\|_{L^{2}(d x)}\right) \mapsto \rho_{P}^{-1 / 2} u \in\left(L^{2}\left(\rho_{P}\right),\|\cdot\|_{L^{2}\left(\rho_{P}\right)}\right)$ is an isometry. It turns out that it will be easier to study first the operator $\mathcal{S}$ in order to get the spectal properties of $\mathcal{A}$.

Proposition A.1. The operator $\mathcal{S}$ is a Schrödinger operator: it admits the following expression for all $u \in C_{c}^{2}(\mathbb{R}): \mathcal{S} u=-u^{\prime \prime}+w_{P} u$ with

$$
w_{P}=\frac{1}{2}\left(\frac{1}{2} V^{\prime 2}-V^{\prime \prime}+2 P V^{\prime} \mathcal{H}\left[\rho_{P}\right]-2 P \mathcal{H}\left[\rho_{P}^{\prime}\right]+2 P^{2} \mathcal{H}\left[\rho_{P}\right]^{2}\right)=\frac{1}{2}\left[\left(\ln \rho_{P}\right)^{\prime \prime}+\frac{1}{2}\left(\ln \rho_{P}\right)^{\prime 2}\right]
$$

Furthermore, $w_{P}$ is continuous and we have $w_{P}(x) \underset{\infty}{\sim} \frac{V^{\prime}(x)^{2}}{4} \underset{|x| \rightarrow \infty}{\longrightarrow}+\infty$.
Proof. We compute directly

$$
\begin{aligned}
\frac{\left(\rho_{P}\left(\rho_{P}^{-1 / 2} u\right)^{\prime}\right)^{\prime}}{\rho_{P}} & =\left(\rho_{P}^{-1 / 2} u\right)^{\prime \prime}+\frac{\rho_{P}^{\prime}}{\rho_{P}}\left(\rho_{P}^{-1 / 2} u\right)^{\prime} \\
& =\left(\rho_{P}^{-1 / 2} u^{\prime}-\frac{1}{2} \rho_{P}^{-3 / 2} \rho_{P}^{\prime} u\right)^{\prime}+\rho_{P}^{\prime} \rho_{P}^{-3 / 2} u^{\prime}-\frac{1}{2} \rho_{P}^{-5 / 2}\left(\rho_{P}^{\prime}\right)^{2} u \\
& =\rho_{P}^{-1 / 2} u^{\prime \prime}+\frac{1}{4} \rho_{P}^{-5 / 2}\left(\rho_{P}^{\prime}\right)^{2} u-\frac{1}{2} \rho_{P}^{-3 / 2} \rho_{P}^{\prime \prime} u \\
& =\rho_{P}^{-1 / 2}\left[u^{\prime \prime}+\frac{1}{4} \rho_{P}^{-2}\left(\rho_{P}^{\prime}\right)^{2} u-\frac{1}{2} \rho_{P}^{-1} \rho_{P}^{\prime \prime} u\right] \\
& =\rho_{P}^{-1 / 2}\left(u^{\prime \prime}-\frac{1}{2}\left[\left(\frac{\rho_{P}^{\prime \prime}}{\rho_{P}}\right)-\frac{1}{2}\left(\frac{\rho_{P}^{\prime}}{\rho_{P}}\right)^{2}\right] u\right) \\
& =\rho_{P}^{-1 / 2}\left(u^{\prime \prime}-\frac{1}{2}\left[\left(\ln \rho_{P}\right)^{\prime \prime}+\frac{1}{2}\left(\ln \rho_{P}\right)^{\prime 2}\right] u\right)=\rho_{P}^{-1 / 2}\left(u^{\prime \prime}-w_{P} u\right)
\end{aligned}
$$

Now, using Lemma 2.2 we have

$$
w_{P}=\frac{1}{2}\left(\frac{1}{2} V^{\prime 2}-V^{\prime \prime}+2 P V^{\prime} \mathcal{H}\left[\rho_{P}\right]-2 P \mathcal{H}\left[\rho_{P}^{\prime}\right]+2 P^{2} \mathcal{H}\left[\rho_{P}\right]^{2}\right)
$$

Notice that $\mathcal{H}\left[\rho_{P}^{\prime}\right]$ and $\mathcal{H}\left[\rho_{P}\right]$ are bounded since they belong to $H^{1}(\mathbb{R})$, as we showed in Lemma 2.2 that $\rho_{P}$ is $H^{2}(\mathbb{R})$. Along with Assumption 1.1 iiii) and Lemma 2.3. we deduce $w_{P}(x) \underset{\infty}{\sim} \frac{1}{4} V^{\prime 2}(x)$.

