# Nonparametric estimation of the expected discounted penalty function in the compound Poisson model 

Florian Dussap

## To cite this version:

Florian Dussap. Nonparametric estimation of the expected discounted penalty function in the compound Poisson model. 2021. hal-03212783v1

HAL Id: hal-03212783
https://hal.science/hal-03212783v1
Preprint submitted on 30 Apr 2021 (v1), last revised 17 Jan 2022 (v2)

HAL is a multi-disciplinary open access archive for the deposit and dissemination of scientific research documents, whether they are published or not. The documents may come from teaching and research institutions in France or abroad, or from public or private research centers.

L'archive ouverte pluridisciplinaire HAL, est destinée au dépôt et à la diffusion de documents scientifiques de niveau recherche, publiés ou non, émanant des établissements d'enseignement et de recherche français ou étrangers, des laboratoires publics ou privés.

# Nonparametric estimation of the expected discounted penalty function in the compound Poisson model 

Florian Dussap*<br>Université de Paris, CNRS, MAP5 UMR 8145, F-75006 Paris, France


#### Abstract

We propose a nonparametric estimator of the expected discounted penalty function function in the compound Poisson risk model. We use a projection estimator on the Laguerre basis and we compute the coefficients using Plancherel theorem. We provide an upper bound on the MISE of our estimator, and we show it achieves parametric rates of convergence without needing a bias-variance compromise. Moreover, we compare our estimator with the Laguerre deconvolution method. We compute an upper bound of the MISE of the Laguerre deconvolution estimator and we compare it on Sobolev-Laguerre spaces with our estimator. Finally, we compare these estimators on simulated data.

MSC-2020: 62G05, 62P05, 91G70. Keywords: nonparametric estimator, rate of estimation, classical risk model, Gerber-Shiu function, Laguerre series expansion, Fourier inversion.


## 1 Introduction

### 1.1 The statistical problem

We consider the classical risk model (compound Poisson model) for the risk reserve process $\left(U_{t}\right)_{t \geqslant 0}$ of an insurance company:

$$
\begin{equation*}
U_{t}=u+c t-\sum_{i=1}^{N_{t}} X_{i}, \quad t \geqslant 0 \tag{1}
\end{equation*}
$$

where $u \geqslant 0$ is the initial capital; $c>0$ is the premium rate; the claim number process $\left(N_{t}\right)$ is a homogeneous Poisson process with intensity $\lambda$; the claim sizes $\left(X_{i}\right)$ are positive and i.i.d. with density $f$ and mean $\mu$, independent of $\left(N_{t}\right)$. We denote by $\tau(u)$ the ruin time:

$$
\tau(u):=\inf \left\{t \geqslant 0 \mid \sum_{i=1}^{N_{t}} X_{i}-c t>u\right\} \in \mathbb{R}_{+} \cup\{+\infty\}
$$

and we make the following assumption to ensure that $\tau(u)$ is not almost surely finite.
Assumption 1 (safety loading condition). Let $\theta:=\frac{\lambda \mu}{c}$, we assume that $\theta<1$.
To study simultaneously the ruin time, the deficit at ruin, and the surplus level before the ruin, Gerber and Shiu 1998 introduced the function:

$$
\begin{equation*}
\phi(u):=\mathbb{E}\left[\mathrm{e}^{-\delta \tau(u)} w\left(U_{\tau(u)-},\left|U_{\tau(u)}\right|\right) \mathbf{1}_{\tau(u)<+\infty}\right], \tag{2}
\end{equation*}
$$

where $\delta \geqslant 0$, and $w$ is a non-negative function of the surplus before the ruin and the deficit at ruin. This function is called the expected discounted penalty function, but it will also be referred to as the Gerber-Shiu function in the following. Several quantities of interest can be put in the form (2), for example:

[^0]1. if $\delta=0$ and $w(x, y)=1$, then $\phi(u)$ is the probability of ruin;
2. if $\delta>0$ and $w(x, y)=1$, then $\phi(u)$ is the Laplace transform of $\tau(u)$;
3. if $\delta=0$ and $w(x, y)=x+y$, then $\phi(u)$ is the expected jump size causing the ruin.

For more information concerning the compound Poisson model and the Gerber-Shiu function, see Asmussen and Albrecher 2010.

Observations and goal In all this article, we suppose that the premium rate $c$ is known, but the parameters of the aggregate claims process are not, that is $(\lambda, \mu, f)$ is unknown. We suppose we have observed the process $\left(U_{t}\right)_{t \geqslant 0}$ during the interval $[0, T]$, with $T>0$ fixed, so we have access to the number of claims and their size. Our goal is to recover the Gerber-Shiu function from the observations $\left(N_{T} ; X_{1}, \ldots, X_{N_{T}}\right)$.

Several authors have considered the problem of estimating the Gerber-Shiu function using nonparametric methods. The first articles had an asymptotic approach: Frees [1986, Croux and Veraverbeke [1990], Pitts 1994, Politis 2003, and Masiello 2014 constructed nonparametric estimators of the ruin probability, and established their consistency and their asymptotic normality.

Concerning non-asymptotic approaches, Mnatsakanov, Ruymgaart, and Ruymgaart 2008 introduced a regularized Laplace inversion method to estimate the ruin probability in the compound Poisson model. Shimizu 2011, 2012] then extended this method to the estimation of the Ger-ber-Shiu function in more general risk models. However, this method suffers from poor rates of convergence, and numerical difficulties to compute the estimator.

In their paper, Zhang and Su [2018] introduced a projection estimator on the Laguerre basis to overcome these drawbacks. The choice of this basis is motivated by the work of Comte, Cuenod, Pensky, and Rozenholc 2017, where the properties of the Laguerre functions relative to the convolution product are used to solve a Laplace deconvolution problem. The same method was then used in more general risk model: Zhang and Su 2019] estimate the Gerber-Shiu function in a Lévy risk model, where the aggregate claims is a pure-jump Lévy process; Su, Yong, and Zhang 2019b estimate the Gerber-Shiu function in the compound Poisson model perturbed by a Brownian motion; and Su, Shi, and Wang 2019a study the model where both the income and the aggregate claims are compound Poisson processes. Recently, Su and Yu 2020 showed the Laguerre projection estimator of the Gerber-Shiu function in the compound Poisson model is pointwise asymptotically normal in the case $\delta=0$.

In this paper, we construct an estimator of the Gerber-Shiu function (2) in the compound Poisson model (11). This estimator is also a projection estimator on the Laguerre basis, but we compute the coefficients using Plancherel theorem instead of using a Laguerre deconvolution method. We emphasize that our estimator achieves parametric rates of convergence for the MISE regardless of the regularity of the Gerber-Shiu function, and without needing to find a compromise between the bias and the variance.

We also improve the previous results concerning the Laguerre deconvolution method. Previous rates were given in $\mathcal{O}_{\mathrm{P}}$, and we propose a non asymptotic bound on the MISE (Mean Integrated Squared Error) of the estimator. To achieve this goal, we introduce two modified versions of the Laguerre deconvolution estimator: the first one depends on a truncation parameter, whereas the second one does not but it is only defined in the case $\delta=0$.

To control the MISE of the second version of the Laguerre deconvolution estimator, we had to prove that the primitives of the Laguerre functions were uniformly bounded (see Lemma 3.3). This result is interesting in itself, but the proof is based on the study of the properties of the ODE's satisfied by Laguerre polynomials, and has little to do with this article's main subject. The interested reader can find all the details in Appendix B

Outline of the paper In the remaining part of this section, we introduce the notations and we give preliminary results on the Gerber-Shiu function. In Section 2 we construct our estimator and we study its MISE. In Section 3, we introduce two modified version of the Laguerre deconvolution estimator and we study their MISE. In Section 4 we compute convergence rates of the different
estimators we considered on Sobolev-Laguerre spaces and also in the case where the claim sizes are exponentially distributed. In Section 5, we compare numerically the estimators on simulated data. We gathered all the proofs in Section 7

### 1.2 Notations and preliminaries on the Gerber-Shiu function

We use the following notations in the paper:

- $\mathbb{N}:=\{0,1,2,3, \ldots\}, \mathbb{N}^{*}:=\mathbb{N} \backslash\{0\}, \mathbb{R}_{+}:=[0,+\infty)$, and $\mathbb{T}:=\{z \in \mathbb{C}| | z \mid=1\}$.
- " $x \lesssim y$ " means $x \leqslant C y$ for an absolute constant $C>0$.
- $x \wedge y:=\min (x, y)$ and $x \vee y:=\max (x, y)$.
- $\mathcal{L} f(s):=\int_{0}^{+\infty} \mathrm{e}^{-s x} f(x) \mathrm{d} x$ is the Laplace transform of $f$.
- $\mathcal{F} \psi(\omega):=\int_{\mathbb{R}} \mathrm{e}^{\mathrm{i} \omega x} \psi(x) \mathrm{d} x$ is the Fourier transform of $\psi$.
- $\operatorname{Leb}(A)$ is the Lebesgue measure of the set $A$.
- $\left\|\mathbf{A}_{m}\right\|_{\mathrm{op}}:=\sup _{x \in \mathbb{R}^{m} \backslash\{0\}} \frac{\left\|\mathbf{A}_{m} x\right\|_{\ell^{2}}}{\|x\|_{\ell^{2}}}$ is the $\ell^{2}$-operator norm of the matrix $\mathbf{A}_{m} \in \mathbb{R}^{m \times m}$.

The key result to estimate the Gerber-Shiu function is the following theorem.
Theorem 1.1 Gerber and Shiu (1998). Under Assumption 1. the Gerber-Shiu function satisfies the equation:

$$
\begin{equation*}
\phi=\phi * g+h . \tag{3}
\end{equation*}
$$

where $g$ and $h$ are given by:

$$
g(x):=\frac{\lambda}{c} \int_{x}^{+\infty} \mathrm{e}^{-\rho_{\delta}(y-x)} f(y) \mathrm{d} y, \quad h(u):=\frac{\lambda}{c} \int_{u}^{+\infty} \mathrm{e}^{-\rho_{\delta}(x-u)}\left(\int_{x}^{+\infty} w(x, y-x) f(y) \mathrm{d} y\right) \mathrm{d} x(4)
$$

and where $\rho_{\delta}$ is the (unique) non-negative solution of the so-called Lundberg equation:

$$
\begin{equation*}
c s-\lambda(1-\mathcal{L} f(s))=\delta \tag{5}
\end{equation*}
$$

Remark 1.2. Since $\mathcal{L} f(s) \in[0,1]$, it is easy to see that $\rho_{\delta} \in\left[\frac{\delta}{c}, \frac{\delta+\lambda}{c}\right]$. Moreover, we know that $\rho_{\delta}=0$ when $\delta=0$.

We need to ensure that $\phi, g$ and $h$ belong to $\mathrm{L}^{2}\left(\mathbb{R}_{+}\right)$in order to use a projection estimator. We see that $\sup _{x} g(x) \leqslant \sup _{x} \frac{\lambda}{c} \mathbb{P}[X>x] \leqslant \frac{\lambda}{c}$ and $\int_{0}^{\infty} g(x) \mathrm{d} x \leqslant \frac{\lambda}{c} \int_{0}^{\infty} \mathbb{P}[X>x] \mathrm{d} x=\theta$, hence $g \in \mathrm{~L}^{1}\left(\mathbb{R}_{+}\right) \cap \mathrm{L}^{\infty}\left(\mathbb{R}_{+}\right)$, therefore $g \in \mathrm{~L}^{2}\left(\mathbb{R}_{+}\right)$. To ensure that $h \in \mathrm{~L}^{c} \int^{c}\left(\mathbb{R}_{+}\right)$we make the following assumption.

Assumption 2. $\int_{0}^{\infty}(1+x)\left(\int_{x}^{\infty} w(x, y-x) f(y) \mathrm{d} y\right) \mathrm{d} x$ is finite.
Then:

$$
\begin{aligned}
\sup _{u>0} h(u) & \leqslant \frac{\lambda}{c} \int_{0}^{\infty}\left(\int_{x}^{\infty} w(x, y-x) f(y) \mathrm{d} y\right) \mathrm{d} x \\
\int_{0}^{\infty} h(u) \mathrm{d} u & \leqslant \frac{\lambda}{c} \int_{0}^{\infty} x\left(\int_{x}^{\infty} w(x, y-x) f(y) \mathrm{d} y\right) \mathrm{d} x
\end{aligned}
$$

Hence, $h$ belongs to $L^{1}\left(\mathbb{R}_{+}\right) \cap \mathrm{L}^{\infty}\left(\mathbb{R}_{+}\right)$, so $h \in \mathrm{~L}^{2}\left(\mathbb{R}_{+}\right)$. Integrating Equation (3) yields:

$$
\|\phi\|_{\mathrm{L}^{1}}=\int_{0}^{\infty} \phi(u) \mathrm{d} u=\frac{\int_{0}^{\infty} h(u) \mathrm{d} u}{1-\int_{0}^{\infty} g(x) \mathrm{d} x}
$$

which is finite under Assumption 1 since $\int_{0}^{\infty} g(x) \mathrm{d} x \leqslant \theta<1$. Since $\phi$ belongs to $\mathrm{L}^{1}\left(\mathbb{R}_{+}\right)$and $g$ belongs to $\mathrm{L}^{2}\left(\mathbb{R}_{+}\right)$, their convolution product belongs to $\mathrm{L}^{2}\left(\mathbb{R}_{+}\right)$, hence $\phi=\phi * g+h$ belongs to $\mathrm{L}^{2}\left(\mathbb{R}_{+}\right)$as well.

