# Minimax estimation of the Wigner function in quantum homodyne tomography with ideal detectors 

## Citation for published version (APA):

Guta, M. I., \& Artiles, L. M. (2006). Minimax estimation of the Wigner function in quantum homodyne tomography with ideal detectors. (Report Eurandom; Vol. 2006030). Eurandom.

## Document status and date:

Published: 01/01/2006

## Document Version:

Publisher's PDF, also known as Version of Record (includes final page, issue and volume numbers)

## Please check the document version of this publication:

- A submitted manuscript is the version of the article upon submission and before peer-review. There can be important differences between the submitted version and the official published version of record. People interested in the research are advised to contact the author for the final version of the publication, or visit the DOI to the publisher's website.
- The final author version and the galley proof are versions of the publication after peer review.
- The final published version features the final layout of the paper including the volume, issue and page numbers.
Link to publication


## General rights

Copyright and moral rights for the publications made accessible in the public portal are retained by the authors and/or other copyright owners and it is a condition of accessing publications that users recognise and abide by the legal requirements associated with these rights.

- Users may download and print one copy of any publication from the public portal for the purpose of private study or research.
- You may not further distribute the material or use it for any profit-making activity or commercial gain
- You may freely distribute the URL identifying the publication in the public portal.

If the publication is distributed under the terms of Article 25fa of the Dutch Copyright Act, indicated by the "Taverne" license above, please follow below link for the End User Agreement:
www.tue.nl/taverne

## Take down policy

If you believe that this document breaches copyright please contact us at:
openaccess@tue.nl
providing details and we will investigate your claim.

# Minimax estimation of the Wigner function in quantum homodyne tomography with ideal detectors 

Mădălin Guţă * and Luis Artiles ${ }^{\dagger}$


#### Abstract

We estimate the quantum state of a light beam from results of quantum homodyne measurements performed on identically prepared pulses. The state is represented through the Wigner function, a "quasiprobability density" on $\mathbb{R}^{2}$ which may take negative values and must respect intrinsic positivity constraints imposed by quantum physics. The data consists of $n$ i.i.d. observations from a probability density equal to the Radon transform of the Wigner function. We construct an estimator for the Wigner function, and prove that it is minimax efficient for the pointwise risk over a class of infinitely differentiable functions. A similar result was previously derived by Cavalier in the context of positron emission tomography. Our work extends this result to the space of smooth Wigner functions, which is the relevant parameter space for quantum homodyne tomography.


## 1 Introduction

The phenomena occurring at the interface between the microscopic and the macroscopic worlds have an intrinsic probabilistic nature. When measuring an atom or a laser pulse we obtain a random result whose probability distribution is determined by the state, or preparation of the quantum system. For example, if we count the number of photons coming from a laser source we observe a Poisson distributed random variable with mean equal to the intensity of the laser. The statistical inverse problem of inferring the state from results of measurements on many identically prepared systems, is called quantum state estimation. Recently it has become possible to apply such a method to the reconstruction of the quantum state of a light beam. The measurement technique is called Quantum Homodyne Tomography [24] and is used to confirm the creation of new and exotic quantum states of light such as squeezed states [3], single-photon-added coherent states [26] and Schrödinger cat states [21]. As experiments become more and more complicated, the costs - in terms of money and time - of running a measurement rise, and one needs to apply more sophisticated statistical techniques to reconstruct the state from a limited number of data. This paper makes a step in this direction by providing minimax convergence rates for a class of physical states.

The object to be estimated is a real function of two variables called the Wigner function [25], which can be seen as a joint density of the electric and magnetic fields of the laser beam. However, since in quantum mechanics we cannot measure both electric and magnetic fields simultaneously, this function is in general not a probability density but has many features in common with the latter, for example the marginals along any direction are bona-fide probability densities. In quantum optics the Wigner function is a preferred representation of the quantum state [10] because many interesting quantum patterns such as squeezing or oscillation between negative and positive values, can be easily identified from the shape of the function. The estimation methods used by the physicists involve a number of ad-hoc approximations, binnings and truncations, making it difficult to verify the reliability of the procedure. Moreover, the quantum features in which the experimenter is interested may be washed out in the resulting estimator.

[^0]From the statistical viewpoint, we deal with an ill posed inverse problem which is closely related to the problem of estimating a bivariate probability density in Positron Emission Tomography. In both cases the parameter of interest is a density and the data consists of independent identically distributed samples from the Radon Transform [7] of that density, with uniformly distributed angles. However as we mentioned above, the Wigner function is not necessarily positive but must satisfy other positivity conditions dictated by the laws of quantum mechanics. This means that some of the existing statistical results are not directly applicable in this case, and the problem of estimating the Wigner function should be studied separately in order to identify the specific "quantum features" that could be exploited in designing new estimators. The problem of state estimation in Quantum Homodyne Tomography was investigated in the non-parametric setting in [1]. That paper provides sufficient conditions for consistency of various estimators for two different representations of the state, namely the density matrix and the Wigner function. The estimation of the Wigner function was further investigated in [4] in a set-up which takes into account the detection losses occurring in the measurement, leading to an additional Gaussian noise in the measurement data. Although both [4] and the present paper consider a family of very smooth Wigner functions, the estimation techniques are very different. In [4] the bias dominates the variance due to the presence of the Gaussian noise, and for deriving the lower bound it is enough to consider a "worst family" consisting of just two states. In this paper, the variance is dominating and the lower bound is based on a one parameter family of states. The techniques that we use here are very similar to the one of Cavalier [5] who considers the same tomography problem in the context of smooth probability densities rather than Wigner functions. For the problem of Positron Emission Tomography and the related regression problem of X-ray Tomography we refer to $[17,18,6]$ and the references therein. For an introduction to Quantum Statistical Inference we refer to $[2,1]$.
In Section 2 we give a short introduction to quantum mechanics, present some properties of Wigner functions and the relation to the density matrix which will play an important role in the proof of the lower bound. The statistical set-up and the main results are presented in Section 3. Following [5], we identify a class of very smooth Wigner functions

$$
\begin{equation*}
\mathcal{W}(\beta, L)=\left\{W: \frac{1}{4 \pi^{2}} \int_{\mathbb{R}^{2}}|\widetilde{W}(w)|^{2} \exp (2 \beta|w|) d w \leq L\right\} \tag{1.1}
\end{equation*}
$$

where $\beta$ and $L$ are positive constants. From the physical point of view, all the states which have been produced in the lab up to date belong to such a class, and a more detailed argument can be found in [4], as to why this assumption is realistic.

We consider a family of estimators depending on a bandwidth which is chosen according to the parameters of the class. The upper bound for the pointwise risk is proven in Theorem 1 and has the same almost parametric expression as the bound derived in [5]. For the lower bound we consider a family of Wigner functions which is different from the worst parametric family of probability densities of [5]. The latter cannot be used in our situation since it does not correspond to Wigner functions of physical states. Thus the main novelty of the paper is the derivation of the lower bound in the physical context of Quantum Homodyne Tomography rather than that of Positron Emission Tomography.

The technical Lemmas are grouped in Section 4.

## 2 Quantum Homodyne Tomography and the Wigner function

In this section we briefly present some basic notions of quantum mechanics, the mathematical set-up of Quantum Homodyne Tomography, and some properties of Wigner functions. More background information on quantum statistical inference can be found in the review paper [2] and the textbooks [14] and [15]. Quantum Homodyne Tomography is discussed in detail in [1, 4].

The state of a quantum system encodes all necessary information for computing the probability distribution of results of any given measurement. Mathematically, the state is described by a density matrix, which is a compact operator $\rho$ on a complex Hilbert space $\mathcal{H}$ having the following properties:

1. Selfadjoint: $\rho=\rho^{*}$, where $\rho^{*}$ is the adjoint of $\rho$.
2. Positive: $\rho \geq 0$, or equivalently $\langle\psi, \rho \psi\rangle \geq 0$ for all $\psi \in \mathcal{H}$.
3. Trace one: $\operatorname{Tr}(\rho)=1$.

Positivity implies that the eigenvalues of $\rho$ are all nonegative and by the last property, they sum up to one. Note that the above requirements parallel properties of probability densities.