Remark A.2. Note that the function $w_{P}$ need not be positive on $\mathbb{R}$. In fact, neglecting the terms involving the Hilbert transforms of $\rho_{P}$ and $\rho_{P}^{\prime}$, $w_{P}$ would only be positive outside of a compact set. However, using the positivity of $\mathcal{A}$, which will be shown further in the article, we can show that the operator $-u^{\prime \prime}+w_{P} u$ is itself positive on $L^{2}(d x)$. It can also be checked that, by integration by partss, $\mathcal{S}$ is symmetric on $\mathcal{C}_{c}^{2}(\mathbb{R})$ with the inner product of $L^{2}(d x)$.

We now introduce an extension of $\mathcal{S}$ by defining its associated bilinear form.
Definition A. 3 (Quadratic form associated to $\mathcal{S}$ ).
Let $\alpha \geqslant 0$ such that $w_{P}+\alpha \geqslant 1$. We define the quadratic form associated to $\mathcal{S}+\alpha I$, defined for all $u \in \mathcal{C}_{c}^{2}(\mathbb{R})$ by

$$
q_{\alpha}(u, u):=\int_{\mathbb{R}} u^{\prime 2} d x+\int_{\mathbb{R}} u^{2}\left(w_{P}+\alpha\right) d x
$$

This quadratic form can be extended to a larger domain denoted by $Q(\mathcal{S}+\alpha I)$, called the form domain of the operator $\mathcal{S}+\alpha I$. By the theory of Schrödinger operators, it is well-known (see Dav96] [Theorem 8.2.1]) that such a domain is given by

$$
Q(\mathcal{S}+\alpha I)=\left\{u \in H^{1}(\mathbb{R}), u\left(w_{P}+\alpha\right)^{1 / 2} \in L^{2}(\mathbb{R})\right\}=\left\{u \in H^{1}(\mathbb{R}), u V^{\prime} \in L^{2}(\mathbb{R})\right\}=: H_{V^{\prime}}^{1}(\mathbb{R})
$$

The space $H_{V^{\prime}}^{1}(\mathbb{R})$ can be seen to be the completion under the norm $q_{\alpha}$ of $\mathcal{C}_{c}^{\infty}$. Now that the quadratic form associated to $\mathcal{S}+\alpha I$ has been extended to its form domain, it is possible to go back to the operator and extend it by its Friedrichs extension.

Theorem A. 4 (Friedrichs extension of $\mathcal{S}+\alpha I$ ).
There exists an extension $(S+\alpha I)_{F}$ of the operator $\mathcal{S}+\alpha I$, called the Friedrichs extension of $\mathcal{S}+\alpha I$ defined on $H_{V^{\prime}}^{2}(\mathbb{R}):=\left\{u \in H^{2}(\mathbb{R}), u V^{\prime} \in L^{2}(\mathbb{R})\right\}$.

Proof. We denote

$$
\begin{aligned}
& D\left((S+\alpha I)_{F}\right) \\
& \quad=\left\{v \in H_{V^{\prime}}^{1}(\mathbb{R}), u \in H_{V^{\prime}}^{1}(\mathbb{R}) \longmapsto q_{\alpha}(u, v) \text { can be extended to a continuous linear form on } L^{2}(\mathbb{R})\right\}
\end{aligned}
$$

If $v \in D\left((S+\alpha I)_{F}\right)$, by Riesz's theorem there exists a unique $f_{v} \in L^{2}(\mathbb{R})$ such that $q_{\alpha}(u, v)=$ $\left\langle u, f_{v}\right\rangle_{L^{2}(d x)}$ holds for all $u \in L^{2}(\mathbb{R})$ and we can set $(S+\alpha I)_{F} v:=f_{v}$. Note that it is indeed a way of extending $\mathcal{S}+\alpha I$ since for all $u, v \in \mathcal{C}_{c}^{2}(\mathbb{R}), q_{\alpha}(u, v)=\langle u,(\mathcal{S}+\alpha I) v\rangle_{L^{2}(d x)}$.