## 2 The Laguerre-Fourier estimator

We use the Laguerre functions $\left(\psi_{k}\right)_{k \in \mathbb{N}}$ as an orthonormal basis of $\mathrm{L}^{2}\left(\mathbb{R}_{+}\right)$:

$$
\begin{equation*}
\forall x \in \mathbb{R}_{+}, \psi_{k}(x):=\sqrt{2} L_{k}(2 x) \mathrm{e}^{-x}, \quad L_{k}(x):=\sum_{j=0}^{k}\binom{k}{j} \frac{(-x)^{j}}{j!} . \tag{6}
\end{equation*}
$$

We choose this basis for several reasons: the support of the Laguerre functions is $\mathbb{R}_{+}$, which is well suited since the functions we want to estimate are defined on $\mathbb{R}_{+}$; exponential functions (and more broadly mixture of gamma functions, see the proof of Lemma 3.9 in Mabon 2017]) have an exponentially small bias in this basis, which is interesting since the claim sizes distribution is often chosen in these families of distribution; and the Fourier transform of the Laguerre function is known explicitly:

$$
\forall \omega \in \mathbb{R}, \quad \mathcal{F} \psi_{k}(\omega)=(-1)^{k} \sqrt{2} \frac{(1+\mathrm{i} \omega)^{k}}{(1-\mathrm{i} \omega)^{k+1}}
$$

We denote by $\left(a_{k}\right)_{k \geqslant 0}$ the Laguerre coefficients of $\phi$. If $m \in \mathbb{N}^{*}$, we denote by $\phi_{m}$ the projection of $\phi$ on the subspace of $L^{2}\left(\mathbb{R}_{+}\right)$spanned by the first $m$ Laguerre functions $\psi_{0}, \ldots, \psi_{m-1}$, that is:

$$
\phi_{m}:=\sum_{k=0}^{m-1} a_{k} \psi_{k}, \text { with } a_{k}=\left\langle\phi, \psi_{k}\right\rangle .
$$

The Laguerre coefficients of $\phi$ can be computed using Plancherel theorem:

$$
a_{k}=\left\langle\phi, \psi_{k}\right\rangle=\frac{1}{2 \pi}\left\langle\mathcal{F} \phi, \mathcal{F} \psi_{k}\right\rangle .
$$

Taking the Fourier transform in equation (3), we see that $\mathcal{F} \phi=\frac{\mathcal{F} h}{1-\mathcal{F g}}$. Let $\hat{g}, \hat{h} \in \mathrm{~L}^{2}\left(\mathbb{R}_{+}\right)$be some estimators of $g$ and $h$ (we provide these estimators later in equation 12), we estimate the coefficients of $\phi$ by:

$$
\begin{equation*}
\hat{a}_{k}:=\frac{1}{2 \pi}\left\langle\frac{\mathcal{F} \hat{h}}{1-\widehat{\mathcal{F} g}}, \mathcal{F} \psi_{k}\right\rangle, \tag{7}
\end{equation*}
$$

where $\widehat{\mathcal{F} g}:=(\mathcal{F} \hat{g}) \mathbf{1}_{|\mathcal{F} \hat{g}| \leqslant \theta_{0}}$ for some truncation parameter $\theta_{0}<1$. The estimator of $\phi$ is then:

$$
\hat{\phi}_{m_{1}}:=\sum_{k=0}^{m_{1}-1} \hat{a}_{k} \psi_{k}
$$

where $m_{1}$ is the dimension of the projection space.
Proposition 2.1. Under Assumptions 1 and 2, if $\theta<\theta_{0}$, we have:

$$
\left\|\phi-\hat{\phi}_{m_{1}}\right\|_{\mathrm{L}^{2}}^{2} \leqslant\left\|\phi-\phi_{m_{1}}\right\|_{\mathrm{L}^{2}}^{2}+\frac{2}{\left(1-\theta_{0}\right)^{2}}\|h-\hat{h}\|_{\mathrm{L}^{2}}^{2}+\frac{2\|h\|_{\mathrm{L}^{1}}^{2}}{\left(1-\theta_{0}\right)^{2}(1-\theta)^{2}}\left(1+\frac{\|g\|_{\mathrm{L}^{1}}^{2}}{\left(\theta_{0}-\theta\right)^{2}}\right)\|g-\hat{g}\|_{\mathrm{L}^{2}}^{2} .
$$

Remark 2.2. We emphasize the fact that this result is proven using only two properties: the function $\phi$ satisfies the equation (3) and $\theta_{0}>\theta>\|g\|_{\mathrm{L}^{1}}$. Hence, it can be applied to other problems where the target function satisfies an equation of the form (3). For example, it is the case in Zhang and Su 2019, Su et al. 2019a and Su et al. 2019b.

We now need to provide good estimators of $g$ and $h$. We choose to estimate them by projection on the Laguerre basis too. Let $\left(b_{k}\right)_{k \geqslant 0}$ and $\left(c_{k}\right)_{k \geqslant 0}$ be the coefficients of $g$ and $h$, that is:

$$
\begin{array}{ll}
g=\sum_{k=0}^{+\infty} b_{k} \psi_{k}, & \text { with } \quad b_{k}:=\left\langle g, \psi_{k}\right\rangle, \\
h=\sum_{k=0}^{+\infty} c_{k} \psi_{k}, & \text { with } \quad c_{k}:=\left\langle h, \psi_{k}\right\rangle . \tag{9}
\end{array}
$$

By Fubini's theorem and using equation (4):

$$
\begin{aligned}
b_{k} & =\int_{0}^{+\infty} g(x) \psi_{k}(x) \mathrm{d} x \\
& =\frac{\lambda}{c} \int_{0}^{+\infty}\left(\int_{x}^{+\infty} \mathrm{e}^{-\rho_{\delta}(y-x)} f(y) \mathrm{d} y\right) \psi_{k}(x) \mathrm{d} x \\
& =\frac{\lambda}{c} \int_{0}^{+\infty}\left(\int_{0}^{y} \mathrm{e}^{-\rho_{\delta}(y-x)} \psi_{k}(x) \mathrm{d} x\right) f(y) \mathrm{d} y \\
& =\frac{\lambda}{c} \mathbb{E}\left[\int_{0}^{X} \mathrm{e}^{-\rho_{\delta}(X-x)} \psi_{k}(x) \mathrm{d} x\right]
\end{aligned}
$$

The same calculation for $c_{k}$ yields:

$$
c_{k}=\frac{\lambda}{c} \mathbb{E}\left[\int_{0}^{X}\left(\int_{u}^{X} \mathrm{e}^{-\rho_{\delta}(x-u)} w(x, X-x) \mathrm{d} x\right) \psi_{k}(u) \mathrm{d} u\right] .
$$

We estimate these coefficients by empirical means. However, we first need to estimate $\rho_{\delta}$. Since $\rho_{\delta}$ is the non-negative solution of the Lundberg equation (5), we estimate it by $\hat{\rho}_{\delta}$ the non-negative solution of the empirical Lundberg equation:

$$
c s-\hat{\lambda}(1-\widehat{\mathcal{L} f}(s))=\delta,
$$

where $\hat{\lambda}:=\frac{N_{T}}{T}$ and $\widehat{\mathcal{L} f}(s):=\frac{1}{N_{T}} \sum_{i=1}^{N_{T}} \mathrm{e}^{-s X_{i}}$. When $\delta=0$ we know that $\rho_{\delta}=0$ so we do not need to estimate it, thus we set $\hat{\rho}_{0}=0$. The estimated coefficients of $g$ and $h$ are:

$$
\begin{align*}
& \hat{b}_{k}=\frac{1}{c T} \sum_{i=1}^{N_{T}} \int_{0}^{X_{i}} \mathrm{e}^{-\hat{\rho}_{\delta}\left(X_{i}-x\right)} \psi_{k}(x) \mathrm{d} x  \tag{10}\\
& \hat{c}_{k}=\frac{1}{c T} \sum_{i=1}^{N_{T}} \int_{0}^{X_{i}}\left(\int_{u}^{X_{i}} \mathrm{e}^{-\hat{\rho}_{\delta}(x-u)} w\left(x, X_{i}-x\right) \mathrm{d} x\right) \psi_{k}(u) \mathrm{d} u \tag{11}
\end{align*}
$$

and the estimators of $g$ and $h$ are:

$$
\begin{equation*}
\hat{g}_{m_{2}}:=\sum_{k=0}^{m_{2}-1} \hat{b}_{k} \psi_{k}, \quad \hat{h}_{m_{3}}:=\sum_{k=0}^{m_{3}-1} \hat{c}_{k} \psi_{k} \tag{12}
\end{equation*}
$$

where $m_{2}$ and $m_{3}$ are the dimensions of the projection spaces. As we did for $\phi$, we denote by $g_{m_{2}}$ and $h_{m_{3}}$ the projections of $g$ and $h$ on the subspaces $\operatorname{Span}\left(\psi_{0}, \ldots, \psi_{m_{2}-1}\right)$ and $\operatorname{Span}\left(\psi_{0}, \ldots, \psi_{m_{3}-1}\right)$.
Remark 2.3. The dimensions $m_{1}, m_{2}, m_{3}$ do not have to be the same for the estimation of $\phi, g$ and $h$. In practice, we will choose different dimensions.

In order to give a bound on the mean integrated squared error of our estimators $\hat{g}_{m_{2}}$ and $\hat{h}_{m_{3}}$, we need to make an additional assumption.
Assumption 3. Let $W(X):=\int_{0}^{X}\left(\int_{u}^{X} w(x, X-x) \mathrm{d} x\right)^{2} \mathrm{~d} u$. If $\delta=0$, we assume that $\mathbb{E}[W(X)]$ is finite, and if $\delta>0$, we assume that $\mathbb{E}\left[W(X)^{2}\right]$ is finite.
Remark 2.4. Assume that $w(x, y) \lesssim(1+x)^{k_{1}}(1+y)^{k_{2}}$ for some $k_{1}, k_{2} \geqslant 0$. Then Assumption 2 holds if $X$ has a finite moment of order $2+\left(k_{1}+k_{2}\right)$ and Assumption 3 holds if $X$ has a finite moment of order $3+2\left(k_{1}+k_{2}\right)$ if $\delta=0$, and a finite moment of order $6+4\left(k_{1}+k_{2}\right)$ if $\delta>0$.

Theorem 2.5. Under Assumptions 1, 2 and 3, if $\delta=0$ then it holds:

$$
\begin{aligned}
& \mathbb{E}\left\|g-\hat{g}_{m_{2}}\right\|_{\mathrm{L}^{2}}^{2} \leqslant\left\|g-g_{m_{2}}\right\|_{\mathrm{L}^{2}}^{2}+\frac{\lambda}{c^{2} T} \mathbb{E}[X] \\
& \mathbb{E}\left\|h-\hat{h}_{m_{3}}\right\|_{\mathrm{L}^{2}}^{2} \leqslant\left\|h-h_{m_{3}}\right\|_{\mathrm{L}^{2}}^{2}+\frac{\lambda}{c^{2} T} \mathbb{E}[W(X)]
\end{aligned}
$$

and if $\delta>0$ then it holds:

$$
\begin{aligned}
& \mathbb{E}\left\|g-\hat{g}_{m_{2}}\right\|_{\mathrm{L}^{2}}^{2} \leqslant\left\|g-g_{m_{2}}\right\|_{\mathrm{L}^{2}}^{2}+\frac{C(\lambda)}{c^{2} T}\left(\mathbb{E}[X]+\frac{\mathbb{E}\left[X^{2}\right]^{\frac{1}{2}}}{(1-\theta)^{2} \delta^{2}}\right) \\
& \mathbb{E}\left\|h-\hat{h}_{m_{3}}\right\|_{\mathrm{L}^{2}}^{2} \leqslant\left\|h-h_{m_{3}}\right\|_{\mathrm{L}^{2}}^{2}+\frac{C(\lambda)}{c^{2} T}\left(\mathbb{E}[W(X)]+\frac{\mathbb{E}\left[W(X)^{2}\right]^{\frac{1}{2}}}{(1-\theta)^{2} \delta^{2}}\right),
\end{aligned}
$$

where $C(\lambda)$ is a $\mathcal{O}\left(\lambda^{2}\right)$.
Remark 2.6. The variance terms do not depend on $m_{2}$ and $m_{3}$, so no compromise between the bias and the variance is needed: we just have to take $m_{2}$ and $m_{3}$ as large as possible such that the bias is smaller than $1 / T$. See Section 4 for a discussion concerning the choice of $m_{2}$ and $m_{3}$ when the functions $g$ and $h$ belong to a Sobolev-Laguerre space.