Since $\rho$ is a selfadjoint operator, it can be diagonalized, i.e. it has the form

$$
\begin{equation*}
\rho=\sum_{i=1}^{\operatorname{dim} \mathcal{H}} \lambda_{i} \rho_{i} \tag{2.1}
\end{equation*}
$$

where $\rho_{i}$ is the projection onto the one dimensional space generated by the eigenvector $e_{i} \in \mathcal{H}$ of $\rho$ corresponding to the eigenvalue $\lambda_{i}$, i.e., $\rho e_{i}=\lambda_{i} e_{i}$. With respect to a fixed orthonormal basis $\left\{\psi_{i}\right\}_{i \geq 1}$ in $\mathcal{H}$, the operator $\rho$ can be represented as a matrix with elements $\rho_{i, j}=\left\langle\psi_{i}, \rho \psi_{j}\right\rangle$.
Let us consider now the following problem. We are given a quantum system prepared in an unknown state $\rho$ and we would like to determine $\rho$. In order to obtain information about the system we have to measure its properties. The laws of quantum mechanics say that for any given measurement with space of outcomes given by the measure space $(\mathcal{X}, \Sigma)$, the result of the measurement performed on a system prepared in state $\rho$ is random and has probability distribution $\mathbb{P}_{\rho}$ over $(\Omega, \Sigma)$ such that the map

$$
\begin{equation*}
\rho \mapsto \mathbb{P}_{\rho}, \tag{2.2}
\end{equation*}
$$

is affine, i.e. it maps convex combinations of states into the corresponding convex combination of probability distributions. This has a natural interpretation: a system can be prepared in a mixture $\lambda \rho_{1}+(1-\lambda) \rho_{2}$ of states by randomly choosing the preparation procedure according to the individual states $\rho_{1}$ and $\rho_{2}$ with probabilities $\lambda$ and respectively $1-\lambda$. The distribution of the results will then reflect this randomized preparation as well.
The simplest measurement is that of an observable such as energy, position, spin. Any given observables is described by some selfadjoint operator $\mathbf{X}$ on the Hilbert space $\mathcal{H}$, and supposing for simplicity that it has a discrete spectrum, it can be written in the diagonal form

$$
\begin{equation*}
\mathbf{X}=\sum_{i=1}^{\operatorname{dim} \mathcal{H}} x_{a} \mathbf{P}_{a} \tag{2.3}
\end{equation*}
$$

with $x_{a} \in \mathbb{R}$ eignevalues, and $\mathbf{P}_{a}$ one dimensional projections onto the eigenvectors of $\mathbf{X}$. The result $X$ of the measurement of the observable $\mathbf{X}$ for a preparation given by the state $\rho$, is a random element of the set $\Omega=\left\{x_{1}, x_{2}, \ldots\right\}$ of eigenvalues of $\mathbf{X}$ and has the probability distribution

$$
\begin{equation*}
\mathbb{P}_{\rho}\left[X=x_{a}\right]=\operatorname{Tr}\left(\mathbf{P}_{a} \rho\right) . \tag{2.4}
\end{equation*}
$$

This measurement will give statistical information about the diagonal elements of the density matrix $\rho$ with respect to the eigenbasis of $\mathbf{X}$, and it suggests that in order to estimate all matrix elements of $\rho$ one would have to probe the system from a number of "directions" by performing different measurements on identically prepared systems.
This framework can be generalized to the case of infinite dimensional Hilbert space, and measurements with outcomes in arbitrary measure spaces. In the next section we will see that an infinite dimensional density matrix can be estimated by measuring a randomly chosen observable from a continuous family of non-commuting observables.

### 2.1 Quantum homodyne tomography and the Wigner function

An important example of quantum system is the monochromatic light in a cavity, described by density matrices on the Hilbert space of complex valued square integrable functions on the real line, $L_{2}(\mathbb{R})$. A distinguished orthonormal basis of this space is given by the vectors

$$
\begin{equation*}
\psi_{k}(x)=H_{k}(x) e^{-x^{2} / 2}, \quad k=0,1,2 \ldots \tag{2.5}
\end{equation*}
$$

where $H_{k}$ are the Hermite polynomials normalized such that $\psi_{k}(x)$ is a unit vector representing the pure state of $k$ photons inside the cavity. We will denote the matrix elements of $\rho$ with respect to this basis by $\rho_{i, j}$. Notice that the diagonal of the density matrix is a probability distribution over the nonnegative numbers $p_{k}=\rho_{k, k}$. This is the probability distribution of results when measuring the number of photons in the cavity prepared in state $\rho$. Clearly this distribution does not contain information about the off-diagonal elements of $\rho$ thus it is not sufficient for identifying the state of the system. This is a typical situation in state estimation and one has to devise experiments in which the quantum systems are looked at subsequently from "different directions", a broadly described methodology which in the physics literature goes by the name of quantum tomography. Quantum Homodyne Tomography is one such measurement method which is frequently used in quantum optics at the estimation of the quantum state of light [3, 21, 26]. A detailed description of the measurement set-up can be found in $[1,4]$. Here we present the statistical problem associated to this measurement.
We observe i.i.d. random variables $\left(X_{1}, \Phi_{1}\right), \ldots,\left(X_{n}, \Phi_{n}\right)$ with values in $\mathbb{R} \times[0,2 \pi]$ and distribution $\mathbb{P}_{\rho}$ whose density with respect to the measure $(2 \pi)^{-1} d \phi \times d x$ is given by

$$
\begin{equation*}
p_{\rho}(x, \phi)=\sum_{j, k=0}^{\infty} \rho_{j, k} \psi_{k}(x) \psi_{j}(x) e^{-i(j-k) \phi} \tag{2.6}
\end{equation*}
$$

Since $\rho$ is a positive definite matrix of trace 1 , and $\psi_{j}$ form an orthonormal basis, it follows that $p_{\rho}$ is a probability density: real, nonnegative, integrates to 1 . The data ( $X_{\ell}, \Phi_{\ell}$ ), $\ell=1, \ldots, n$, come from independent measurements on identically prepared pulses of light whose state is completely encoded in the matrix $\rho$.

For each of the systems independently, we repeat the following experimental procedure: we first choose an angle $\Phi$ uniformly distributed over $[0,2 \pi]$ and then measure a certain observable $\mathbf{X}_{\Phi}$ called quadrature, obtaining a real valued result with probability density $p_{\rho}(x, \phi)$ for any fixed value $\phi$ taken by the angle $\Phi$. The quadrature is defined as the linear combination $X_{\phi}:=\cos \phi \mathbf{Q}+\sin \phi \mathbf{P}$, where $\mathbf{Q}$ and $\mathbf{P}$ are the electric and magnetic fields of the light beam given by the selfadjoint operators

$$
\mathbf{Q} \psi(x)=x \psi(x), \quad \text { and } \quad \mathbf{P} \psi(x)=-i \frac{d \psi}{d x}
$$

The characteristic functions of the densities $p_{\rho}(x, \phi)$ can be put together to define a function of two variables

$$
\begin{equation*}
\widetilde{W}_{\rho}(u, v):=\operatorname{Tr}\left(\rho \exp \left(-i t \mathbf{X}_{\phi}\right)\right)=\mathcal{F}_{1}\left[p_{\rho}(\cdot, \phi)\right](t), \tag{2.7}
\end{equation*}
$$

where we have used the polar coordinates $(u, v)=(t \cos \phi, t \sin \phi)$, and $\mathcal{F}_{1}$ is the Fourier transform with respect to the first variable, for fixed $\phi$. Note that our convention for defining the Fourier transform and its inverse are the following

$$
\mathcal{F}[f](t)=\int_{-\infty}^{\infty} f(x) e^{-i x t} d x, \quad \mathcal{F}^{-1}[g](x)=\frac{1}{2 \pi} \int_{-\infty}^{\infty} g(t) e^{i x t} d t
$$

Equivalently, we can write

$$
\begin{equation*}
\widetilde{W}_{\rho}(u, v)=\operatorname{Tr}(\rho \exp (-i u \mathbf{Q}-i v \mathbf{P})) \tag{2.8}
\end{equation*}
$$

which resembles a characteristic function of a bivariate probability density, namely the joint density of $\mathbf{Q}$ and $\mathbf{P}$. However, since the operators $\mathbf{Q}$ and $\underset{\mathbf{P}}{ }$ do not commute with each other, we cannot speak of their joint probability distribution and the function $\widetilde{W}_{\rho}(u, v)$ is in general not the characteristic function of a probability density but rather of the so called Wigner function

$$
\begin{equation*}
W_{\rho}(q, p):=\mathcal{F}_{2}^{-1}\left[\widetilde{W}_{\rho}\right](q, p), \tag{2.9}
\end{equation*}
$$

a "quasi-distribution" which may take negative values but whose marginals are bona-fide probability densities. As we will see below, the Wigner function $W_{\rho}$ is in one to one correspondence with the density matrix $\rho$, and in quantum optics one frequently uses the Wigner function as an alternative representation of the quantum state, having the advantage that it illustrates important "quantum features" such as squeezing and negative oscillations. From (2.7) and (2.9) we deduce that the probability density of the data $p_{\rho}(x, \phi)$ is the Radon transform of the Wigner function

$$
\mathcal{R}\left[W_{\rho}\right](x, \phi)=\int_{-\infty}^{\infty} W_{\rho}(x \cos \phi-t \sin \phi, x \sin \phi+t \cos \phi) d t
$$

adding Quantum Homodyne Tomography to a long the list of applications ranging from computerized tomography to astronomy and geophysics [7].
Another important feature of the Wigner function is that it can be used as a computational tool: for any selfadjoint operator $\mathbf{X}$ there exists a function $W_{\mathbf{X}}$ from $\mathbb{R}^{2}$ to $\mathbb{R}$ such that the expectation of $\mathbf{X}$ is given by

$$
\begin{equation*}
\operatorname{Tr}(\mathbf{X} \rho)=2 \pi \iint W_{\mathbf{X}}(q, p) W_{\rho}(q, p) d q d p \tag{2.10}
\end{equation*}
$$