We want to show that $D\left((S+\alpha I)_{F}\right)=H_{V^{\prime}}^{2}(\mathbb{R})$. Let $f \in D\left((S+\alpha I)_{F}\right)$ and $g:=(S+\alpha I)_{F} f \in L^{2}(\mathbb{R})$. By definition of $q_{\alpha}$, for all $u \in \mathcal{C}_{c}^{2}(\mathbb{R})$ :

$$
\int_{\mathbb{R}} g u d x=\int_{\mathbb{R}} f^{\prime} u^{\prime} d x+\int_{\mathbb{R}}\left(w_{P}+\alpha\right) f u d x=-\int_{\mathbb{R}} f u^{\prime \prime} d x+\int_{\mathbb{R}}\left(w_{P}+\alpha\right) f u d x
$$

Therefore in the sense of distributions, we get $-f^{\prime \prime}=g-\left(w_{P}+\alpha\right) f$ which is a function in $L^{2}(\mathbb{R})$, hence $f \in H_{V^{\prime}}^{2}(\mathbb{R})$. Conversely, if $f \in H_{V^{\prime}}^{2}(\mathbb{R})$, it is possible to extend $u \mapsto q_{\alpha}(f, u)$ to a continuous linear form on $L^{2}(\mathbb{R})$ by

$$
u \mapsto-\int_{\mathbb{R}} u f^{\prime \prime} d x+\int_{\mathbb{R}} u f\left(w_{P}+\alpha\right) d x
$$

which concludes the fact that $D\left((S+\alpha I)_{F}\right)=H_{V^{\prime}}^{2}(\mathbb{R})$.
In the following, we will deal only with $(\mathcal{S}+\alpha I)_{F}: H_{V^{\prime}}^{2}(\mathbb{R}) \longrightarrow L^{2}(\mathbb{R})$ and denote it $(\mathcal{S}+\alpha I)$.
Remark A.5. Note that in the previous proof, the application of Riesz's theorem doesn't allow to say that $(\mathcal{S}+\alpha I): v \in\left(H_{V^{\prime}}^{2}(\mathbb{R}),\|\cdot\|_{q_{\alpha}}\right) \mapsto f_{v} \in\left(L^{2}(\mathbb{R}),\|\cdot\|_{L^{2}(d x)}\right)$, where $\|\cdot\|_{q_{\alpha}}$ stands for the norm associated to the bilinear positive definite form $q_{\alpha}$, is continuous. It can be seen by the fact that $v \in\left(D(S+\alpha I),\|\cdot\|_{q_{\alpha}}\right) \mapsto q(., v) \in\left(L^{2}(\mathbb{R})^{\prime},\|\cdot\|_{L^{2}(d x)^{\prime}}\right)$, where $L^{2}(\mathbb{R})^{\prime}$ stands for the topological dual of $L^{2}(\mathbb{R})$ equipped with its usual norm, is not continuous. Indeed the $\|\cdot\|_{q_{\alpha}}$ norm doesn't control the second derivative of $v$ and hence doesn't provide any module of continuity for the $L^{2}(\mathbb{R})$-extended linear form $q(., v)$.

Theorem A. 6 (Inversion of $\mathcal{S}+\alpha I$ ).
For every $f \in L^{2}(\mathbb{R})$, there exists a unique $u \in H_{V^{\prime}}^{2}(\mathbb{R})$ such that $(\mathcal{S}+\alpha I) u=f$. Furthermore, the map $(S+\alpha I)^{-1}$ is continuous from $\left(L^{2}(\mathbb{R}),\|\cdot\|_{L^{2}(d x)}\right)$ to $\left(H_{V^{\prime}}^{2}(\mathbb{R}),\|\cdot\|_{q_{\alpha}}\right)$.