Let $m_{1}, m_{2}, m_{3} \in \mathbb{N}^{*}$, we estimate $g$ by $\hat{g}_{m_{2}}$ and $h$ by $\hat{h}_{m_{3}}$. We plug these estimators in (7) and we estimate $\phi$ by:

$$
\hat{\phi}_{m_{1}, m_{2}, m_{3}}:=\sum_{k=0}^{m_{1}-1}\left\langle\frac{\mathcal{F} \hat{h}_{m_{3}}}{1-\widehat{\mathcal{F} g_{m_{2}}}}, \psi_{k}\right\rangle \psi_{k}
$$

with $\widehat{\mathcal{F} g_{m_{2}}}:=\mathcal{F} \hat{g}_{m_{2}} \mathbf{1}_{\left|\mathcal{F} \hat{g}_{m_{2}}\right| \leqslant \theta_{0}}$. Combining Proposition 2.1 with Theorem 2.5, we obtain:
Corollary 2.7. Under Assumptions 1, 2 and 3, if $\theta<\theta_{0}$ then it holds:

$$
\mathbb{E}\left\|\phi-\hat{\phi}_{m_{1}, m_{2}, m_{3}}\right\|_{\mathrm{L}^{2}}^{2} \leqslant\left\|\phi-\phi_{m_{1}}\right\|_{\mathrm{L}^{2}}^{2}+\frac{C}{\left(1-\theta_{0}\right)^{2}}\left(\left\|g-g_{m_{2}}\right\|_{\mathrm{L}^{2}}^{2}+\left\|h-h_{m_{3}}\right\|_{\mathrm{L}^{2}}^{2}+\frac{1}{c T}\right)
$$

where $C$ is a constant depending on $\lambda, c, \theta,\|g\|_{\mathrm{L}^{1}},\|h\|_{\mathrm{L}^{1}}, \mathbb{E}[X], \mathbb{E}[W(X)]$ and $\theta_{0}-\theta$; and also $\delta$, $\mathbb{E}\left[X^{2}\right], \mathbb{E}\left[W(X)^{2}\right]$ if $\delta>0$.

We want to compare our estimator with the Laguerre deconvolution method. However, there is no result on the MISE of this method for estimating the Gerber-Shiu function, so we study it in the next section.

## 3 The Laguerre deconvolution estimator

For the Laguerre deconvolution method, we need an additional assumption on the coefficients of $g$.

Assumption 4. The sequence $\left(b_{k+1}-b_{k}\right)_{k \geqslant 0}$ belongs to $\ell^{1}(\mathbb{N})$.
The reason why the Laguerre basis is well suited for deconvolution on $\mathbb{R}_{+}$is the following relation satisfied by the Laguerre functions:

$$
\forall k, j \in \mathbb{N}, \quad \psi_{k} * \psi_{j}=\frac{1}{\sqrt{2}}\left(\psi_{k+j}-\psi_{k+j+1}\right)
$$

see formula 22.13.14 in Abramowitz and Stegun 1972. The reader interested in the use of the Laguerre basis for deconvolution problems can read Mabon 2017. Expanding the renewal equation (3) on the Laguerre basis, one easily obtains the following relation between the coefficients of $\phi, g$ and $h$ :

$$
\forall k \in \mathbb{N}, \quad a_{k}=(\beta * a)_{k}+c_{k},
$$

where the sequence $\left(\beta_{k}\right)_{k \geqslant 0}$ is defined by $\beta_{0}:=\frac{b_{0}}{\sqrt{2}}$ and $\beta_{k}:=\frac{b_{k}-b_{k-1}}{\sqrt{2}}$ for $k \geqslant 1$. This relation can be written in a matrix form: if $\mathbf{a}_{m}:=\left(a_{0}, \ldots, a_{m-1}\right)^{T}$ and $\mathbf{c}_{m}:=\left(c_{0}, \ldots, c_{m-1}\right)^{T}$ are the vectors of the $m$ first coefficients of $\phi$ and $h$, then it holds:

$$
\begin{equation*}
\mathbf{A}_{m} \times \mathbf{a}_{m}=\mathbf{c}_{m} \Longleftrightarrow \mathbf{a}_{m}=\mathbf{A}_{m}^{-1} \times \mathbf{c}_{m} \tag{13}
\end{equation*}
$$

where $\mathbf{A}_{m}$ is the lower triangular Toeplitz matrix:

$$
\forall i, j \in\{0, \ldots, m-1\}, \quad\left(\mathbf{A}_{m}\right)_{i, j}:= \begin{cases}1-\frac{1}{\sqrt{2}} b_{0} & \text { if } i=j  \tag{14}\\ \frac{1}{\sqrt{2}}\left(b_{i-j-1}-b_{i-j}\right) & \text { if } i>j \\ 0 & \text { else }\end{cases}
$$

This matrix is invertible if and only if $1-\frac{b_{0}}{\sqrt{2}} \neq 1$, which is the case because $\frac{b_{0}}{\sqrt{2}} \leqslant \theta<1$ under Assumption 1 .

Lemma 3.1. Under Assumption 4, we have $\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}} \leqslant \frac{2}{1-\theta}$ for all $m \in \mathbb{N}^{*}$.
This lemma is borrowed from Zhang and Su 2018 (Lemma 4.3 in their article). There were missing elements in their proof, so we give a new proof of this lemma, for the sake of completeness.

The naive Laguerre deconvolution estimator consists in estimating the matrix $\mathbf{A}_{m}$ and the coefficients $\mathbf{c}_{m}$ in (13), to obtain an estimation of the coefficients of $\phi$. More precisely, the matrix $\mathbf{A}_{m}$ is estimated by plugging $\hat{b}_{k}$, defined by (10), in 14):

$$
\forall i, j \in\{0, \ldots, m-1\}, \quad\left(\hat{\mathbf{A}}_{m}\right)_{i, j}:= \begin{cases}1-\frac{1}{\sqrt{2}} \hat{b}_{0} & \text { if } i=j  \tag{15}\\ \frac{1}{\sqrt{2}}\left(\hat{b}_{i-j-1}-\hat{b}_{i-j}\right) & \text { if } i>j \\ 0 & \text { else }\end{cases}
$$

This matrix is invertible if and only if $\frac{\hat{b}_{0}}{\sqrt{2}} \neq 1$, which is almost surely the case since $\frac{\hat{b}_{0}}{\sqrt{2}}=$ $\frac{1}{c T} \sum_{i=1}^{N_{T}}\left(1-\mathrm{e}^{-X_{i}}\right)$ is a continuous random variable. The coefficients of $\phi$ are estimated by:

$$
\begin{equation*}
\hat{\mathbf{a}}_{m}^{\mathrm{Lag}_{0}}:=\hat{\mathbf{A}}_{m}^{-1} \times \hat{\mathbf{c}}_{m}, \tag{16}
\end{equation*}
$$

where $\hat{\mathbf{c}}_{m}:=\left(\hat{c}_{0}, \ldots, \hat{c}_{m-1}\right)^{T}$. Under Assumptions 1, 2, 3, and 4, Zhang and Su 2018 show that if $\mathbb{E}\left[X^{2}\right]$ is finite and if $m=o(T)$, then $\left\|\phi-\hat{\phi}_{m}^{\mathrm{Lag}_{0}}\right\|_{\mathrm{L}^{2}}^{2} \leqslant\left\|\phi-\phi_{m}\right\|_{\mathrm{L}^{2}}^{2}+\mathcal{O}_{\mathrm{P}}\left(\frac{m}{T}\right)$.

In the following, we propose two ways inspired by Comte and Mabon 2017 to estimate the Gerber-Shiu function, using the Laguerre deconvolution method. To obtain a non asymptotic result on the MISE of the estimator, a cutoff is required when inverting the matrix $\hat{\mathbf{A}}_{m}$. Let $\theta_{0}<1$ be a truncation parameter, we estimate $\mathbf{A}_{m}^{-1}$ by:

$$
\widetilde{\mathbf{A}}_{m, 1}^{-1}:=\hat{\mathbf{A}}_{m}^{-1} \mathbf{1}_{\Delta_{m}^{1}} \quad \text { where } \quad \Delta_{m}^{1}:=\left\{\left\|\hat{\mathbf{A}}_{m}^{-1}\right\|_{\mathrm{op}} \leqslant \frac{2}{1-\theta_{0}}\right\}
$$

and we estimate the coefficients $\mathbf{a}_{m}$ by $\hat{\mathbf{a}}_{m}^{\mathrm{Lag}_{1}}:=\widetilde{\mathbf{A}}_{m, 1}^{-1} \times \hat{\mathbf{c}}_{m}$.
Theorem 3.2. Under Assumptions 1, 2, 3, and 4, if $\theta<\theta_{0}$ then it holds:

$$
\forall m \in \mathbb{N}^{*}, \quad \mathbb{E}\left\|\phi-\hat{\phi}_{m}^{\mathrm{Lag}_{1}}\right\|_{\mathrm{L}^{2}}^{2} \leqslant\left\|\phi-\phi_{m}\right\|_{\mathrm{L}^{2}}^{2}+C \frac{m}{T},
$$

where $C$ is a constant depending on $\lambda, c, \theta, \mathbb{E}[X], \mathbb{E}[W(X)]$ and $\theta_{0}-\theta$; and also $\delta, \mathbb{E}\left[X^{2}\right]$, $\mathbb{E}\left[W(X)^{2}\right]$ if $\delta>0$.

We propose a second way to estimate $\phi$ using the Laguerre deconvolution method, in the case $\delta=0$. It avoids the use of a truncation parameter $\theta_{0}$, but at the expense of an extra $\log (m)$ factor in the upper bound of the MISE, and it uses an additional independence assumption. We estimate the Laguerre coefficients of $g$ by 10 , that is in this case:

$$
\hat{b}_{k}:=\frac{1}{c T} \sum_{i=1}^{N_{T}} \Psi_{k}\left(X_{i}\right),
$$

where $\Psi_{k}(x):=\int_{0}^{x} \psi_{k}(t) \mathrm{d} t$. The matrix $\mathbf{A}_{m}$ is still estimated by 15.

Lemma 3.3. The functions $\Psi_{k}(x):=\int_{0}^{x} \psi_{k}(t) \mathrm{d} t$ are uniformly bounded.
This lemma is a technical result interesting in itself and we prove it in Appendix B Using this theorem, we can control the risk of $\hat{\mathbf{A}}_{m}$ in operator norm.

Proposition 3.4. If $p \geqslant 1$ and $\log m \geqslant p$, then it holds:

$$
\mathbb{E}\left[\left\|\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right\|_{\mathrm{op}}^{2 p}\right] \leqslant C(p, \lambda) \mu^{p}\left(\frac{m \log m}{c T}\right)^{p}+C(p)\left(\frac{m \log m}{c T}\right)^{2 p}
$$

where $C(p, \lambda)$ is a $\mathcal{O}\left(\lambda^{p}\right)$, and $C(p)$ is a constant depending on $p$.
This time, we estimate the inverse of the matrix $\mathbf{A}_{m}$ by:

$$
\widetilde{\mathbf{A}}_{m, 2}^{-1}:=\hat{\mathbf{A}}_{m}^{-1} \mathbf{1}_{\Delta_{m}^{2}}, \quad \text { where } \quad \Delta_{m}^{2}:=\left\{\left\|\hat{\mathbf{A}}_{m}^{-1}\right\|_{\mathrm{op}}^{2} \leqslant \frac{c T}{m \log m}\right\}
$$

we estimate the coefficients of $\phi$ by $\hat{\mathbf{a}}_{m}^{\mathrm{Lag}_{2}}:=\widetilde{\mathbf{A}}_{m, 2}^{-1} \times \hat{\mathbf{c}}_{m}$, and we estimate $\phi$ by:

$$
\hat{\phi}_{m}^{\mathrm{Lag}_{2}}:=\sum_{k=0}^{m-1} \hat{a}_{k}^{\mathrm{Lag}_{2}} \psi_{k}
$$

To provide an upper bound on the MISE of $\hat{\phi}_{m}^{\mathrm{Lag}_{2}}$, we need $\widetilde{\mathbf{A}}_{m, 2}^{-1}$ and $\hat{\mathbf{c}}_{m}$ to be independent. For this reason, we assume that we have a second observation set $\left\{N_{T}^{\prime} ; X_{1}^{\prime}, \ldots, X_{N_{T}^{\prime}}^{\prime}\right\}$ identical in law but independent from the main on $\underbrace{1}$. We use this second set to estimate $\widetilde{\mathbf{A}}_{m, 2}^{-1}$.
Theorem 3.5. We assume that $\delta=0$. Under Assumptions 1, 2, 3 and 4. if $m \log m \leqslant c T$ then it holds:

$$
\mathbb{E}\left\|\phi-\hat{\phi}_{m}^{\mathrm{Lag}_{2}}\right\|_{\mathrm{L}^{2}}^{2} \leqslant\left\|\phi-\phi_{m}\right\|_{\mathrm{L}^{2}}^{2}+\frac{C(\lambda)}{c T(1-\theta)^{2}}\left(\frac{\mathbb{E}[W(X)]}{c}+\|\phi\|_{\mathrm{L}^{2}}^{2}\left(\mu+\mu^{2}\right) m \log m\right)+\mathcal{O}\left(\frac{1}{T^{2}}\right)
$$

with $C(\lambda)=\mathcal{O}\left(\lambda \vee \lambda^{2}\right)$.
Remark 3.6. Contrary to the Laguerre-Fourier method, there is only one bias term with the Laguerre deconvolution method. However, the variance term is more complicated and a biasvariance compromise is needed. It leads to non-parametric rates of convergence.