In particular, the correspondence between the density matrix $\rho$ and the Wigner function $W_{\rho}$ is an $L_{2}$ isometry up to a constant

$$
\begin{equation*}
\left\|W_{\rho}-W_{\tau}\right\|_{2}^{2}:=\iint\left|W_{\rho}(q, p)-W_{\tau}(q, p)\right|^{2} d p d q=\frac{1}{2 \pi} \sum_{i, j=0}^{\infty}\left|\rho_{i, j}-\tau_{i, j}\right|^{2}:=\frac{1}{2 \pi}\|\rho-\tau\|_{2}^{2} \tag{2.11}
\end{equation*}
$$

The space of Wigner function has an overlap with that of probability distributions. For example, all Gaussian densities which are bounded from above by $1 / \pi$ are Wigner functions, while the rest do not correspond to physical states. This is due to the celebrated Heisenberg's uncertainty relations which say that the non-commuting observables $\mathbf{P}$ and $\mathbf{Q}$ cannot have probability distributions such that the product of their variances is smaller than $\frac{1}{4}$. In general, a Wigner function cannot be too "peaked":

$$
\begin{equation*}
\left|W_{\rho}(q, p)\right| \leq \frac{1}{\pi}, \quad \text { for all }(q, p) \in \mathbb{R}^{2} \tag{2.12}
\end{equation*}
$$

Some examples of quantum states which can be created in laboratory are given in Table 1 of [1]. Typically, the corresponding Wigner functions have a Gaussian tail but need not be positive. For example the state of one-photon in the cavity is described by the density matrix with $\rho_{1,1}=1$ and all other elements zero which is equal to the orthogonal projection onto the vector $\psi_{1}$. The corresponding Wigner function is

$$
W(q, p)=\frac{1}{\pi}\left(2 q^{2}+2 p^{2}-1\right) \exp \left(-q^{2}-p^{2}\right)
$$

In conclusion, although we deal with a problem which is similar to that of Positron Emission Tomography, the parameter space is different from the space of probability densities and special techniques have to be developed for this situation.

## 3 The main results

Our problem is that of estimating the Wigner function $W_{\rho}(z)$, defined on the plane i.e. $z=(q, p) \in \mathbb{R}^{2}$. We consider the class $\mathcal{W}(\beta, L)$ of very smooth Wigner functions whose Fourier transform satisfies

$$
\begin{equation*}
\frac{1}{4 \pi^{2}} \int_{\mathbb{R}^{2}}|\widetilde{W}(w)|^{2} \exp (2 \beta|w|) d w \leq L \tag{3.1}
\end{equation*}
$$

for $\beta$ and $L$ positive constants. Such classes appeared in the statistical literature in [16], and we subsequently used in the context of density estimation [12], functional estimation, regression problems [11], and tomography [5]. In [4] we have argued that from the physical point of view it is natural to assume that the Wigner function of a state which can be created in the lab belongs to such a class.

We use a kernel-type estimator based on the following function called a band-limited filter

$$
\begin{equation*}
K_{\delta_{n}}(u)=\frac{1}{4 \pi} \int_{-\delta_{n}}^{\delta_{n}} r e^{i r u} d r . \tag{3.2}
\end{equation*}
$$

This filter has already been used in the context of tomography [20, 18, 5], and its Fourier transform is

$$
\begin{equation*}
\widetilde{K}_{\delta_{n}}(t)=\frac{1}{2}|t| I_{\delta_{n}}(t) \tag{3.3}
\end{equation*}
$$

where $I_{\delta_{n}}$ is the indicator function of $\left\{t:|t| \leq 1 / \delta_{n}\right\}$. The estimator that we use is the following

$$
\begin{equation*}
\widehat{W}_{n}(z)=\frac{1}{n} \sum_{i=1}^{n} K_{\delta_{n}}\left(\left[z, \Phi_{i}\right]-X_{i}\right), \tag{3.4}
\end{equation*}
$$

where $[z, \phi]:=q \cos \phi+p \sin \phi$, and $\left(X_{i}, \Phi_{i}\right)$ are i.i.d. observations with density $p_{\rho}(x, \phi)=\mathcal{R}\left[W_{\rho}\right](x, \phi)$.
Following [20] we define the dual operator $\mathcal{R}^{\#}$ on $L_{1}(\mathbb{R} \times[0, \pi])$ by

$$
\begin{equation*}
\mathcal{R}^{\#}[h](z)=\int_{0}^{2 \pi} h([z, \phi], \phi) d \phi . \tag{3.5}
\end{equation*}
$$

Then

$$
\begin{equation*}
\mathcal{R}^{\#} \mathcal{R}[W](z)=\int_{0}^{2 \pi} \mathcal{R}[W]([z, \phi], \phi) d \phi \tag{3.6}
\end{equation*}
$$

represents the integrals of $W$ over all lines passing through the point $x$. Note that in general $R^{\#} \mathcal{R}[W](z) \geq 0$ for all Wigner functions $W$ and all $z \in \mathbb{R}^{2}$, but there exist states for which $\mathcal{R} \# \mathcal{R}[W](0)=0$, for example the number states $\psi_{k}$ with $k$ odd. In [5] it is assumed that the probability distributions $f$ to be estimated are strictly positive which implies that $\mathcal{R}^{\#} \mathcal{R}[f](z)>0$ for all $z$. For the upper bound we will assume that the latter condition holds.

### 3.1 The upper bound

Theorem 1 For any $W \in \mathcal{W}(\beta, L)$ and any fixed $z \in \mathbb{R}^{2}$ such that $R^{\#} \mathcal{R}[W](z)>0$ we have

$$
\begin{equation*}
\mathbb{E}\left[\left(\widehat{W}_{n}(z)-W(z)\right)^{2}\right]=C^{*} \mathcal{R}^{\#} \mathcal{R}[W](z) \times \frac{(\log n)^{3}}{n}(1+o(1)), \quad \text { as } n \rightarrow \infty \tag{3.7}
\end{equation*}
$$

where $C^{*}=\frac{\pi}{3(4 \pi \beta)^{3}}$.
Proof. We will provide only the main steps of the proof pointing out where the assumption on the class of Wigner function plays a significant role. For a more detailed proof of the bounds for the class of probability densities in $\mathcal{A}(\beta, L)$ we refer to [5]. The risk can be decomposed in two parts, the bias and the stochastic part

$$
\begin{align*}
\mathbb{E}\left[\left(\widehat{W}_{n}(z)-W(z)\right)^{2}\right] & =\left(\mathbb{E}\left(\widehat{W}_{n}(z)\right)-W(z)\right)^{2}+\mathbb{E}\left[\left(\widehat{W}_{n}(z)-\mathbb{E}\left(\widehat{W}_{n}(z)\right)\right)^{2}\right] \\
& :=b_{n}^{2}(z)+\sigma_{n}^{2}(z) \tag{3.8}
\end{align*}
$$

On the one hand, using property (3.3) of the kernel and the inverse Fourier transform, the bias can be written as

$$
\begin{equation*}
b_{n}(z)=\frac{1}{(2 \pi)^{2}} \int \widetilde{W}(w) I\left(|w|>\frac{1}{\delta_{n}}\right) e^{-i\langle w, z\rangle} d w . \tag{3.9}
\end{equation*}
$$

By using Cauchy-Schwarz and the fact that $W \in \mathcal{W}(\beta, L)$ we get

$$
\begin{equation*}
b_{n}^{2}(z) \leq \frac{c_{1}}{\delta_{n}} e^{-2 \beta / \delta_{n}}(1+o(1)), \quad \text { as } \delta_{n} \rightarrow 0, \tag{3.10}
\end{equation*}
$$

for some constant $c_{1}$. With the choice $1 / \delta_{n}=\log n /(2 \beta)$ the bias upper bound becomes

$$
\begin{equation*}
b_{n}^{2}(z) \leq c_{2} \frac{\log n}{n}(1+o(1)), \quad \text { as } n \rightarrow \infty \tag{3.11}
\end{equation*}
$$