Proof. Let $f \in L^{2}(\mathbb{R})$, the map $u \longmapsto\langle u, f\rangle_{L^{2}(d x)}$ is continuous on $\left(H_{V^{\prime}}^{1}(\mathbb{R}),\|\cdot\|_{q_{\alpha}}\right)$ which is a Hilbert space. Therefore by Riesz's theorem, there exists a unique $v_{f} \in H_{V^{\prime}}^{1}(\mathbb{R})$ such that for all $u \in H_{V^{\prime}}^{1}(\mathbb{R})$, $\langle f, u\rangle_{L^{2}(d x)}=q_{\alpha}\left(v_{f}, u\right)$ from which we deduce that, in the sense of distributions, $f=-v_{f}^{\prime \prime}+\left(w_{P}+\alpha\right) v_{f}$ which implies that $v_{f} \in H_{V^{\prime}}^{2}(\mathbb{R})$. Since $v_{f} \in D(S+\alpha I)$, we have then for all $u \in L^{2}(\mathbb{R}),\langle f, u\rangle_{L^{2}(d x)}=$ $q_{\alpha}\left(v_{f}, u\right)=\left\langle(\mathcal{S}+\alpha) v_{f}, u\right\rangle_{L^{2}(d x)}$, hence $(\mathcal{S}+\alpha I) v_{f}=f$. Finally, by Riesz's theorem, $f \in L^{2}(\mathbb{R}) \mapsto v_{f} \in$ $H_{V^{\prime}}^{1}(\mathbb{R})$ is continuous hence so is $(\mathcal{S}+\alpha I)^{-1}$.

Remark A.7. It would be tempting to use Banach's isomorphism theorem to say that since $(\mathcal{S}+\alpha I)^{-1}$ is bijective and continuous, so must be $\mathcal{S}+\alpha I$. But since $\left(H_{V^{\prime}}^{2}(\mathbb{R}),\|\cdot\|_{q_{\alpha}}\right)$ is not a Banach space (it's not closed in $H_{V^{\prime}}^{1}(\mathbb{R})$ ) we can't apply it.

We are now able to diagonalize the resolvent of $\mathcal{S}$.
Theorem A. 8 (Diagonalization of $\left.(\mathcal{S}+\alpha I)^{-1}\right)$.
There exists a complete orthonormal set $\left(\psi_{n}\right)_{n \geqslant 0}$ of $L^{2}(d x)$ (meaning that

$$
\overline{\operatorname{span}\left\{\psi_{n}, n \geqslant 0\right\}} \|^{\|\cdot\|_{L^{2}(d x)}}=L^{2}(d x)
$$

and $\left.\left\langle\psi_{i}, \psi_{j}\right\rangle_{L^{2}(d x)}=\delta_{i, j}\right)$, where each $\psi_{n} \in H_{V^{\prime}}^{2}$ and $\left(\mu_{n}(\alpha)\right)_{n \geqslant 0} \in[0,1]^{\mathbb{N}}$ with $\mu_{n}(\alpha) \underset{N \rightarrow \infty}{\longrightarrow} 0$ such that $(\mathcal{S}+\alpha I)^{-1} \psi_{n}=\mu_{n}(\alpha) \psi_{n}$ for all $n \geqslant 0$. We also have $\left\|(\mathcal{S}+\alpha I)^{-1}\right\|_{\mathcal{L}\left(L^{2}(d x)\right)} \leqslant 1$.
Proof. By Proposition A.1, $w_{P}+\alpha$ is continuous and goes to infinity at infinity. By Rellich criterion RS78 [see Theorem XIII.65], the unit ball of $H_{V^{\prime}}^{2}(\mathbb{R})$, ie the set

$$
\left\{u \in H_{V^{\prime}}^{2}(\mathbb{R}), \int_{\mathbb{R}} u^{\prime 2}+\int_{\mathbb{R}}\left(w_{P}+\alpha\right) u^{2} \leqslant 1\right\}
$$