## 4 Convergence rates of the Laguerre estimators

### 4.1 Sobolev-Laguerre spaces

To study the bias of a function in the Laguerre basis, we consider the Sobolev-Laguerre spaces. These functional spaces have been introduced by Bongioanni and Torrea 2009 to study the Laguerre operator. The connection with the Laguerre coefficients was established later by Comte and Genon-Catalot 2015.

Definition 4.1. For $s>0$, we define the Sobolev-Laguerre ball of radius $L>0$ and regularity $s$ as:

$$
\mathrm{W}^{s}\left(\mathbb{R}_{+}, L\right):=\left\{v \in \mathrm{~L}^{2}\left(\mathbb{R}_{+}\right) \mid \sum_{k=0}^{+\infty}\left\langle v, \psi_{k}\right\rangle^{2} k^{s} \leqslant L\right\}
$$

and we define the Sobolev-Laguerre space of regularity $s$ as $\mathrm{W}^{s}\left(\mathbb{R}_{+}\right):=\bigcup_{L>0} \mathrm{~W}^{s}\left(\mathbb{R}_{+}, L\right)$.
By Proposition 7.2 in Comte and Genon-Catalot 2015, when $s$ is a natural number, $v \in \mathrm{~W}^{s}\left(\mathbb{R}_{+}\right)$if and only if $v$ is $(s-1)$ times differentiable, $v^{(s-1)}$ is absolutely continuous, and for all $0 \leqslant k \leqslant s-1$ we have $x^{\frac{k+1}{2}} \sum_{j=0}^{k+1}\binom{k+1}{j} v^{(j)} \in \mathrm{L}^{2}\left(\mathbb{R}_{+}\right)$.

[^1]We are interested in the Sobolev-Laguerre spaces because of the following observation. If $v$ belongs to a Sobolev-Laguerre ball $\mathrm{W}^{s}\left(\mathbb{R}_{+}, L\right)$, then its bias is controlled by:

$$
\left\|v-v_{m}\right\|_{\mathrm{L}^{2}}^{2}=\sum_{k=m}^{+\infty}\left\langle v, \psi_{k}\right\rangle^{2}=\sum_{k=m}^{+\infty}\left\langle v, \psi_{k}\right\rangle^{2} k^{s} k^{-s} \leqslant L m^{-s} .
$$

Combining this upper bound on the bias term with Corollary 2.7 and Theorems 3.2 and 3.5 , we obtain convergence rates for the Laguerre-Fourier estimator and the Laguerre deconvolution estimators, on Sobolev-Laguerre spaces.

Theorem 4.2. Under Assumptions 1, 2 and 3, if $\theta<\theta_{0}$ and if $\phi \in \mathrm{W}^{s_{1}}\left(\mathbb{R}_{+}\right), g \in \mathrm{~W}^{s_{2}}\left(\mathbb{R}_{+}\right)$and $h \in \mathrm{~W}^{s_{3}}\left(\mathbb{R}_{+}\right)$, then choosing $m_{i}>(c T)^{\frac{1}{1+s_{i}}}$ for all $i \in\{1,2,3\}$ yields:

$$
\mathbb{E}\left\|\phi-\hat{\phi}_{m_{1}, m_{2}, m_{3}}\right\|_{\mathrm{L}^{2}}^{2}=\mathcal{O}\left(\frac{1}{c T}\right) .
$$

Remark 4.3. If $\phi, g$ and $h$ belong to some Sobolev-Laguerre spaces, we can just choose $m_{1}=m_{2}=$ $m_{3}=\lceil c T\rceil$ and obtain the parametric rate $\mathcal{O}\left(\frac{1}{c T}\right)$ for the Laguerre-Fourier estimator.

Theorem 4.4. We make Assumptions 1, 2 3 and 4 and we assume that $\phi \in \mathrm{W}^{s}\left(\mathbb{R}_{+}\right)$.

1. If $\theta<\theta_{0}$, then choosing $m_{\mathrm{opt}} \propto(c T)^{\frac{1}{1+s}}$ yields:

$$
\mathbb{E}\left\|\phi-\hat{\phi}_{m_{\mathrm{opt}}}^{\mathrm{Lag}_{1}}\right\|_{\mathrm{L}^{2}}^{2}=\mathcal{O}\left((c T)^{-\frac{s}{1+s}}\right)
$$

2. If $\delta=0$, then choosing $m_{\mathrm{opt}} \propto(c T)^{\frac{1}{1+s}}$ yields:

$$
\mathbb{E}\left\|\phi-\hat{\phi}_{m_{\mathrm{opt}}}^{\mathrm{Lag}_{2}}\right\|_{\mathrm{L}^{2}}^{2}=\mathcal{O}\left((c T)^{-\frac{s}{1+s}} \log (c T)\right)
$$

Remark 4.5. The Fourier-Laguerre estimator and the Laguerre deconvolution estimator $\hat{\phi}_{m}^{\mathrm{Lag}_{1}}$ both depend on a truncation parameter $\theta_{0}$ that needs to be chosen such that $\theta<\theta_{0}$. We see two ways to ensure that.

1. We can assume that we know some $\theta_{0}<1$ such that $\theta<\theta_{0}$. Then our convergence rates are those of Theorems 4.2 and 4.4
2. We can choose $\theta_{0}=1-(\log T)^{1 / 2}$. Then for $T$ large enough (more precisely $T>\mathrm{e}^{(1-\theta)^{2}}$ ), our convergence rates of are those of Theorems 4.2 and 4.4 multiplied by $\log (T)$.

In our simulations, we chose the first way.

### 4.2 The exponential case

In this section, we want to compute the convergence rate of the estimators, in the exponential case: $X \sim \operatorname{Exp}(1 / \mu)$. This distribution is often considered in risk theory and closed forms of the Gerber-Shiu function are available in this case. Indeed, the Gerber-Shiu functions we are interested in (see Section 1) are given by:

$$
\phi(u)= \begin{cases}\theta \exp \left(-\frac{1-\theta}{\mu} u\right) & \text { (ruin probability), }  \tag{17}\\ \frac{\theta}{1+\mu \rho_{\delta}} \exp \left(-\left[\frac{1-\theta}{\mu}+\rho_{\delta}-\frac{\delta}{c}\right] u\right) & \text { (Laplace transform of the ruin time), } \\ \mu(1+2 \theta) \exp \left(-\frac{1-\theta}{\mu} u\right)-\mu \exp \left(-\frac{u}{\mu}\right) & \text { (jump size causing the ruin). }\end{cases}
$$

These formulas are obtained by Laplace inversion, see Asmussen and Albrecher 2010, chapter XII. We use the following lemma to compute the bias terms of the functions $\phi, g$ and $h$.

Lemma 4.6. Let $C, \gamma$ be positive numbers and let $F(x)=C \exp (-\gamma x) \mathbf{1}_{\mathbb{R}_{+}}(x)$. The Laguerre coefficients of $F$ are given by:

$$
\left\langle F, \psi_{k}\right\rangle=\frac{C \sqrt{2}}{\gamma+1}\left(\frac{\gamma-1}{\gamma+1}\right)^{k}
$$

Hence if $m \geqslant 0$ we have:

$$
\sum_{k=m}^{+\infty}\left\langle F, \psi_{k}\right\rangle^{2}=\frac{C^{2}}{2 \gamma}\left(\frac{\gamma-1}{\gamma+1}\right)^{2 m}
$$

Proposition 4.7. If the density of $X$ is $f(x)=\frac{1}{\mu} \mathrm{e}^{-\frac{x}{\mu}}$, then the bias term of $\phi$ is given by:

$$
\left\|\phi-\phi_{m}\right\|_{\mathrm{L}^{2}}^{2} \leqslant L \mathrm{e}^{-r m}
$$

with $L$ and $r$ given by:

1. Ruin probability: $L=\frac{\theta^{2}}{2 \gamma}, r=2 \log \left|\frac{\gamma+1}{\gamma-1}\right|$ and $\gamma=\frac{1-\theta}{\mu}$.
2. Laplace transform of the ruin time: $L=\frac{\theta^{2}}{2 \gamma\left(1+\mu \rho_{\delta}\right)^{2}}, r=2 \log \left|\frac{\gamma+1}{\gamma-1}\right|$ and $\gamma=\frac{1-\theta}{\mu}+\rho_{\delta}-\frac{\delta}{c}$.
3. Jump size causing the ruin: $L=\mu^{3} \frac{(1+2 \theta)^{2}}{1-\theta}, r=2 \log \left(\left|\frac{1-\theta+\mu}{1-\theta-\mu}\right| \wedge\left|\frac{1+\mu}{1-\mu}\right|\right)$.

Combining this Proposition with Theorems 3.2 and 3.5 we easily obtain convergence rates for the Gerber-Shiu functions we are interested in.

Theorem 4.8. We assume that the density of $X$ is $f(x)=\frac{1}{\mu} \mathrm{e}^{-\frac{x}{\mu}}$, we make Assumptions 1. 2, ,3 and 4, and we assume that the bias term of $\phi$ decreases as:

$$
\left\|\phi-\phi_{m}\right\|_{\mathrm{L}^{2}}^{2} \leqslant L \mathrm{e}^{-r m}
$$

1. If $\theta<\theta_{0}$, then choosing $m_{\mathrm{opt}}=\left\lceil\frac{1}{r} \log (c T)\right\rceil$ yields:

$$
\mathbb{E}\left\|\phi-\hat{\phi}_{m_{\mathrm{opt}}}^{\mathrm{Lag}_{1}}\right\|^{2}=\mathcal{O}\left(\frac{\log (c T)}{c T}\right)
$$

2. If $\delta=0$, then choosing $m_{\mathrm{opt}}=\left\lceil\frac{1}{r} \log (c T)\right\rceil$ yields:

$$
\mathbb{E}\left\|\phi-\hat{\phi}_{m_{\mathrm{opt}}}^{\mathrm{Lag}_{2}}\right\|_{\mathrm{L}^{2}}^{2}=\mathcal{O}\left(\frac{\log (c T) \log \log (c T)}{c T}\right)
$$

For the Laguerre-Fourier estimator, we also need to know the decreasing rate of the bias term of $g$ and $h$. For the ruin probability, the Laplace transform of $\tau$, and the jump size causing the ruin, direct calculations show that $g$ and $h$ are given by a positive multiple of $\mathrm{e}^{-x / \mu}$. Thus, Lemma 4.6 yields that their bias term is less than $\exp \left(-r^{\prime} m\right)$, with $r^{\prime}:=2 \log \left|\frac{1+\mu}{1-\mu}\right|$. Together with Corollary 2.7, we obtain the convergence rates of the Laguerre-Fourier estimator.
Theorem 4.9. If the density of $X$ is $f(x)=\frac{1}{\mu} \mathrm{e}^{-\frac{x}{\mu}}$, under Assumptions 1 , 2and 3, if $\theta<\theta_{0}$ and if the bias term of $\phi$ decreases as:

$$
\left\|\phi-\phi_{m}\right\|_{\mathrm{L}^{2}}^{2} \leqslant L \mathrm{e}^{-r m}
$$

then choosing $m_{1}>\left\lceil\frac{1}{r} \log (c T)\right\rceil$ and $m_{2}, m_{3}>\left\lceil\frac{1}{r^{\prime}} \log (c T)\right\rceil$ with $r^{\prime}:=2 \log \left|\frac{1+\mu}{1-\mu}\right|$ yields:

$$
\mathbb{E}\left\|\phi-\hat{\phi}_{m_{1}, m_{2}, m_{3}}\right\|_{\mathrm{L}^{2}}^{2}=\mathcal{O}\left(\frac{1}{c T}\right)
$$

## 5 Numerical comparison

In this section, we compare the performance on simulated data of the Laguerre-Fourier estimator (see Section 2) and the Laguerre deconvolution estimators (see Section 3). We consider three types of Gerber-Shiu functions:

1. $\phi^{(1)}(u)=\mathbb{P}[\tau(u)<\infty]$ the ruin probability;
2. $\phi^{(2)}(u)=\mathbb{E}\left[\left(U_{\tau(u)-}+\left|U_{\tau(u)}\right|\right) \mathbf{1}_{\tau(u)<\infty}\right]$ the expected claim size causing the ruin;
3. $\phi^{(3)}(u)=\mathbb{E}\left[\mathrm{e}^{-\delta \tau(u)} \mathbf{1}_{\tau(u)<\infty}\right]$ the Laplace transform of the ruin time, for $\delta=0.1$.