On the other hand, the variance is equal to

$$
\begin{align*}
\sigma_{n}^{2}(z)= & \frac{1}{n} \operatorname{Var} K_{\delta_{n}}([z, \Phi]-X)= \\
& \frac{1}{n} \mathbb{E}\left[K_{\delta_{n}}^{2}([z, \Phi]-X)\right]-\frac{1}{n}\left(\mathbb{E}\left[K_{\delta_{n}}([z, \Phi]-X)\right]\right)^{2}, \tag{3.12}
\end{align*}
$$

where $(X, \Phi)$ is a random variable with probability density $p_{\rho}(x, \phi)=\mathcal{R}\left[W_{\rho}\right](x, \phi)$. The second term is

$$
\frac{1}{n}\left(\mathbb{E}\left[K_{\delta_{n}}([z, \Phi]-X)\right]\right)^{2}=\frac{1}{n}\left(\frac{1}{2 \pi} \int_{0}^{2 \pi} \int_{\mathbb{R}} p_{\rho}(x, \phi) K_{\delta_{n}}([z, \phi]-x) d \phi d x\right)^{2}
$$

By Parseval equality and (3.3) we can pass to the Fourier transformed expressions and then bound from above using Cauchy-Schwarz inequality:

$$
\begin{equation*}
\frac{1}{n}\left(\mathbb{E}\left[K_{\delta_{n}}([z, \Phi]-X)\right]\right)^{2} \leq \frac{1}{(2 \pi)^{4} n} \int I_{\delta_{n}}(|w|) e^{-2 \beta|w|} d w \times \int|\widetilde{W}(w)|^{2} e^{2 \beta|w|} d w=O(1 / n), \tag{3.13}
\end{equation*}
$$

with $O(1 / n)$ uniformly with respect to $W \in \mathcal{W}(\beta, L)$.
The first term is

$$
\begin{equation*}
\mathbb{E}\left[K_{\delta_{n}}^{2}([z, \Phi]-X)\right]=\frac{1}{2 \pi} \int_{0}^{\pi} \int_{\mathbb{R}} K_{\delta_{n}}^{2}([z, \phi]-y) p_{\rho}(y, \phi) d \phi d y \tag{3.14}
\end{equation*}
$$

Now denote

$$
\begin{equation*}
G(u)=\frac{3}{\pi}\left(\int_{0}^{1} r \cos (u r) d r\right)^{2} \tag{3.15}
\end{equation*}
$$

and let $G_{\delta}(u)=(1 / \delta) G(u / \delta)$. We have

$$
\begin{equation*}
\int_{\mathbb{R}} K_{\delta_{n}}^{2}([z, \phi]-y) p_{\rho}(y, \phi) d y=\frac{1}{6(2 \pi)^{2}}\left(\frac{1}{\delta_{n}}\right)^{3}\left(G_{\delta_{n}} * \mathcal{R}[W](\cdot, \phi)\right)([z, \phi]) \tag{3.16}
\end{equation*}
$$

Using Lemma 4 of [5], we have that as $\delta \rightarrow 0$,

$$
\int_{0}^{\pi}\left(G_{\delta} * \mathcal{R}[W](\cdot, \phi)\right) d \phi=\mathcal{R}^{\#} \mathcal{R}[W](z)(1+o(1))+O\left(\delta^{1 / 3}\right)
$$

In proving this we need an estimate of the smoothness of $\mathcal{R}[W](x, \phi)$ as describe in Lemma 2 which gives a similar result to that of Lemma 3 of [5]. Thus, the second term of the variance can be written as

$$
\mathbb{E}\left[K_{\delta_{n}}^{2}([z, \Phi]-X)\right]=\frac{1}{6(2 \pi)^{2}}\left(\frac{1}{\delta_{n}}\right)^{3} \mathcal{R}^{\#} \mathcal{R}[W](z)(1+o(1))+O\left(\left(\frac{1}{\delta_{n}}\right)^{3-1 / 3}\right),
$$

as $\delta_{n} \rightarrow 0$. Thus as $n \rightarrow \infty$

$$
\begin{aligned}
\sigma_{n}^{2}(z) & \leq C^{*} \mathcal{R}^{\#} \mathcal{R}[W](z) \frac{(\log n)^{3}}{n}(1+o(1))+O\left(\frac{(\log n)^{3-1 / 3}}{n}\right) \\
& =\frac{\log ^{3} n}{n}\left(C^{*} \mathcal{R}^{\#} \mathcal{R}[W](z)(1+o(1))+O\left(\log ^{-1 / 3} n\right)\right)
\end{aligned}
$$

If $R^{\#} \mathcal{R}[W](z)>0$ then we obtain the claimed constant. Notice that if $R^{\#} \mathcal{R}[W](z)=0$ then the convergence rate is faster than $\frac{\log ^{3} n}{n}$.

### 3.2 The lower bound

In order to prove a lower bound result we consider the slightly modified class of Wigner functions

$$
\begin{equation*}
\mathcal{W}\left(\beta, L, \alpha_{n}\right)=\left\{W \in \mathcal{W}(\beta, L): \mathcal{R}^{\#} \mathcal{R}[W](z) \geq \alpha_{n}\right\} \tag{3.17}
\end{equation*}
$$

for a sequence $\alpha_{n}$ such that $\lim _{n \rightarrow \infty} \alpha_{n}=0$ and $\lim _{n \rightarrow \infty}\left(\alpha_{n}(\log n)^{1 / 3}\right)=\infty$. Let us denote

$$
\begin{equation*}
r_{n}(W, z)=\left(C^{*} \mathcal{R}^{\#} \mathcal{R}[W](z) \frac{(\log n)^{3}}{n}\right)^{1 / 2} \tag{3.18}
\end{equation*}
$$

Theorem 2 For any fixed $z \in \mathbb{R}^{2}$, we have

$$
\liminf _{n \rightarrow \infty} \inf _{\widehat{W}_{n}} \sup _{W \in \mathcal{W}\left(\beta, L, \alpha_{n}\right)} \mathbb{E}\left[\left(\frac{\widehat{W}_{n}(z)-W(z)}{r_{n}(W, z)}\right)^{2}\right] \geq 1
$$

where $\inf _{\widehat{W}_{n}}$ denotes the infimum over all estimators of $W(z)$.
Proof. The proof is based on the standard procedure of building a hardest parametric subfamily for the class $\mathcal{W}(\beta, L)$ of Wigner functions of the form

$$
\begin{equation*}
W_{c}=W_{\alpha}+c g_{a} \tag{3.19}
\end{equation*}
$$

where $c$ is a parameter in a neighborhood of the origin, $W_{\alpha}$ and $g_{a}$ are functions which will be defined shortly. The essential point of the proof is that the family of probability densities $f_{0}+c g_{a}$ used in [5] is not always contained in our parameter space consisting of Wigner functions. For illustration we will show that for some parameters $\beta$, the function $f_{0}$ defined in equation (29) of [5] is not the Wigner function of a quantum state. Indeed, suppose for the moment that this was the case, i.e. $f_{0}=W_{\rho}$ for some state $\rho$. Then by the rotation symmetry of $f_{0}$, the density matrix $\rho$ must be diagonal and

$$
p_{\rho}(x, \phi)=\sum_{k=0}^{\infty} \rho_{k, k} \psi_{k}^{2}(x) .
$$

Using the inequality [8] $\left\|\psi_{k}\right\|_{\infty} \leq k$, where $k$ is a constant whose value is slightly bigger that 1 , and the fact that $\sum \rho_{k, k}=1$, we find that $\left\|p_{\rho}\right\|_{\infty} \leq k$. However, the Radon transform of the Wigner function $W_{\rho}$ is

$$
p_{\rho}(x, \phi)=\mathcal{R}\left[W_{\rho}\right](x, \phi)=\mathcal{R}\left[f_{0}\right](x, \phi)=\frac{\beta}{\pi} \frac{1}{x^{2}+\beta^{2}}
$$

which violates the above bound for $\beta<1 /(\pi k)$.

We thus define a parametric subfamily of $\mathcal{W}(\beta, L)$ which is a suitable modification of the family considered in [5] in order to cope with this problem.