considered as a subset of $L^{2}(\mathbb{R})$ is relatively compact in $\left(L^{2}(d x),\|\cdot\|_{L^{2}(d x)}\right)$. Hence, we can conclude that the injection $\iota:\left(H_{V^{\prime}}^{2}(\mathbb{R}),\|\cdot\|_{q_{\alpha}}\right) \longrightarrow\left(L^{2}(d x),\|\cdot\|_{L^{2}(d x)}\right)$ is a compact operator. Since $(S+\alpha I)^{-1}$ : $\left(L^{2}(d x),\|\cdot\|_{L^{2}(d x)}\right) \longrightarrow\left(H_{V^{\prime}}^{2}(\mathbb{R}),\|\cdot\|_{q_{\alpha}}\right)$ is continuous then $(\mathcal{S}+\alpha I)^{-1}$ is compact from $\left(L^{2}(d x),\|\cdot\|_{L^{2}(d x)}\right)$ to itself. The fact that $(\mathcal{S}+\alpha I)^{-1}$ is self-adjoint and positive allows us to apply the spectral theorem to obtain $\left(\mu_{n}(\alpha)\right)_{n \geqslant 0}$ positive eigenvalues verifying $\mu_{n}(\alpha) \underset{N \rightarrow \infty}{\longrightarrow} 0$ by compactness and a Hilbertian basis $\left(\psi_{n}\right)_{n \geqslant 0} \in L^{2}(\mathbb{R})^{\mathbb{N}}$, such that for all $n \geqslant 0,(\mathcal{S}+\alpha I)^{-1} \psi_{n}=\mu_{n}(\alpha) \psi_{n}$. It is then easy to see that for all $n, \psi_{n} \in H_{V^{\prime}}^{2}(\mathbb{R})$ since they belong to the range of $(\mathcal{S}+\alpha I)^{-1}$. Finally, since for all $\phi \in L^{2}(\mathbb{R})$, $\langle(\mathcal{S}+\alpha I) \phi, \phi\rangle_{L^{2}(d x)} \geqslant\|\phi\|_{L^{2}(d x)}^{2}$, the spectrum of $(\mathcal{S}+\alpha I)^{-1}$ is contained in $[0,1]$. It allows us to conclude that $\left\|\mid(\mathcal{S}+\alpha I)^{-1}\right\| \|_{L^{2}(d x)} \leqslant 1$.

It is now straightforward to see how to extend $\mathcal{A}=\rho_{P}^{-1 / 2} \mathcal{S} \rho_{P}^{1 / 2}$ on $H_{V^{\prime}}^{2}\left(\rho_{P}\right):=\rho_{P}^{-1 / 2} H_{V^{\prime}}^{2}(\mathbb{R})$ equipped with the norm $\|\cdot\|_{q_{\alpha}, \rho_{P}}$ to $\left(L^{2}\left(\rho_{P}\right),\|\cdot\|_{L^{2}\left(\rho_{P}\right)}\right)$. The norm $\|\cdot\|_{q_{\alpha}, \rho_{P}}$ is defined for all $u \in H^{2}\left(\rho_{P}\right)$ by

$$
\|u\|_{q_{\alpha}, \rho_{P}}=\int_{\mathbb{R}} u^{\prime 2} \rho_{P} d x+\int_{\mathbb{R}} u^{2}\left(w_{P}+\alpha\right) \rho_{P} d x
$$

It is easy to see that $(\mathcal{A}+\alpha I)^{-1}$ is continuous. We stress the fact that $H_{V^{\prime}}^{2}\left(\rho_{P}\right) \neq\left\{u \in H^{2}\left(\rho_{P}\right), u V^{\prime} \in\right.$ $\left.L^{2}\left(\rho_{P}\right)\right\}$. Indeed if $u \in H_{V^{\prime}}^{2}(\mathbb{R})$, even though $u \rho_{P}^{-1 / 2}$ and its derivative belong to $L^{2}\left(\rho_{P}\right),\left(\rho_{P}^{-1 / 2} u\right)^{\prime \prime} \notin$ $L^{2}\left(\rho_{P}\right)$. The reader can check that for such a function to be in $L^{2}\left(\rho_{P}\right)$, it would be necessary to have that $u^{2} V^{\prime 4} \in L^{2}(d x)$ which is not the case.

Remark A.9. The kernel of $\mathcal{A}$ is generated by the function $\widetilde{1}$. Indeed if $\phi \in H_{V^{\prime}}^{2}\left(\rho_{P}\right)$ is in the kernel of $\mathcal{A}$ then

$$
0=-\frac{\left(\phi^{\prime} \rho_{P}\right)^{\prime}}{\rho_{P}} \Rightarrow \exists c \in \mathbb{R}, \phi^{\prime}=\frac{c}{\rho_{P}}
$$