We also consider 3 sets of parameters:

1. $X$ follows an exponential distribution, $\lambda=1, \mu=1, c=1.5$. In this setting, $\theta \approx 0.67$.
2. $X$ follows an exponential distribution, $\lambda=1.25, \mu=2, c=3$. In this setting, $\theta \approx 0.83$
3. $X$ follows a $\Gamma\left(2, \frac{\mu}{2}\right)$ distribution, $\lambda=1.25, \mu=2, c=3$. In this setting, $\theta \approx 0.83$.

Model selection Each estimator we consider depends on one or several parameters that need to be chosen. The Laguerre-Fourier estimator and the Laguerre deconvolution estimator $\hat{\phi}_{m}^{\mathrm{Lag}_{1}}$ depend on a truncation parameter $\theta_{0}$, which needs to be chosen such that $\theta<\theta_{0}$. We choose $\theta_{0}=0.95$ in our simulations.

- The Laguerre-Fourier estimator depends on four parameters: $m_{1}, m_{2}$ and $m_{3}$, the dimensions of the projection spaces for the functions $\phi, g$ and $h$, and $\theta_{0}$ the truncation parameter in the estimation of $\widehat{\mathcal{F} g}$. As said in Remark 4.3. we can choose $m_{1}=m_{2}=m_{3}=\lceil c T\rceil$, no selection procedure is required. Still, we propose a model reduction procedure for the choice of $m_{2}$ and $m_{3}$, that we describe in Appendix A
- The naive Laguerre deconvolution estimator $\hat{\phi}_{m}^{\mathrm{Lag}_{0}}$, defined by (16), depends on one parameter: $m$, the dimension of the projection space for $\phi$. However, there is no model selection procedure for $m$. In their numerical section, Zhang and Su 2018 only consider (as we do) Gerber-Shiu functions with exponential decay; hence the bias term also decays with exponential rate. Using this fact, they chose $m=\left\lfloor 5 T^{1 / 10}\right\rfloor$. We make the same choice in our simulations and we write $\hat{\phi}^{\mathrm{ZS}}$ this estimator.
- The Laguerre deconvolution estimators $\hat{\phi}_{m}^{\mathrm{Lag}_{1}}$ and $\hat{\phi}_{m}^{\mathrm{Lag}_{2}}$ also depend on $m$. For $i \in\{1,2\}$, we choose $\hat{m}^{\mathrm{Lag}_{i}}$ as the minimizer of a penalized criterion:

$$
\begin{equation*}
\hat{m}^{\mathrm{Lag}_{i}} \in \underset{m \in \mathcal{M}_{i}}{\arg \min }\left\{-\left\|\hat{\phi}_{m}^{\mathrm{Lag}_{i}}\right\|_{\mathrm{L}^{2}}^{2}+\kappa_{i} \operatorname{pen}_{i}(m)\right\} \tag{18}
\end{equation*}
$$

where the model collections are:

$$
\begin{gathered}
\mathcal{M}_{1}:=\left\{1 \leqslant m \leqslant M \left\lvert\,\left\|\hat{\mathbf{A}}_{m}^{-1}\right\|_{\mathrm{op}} \leqslant \frac{1}{1-\theta_{0}}\right.\right\} \\
\mathcal{M}_{2}:=\left\{1 \leqslant m \leqslant M \left\lvert\,\left\|\hat{\mathbf{A}}_{m}^{-1}\right\|_{\mathrm{op}}^{2} \leqslant \frac{c T}{m \log (m)}\right.\right\}
\end{gathered}
$$

with $M=\lceil c T\rceil \wedge 500$ (we do not compute more than 500 coefficients, because of computation time).
In the following, if $F(X)$ is a function of $X$, we write $\overline{F(X)}:=\frac{1}{N_{T}} \sum_{i=1}^{N_{T}} F\left(X_{i}\right)$ its empirical mean from the sample $\left\{X_{1}, \ldots, X_{N_{T}}\right\}$. For the penalty terms, we choose empirical versions of the variance terms in Theorems 3.2 and 3.5 .

$$
\begin{gathered}
\operatorname{pen}_{1}(m):=\frac{1}{\left(1-\theta_{0}\right)^{2}}\left(\left\|\hat{\phi}_{m}^{\mathrm{Lag}_{1}}\right\|_{\mathrm{L}^{2}}^{2} m \hat{V}_{g}+\hat{V}_{h}\right) \\
\operatorname{pen}_{2}(m):=\left(\hat{\lambda} \vee \hat{\lambda}^{2}\right)\left(\frac{\overline{W(X)}}{c}+\left\|\hat{\phi}_{m}^{\mathrm{Lag}_{1}}\right\|_{\mathrm{L}^{2}}^{2}\left(\bar{X} \vee \bar{X}^{2}\right) m \log (m)\right),
\end{gathered}
$$

with:

$$
\hat{V}_{g}:=\left\{\begin{array}{ll}
\frac{\hat{\lambda}}{c^{2} T} \bar{X} & \text { if } \delta=0, \\
\frac{\hat{\lambda}^{2}}{c^{2} T}\left(\bar{X}+\frac{\left(\overline{X^{2}}\right)^{1 / 2}}{\delta\left(1-\theta_{0}\right)^{2}}\right) & \text { if } \delta>0,
\end{array} \quad \hat{V}_{h}:= \begin{cases}\frac{\hat{\lambda}}{c^{2} T} \overline{W(X)} & \text { if } \delta=0, \\
\frac{\hat{\lambda}^{2}}{c^{2} T}\left(\overline{W(X)}+\frac{\left(\overline{W(X)^{2}}\right)^{1 / 2}}{\delta\left(1-\theta_{0}\right)^{2}}\right) & \text { if } \delta>0 .\end{cases}\right.
$$

The constants $\kappa_{1}$ and $\kappa_{2}$ are constants we calibrated on preliminary simulations. We choose:
$-\kappa_{1}=0.01, \kappa_{2}=0.01$ for the ruin probability;
$-\kappa_{1}=0.1, \kappa_{2}=1$ for the expected claim size causing the ruin;
$-\kappa_{1}=10^{-8}$ for the Laplace transform of the ruin time, $\delta=0.1$.
There is no constant $\kappa_{2}$ in the last case because the Laguerre deconvolution estimator $\hat{\phi}_{m}^{\mathrm{Lag}_{2}}$ is defined only if $\delta=0$.
We write $\hat{\phi}^{\mathrm{Lag}_{1}}:=\hat{\phi}_{\hat{m}^{\mathrm{Lag}_{1}}}^{\mathrm{Lag}_{1}}$ and $\hat{\phi}^{\mathrm{Lag}_{2}}:=\hat{\phi}_{\hat{m}^{\mathrm{Lag}_{2}}}^{\mathrm{Lag}_{2}}$ in the following.
MISE calculation We compare the estimators by looking at their Mean Integrated Squared Error (MISE): $\mathbb{E}\|\phi-\hat{\phi}\|_{\mathrm{L}^{2}}^{2}$. We compute the norm $\|\cdot\|_{\mathrm{L}^{2}}$ using Romberg's method, and we compute the expectation by an empirical mean over $n=200$ paths of the process $\left(U_{t}\right)_{t \in[0, T]}$. We also compute a $95 \%$ confidence interval for the MISE. We have two goals in this section:

1. To compare the performance of our Laguerre-Fourier estimator with the Laguerre deconvolution estimators.
2. To see if the model selection procedures (18) for the Laguerre deconvolution estimators lead to the same performance than the naive choice $m=\left\lfloor 5 T^{1 / 10}\right\rfloor$.

Results We display our results in Tables 1, 22 and 3
Concerning the estimation of the ruin probability (Table 1), we see that all the estimators perform well with the first set of parameter (exponential distribution, $\theta=0.67$ ). However, with the two other sets of parameters (exponential distribution and Gamma(2) distribution, $\theta=0.83$ ), the difference is clear: the Laguerre-Fourier estimator has the smallest risk, followed by the estimator of Zhang and Su 2018], and the Laguerre deconvolution estimators come last. We notice that $\hat{\phi}^{\mathrm{Lag}_{2}}$ seems to be better than $\hat{\phi}^{\mathrm{Lag}_{1}}$ in this case.

Concerning the estimation of the expected jump size causing the ruin (Table 22, the difference is even clearer. With the first set of parameters, we see that the Laguerre-Fourier is better for small sample size $\left(\mathbb{E}\left[N_{T}\right]=100\right)$, but equivalent to the other estimators for larger sample sizes. We also notice that the estimator $\hat{\phi}^{\mathrm{ZS}}$ and $\hat{\phi}^{\mathrm{Lag}_{2}}$ have the same risk. With the two other sets of parameters, we find again that the Laguerre-Fourier estimator is better than the estimator $\hat{\phi}^{\mathrm{ZS}}$, which is better than the Laguerre deconvolution estimators. This time, we see that $\hat{\phi}^{\text {Lag }_{1}}$ has better performances than $\hat{\phi}^{\mathrm{Lag}_{2}}$.

Concerning the estimation of Laplace transform of the ruin time (Table 3), we see no difference between the MISE of the Laguerre-Fourier estimator and the Laguerre deconvolution estimators.

To conclude, we can say that our Laguerre-Fourier estimator has better performances than the Laguerre deconvolution estimators on simulated data, even in the exponential case where they have theoretically the same MISE (up to a log factor). Furthermore, the Laguerre deconvolution estimators with the model selection procedure (18) fail to match the performance of the estimator of Zhang and Su 2018], for which we choose the parameter $m$ knowing the bias decay rate of $\phi$, in most cases.

## 6 Conclusion

Using a projection estimator on the Laguerre basis, and computing the coefficients with Fourier transforms, we constructed an estimator of the Gerber-Shiu function that achieves parametric rates of convergence, without needing a model selection procedure. It is worth noticing that our