Construction of $W_{\alpha}$. Consider the Mehler formula, (see [8], 10.13.22)

$$
\begin{equation*}
\sum_{n=0}^{\infty} z^{n} \frac{1}{\sqrt{\pi} n!2^{n}} H_{n}(x)^{2} e^{-x^{2}}=\frac{1}{\sqrt{\pi\left(1-z^{2}\right)}} \exp \left(-x^{2} \frac{1-z}{1+z}\right) . \tag{3.20}
\end{equation*}
$$

Integrating both terms with $f_{\alpha}(z)=\alpha(1-z)^{\alpha}$ we get

$$
\begin{equation*}
p_{\alpha}(x, \phi):=\sum_{n=0}^{\infty} \psi_{n}(x)^{2} \int_{0}^{1} f_{\alpha}(z) z^{n} d z=\int_{0}^{1} \frac{f_{\alpha}(z)}{\sqrt{\pi\left(1-z^{2}\right)}} \exp \left(-x^{2} \frac{1-z}{1+z}\right) d z \tag{3.21}
\end{equation*}
$$

The Fourier transform of $p_{\alpha}$ is

$$
\begin{equation*}
\widetilde{W}_{\alpha}(w):=\mathcal{F}\left[p_{\alpha}\right](w)=\int_{0}^{1} \frac{f_{\alpha}(z)}{1-z} \exp \left(-|w|^{2} \frac{1+z}{4(1-z)}\right) d z \tag{3.22}
\end{equation*}
$$

Notice that the normalization condition $\int p_{\alpha}=1$ is equivalent to $\widetilde{W}_{\alpha}(0)=1$ which is satisfied for the chosen functions $f_{\alpha}$, thus $p_{\alpha}$ is a probability density corresponding to a diagonal density matrix $\rho^{\alpha}$ with elements

$$
\begin{equation*}
\rho_{k, k}^{\alpha}=\int_{0}^{1} z^{k} f_{\alpha} d z \tag{3.23}
\end{equation*}
$$

We denote by $W_{\alpha}$ the Wigner function whose Fourier transform is defined in equation (3.22) with $\alpha>0$ a parameter to be fixed later. This function is considered in a more general form in [4], and corresponds to $f_{\alpha}^{\epsilon}$ for $\epsilon=0$.

Construction of $g_{a}$. Let [5]

$$
\begin{equation*}
H_{a}(s)=\frac{1}{(2 \pi)^{2}} \int_{0}^{\infty} \frac{r}{1+a^{-1} \sinh ^{2} \beta r} \cos (s r) d r, \quad s \in \mathbb{R} \tag{3.24}
\end{equation*}
$$

where $a>0$ is a parameter which will depend on $n$ as $a=a_{n}=n^{\eta}$ with $0<\eta<1$. The Fourier transform of $g_{a}$ is

$$
\begin{equation*}
\widetilde{H}_{a}(t)=\frac{1}{4 \pi} \frac{t}{1+a^{-1} \sinh ^{2} \beta t}, \quad t \in \mathbb{R} \tag{3.25}
\end{equation*}
$$

Let

$$
\begin{equation*}
g_{a}(z)=\frac{1}{2(2 \pi)^{3}} \int_{\mathbb{R}^{2}} \frac{|w|}{1+a^{-1} \sinh ^{2} \beta|w|} \cos (\langle z, w\rangle) d w, \quad z \in \mathbb{R}^{2} \tag{3.26}
\end{equation*}
$$

and its Fourier transform

$$
\begin{equation*}
\widetilde{g}_{a}(w)=\frac{1}{4 \pi} \frac{|w|}{1+a^{-1} \sinh ^{2} \beta|w|}, \quad w \in \mathbb{R}^{2} \tag{3.27}
\end{equation*}
$$

Now, let us consider the family

$$
\begin{equation*}
W_{c}=W_{\alpha}+c g_{a} \tag{3.28}
\end{equation*}
$$

where the real parameter $c$ satisfies

$$
\begin{equation*}
|c| \leq C_{a}=\frac{q}{\sqrt{a}(\log a)^{3 / 2}} \tag{3.29}
\end{equation*}
$$

with $q>0$ a sufficiently small constant. By translating with $z$ in $\mathbb{R}^{2}$ we obtain our hardest family for pointwise estimation at the point $z$

$$
\begin{equation*}
W_{c}^{z}(\zeta)=W_{c}(\zeta-z) \tag{3.30}
\end{equation*}
$$

We will check now that $W_{c}^{z}$ belongs to the class $\mathcal{W}(\beta, L)$ for an appropriate choice of $\alpha$ in $W_{\alpha}$, which means that $W_{c}^{z}$ is a Wigner function and

$$
\begin{equation*}
\frac{1}{4 \pi^{2}} \int\left|\widetilde{W}_{c}^{z}(w)\right|^{2} e^{2 \beta|w|} d w=\frac{1}{4 \pi^{2}} \int\left|\widetilde{W}_{c}(w)\right|^{2} e^{2 \beta|w|} d w \leq L \tag{3.31}
\end{equation*}
$$

By Lemma 5 we have $W_{\alpha} \in \mathcal{W}(\beta, L / 4)$ for a small enough $\alpha>0$, and from Lemma 5 of [5] we have that

$$
\begin{equation*}
\frac{c^{2}}{4 \pi^{2}} \int\left|\tilde{g}_{a}(w)\right|^{2} e^{2 \beta|w|} d w \leq L / 4 \tag{3.32}
\end{equation*}
$$

for all $c \leq C_{a}$. The last two conditions together imply (3.31).
Furthermore, $W_{c}^{z}$ has to be a Wigner function. As translations in the plane transform Wigner functions into Wigner functions, we need only to show this for $W_{c}$. This means that there exist a family of density matrices $\rho_{c}$ such that their corresponding Wigner functions are $W_{c}$. The invariance of $W_{c}$ under rotations in the plane translates into the fact that $\rho^{c}$ has all off-diagonal elements equal to zero, and thus we only need to show that all its diagonal elements are positive and add up to one. The relation between the diagonal matrix elements and the Wigner function is [19]

$$
\begin{equation*}
\rho_{k, k}^{c}=\frac{1}{2 \pi} \int_{\mathbb{R}^{2}} e^{|w|^{2} / 4} L_{k}\left(|w|^{2} / 2\right) \widetilde{W}_{c}(w) d w \tag{3.33}
\end{equation*}
$$

where $L_{k}$ are the Laguerre polynomials. By linearity we have

$$
\begin{equation*}
\rho_{k, k}^{c}=\rho_{k, k}^{\alpha}+c \tau_{k, k}^{a} \tag{3.34}
\end{equation*}
$$

where

$$
\begin{equation*}
\tau_{k, k}^{a}=\int_{0}^{\infty} t e^{-t^{2} / 4} L_{k}\left(t^{2} / 2\right) \widetilde{H}_{a}(t) d t \tag{3.35}
\end{equation*}
$$

and

$$
\begin{equation*}
\rho_{k, k}^{\alpha}=\alpha \int_{0}^{1} z^{k}(1-z)^{\alpha} . \tag{3.36}
\end{equation*}
$$

Corroborating the result shown in Lemma 4

$$
\tau_{k k}^{a}=O\left(k^{-5 / 4}(\log a)^{4}\right)
$$

as $a, k \rightarrow \infty$, with that of Lemma 2 in [4]

$$
\rho_{k, k}^{\alpha} \sim k^{-(1+\alpha)}
$$

we conclude that if $\alpha<1 / 4$, then $\rho_{k, k}^{c} \geq 0$ for all $k \geq 0$ and $|c| \leq C_{a}$ for $a$ sufficiently large.
Now we can use the fact that for the family of translated functions, as defined in $(3.30), \mathcal{R}^{\#} \mathcal{R}\left[W^{z}\right](z)=$ $\mathcal{R}^{\#} \mathcal{R}[W](0)$. Indeed, using (2.7) we get

$$
\begin{align*}
\mathcal{R}\left[W^{z}\right](x, \phi) & =\frac{1}{2 \pi} \int \widetilde{W}^{z}(t \cos \phi, t \sin \phi) e^{i t x} d t \\
& =\frac{1}{2 \pi} \int \widetilde{W}(t \cos \phi, t \sin \phi) e^{-i t[z, \phi]} e^{i t x} d t \\
& =\frac{1}{2 \pi} \int \widetilde{W}(t \cos \phi, t \sin \phi) e^{i t(x-[z, \phi])} d t \\
& =\mathcal{R}[W](x-[z, \phi], \phi) . \tag{3.37}
\end{align*}
$$

Now, from definition of $\mathcal{R}^{\#}$ we get

$$
\begin{aligned}
\mathcal{R}^{\#} \mathcal{R}\left[W^{z}\right](z) & =\int_{0}^{2 \pi} \mathcal{R}\left[W^{z}\right]([z, \phi], \phi) d \phi \\
& =\int_{0}^{2 \pi} \mathcal{R}[W](0, \phi) d \phi=\mathcal{R}^{\#} \mathcal{R}[W](0) .
\end{aligned}
$$

Thus, $R_{c}(z):=\mathcal{R}^{\#} \mathcal{R}\left[W_{c}^{z}\right](z)=\mathcal{R}^{\#} \mathcal{R}\left[W_{\alpha}\right](0)+c \mathcal{R}^{\#} \mathcal{R}\left[g_{a}\right](0)$ for any $z$. Given in our case $W_{\alpha}$ and $g_{a}$ are invariant under rotations we obtain

$$
\begin{align*}
R_{0} & :=\mathcal{R}^{\#} \mathcal{R}\left[W_{\alpha}\right](0)=\int_{0}^{2 \pi} \mathcal{R}\left[W_{\alpha}\right](0, \phi) d \phi=\int_{0}^{2 \pi} p_{\alpha}(0, \phi) d \phi \\
& =2 \sqrt{\pi} \alpha \int_{0}^{1} \frac{(1-z)^{\alpha-1 / 2}}{(1+z)^{1 / 2}} d z>0 \tag{3.38}
\end{align*}
$$