But since $\phi^{\prime}$ is in $L^{2}\left(\rho_{P}\right)$ then $c=0$ which implies that $\phi$ is constant. We must restrict $\mathcal{A}$ to the orthogonal of $\operatorname{Ker} \mathcal{A}$ with respect to the inner product of $L^{2}\left(\rho_{P}\right)$, ie

$$
H_{V^{\prime}, 0}^{2}\left(\rho_{P}\right):=\left\{u \in H_{V^{\prime}}^{2}\left(\rho_{P}\right) \mid \int_{\mathbb{R}} u \rho_{P}=0\right\}
$$

Doing so makes $\mathcal{A}$ injective.
Before inverting $\mathcal{A}$, we need the following lemma:
Lemma A.10. The following equality holds

$$
(\mathcal{A}+\alpha I)\left(H_{V^{\prime}, 0}^{2}\left(\rho_{P}\right)\right)=L_{0}^{2}\left(\rho_{P}\right):=\left\{u \in L^{2}\left(\rho_{P}\right), \int_{\mathbb{R}} u \rho_{P} d x=0\right\}
$$

Proof. Let $\phi=\widetilde{c}$ for $c \in \mathbb{R},(\mathcal{A}+\alpha I) \phi=\widetilde{\alpha c}$ then $(\mathcal{A}+\alpha I)(\mathbb{R} . \widetilde{1})=\mathbb{R} \tilde{1}$. Hence since $\mathcal{A}+\alpha I$ is self-adjoint with respect to the inner product of $L^{2}\left(\rho_{P}\right)$ and that $\mathbb{R} \widetilde{1}$ is stable by $\mathcal{A}+\alpha I$, then $(\mathcal{A}+\alpha I)\left((\mathbb{R} . \widetilde{1})^{\perp} \cap\right.$ $\left.H_{V^{\prime}}^{2}\left(\rho_{P}\right)\right) \subset(\mathbb{R} . \widetilde{1})^{\perp}$. For the converse, let $u \in(\mathbb{R} . \widetilde{1})^{\perp}$, since $\mathcal{A}+\alpha I$ is bijective, there exists $v \in H_{V^{\prime}}^{2}\left(\rho_{P}\right)$ such that $u=(\mathcal{A}+\alpha I) v$. For all $w \in \mathbb{R} . \widetilde{1}$,

$$
0=\langle u, w\rangle_{L^{2}\left(\rho_{P}\right)}=\langle(\mathcal{A}+\alpha I) v, w\rangle_{L^{2}\left(\rho_{P}\right)}=\langle v,(\mathcal{A}+\alpha I) w\rangle_{L^{2}\left(\rho_{P}\right)}
$$

Hence $v \in[(\mathcal{A}+\alpha I)(\mathbb{R} \widetilde{1})]^{\perp}=\mathbb{R} \widetilde{1}^{\perp}$ and so $(\mathbb{R} . \widetilde{1})^{\perp} \subset(\mathcal{A}+\alpha I)\left((\mathbb{R} . \widetilde{1})^{\perp}\right)$.
It is easy to see that $L_{0}^{2}\left(\rho_{P}\right)$ is a closed subset of $L^{2}\left(\rho_{P}\right)$ as it is the kernel of the linear form $\phi \in L^{2}\left(\rho_{P}\right) \mapsto\langle\phi, \widetilde{1}\rangle_{L^{2}\left(\rho_{P}\right)}$, making it a Hilbert space.
Proposition 1.11 (Diagonalization and invertibility of $\mathcal{A}$ ). There exists a complete orthonormal set of $\left(L_{0}^{2}\left(\rho_{P}\right),\langle., .\rangle_{L^{2}\left(\rho_{P}\right)}\right),\left(\phi_{n}\right)_{n \in \mathbb{N}} \in H_{V^{\prime}, 0}^{2}\left(\rho_{P}\right)^{\mathbb{N}}$ such that $\mathcal{A} \phi_{n}=\lambda_{n} \phi_{n}$ (meaning that

$$
\overline{\operatorname{span}\left\{\phi_{n}, n \geqslant 0\right\}}\|\cdot\|_{L^{2}\left(\rho_{P}\right)}=L_{0}^{2}\left(\rho_{P}\right)
$$

and $\left.\left\langle\phi_{i}, \phi_{j}\right\rangle_{L^{2}\left(\rho_{P}\right)}=\delta_{i, j}\right)$. Furthermore, $\mathcal{A}: H_{V^{\prime}, 0}^{2}\left(\rho_{P}\right) \longrightarrow L_{0}^{2}\left(\rho_{P}\right):=\left\{u \in L^{2}\left(\rho_{P}\right), \int_{\mathbb{R}} u \rho_{P} d x=0\right\}$ is bijective, $\mathcal{A}^{-1}$ is continuous when considered as an operator of $L_{0}^{2}\left(\rho_{P}\right)$.