| Parameters | Estimator | $\mathbb{E}\left[N_{T}\right]=100$ | $\mathbb{E}\left[N_{T}\right]=200$ | $\mathbb{E}\left[N_{T}\right]=400$ |
| :---: | :---: | :---: | :---: | :---: |
| $\begin{gathered} X \sim \operatorname{Exp}(1) \\ \lambda=1 \\ c=1.5 \\ \theta=0.67 \end{gathered}$ | LagFou | $\begin{gathered} 0.14 \\ {[0.07,0.21]} \\ m_{1}=150 \\ \hline \end{gathered}$ | $\begin{gathered} 0.053 \\ {[0.039,0.067]} \\ m_{1}=300 \\ \hline \end{gathered}$ | $\begin{gathered} 0.022 \\ {[0.017,0.027]} \\ m_{1}=500 \\ \hline \end{gathered}$ |
|  | ZS | $\begin{gathered} 0.23 \\ {[0.02,0.44]} \\ m=8 \end{gathered}$ | $\begin{gathered} 0.053 \\ {[0.039,0.068]} \\ m=9 \end{gathered}$ | $\begin{gathered} 0.022 \\ {[0.017,0.028]} \\ m=10 \end{gathered}$ |
|  | LagDec1 | $\begin{gathered} 0.25 \\ {[0.01,0.48]} \\ \hat{m}=3.3 \\ \hline \end{gathered}$ | $\begin{gathered} 0.055 \\ {[0.042,0.069]} \\ \hat{m}=3.8 \end{gathered}$ | $\begin{gathered} 0.024 \\ {[0.019,0.029]} \\ \hat{m}=4.2 \end{gathered}$ |
|  | LagDec2 | $\begin{gathered} 0.23 \\ {[0.02,0.45]} \\ \hat{m}=6.0 \end{gathered}$ | $\begin{gathered} 0.053 \\ {[0.039,0.068]} \\ \hat{\hat{m}}=6.5 \end{gathered}$ | $\begin{gathered} 0.023 \\ {[0.017,0.028]} \\ \hat{m}=7.0 \end{gathered}$ |
| $\begin{gathered} X \sim \operatorname{Exp}(1 / 2) \\ \lambda=1.25 \\ c=3 \\ \theta=0.83 \end{gathered}$ | LagFou | $\begin{gathered} 0.95 \\ {[0.80,1.09]} \\ m_{1}=240 \\ \hline \end{gathered}$ | $\begin{gathered} 0.67 \\ {[0.53,0.80]} \\ m_{1}=480 \\ \hline \end{gathered}$ | $\begin{gathered} \hline \hline 0.43 \\ {[0.31,0.55]} \\ m_{1}=500 \\ \hline \end{gathered}$ |
|  | ZS | $\begin{gathered} 1.57 \\ {[1.26,1.89]} \\ m=8 \end{gathered}$ | $\begin{gathered} 1.02 \\ {[0.74,1.30]} \\ m=9 \end{gathered}$ | $\begin{gathered} 0.54 \\ {[0.46,0.61]} \\ m=9 \end{gathered}$ |
|  | LagDec1 | $\begin{gathered} 8.51 \\ {[5.07,11.95]} \\ \hat{\hat{m}}=11.0 \end{gathered}$ | $\begin{gathered} 3.82 \\ {[1.57,6.07]} \\ \hat{\hat{m}}=12.8 \end{gathered}$ | $\begin{gathered} 0.72 \\ {[0.36,1.07]} \\ \hat{m}=14.4 \end{gathered}$ |
|  | LagDec2 | $\begin{gathered} 2.96 \\ {[2.11,3.82]} \\ \hat{m}=14.2 \\ \hline \end{gathered}$ | $\begin{gathered} 1.94 \\ {[1.08,2.80]} \\ \hat{m}=18.1 \\ \hline \end{gathered}$ | $\begin{gathered} 0.64 \\ {[0.36,0.92]} \\ \hat{m}=21.8 \end{gathered}$ |
| $\begin{gathered} X \sim \Gamma(2,1 / 4) \\ \lambda=1.25 \\ c=3 \\ \theta=0.83 \end{gathered}$ | LagFou | $\begin{gathered} 0.64 \\ {[0.52,0.77]} \\ m_{1}=240 \\ 1.77 \end{gathered}$ | $\begin{gathered} 0.46 \\ {[0.37,0.56]} \\ m_{1}=480 \\ \hline \end{gathered}$ | $\begin{gathered} 0.30 \\ {[0.22,0.38]} \\ m_{1}=500 \\ \hline \end{gathered}$ |
|  | ZS | $\begin{gathered} 1.77 \\ {[1.07,2.47]} \\ m=8 \end{gathered}$ | $\begin{gathered} 0.62 \\ {[0.45,0.78]} \\ m=9 \end{gathered}$ | $\begin{gathered} 0.30 \\ {[0.23,0.36]} \\ m=9 \end{gathered}$ |
|  | LagDec1 | $\begin{gathered} 7.22 \\ {[4.25,10.19]} \\ \hat{\hat{m}}=9.3 \end{gathered}$ | $\begin{gathered} 1.71 \\ {[0.87,2.56]} \\ \hat{\vec{m}}=10.6 \end{gathered}$ | $\begin{gathered} 0.45 \\ {\left[\frac{0.21,0.70]}{\hat{m}}=11.6\right.} \end{gathered}$ |
|  | LagDec2 | $\begin{gathered} 2.71 \\ {[1.80,3.62]} \\ \hat{\hat{m}}=12.1 \end{gathered}$ | $\begin{gathered} 1.07 \\ {\left[\frac{0.66,1.47]}{\hat{m}}=15.7\right.} \end{gathered}$ | $\begin{gathered} 0.41 \\ {\left[\frac{0.22,0.60]}{\hat{m}}=18.0\right.} \end{gathered}$ |

Table 1: Ruin Probability. For two sets of parameters, we compare the three estimators of the ruin probability: the Laguerre-Fourier estimator (LagFou), the estimator of Zhang and Su [2018] (ZS), and the Laguerre deconvolution estimators (LagDec1 and LagDec2). In each case, we display the estimation of the MISE over 200 samples with a $95 \%$ confidence interval and the model used ( $\overline{\hat{m}}$ is the mean selected model in the case of the Laguerre deconvolution estimators).

| Parameters | Estimator | $\mathbb{E}\left[N_{T}\right]=100$ | $\mathbb{E}\left[N_{T}\right]=200$ | $\mathbb{E}\left[N_{T}\right]=400$ |
| :---: | :---: | :---: | :---: | :---: |
| $\begin{gathered} X \sim \operatorname{Exp}(1) \\ \lambda=1 \\ c=1.5 \\ \theta=0.67 \end{gathered}$ | LagFou | $\begin{gathered} 1.71 \\ {[1.09,2.32]} \\ m_{1}=150 \end{gathered}$ | $\begin{gathered} 0.60 \\ {[0.46,0.73]} \\ m_{1}=300 \\ \hline \end{gathered}$ | $\begin{gathered} 0.34 \\ {[0.27,0.40]} \\ m_{1}=500 \\ \hline \end{gathered}$ |
|  | ZS | $\begin{gathered} 1.80 \\ {[1.07,2.53]} \\ m=8 \end{gathered}$ | $\begin{gathered} 0.62 \\ {[0.47,0.77]} \\ m=9 \end{gathered}$ | $\begin{gathered} 0.34 \\ {[0.27,0.41]} \\ m=10 \end{gathered}$ |
|  | LagDec1 | $\begin{gathered} 1.41 \\ {[1.19,1.64]} \\ \hat{m}=1.9 \end{gathered}$ | $\begin{gathered} 0.84 \\ {[0.74,0.93]} \\ \hat{m}=2.2 \end{gathered}$ | $\begin{gathered} 0.44 \\ {[0.39,0.50]} \\ \hat{m}=2.8 \end{gathered}$ |
|  | LagDec2 | $\begin{gathered} 1.86 \\ {[1.11,2.61]} \\ \hat{m}=3.7 \end{gathered}$ | $\begin{gathered} 0.64 \\ {[0.49,0.78]} \\ \hat{\vec{m}}=4.1 \end{gathered}$ | $\begin{gathered} 0.35 \\ {[0.28,0.42]} \\ \hat{\vec{m}}=4.6 \end{gathered}$ |
| $\begin{gathered} X \sim \operatorname{Exp}(1 / 2) \\ \lambda=1.25 \\ c=3 \\ \theta=0.83 \end{gathered}$ | LagFou | 46.2 $[30.0,62.3]$ $m_{1}=240$ | $\begin{gathered} \hline \hline 28.1 \\ {[21.7,34.5]} \\ m_{1}=480 \\ \hline \end{gathered}$ | $\begin{gathered} \hline \hline 20.5 \\ {[15.1,25.8]} \\ m_{1}=500 \\ \hline \end{gathered}$ |
|  | ZS | $\begin{gathered} 96.3 \\ {[62.4,130.2]} \\ m=8 \end{gathered}$ | $\begin{gathered} 48.0 \\ {[31.3,64.6]} \\ m=9 \end{gathered}$ | $\begin{gathered} 27.9 \\ {[23.0,32.7]} \\ m=9 \end{gathered}$ |
|  | LagDec1 | $\begin{gathered} 77.5 \\ {[71.6,83.5]} \\ \hat{m}=3.0 \\ \hline \end{gathered}$ | $\begin{gathered} 56.3 \\ {[45.1,67.4]} \\ \hat{m}=4.9 \\ \hline \end{gathered}$ | $\begin{gathered} 38.9 \\ {[29.7,48.1]} \\ \hat{m}=7.1 \\ \hline \end{gathered}$ |
|  | LagDec2 | $\begin{gathered} 197.7 \\ {[115.1,280.3]} \\ \hat{m}=9.7 \end{gathered}$ | $\begin{gathered} 96.7 \\ {[47.4,146.0]} \\ \hat{m}=12.5 \end{gathered}$ | $\begin{gathered} 48.5 \\ {[27.5,69.6]} \\ \hat{\vec{m}}=14.6 \end{gathered}$ |
| $\begin{gathered} X \sim \Gamma(2,1 / 4) \\ \lambda=1.25 \\ c=3 \\ \theta=0.83 \end{gathered}$ | LagFou | $\begin{gathered} \hline \hline 11.7 \\ {[9.4,14.0]} \\ m_{1}=240 \end{gathered}$ | $\begin{gathered} 9.2 \\ {[7.5,10.9]} \\ m_{1}=480 \end{gathered}$ | $\begin{gathered} \hline \hline 6.2 \\ {[4.1,8.3]} \\ m_{1}=500 \\ \hline \end{gathered}$ |
|  | ZS | $\begin{gathered} 18.5 \\ {[12.0,25.0]} \\ m=8 \end{gathered}$ | $\begin{gathered} 13.6 \\ {[10.1,17.1]} \\ m=9 \end{gathered}$ | $\begin{gathered} 5.9 \\ {[4.6,7.2]} \\ m=9 \end{gathered}$ |
|  | LagDec1 | $\begin{gathered} 19.2 \\ {\left[\frac{18.2,20.2]}{\hat{m}}=2.8\right.} \end{gathered}$ | $\begin{gathered} 15.0 \\ {\left[\begin{array}{c} 13.1,17.0] \\ \hat{m} \\ \hline \end{array} .3 .3\right.} \end{gathered}$ | 8.4 $[7.1,9.7]$ $\hat{m}=5.9$ |
|  | LagDec2 | $\begin{gathered} 28.2 \\ {[19.4,37.1]} \\ \hat{m}=7.8 \\ \hline \end{gathered}$ | $\begin{gathered} 24.6 \\ {[16.5,32.7]} \\ \hat{\hat{m}}=10.0 \\ \hline \end{gathered}$ | $\begin{gathered} 8.3 \\ {[5.6,11.1]} \\ \hat{m}=11.3 \end{gathered}$ |

Table 2: Expected claim size causing the ruin. For two sets of parameters, we compare the three estimators of the expected jump size causing the ruin: the Laguerre-Fourier estimator (LagFou), the estimator of Zhang and Su 2018 (ZS), and the Laguerre deconvolution estimators (LagDec1 and LagDec2). In each case, we display the estimation of the MISE over 200 samples with a $95 \%$ confidence interval and the model used ( $\overline{\hat{m}}$ is the mean selected model in the case of the Laguerre deconvolution estimators).

| Parameters | Estimator | $\mathbb{E}\left[N_{T}\right]=100$ | $\mathbb{E}\left[N_{T}\right]=200$ | $\mathbb{E}\left[N_{T}\right]=400$ |
| :---: | :---: | :---: | :---: | :---: |
| $\begin{gathered} X \sim \operatorname{Exp}(1) \\ \lambda=1 \\ c=1.5 \\ \theta=0.67 \end{gathered}$ | LagFou | $\begin{gathered} 2.50 \\ {[1.91,3.09]} \\ m_{1}=150 \end{gathered}$ | $\begin{gathered} 1.09 \\ {[0.87,1.31]} \\ m_{1}=300 \end{gathered}$ | $\begin{gathered} 0.64 \\ {[0.52,0.77]} \\ m_{1}=500 \end{gathered}$ |
|  | ZS | $\begin{gathered} 2.50 \\ {[1.93,3.07]} \\ m=8 \end{gathered}$ | $\begin{gathered} 1.10 \\ {[0.88,1.33]} \\ m=9 \end{gathered}$ | $\begin{gathered} 0.66 \\ {[0.53,0.79]} \\ m=10 \end{gathered}$ |
|  | LagDec1 | $\begin{gathered} 2.52 \\ {[1.95,3.08]} \\ \hat{\hat{m}}=4.2 \end{gathered}$ | $\begin{gathered} 1.11 \\ {[0.89,1.34]} \\ \hat{m}=4.6 \end{gathered}$ | $\begin{gathered} 0.67 \\ {[0.54,0.80]} \\ \hat{m}=4.9 \end{gathered}$ |
| $\begin{gathered} X \sim \operatorname{Exp}(1 / 2) \\ \lambda=1.25 \\ c=3 \\ \theta=0.83 \end{gathered}$ | LagFou | $\begin{gathered} \hline \hline 11.81 \\ {[9.26,14.36]} \\ m_{1}=240 \\ \hline \end{gathered}$ | $\begin{gathered} \hline \hline 5.60 \\ {[4.49,6.72]} \\ m_{1}=480 \\ \hline \end{gathered}$ | $\begin{gathered} 2.51 \\ {[2.04,2.98]} \\ m_{1}=500 \\ \hline \end{gathered}$ |
|  | ZS | $\begin{gathered} 12.47 \\ {[10.53,14.41]} \\ m=8 \end{gathered}$ | $\begin{gathered} 6.13 \\ {[5.07,7.19]} \\ m=9 \end{gathered}$ | $\begin{gathered} 3.30 \\ {[2.86,3.75]} \\ m=9 \end{gathered}$ |
|  | LagDec1 | $\begin{gathered} 13.22 \\ {[10.61,15.84]} \\ \hat{m}=10.2 \\ \hline \end{gathered}$ | $\begin{gathered} 5.82 \\ {[4.57,7.06]} \\ \hat{m}=11.2 \end{gathered}$ | $\begin{gathered} 2.65 \\ {[2.15,3.14]} \\ \hat{m}=12.4 \end{gathered}$ |
| $\begin{gathered} X \sim \Gamma(2,1 / 4) \\ \lambda=1.25 \\ c=3 \\ \theta=0.83 \end{gathered}$ | LagFou | $\begin{gathered} \hline \hline 10.26 \\ {[8.08,12.45]} \\ m_{1}=240 \\ \hline \end{gathered}$ | $\begin{gathered} \hline 4.09 \\ {[3.28,4.91]} \\ m_{1}=480 \\ \hline \end{gathered}$ | $\begin{gathered} 2.01 \\ {[1.57,2.46]} \\ m_{1}=500 \\ \hline \end{gathered}$ |
|  | ZS | $\begin{gathered} 9.76 \\ {[8.07,11.45]} \\ m=8 \end{gathered}$ | $\begin{gathered} 4.16 \\ {[3.36,4.96]} \\ m=9 \end{gathered}$ | $\begin{gathered} 2.20 \\ {[1.77,2.63]} \\ m=9 \end{gathered}$ |
|  | LagDec1 | $\begin{gathered} 10.39 \\ {[8.39,12.38]} \\ \hat{m}=8.8 \\ \hline \end{gathered}$ | $\begin{gathered} 4.15 \\ {[3.29,5.00]} \\ \hat{m}=9.5 \end{gathered}$ | $\begin{gathered} 2.05 \\ {[1.59,2.52]} \\ \hat{m}=10.3 \end{gathered}$ |