For the second term, using eq. (3.24), Lemma 5 in [5], and the definition of $C_{a}$

$$
\begin{equation*}
\sup _{|c| \leq C_{a}}\left|c \mathcal{R}^{\#} \mathcal{R}\left[g_{a}\right](0)\right|=\sup _{|c| \leq C_{a}}\left|2 \pi c H_{a}(0)\right|=o(1) \quad \text { as } \quad a \rightarrow \infty \tag{3.39}
\end{equation*}
$$

We conclude that $R_{c}=R_{0}(1+o(1)) \geq \alpha_{n}$, for $n \rightarrow \infty$ and thus, for $a$ large enough, $W_{c}^{z} \in \mathcal{W}\left(\beta, L, \alpha_{n}\right)$. Moreover, from (3.18),

$$
\begin{equation*}
r_{n}\left(W_{c}^{z}, z\right)^{2}=r_{n}\left(W_{\alpha}, 0\right)^{2}(1+o(1))=C^{*} R_{0} \frac{(\log n)^{3}}{n}(1+o(1)) \tag{3.40}
\end{equation*}
$$

The rest of the proof is based on the Van Trees inequality [23, 9] and follows along the lines of [5]. The main difference is in the proof of Lemma 6 where the Fisher information of the family of densities $\frac{1}{2 \pi} \mathcal{R}\left[W_{c}^{z}\right]$ is approximated. Take a continuously differentiable probability density, $\lambda_{0}(c)$, defined on the interval $[-1,1]$, such that $\lambda_{0}(-1)=\lambda_{0}(1)=0$, with a finite Fisher information $I_{0}$. The new density $\lambda(c)=\lambda_{a}(c)=C_{a}^{-1} \lambda_{0}\left(C_{a}^{-1} c\right)$ is a prior density with finite Fisher information $I(\lambda)=I_{0} C_{a}^{-2}$. Finally let us define $I(c)$ the Fisher information of the family of densities $\frac{1}{2 \pi} \mathcal{R}\left[W_{c}^{z}\right]$. Using the Van Trees inequality

$$
\begin{align*}
\inf _{\widehat{W}_{n}} \sup _{W \in \mathcal{W}\left(\beta, L, \alpha_{n}\right)} \mathbb{E}_{W}\left[\left(\widehat{W}_{n}(z)-W(z)\right)^{2}\right] & \geq \inf _{\widehat{W}_{n}} \sup _{|c|<C_{a}} \mathbb{E}_{W_{c}^{z}}\left[\left(\widehat{W}_{n}(z)-W_{c}^{z}(z)\right)^{2}\right] \\
\geq & \inf _{\widehat{W}_{n}} \int_{-C_{a}}^{C_{a}} \mathbb{E}_{W_{c}^{z}}\left[\left(\widehat{W}_{n}(z)-W_{c}^{z}(z)\right)^{2}\right] \lambda_{a}(c) d c \tag{3.41}
\end{align*} \frac{\left(\int_{-C_{a}}^{C_{a}}\left(\partial W_{c}^{z}(z) / \partial c\right) \lambda_{a}(c) d c\right)^{2}}{n \int_{-C_{a}}^{C_{a}} I(c) \lambda_{a}(c) d c+I(\lambda)} .
$$

Now, from equation (3.26) and Lemma 5 of [5], we get that $\partial W_{c}^{z}(z) / \partial c=g_{a}(0)=C^{*}(\log a)^{3}(1+o(1))$. By inserting the expression of $I(c)$ from Lemma 6, and using (3.40) and $a=a_{n}=n^{\eta}$ we get the lower bound

$$
\begin{aligned}
\inf _{\widehat{W}_{n}} \sup _{|c|<C_{a}} \mathbb{E}\left[\left(\frac{\widehat{W}_{n}(z)-W_{c}^{z}(z)}{r\left(W_{c}^{z}, z\right)}\right)^{2}\right] \quad & \geq \frac{\left(\int_{-C_{a}}^{C_{a}}\left(\partial W_{c}^{z}(z) / \partial c\right) \lambda_{c} d c\right)^{2}}{n \int_{-C_{a}}^{C_{a}} I(c) \lambda_{c} d c+I(\lambda)} \cdot\left(r_{n}\left(W_{c}^{z}, z\right)(1+o(1))\right)^{-2} \\
& =\frac{\left(C^{*} \eta^{3}(\log n)^{3}\right)^{2} r_{n}\left(W_{\alpha}, 0\right)^{-2}(1+o(1))}{C^{*} \eta^{3} R_{0}^{-1} n(\log n)^{3}(1+o(1))+q^{-2} I_{0} \eta^{3} n^{\eta}(\log n)^{3}}
\end{aligned}
$$

By letting $n \rightarrow \infty$ followed by $\eta \rightarrow 1$ we finally obtain

$$
\liminf _{n \rightarrow \infty} \inf _{\widehat{W}_{n}} \sup _{W \in \mathcal{W}\left(\beta, L, \alpha_{n}\right)} \mathbb{E}\left[\left(\frac{\widehat{W}_{n}(z)-W(z)}{r(W, z)}\right)^{2}\right] \geq 1 .
$$

## 4 Technical Lemmas

Lemma 1 Let $W \in \mathcal{W}(\beta, L)$. Then the following inequalities hold

$$
\begin{align*}
|\widetilde{W}(\omega)| & \leq 1, \quad \forall \omega \in \mathbb{R}^{2},  \tag{4.1}\\
\int_{\mathbb{R}^{2}}|\widetilde{W}(w)| d w & \leq Q  \tag{4.2}\\
\sup _{W \in \mathcal{W}(\beta, L)} \mathcal{R}^{\#} \mathcal{R}[W](z) & \leq 1+\frac{Q}{2 \pi} \tag{4.3}
\end{align*}
$$

where $Q$ is a constant depending only on $\beta$ and $L$.
Proof. As $W=W_{\rho}$ for some density matrix $\rho$ we use the definition of $\widetilde{W}_{\rho}$ to obtain

$$
\left|\widetilde{W}_{\rho}(u, v)\right|=\left|\operatorname{Tr}\left(\rho e^{-i u \mathbf{Q}-i v \mathbf{P}}\right)\right| \leq\|\rho\|_{1}\left\|e^{-i u \mathbf{Q}-i v \mathbf{P}}\right\|=1,
$$

where we have used the normalization of the density matrix $\rho$ and the fact that $\|U\|=1$ for any unitary operator $U$. The inequality (4.2) is a direct consequence of the definition of $\mathcal{W}(\beta, L)$ and the Cauchy-Schwarz inequality. Now using (2.7) and the previous inequalities we get

$$
\begin{aligned}
\left|\mathcal{R}^{\#} \mathcal{R}[W](z)\right| & =\left|\int_{0}^{2 \pi} \mathcal{R}[W]([z, \phi], \phi) d \phi\right|= \\
& \left|\frac{1}{2 \pi} \int_{0}^{2 \pi} d \phi \int d t \widetilde{W}(t \cos \phi, t \sin \phi) e^{i t[z, \phi]}\right| \leq \\
& \frac{1}{2 \pi} \int_{0}^{2 \pi} d \phi \int d t|\widetilde{W}(t \cos \phi, t \sin \phi)| \leq \\
& \frac{1}{2 \pi}\left(\int_{|\omega| \geq 1} d \omega|\widetilde{W}(\omega)|+\int_{0}^{2 \pi} d \phi \int_{0}^{1} d t|\widetilde{W}(t \cos \phi, t \sin \phi)|\right) \leq \\
& \frac{1}{2 \pi}(Q+2 \pi) .
\end{aligned}
$$

Lemma 2 For all $W \in \mathcal{W}(\beta, L)$ we have

$$
\int_{0}^{2 \pi}|\mathcal{R}[W](x, \phi)-\mathcal{R}[W](y, \phi)| d \phi \leq \frac{Q}{\pi}|x-y|
$$

with $Q$ the constant depending on $\beta$ and $L$ defined in Lemma 1.

Proof. By (2.7) and Lemma 1 we have

$$
\begin{aligned}
& \int_{0}^{2 \pi}|\mathcal{R}[W](x, \phi)-\mathcal{R}[W](y, \phi)| d \phi \leq \\
& \frac{1}{2 \pi} \int_{0}^{2 \pi} \int\left|e^{i t x}-e^{i t y}\right||\widetilde{W}(t \cos \phi, t \sin \phi)| d \phi d t \leq \frac{|x-y|}{\pi} \int_{\mathbb{R}^{2}}|\widetilde{W}(w)| d w \leq \frac{Q}{\pi}|x-y|
\end{aligned}
$$

where we have used that $\left|e^{i t x}-e^{i t y}\right| \leq 2|t||x-y|$.