Proof. Since $(\mathcal{S}+\alpha I)^{-1}$ considered as an operator of $L^{2}(d x)$, is compact so is $(\mathcal{A}+\alpha I)^{-1}$ on $L^{2}\left(\rho_{P}\right)$ and since $\mathcal{A}$ is self-adjoint, by the spectral theorem, $(\mathcal{A}+\alpha I)^{-1}$ is diagonalizable. With the notations of Theorem A.8 $(\mathcal{A}+\alpha I)^{-1}$ has eigenvalues $\left(\mu_{n}(\alpha)\right)_{n \geqslant 0}$ and corresponding eigenfunctions $\phi_{n}=\rho_{P}^{-1 / 2} \psi_{n} \in$ $H_{V^{\prime}}^{2}\left(\rho_{P}\right)$. Hence for all $n \in \mathbb{N}, \mathcal{A} \phi_{n}=\lambda_{n} \phi_{n}$ with $\lambda_{n}:=\left(\frac{1}{\mu_{n}(\alpha)}-\alpha\right)$. Now,

$$
\lambda_{n}\left\|\phi_{n}\right\|_{L^{2}\left(\rho_{P}\right)}=\int_{\mathbb{R}}\left(\mathcal{A} \phi_{n}\right) \phi_{n} \rho_{P} d x=-\int_{\mathbb{R}}\left(\rho_{P} \phi_{n}^{\prime}\right)^{\prime} \phi_{n}=\int_{\mathbb{R}} \phi_{n}^{\prime 2} \rho_{P} \geqslant 0
$$

Furthermore, the kernel of $\mathcal{A}$ is $\mathbb{R} . \widetilde{1}$, thus the spectrum of $\mathcal{A}$ restricted to $H_{V^{\prime}, 0}^{2}\left(\rho_{P}\right)$ is positive. But since $(\mathcal{A}+\alpha I)^{-1}$ is a compact operator of $L^{2}\left(\rho_{P}\right)$ and that $(\mathcal{A}+\alpha I)$ maps $\mathbb{R} . \widetilde{1}^{\perp}$ to $\mathbb{R}^{1} \widetilde{1}^{\perp}$ with respect to the inner product of $L^{2}\left(\rho_{P}\right)$ (see lemma A.10, then $(\mathcal{A}+\alpha I)^{-1}$ is compact as an operator from $L_{0}^{2}\left(\rho_{P}\right)$ to itself. By Fredholm alternative, for every $\lambda \in \mathbb{R} \lambda \neq 0$, either $(\mathcal{A}+\alpha I)^{-1}-\lambda I$ is bijective
either $\lambda \in S p\left((\mathcal{A}+\alpha I)^{-1}\right)$. These conditions are equivalent to: either $\mathcal{A}+\left(\alpha-\frac{1}{\lambda}\right) I$ is bijective as an operator from $H_{V^{\prime}, 0}^{2}\left(\rho_{P}\right)$ to $L_{0}^{2}\left(\rho_{P}\right)$, either $-\alpha+\frac{1}{\lambda} \in S p(\mathcal{A})$. If we set $\lambda=\frac{1}{\alpha}$ then either $\mathcal{A}$ is bijective either $0 \in S p(\mathcal{A})$, since the latter is wrong then $\mathcal{A}: H_{V^{\prime}, 0}^{2}\left(\rho_{P}\right) \rightarrow L_{0}^{2}\left(\rho_{P}\right)$ is bijective. The spectrum of $\mathcal{A}$ is $\left(\frac{1}{\mu_{n}(\alpha)}-\alpha\right)_{n \geqslant 0} \subset\left(\lambda_{1},+\infty\right) \subset(0,+\infty)$, where $\lambda_{1}$ is the smallest eigenvalue. Hence, we deduce $\left\|\mid \mathcal{A}^{-1}\right\|_{\mathcal{L}\left(L^{2}\left(\rho_{P}\right)\right)} \leqslant \lambda_{1}^{-1}$.

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