Table 3: Laplace transform, $\delta=0.1$. For two sets of parameters, we compare three estimators of the Laplace transform of the ruin time: the Laguerre-Fourier estimator (LagFou), the estimator of Zhang and Su 2018 (ZS) and the Laguerre deconvolution estimator (LagDec1). In each case, we display the estimation of the MISE $\times 10^{2}$ over 200 samples with a $95 \%$ confidence interval and the model used ( $\overline{\hat{m}}$ is the mean selected model in the case of the Laguerre deconvolution estimator).
results are non-asymptotic and concern the MISE of the estimator. In comparison, the Laguerre deconvolution estimators have slower rates of convergence and necessitate a model selection procedure in practice. The better performances of our procedure are confirmed by a numerical study, on simulated data.

Knowing that the Laguerre deconvolution method does not achieve the best rate of convergence in the compound Poisson model is important. Indeed, this method is used to estimate the GerberShiu function in more general models, see Zhang and Su [2019], Su et al. 2019a and Su et al. 2019b. These papers have one thing in common: they all want to estimate a function $\phi$ that satisfies an equation of the form $\phi=\phi * g+h$, with $g$ and $h$ functions that depend on the specificity of each problem. If we applied the procedure described in the beginning of Section 2 , we could obtain an estimator that would achieve the same rate of convergence as the estimators of $g$ and $h$ (see Remark 2.2 . Hence the Laguerre deconvolution method used in these papers is not optimal since a factor $m$ appears in the variance term in the construction step of $\hat{\phi}_{m}$ from $\hat{g}_{m}$ and $\hat{h}_{m}$.

## 7 Proofs

Proof of Proposition 2.1. By Pythagoras theorem, $\left\|\phi-\hat{\phi}_{m}\right\|_{\mathrm{L}^{2}}^{2}=\left\|\phi-\phi_{m}\right\|_{\mathrm{L}^{2}}^{2}+\left\|\phi_{m}-\hat{\phi}_{m}\right\|_{\mathrm{L}^{2}}^{2}$. Let $\Pi_{m}$ be the projector on $\operatorname{Span}\left(\mathcal{F} \psi_{0}, \ldots, \mathcal{F} \psi_{m-1}\right)$. Since $\left\|\mathcal{F} \psi_{k}\right\|^{2}=2 \pi$, we get:

$$
\begin{align*}
\left\|\phi_{m}-\hat{\phi}_{m}\right\|_{\mathrm{L}^{2}}^{2} & =\sum_{k=0}^{m-1}\left(\hat{a}_{k}-a_{k}\right)^{2} \\
& =\frac{1}{(2 \pi)^{2}} \sum_{k=0}^{m-1}\left\langle\frac{\mathcal{F} \hat{h}}{1-\widehat{\mathcal{F} g}}-\frac{\mathcal{F} h}{1-\mathcal{F} g}, \mathcal{F} \psi_{k}\right\rangle^{2} \\
& =\frac{1}{2 \pi}\left\|\Pi_{m}\left(\frac{\mathcal{F} \hat{h}}{1-\widehat{\mathcal{F} g}}-\frac{\mathcal{F} h}{1-\mathcal{F} g}\right)\right\|_{\mathrm{L}^{2}}^{2} \\
& \leqslant \frac{1}{2 \pi}\left\|\frac{\mathcal{F} \hat{h}}{1-\widehat{\mathcal{F} g}}-\frac{\mathcal{F} h}{1-\mathcal{F} g}\right\|_{\mathrm{L}^{2}}^{2} . \tag{19}
\end{align*}
$$

Then since $|\widehat{\mathcal{F} g}| \leqslant \theta_{0}$ by definition, and $|\mathcal{F} g| \leqslant\|g\|_{\mathrm{L}^{1}} \leqslant \theta$, we obtain:

$$
\begin{align*}
\left\|\frac{\mathcal{F} \hat{h}}{1-\widehat{\mathcal{F} g}}-\frac{\mathcal{F} h}{1-\mathcal{F} g}\right\|_{\mathrm{L}^{2}}^{2} & \leqslant 2\left\|\frac{\mathcal{F} \hat{h}}{1-\widehat{\mathcal{F} g}}-\frac{\mathcal{F} h}{1-\widehat{\mathcal{F} g}}\right\|_{\mathrm{L}^{2}}^{2}+2\left\|\frac{\mathcal{F} h}{1-\widehat{\mathcal{F} g}}-\frac{\mathcal{F} h}{1-\mathcal{F} g}\right\|_{\mathrm{L}^{2}}^{2} \\
& \leqslant \frac{2}{\left(1-\theta_{0}\right)^{2}}\|\mathcal{F} \hat{h}-\mathcal{F} h\|_{\mathrm{L}^{2}}^{2}+\frac{2\|h\|_{\mathrm{L}^{1}}^{2}}{\left(1-\theta_{0}\right)^{2}(1-\theta)^{2}}\|\widehat{\mathcal{F} g}-\mathcal{F} g\|_{\mathrm{L}^{2}}^{2} \tag{20}
\end{align*}
$$

To control the last term, we decompose according to the set $\left\{|\mathcal{F} \hat{g}| \leqslant \theta_{0}\right\}$ and its complement:

$$
\|\widehat{\mathcal{F} g}-\mathcal{F} g\|_{\mathrm{L}^{2}}^{2} \leqslant\|\mathcal{F} \hat{g}-\mathcal{F} g\|_{\mathrm{L}^{2}}^{2}+\|g\|_{\mathrm{L}^{1}}^{2} \operatorname{Leb}\left(\left\{|\mathcal{F} \hat{g}|>\theta_{0}\right\}\right) .
$$

Thus if $\theta<\theta_{0}$, then $\left\{|\mathcal{F} \hat{g}|>\theta_{0}\right\} \subseteq\left\{|\mathcal{F} \hat{g}-\mathcal{F} g| \geqslant \theta_{0}-\theta\right\}$, therefore Markov inequality yields:

$$
\begin{equation*}
\|\widehat{\mathcal{F} g}-\mathcal{F} g\|_{\mathrm{L}^{2}}^{2} \leqslant\|\mathcal{F} \hat{g}-\mathcal{F} g\|_{\mathrm{L}^{2}}^{2}+\frac{\|g\|_{\mathrm{L}^{1}}^{2}}{\left(\theta_{0}-\theta\right)^{2}}\|\mathcal{F} \hat{g}-\mathcal{F} g\|_{\mathrm{L}^{2}}^{2} \tag{21}
\end{equation*}
$$

Finally, gathering $\sqrt[192]{20}$ and 21 , and using Plancherel theorem yield the desired result.

### 7.1 Proof of Theorem 2.5

We start with some preliminary lemmas.

Lemma 7.1. Let $Y_{1}, \ldots, Y_{n}$ be i.i.d non-negative random variables. We denote by $\mathcal{L}(s):=\mathbb{E}\left[\mathrm{e}^{-s Y_{1}}\right]$ their Laplace transform and we denote by $\widehat{\mathcal{L}}(s):=\frac{1}{n} \sum_{i=1}^{n} \mathrm{e}^{-s Y_{i}}$ the empirical Laplace transform. Then for $p \geqslant 1$, we have:

$$
\mathbb{E}\left[\sup _{s>0}|\widehat{\mathcal{L}}(s)-\mathcal{L}(s)|^{2 p}\right] \leqslant \frac{p!}{2^{p-1} n^{p}}
$$

Proof. Let $\widehat{F}(x):=\frac{1}{n} \sum_{i=1}^{n} \mathbf{1}_{\left\{Y_{i} \leqslant x\right\}}$ be the empirical c.d.f. of the $Y_{i}$ 's, and let $F(x):=\mathbb{P}[Y \leqslant x]$ be their c.d.f. We notice that for $s>0$ :

$$
\int_{0}^{+\infty} s \mathrm{e}^{-s x} \widehat{F}(x) \mathrm{d} x=\frac{1}{n} \sum_{i=1}^{n} \int_{0}^{+\infty} s \mathrm{e}^{-s x} \mathbf{1}_{\left\{Y_{i} \leqslant x\right\}} \mathrm{d} x=\frac{1}{n} \sum_{i=1}^{n} \mathrm{e}^{-s X_{i}}=: \widehat{\mathcal{L}}(s)
$$

and by the same argument, $\mathcal{L}(s)=\int_{0}^{+\infty} s \mathrm{e}^{-s x} F(x) \mathrm{d} x$. Thus:

$$
\sup _{s>0}|\widehat{\mathcal{L}}(s)-\mathcal{L}(s)| \leqslant \sup _{s>0} \int_{0}^{+\infty} s \mathrm{e}^{-s x}|\widehat{F}(x)-F(x)| \mathrm{d} x \leqslant\|\widehat{F}-F\|_{\infty}
$$

We take the expectation and we get:

$$
\begin{equation*}
\mathbb{E}\left[\sup _{s>0}|\widehat{\mathcal{L}}(s)-\mathcal{L}(s)|^{2 p}\right] \leqslant \mathbb{E}\left[\|\widehat{F}-F\|_{\infty}^{2 p}\right]=2 p \int_{0}^{+\infty} t^{2 p-1} \mathbb{P}\left[\|\widehat{F}-F\|_{\infty} \geqslant t\right] \mathrm{d} t \tag{22}
\end{equation*}
$$

By Massart $1990, \mathbb{P}\left[\sqrt{n}\|\widehat{F}-F\|_{\infty} \geqslant x\right] \leqslant 2 \mathrm{e}^{-2 x^{2}}$, so by setting $t=x / \sqrt{n}$ in 22 , we obtain:

$$
\mathbb{E}\left[\sup _{s>0}|\widehat{\mathcal{L}}(s)-\mathcal{L}(s)|^{2 p}\right] \leqslant \frac{2 p}{n^{p}} \int_{0}^{+\infty} x^{2 p-1} 2 \mathrm{e}^{-2 x^{2}} \mathrm{~d} x=\frac{2 p}{n^{p}} \int_{0}^{+\infty} u^{p-1} \mathrm{e}^{2 u} \mathrm{~d} u=\frac{p!}{n^{p} 2^{p-1}}
$$

Lemma 7.2. Let $Z \sim \mathcal{P}(\lambda)$ and $m_{j}(\lambda):=\mathbb{E}\left[(Z-\lambda)^{j}\right]$ be the $j$-th central moment of $Z$. Then, for all $r \geqslant 2$ we have:

$$
m_{r}(\lambda)=\lambda \sum_{j=0}^{r-2}\binom{r-1}{j} m_{j}(\lambda)
$$

Proof. Let $\mathcal{L}(\lambda, t):=e^{\lambda\left(\mathrm{e}^{t}-t-1\right)}=\mathbb{E}\left[\mathrm{e}^{t(Z-\lambda)}\right]$ and $\varphi(t):=\mathrm{e}^{t}-t-1$. Then $m_{r}(\lambda)=\frac{\partial^{r} \mathcal{L}}{\partial t^{r}}(\lambda, 0)$. By Leibniz's rule:

$$
\frac{\partial^{r} \mathcal{L}}{\partial t^{r}}(\lambda, t)=\frac{\partial^{r-1}}{\partial t^{r-1}}\left(\lambda \varphi^{\prime}(t) \mathcal{L}(\lambda, t)\right)=\lambda \sum_{j=0}^{r-1}\binom{r-1}{j} \frac{\partial^{j} \mathcal{L}}{\partial t^{j}}(\lambda, t) \varphi^{(r-j)}(t)
$$