Lemma 3 For all $0<\alpha \leq 1$ and $|x|>1$ there exist constants $c, C$ depending on $\alpha$, such that

$$
\begin{equation*}
c x^{-(1+2 \alpha)} \leq p_{\alpha}(x) \leq C x^{-(1+2 \alpha)} . \tag{4.4}
\end{equation*}
$$

Moreover there exist constants $C_{1}$ and $C_{2}$ such that the first two derivatives of $p_{\alpha}$ satisfy $\left|p_{\alpha}^{\prime}(x)\right| \leq C_{1} x^{-(2+2 \alpha)}$, and $\left|p_{\alpha}^{\prime \prime}(x)\right| \leq C_{2} x^{-(3+2 \alpha)}$.

Proof We have

$$
p_{\alpha}(x, \phi)=\frac{\alpha}{\sqrt{\pi}} \int_{0}^{1} \frac{(1-z)^{\alpha-1 / 2}}{(1+z)^{1 / 2}} \exp \left(-x^{2} \frac{1-z}{1+z}\right) d z
$$

which by the change of variables $u=x \sqrt{\frac{1-z}{1+z}}$ becomes

$$
\begin{equation*}
p_{\alpha}(x, \phi)=\frac{\alpha 2^{\alpha+1} x}{\sqrt{\pi}} \int_{0}^{x} \frac{u^{2 \alpha}}{\left(u^{2}+x^{2}\right)^{\alpha+1}} \exp \left(-u^{2}\right) d u . \tag{4.5}
\end{equation*}
$$

It can easily be checked that if $|x|>1$, the right hand side is bounded from above by $C x^{-(1+2 \alpha)}$ and from below by $c x^{-(1+2 \alpha)}$, for some positive constants $c, C$.
For the first derivative one can see that

$$
p_{\alpha}^{\prime}(x)=\frac{1}{x} p_{\alpha}(x)-\frac{\alpha(\alpha+1) 2^{\alpha+2} x^{2}}{\sqrt{\pi}} \int_{0}^{x} \frac{u^{2 \alpha}}{\left(u^{2}+x^{2}\right)^{\alpha+2}} d u+\frac{\alpha}{\sqrt{\pi} x} \exp \left(-x^{2}\right)
$$

which can be bounded using the same argument as before to obtain

$$
\left|p_{\alpha}^{(1)}(x)\right| \leq C_{1} x^{-(2+2 \alpha)},
$$

for some $C_{1}$ and $|x|>1$. For the second derivative the procedure is the same.

Lemma 4 Let $\tau^{a}$ be the diagonal matrix with elements

$$
\begin{equation*}
\tau_{k, k}^{a}=\frac{1}{4 \pi} \int_{0}^{\infty} L_{k}\left(t^{2} / 2\right) e^{-t^{2} / 4} \frac{t^{2}}{1+a^{-1} \sinh ^{2} \beta t} d t \tag{4.6}
\end{equation*}
$$

Then

$$
\begin{equation*}
\tau_{k k}^{a}=O\left(k^{-5 / 4}(\log a)^{4}\right) \tag{4.7}
\end{equation*}
$$

Proof. We analyze first the dependance on $a$ for a fixed $k$. The functions $\left\{L_{k}(u) e^{-u / 2}\right\}_{k \geq 0}$ form a orthonormal basis of $L_{2}(\mathbb{R})$. By using Cauchy-Schwarz followed by Lemma 5 from [5] we get

$$
\left|\tau_{k, k}^{a}\right| \leq \frac{1}{4 \pi}\left(\int_{0}^{\infty} \frac{t^{3}}{\left(1+a^{-1} \sinh ^{2} \beta t\right)^{2}} d t\right)^{1 / 2} \leq C(\log a)^{3 / 2}
$$

for some positive constant $C$.

Let now $a$ be fixed and look at the asymptotic behavior of $\tau_{k, k}^{a}$ as $k \rightarrow \infty$. We use the differential equation of the Laguerre polynomials, [13] 8.979:

$$
L_{n}(x)=\frac{1}{n}\left((x-1) L_{n}^{\prime}(x)-x L_{n}^{\prime \prime}(x)\right) .
$$

Thus

$$
\begin{align*}
& \frac{d}{d t} L_{n}\left(t^{2} / 2\right)=t L_{n}^{\prime}\left(t^{2} / 2\right)  \tag{4.8}\\
& \frac{d^{2}}{d t^{2}} L_{n}\left(t^{2} / 2\right)=L_{n}^{\prime}\left(t^{2} / 2\right)+t^{2} L_{n}^{\prime \prime}\left(t^{2} / 2\right) \tag{4.9}
\end{align*}
$$

which implies

$$
\frac{t^{2}}{2} L_{n}^{\prime \prime}\left(t^{2} / 2\right)=\frac{1}{2} \frac{d^{2}}{d t^{2}} L_{n}\left(t^{2} / 2\right)-\frac{1}{2} t^{-1} \frac{d}{d t} L_{n}\left(t^{2} / 2\right)
$$

and

$$
L_{n}\left(t^{2} / 2\right)=\frac{1}{2 n}\left(\left(t^{2}-1\right) t^{-1} \frac{d}{d t} L_{n}\left(t^{2} / 2\right)-\frac{d^{2}}{d t^{2}} L_{n}\left(t^{2} / 2\right)\right) .
$$

Using integration by parts we obtain the formula

$$
\frac{1}{4 \pi} \int_{0}^{\infty} L_{k}\left(t^{2} / 2\right) e^{-t^{2} / 4} \frac{t^{2}}{1+a^{-1} \sinh ^{2} \beta t} d t=\frac{1}{k} \int_{0}^{\infty} L_{k}\left(t^{2} / 2\right) e^{-t^{2} / 4} f(t) d t
$$

where the function $f$ is given by

$$
\frac{P_{1}(t)}{1+a^{-1} \sinh ^{2} \beta t}+\frac{a^{-1}\left(P_{2}(t) \sinh 2 \beta t+P_{3}(t) \cosh 2 \beta t\right)}{\left(1+a^{-1} \sinh ^{2} \beta t\right)^{2}}+\frac{P_{3}(t) a^{-2} \sinh ^{2} 2 \beta t}{\left(1+a^{-1} \sinh ^{2} \beta t\right)^{3}} .
$$

with $P_{i}(t)$ polynomials with degree at most four, whose coefficients do not depend on $a$.
We split the integral into $\int_{0}^{1}$ and $\int_{1}^{\infty}$ and use the following bounds for the behavior of Laguerre polynomials in the two intervals (see [22] Theorem 8.9.12 and Theorem 7.6.4):

$$
\max _{x \in[1, \infty)} e^{-x / 2}\left|L_{n}(x)\right|=O\left(n^{-1 / 4}\right)
$$

and

$$
L_{n}(x)=x^{-1 / 4} O\left(n^{-1 / 4}\right)
$$

uniformly on $(0,1]$. Thus using Lemma 5 of [5], we get that as $k, a \rightarrow \infty$,

$$
\left|\frac{1}{4 \pi} \int_{0}^{\infty} L_{k}\left(t^{2} / 2\right) e^{-t^{2} / 4} \frac{t^{2}}{1+a^{-1} \sinh ^{2} \beta t} d t\right|=O\left((\log a)^{4} k^{-5 / 4}\right)
$$

Lemma 5 For any $(\beta, L)$ there exists an $\alpha>0$ such that $W_{\alpha}$ belongs to the class $\mathcal{W}(\beta, L)$.
Proof. By using Minkowski inequality we get

$$
\begin{aligned}
& \int e^{2 \beta|w|}\left|\widetilde{W}_{\alpha}(w)\right|^{2} d w=2 \pi \int_{0}^{\infty}\left|\int_{0}^{1} \sqrt{r} \frac{f_{\alpha}(z)}{1-z} \exp \left(-r^{2} \frac{1+z}{4(1-z)}+\beta r\right) d z\right|^{2} d r \leq \\
& {\left[\int_{0}^{1}\left(\int_{0}^{\infty} r \exp \left(-r^{2} \frac{1+z}{2(1-z)}+2 \beta r\right) d r\right)^{1 / 2} \frac{f_{\alpha}(z)}{1-z} d z\right]^{2}}
\end{aligned}
$$

The interior integral satisfies the bound

$$
\int_{0}^{\infty} r \exp \left(-r^{2} \frac{1+z}{2(1-z)}+2 \beta r\right) d r \leq C(\beta)(1-z)
$$

for some positive constant $C(\beta)$. Thus

$$
\int e^{2 \beta|w|}\left|\widetilde{W}_{\alpha}(w)\right|^{2} d w \leq C(\beta)\left(\int_{0}^{1} \alpha(1-z)^{\alpha-1 / 2}\right)^{2}=C(\beta)\left(\frac{\alpha}{\alpha+1 / 2}\right)^{2} \rightarrow 0
$$

as $\alpha \rightarrow 0$.

Lemma 6 For $\alpha \leq 1 / 2$, the Fisher information of the family of densities $\mathcal{R} W_{c}^{z}$ satisfies

$$
\begin{equation*}
I(c)=C^{*}(\log a)^{3} R_{0}^{-1}(1+o(1)) \tag{4.10}
\end{equation*}
$$

where $R_{0}$ is defined in (3.38).
Proof. We sketch the proof following the line of [5] and pointing out where the differences appear. After some transformations, the Fisher information of the family can be brought to the form

$$
\begin{equation*}
I(c)=\frac{1}{\pi} \int_{0}^{\pi} d \phi \int \frac{H_{a}^{2}(u)}{\mathcal{R}\left[W_{c}^{z}\right]([z, \phi]-u, \phi)} d u \tag{4.11}
\end{equation*}
$$

By expanding $\mathcal{R}\left[W_{c}^{z}\right]([z, \phi]-u, \phi)^{-1}$ up to the second order and bounding the second derivative one can show that

$$
\begin{align*}
& \left|\int \frac{H_{a}^{2}(u)}{\mathcal{R}\left[W_{c}^{z}\right][(z, \phi]-u, \phi)} d u-\frac{1}{\mathcal{R}\left[W_{c}^{z}\right]([z, \phi], \phi)} \int H_{a}^{2}(u) d u\right| \\
& =O\left(\int u^{2} H_{a}^{2}(u) d u\right)=O\left((\log a)^{2}\right), \quad \text { as } \quad a \rightarrow \infty . \tag{4.12}
\end{align*}
$$

Recall that $\mathcal{R}\left[W_{c}^{z}\right]([z, \phi], \phi)=\mathcal{R}\left[W^{c}\right](0, \phi)$ does not depend on $\phi$, cf. (3.37). Thus we can write

$$
I(c)=\frac{1}{\mathcal{R}\left[W_{c}\right](0, \phi)} \int H_{a}^{2}(u) d u+O\left((\log a)^{2}\right)
$$

From [5] we have

$$
\begin{equation*}
\int H_{a}^{2}(u) d u=\frac{1}{3 \cdot 2 \cdot(4 \pi \beta)^{2}}(\log a)^{3}(1+o(1)) \quad \text { as } \quad a \rightarrow \infty \tag{4.13}
\end{equation*}
$$

From equations (4.11)-(4.13) one obtains the desired result.
The main difference in the proof appears when deriving the bound (4.12). To derive this, one needs a bound of the absolute value of the second derivative

$$
\left|\frac{\partial^{2}}{\partial t^{2}} \frac{1}{\mathcal{R}\left[W_{c}^{z}\right](t, \phi)}\right|_{t=[z, \phi]-u}=O(1), \quad \text { as } \quad a \rightarrow \infty
$$

which is uniform in $u$. We have

$$
\begin{equation*}
\left|\frac{\partial^{2}}{\partial t^{2}} \frac{1}{\mathcal{R}\left[W_{c}^{z}\right](t, \phi)}\right|_{t=[z, \phi]-u} \leq\left|\frac{p_{c}^{\prime \prime}(u, \phi)}{p_{c}(u, \phi)^{2}}\right|+\left|\frac{p_{c}^{\prime}(u, \phi)^{2}}{p_{c}(u, \phi)^{3}}\right| \tag{4.14}
\end{equation*}
$$

where

$$
p_{c}(u, \phi)=\mathcal{R}\left[W_{c}^{z}\right]([z, \phi]-u, \phi)=\mathcal{R}\left[W_{c}\right](-u, \phi)=\mathcal{R}\left[W_{c}\right](u, \phi)=p_{\alpha}(u, \phi)+c H_{a}(u),
$$

and the derivatives are with respect to the first argument. If $u$ is in a compact interval, the two terms on the right hand side of (4.14) are $O(1)$ as $a \rightarrow \infty$ since in that case $\sup _{|c|<C_{a}}\left|c H_{a}(u)\right| \rightarrow 0$, and similarly for its derivatives. Outside the interval, we have that $H_{a}$ and its derivatives are exponentially decreasing, and the dominating term is $p_{\alpha}$. Indeed, according to Lemma $3, p_{\alpha}$ is of order $|u|^{-(1+2 \alpha)}$, and its first two derivatives are bounded from above by $|u|^{-(2+2 \alpha)}$ and $|u|^{-(3+2 \alpha)}$ respectively.

## References

[1] Artiles L., Gill, R. D., Guţă, M., An invitation to Quantum Tomography, J. R. Statist. Soc. B: Methodological, 3, (2005), 22-58.
[2] Barndorff-Nielsen O. E., Gill R. D., Jupp P. E., On quantum statistical inference (with discussion), J. Royal Stat. Soc. B, 65, (2003), 775-816.
[3] Breitenbach G., Schiller S., Mlynek J., Measurement of the quantum state of the squeezed light, Nature, 387, (1997), 471-475.
[4] Butucea C., Guţă, M., Artiles, L., Minimax and adaptive estimation of the Wigner function in quantum homodyne tomography with noisy data (2006), accepted by Ann. Statist.
[5] Cavalier L., Efficient estimation of a density in a problem of tomography, Ann. Statist., 28, (2000), 630-647.
[6] Cavalier L., On the problem of local adaptive estimation in tomography, Bernoulli, 7, (2001), 6378.
[7] Deans S. R., The Radon transform and some of its applications, John Wiley \& Sons, New York (1983).
[8] Erdélyi, A.Higher Transcendental Functions, volume 2, McGraw-Hill (1953).
[9] Gill R. D., Levit B. Y., Applications of the van Trees inequality: a Bayesian Cramer-Rao bound, Bernoulli, 1, (1995), 59-79.
[10] Glauber R. J., Coherent and incoherent states of the radiation field, Phys. Rev., 131, (1963), 2766-2788.
[11] Golubev Y. K., Levit B. Y., Tsybakov A. B., Asymptotically efficient estimation of analytic functions in Gaussian noise, Bernoulli, 2, (1996), 167-181.
[12] Golubev Y. K., Levit, B. Y., Asymptotic efficient estimation for analytical distributions, Math. Methods Statist., 5, (1996), 357-368.
[13] Gradshteyn I. S., Ryzhik I. M., Table of Integrals, Series and Products, Academic Press, New York (1980).
[14] Helstrom C., Quantum Detection and Estimation Theory, Academic Press, New York (1976).
[15] Holevo A. S., Probabilistic and Statistical Aspects of Quantum Theory, North-Holland (1982).
[16] Ibragimov I. A., Hasminkii, R. Z., Estimation of a distribution density, J. Soviet. Math., 25, (1983), 40-57.
[17] Johnstone I. M., Silverman, B. W., Speed of estimation in positron emission tomography and related inverse problems, Ann. Statist., 18, (1990), 251-280.
[18] Korostelëv A. P., Tsybakov A. B., Minimax theory of image reconstruction, volume 82 of Lecture Notes in Statistics, Springer-Verlag, New York (1993).
[19] Leonhardt U., Measuring the Quantum State of Light, Cambridge University Press (1997).
[20] Natterer F., The Mathematics of Computerized Tomography, Wiley, Chichester (1986).
[21] Ourjoumtsev A., Tualle-Brouri R., Laurat J., Grangier P., Generating optical Schrödinger kittens for Quantum Information Processing, Science, 312, (2006), 83.
[22] Szegö G., Orthogonal Polynomials, Americam Mathematical Society, Providence (1975).
[23] van Trees H. L., Detection, Estimation and Modulation Theory, volume 1, Wiley, New-York (1968).
[24] Vogel K., Risken, H., Determination of quasiprobability distributions in terms of probability distributions for the rotated quadrature phase, Phys. Rev. A, 40, (1989), 2847-2849.
[25] Wigner E., On the quantum correction for thermodynamic equations, Phys. Rev., 40, (1932), 749-759.
[26] Zavatta A., Viciani S., Bellini M., Quantum to classical transition with single-photon-added coherent states of light, Science, 306, (2004), 660-662.


[^0]:    *School of Mathematical Sciences, University of Nottingham, University Park, NG7 2RD Nottingham, U.K., madalin.guta@nottingham.ac.uk
    ${ }^{\dagger}$ Eurandom, University of Eindhoven, P.O. Box 513, 5600 MB Eindhoven, The Netherlands, artiles@eurandom.tue.nl
    ${ }^{1}$ AMS 1991 subject classifications. Primary 62G05, 62G20; secondary 62C20.
    Key words and phrases. Nonparametric statistics, Wigner function, Quantum Homodyne Tomography, Kernel estimator, minimax efficient rates