Taking $t=0$ gives the result since $\varphi^{\prime}(0)=0$ and $\varphi^{(k)}(0)=1$ if $k \geqslant 2$.
Corollary 7.3. The central moments $m_{2 r}(\lambda)$ and $m_{2 r+1}(\lambda)$ are polynomials of degree $r$ in $\lambda$.
The next proposition provides an upper bound on the $\mathrm{L}^{p}$-risk of $\hat{\rho}_{\delta}$.
Proposition 7.4. Under Assumption 11, for $p \geqslant 1$, we have:

$$
\mathbb{E}\left[\left(\hat{\rho}_{\delta}-\rho_{\delta}\right)^{2 p}\right] \leqslant \frac{C(p, \lambda)}{c^{2 p}(1-\theta)^{2 p} T^{p}}
$$

where $C(p, \lambda)$ is a $\mathcal{O}\left(\lambda^{p}\right)$.
Proof. By definition, $\rho_{\delta}$ is a solution of the Lundberg equation, so it is a zero of the function:

$$
\ell_{\delta}(s):=c s-(\lambda+\delta)+\lambda \mathcal{L} f(s)
$$

The estimator $\hat{\rho}_{\delta}$ is then a zero of the function:

$$
\hat{\ell}_{\delta}(s):=c s-(\hat{\lambda}+\delta)+\hat{\lambda} \widehat{\mathcal{L} f}(s)
$$

We use a Taylor-Lagrange expansion:

$$
\hat{\ell}_{\delta}\left(\hat{\rho}_{\delta}\right)=0=\ell_{\delta}\left(\rho_{\delta}\right)=\ell_{\delta}\left(\hat{\rho}_{\delta}\right)+\ell_{\delta}^{\prime}(z)\left(\rho_{\delta}-\hat{\rho}_{\delta}\right),
$$

where $z$ is between $\rho_{\delta}$ and $\hat{\rho}_{\delta}$.

$$
\left|\ell_{\delta}^{\prime}(z)\right|=\left|c-\lambda \int_{0}^{+\infty} x \mathrm{e}^{-z x} f(x) \mathrm{d} x\right| \geqslant c-\lambda \int_{0}^{+\infty} x f(x) \mathrm{d} x=c-\lambda \mu>0
$$

under the safety loading condition. Thus:

$$
\begin{aligned}
\left|\rho_{\delta}-\hat{\rho}_{\delta}\right| & \leqslant \frac{1}{c-\lambda \mu}\left|\hat{\ell}_{\delta}\left(\hat{\rho}_{\delta}\right)-\ell_{\delta}\left(\hat{\rho}_{\delta}\right)\right| \\
& =\frac{1}{c(1-\theta)}\left|\hat{\lambda}\left(\widehat{\mathcal{L} f}\left(\hat{\rho}_{\delta}\right)-\mathcal{L}\left(\hat{\rho}_{\delta}\right)\right)+(\hat{\lambda}-\lambda)\left(1-\mathcal{L} f\left(\hat{\rho}_{\delta}\right)\right)\right| \\
& \leqslant \frac{1}{c(1-\theta)}\left(|\hat{\lambda}|\|\widehat{\mathcal{L} f}-\mathcal{L} f\|_{\infty}+2|\hat{\lambda}-\lambda|\right) \\
\mathbb{E}\left[\left(\hat{\rho}_{\delta}-\rho_{\delta}\right)^{2 p}\right] & \leqslant \frac{1}{c^{2 p}(1-\theta)^{2 p}}\left(2^{2 p-1} \mathbb{E}\left[\hat{\lambda}^{2 p}\|\widehat{\mathcal{L} f}-\mathcal{L} f\|_{\infty}^{2 p}\right]+2^{4 p-1} \mathbb{E}\left[|\hat{\lambda}-\lambda|^{2 p}\right]\right) .
\end{aligned}
$$

For the second term, we use Corollary $7.3 \mathbb{E}|\hat{\lambda}-\lambda|^{2 p}=\frac{\mathbb{E}\left|N_{T}-\lambda T\right|^{2 p}}{T^{2 p}}=\frac{\mathcal{O}\left(\lambda^{p}\right)}{T^{p}}$. For the first term, we apply Lemma 7.1 conditional to $N_{T}$ :

$$
\mathbb{E}\left[\hat{\lambda}^{2 p}\|\widehat{\mathcal{L} f}-\mathcal{L} f\|_{\infty}^{2 p}\right] \leqslant \mathbb{E}\left[C(p) \frac{\hat{\lambda}^{2 p}}{N_{T}^{p}}\right]=\frac{C(p)}{T^{p}} \mathbb{E}\left[\hat{\lambda}^{p}\right]=\frac{\mathcal{O}\left(\lambda^{p}\right)}{T^{p}}
$$

Finally:

$$
\mathbb{E}\left[\left(\hat{\rho}_{\delta}-\rho_{\delta}\right)^{2 p}\right] \leqslant \frac{C(p, \lambda)}{c^{2 p}(1-\theta)^{2 p} T^{p}},
$$

with $C(p, \lambda)=\mathcal{O}\left(\lambda^{p}\right)$.
Now, we can prove Theorem 2.5
Proof of Theorem 2.5. By Pythagoras theorem:

$$
\begin{aligned}
& \mathbb{E}\left\|g-\hat{g}_{m_{2}}\right\|_{\mathrm{L}^{2}}^{2}=\left\|g-g_{m_{2}}\right\|_{\mathrm{L}^{2}}^{2}+\mathbb{E}\left\|g_{m_{2}}-\hat{g}_{m_{2}}\right\|_{\mathrm{L}^{2}}^{2}=\left\|g-g_{m_{2}}\right\|_{\mathrm{L}^{2}}^{2}+\sum_{k=0}^{m_{2}-1} \mathbb{E}\left[\left(\hat{b}_{k}-b_{k}\right)^{2}\right], \\
& \mathbb{E}\left\|h-\hat{h}_{m_{3}}\right\|_{\mathrm{L}^{2}}^{2}=\left\|h-h_{m_{3}}\right\|_{\mathrm{L}^{2}}^{2}+\mathbb{E}\left\|h_{m_{3}}-\hat{h}_{m_{3}}\right\|_{\mathrm{L}^{2}}^{2}=\left\|h-h_{m_{3}}\right\|_{\mathrm{L}^{2}}^{2}+\sum_{k=0}^{m_{3}-1} \mathbb{E}\left[\left(\hat{c}_{k}-c_{k}\right)^{2}\right],
\end{aligned}
$$

hence we need to control the variance terms $\sum_{k=0}^{m_{2}-1} \mathbb{E}\left[\left(\hat{b}_{k}-b_{k}\right)^{2}\right]$ and $\sum_{k=0}^{m_{3}-1} \mathbb{E}\left[\left(\hat{c}_{k}-c_{k}\right)^{2}\right]$.
Using equations (4.17) to (4.21) and (4.10) to (4.14) in Zhang and Su 2018], we can obtain equations 28 and 29 below. Still, we give the proofs of these equations for the sake of completeness.

We notice that $\hat{b}_{k}$ and $\hat{c}_{k}$ (defined by (10) and ) can be written as:

$$
\frac{1}{c T} \sum_{i=1}^{N_{T}} \int_{0}^{+\infty} F\left(u, X_{i}, \hat{\rho}_{\delta}\right) \psi_{k}(u) \mathrm{d} u
$$

and that the coefficients $b_{k}$ and $c_{k}$ (defined by (8) and (9)) can be written as:

$$
\mathbb{E}\left[\frac{1}{c T} \sum_{i=1}^{N_{T}} \int_{0}^{+\infty} F\left(u, X_{i}, \rho_{\delta}\right) \psi_{k}(u) \mathrm{d} u\right],
$$

where $F$ is given by:

$$
F(u, X, \rho):= \begin{cases}\mathrm{e}^{-\rho(X-u)} \mathbf{1}_{X>u} & \text { for the coefficients of } g  \tag{23}\\ \int_{u}^{X} \mathrm{e}^{-\rho(X-x)} w(x, X-x) \mathrm{d} x \mathbf{1}_{X>u} & \text { for the coefficients of } h\end{cases}
$$

Thus, we need to give an upper bound on quantities of the form:

$$
\begin{equation*}
V_{m}:=\sum_{k=0}^{m-1} \mathbb{E}\left[\left(\frac{1}{c T} \sum_{i=1}^{N_{T}} \int_{0}^{+\infty} F\left(u, X_{i}, \hat{\rho}_{\delta}\right) \psi_{k}(u) \mathrm{d} u-\mathbb{E}\left[\frac{1}{c T} \sum_{i=1}^{N_{T}} \int_{0}^{+\infty} F\left(u, X_{i}, \rho_{\delta}\right) \psi_{k}(u) \mathrm{d} u\right]\right)^{2}\right] \tag{24}
\end{equation*}
$$

The bound on $V_{m}$ is based on the following decomposition:

$$
\begin{equation*}
\frac{1}{c T} \sum_{i=1}^{N_{T}} \int_{0}^{+\infty} F\left(u, X_{i}, \hat{\rho}_{\delta}\right) \psi_{k}(u) \mathrm{d} u-\mathbb{E}\left[\frac{1}{c T} \sum_{i=1}^{N_{T}} \int_{0}^{+\infty} F\left(u, X_{i}, \rho_{\delta}\right) \psi_{k}(u) \mathrm{d} u\right]=Z_{k}+\Delta_{k} \tag{25}
\end{equation*}
$$

where:

$$
\begin{aligned}
Z_{k} & :=\frac{1}{c T}\left(\sum_{i=1}^{N_{T}} \int_{0}^{+\infty} F\left(u, X_{i}, \hat{\rho}_{\delta}\right) \psi_{k}(u) \mathrm{d} u-\mathbb{E}\left[\sum_{i=1}^{N_{T}} \int_{0}^{+\infty} F\left(u, X_{i}, \hat{\rho}_{\delta}\right) \psi_{k}(u) \mathrm{d} u\right]\right) \\
\Delta_{k} & :=\frac{1}{c T} \sum_{i=1}^{N_{T}} \int_{0}^{+\infty}\left[F\left(u, X_{i}, \hat{\rho}_{\delta}\right)-F\left(u, X_{i}, \rho_{\delta}\right)\right] \psi_{k}(u) \mathrm{d} u .
\end{aligned}
$$

Let us notice that if $\delta=0$, then $\hat{\rho}_{\delta}=\rho_{\delta}=0$, so $\Delta_{k}=0$ and the decomposition reduces to $Z_{k}$.

- Bound on $\sum_{k=0}^{m-1} \mathbb{E}\left[Z_{k}^{2}\right]$. This bound is obtained by a projection argument:

$$
\begin{aligned}
\sum_{k=0}^{m-1} \mathbb{E}\left[Z_{k}^{2}\right] & =\sum_{k=0}^{m-1} \operatorname{Var}\left(\frac{1}{c T} \sum_{i=1}^{N_{T}} \int_{0}^{+\infty} F\left(u, X_{i}, \rho_{\delta}\right) \psi_{k}(u) \mathrm{d} u\right) \\
& =\sum_{k=0}^{m-1} \frac{\lambda}{c^{2} T} \mathbb{E}\left[\left(\int_{0}^{+\infty} F\left(u, X, \rho_{\delta}\right) \psi_{k}(u) \mathrm{d} u\right)^{2}\right] \\
& \leqslant \frac{\lambda}{c^{2} T} \mathbb{E}\left[\int_{0}^{+\infty} F\left(u, X, \rho_{\delta}\right)^{2} \mathrm{~d} u\right],
\end{aligned}
$$

where the last inequality comes from the fact that $\left(\psi_{k}\right)_{k \geqslant 0}$ is an orthonormal basis of $\mathrm{L}^{2}\left(\mathbb{R}_{+}\right)$.
From (23), we see that:

$$
\frac{\lambda}{c^{2} T} \mathbb{E}\left[\int_{0}^{+\infty} F\left(u, X, \rho_{\delta}\right)^{2} \mathrm{~d} u\right] \leqslant \begin{cases}\frac{\lambda}{c^{2} T} \mathbb{E}[X] & \text { for the coefficients of } g  \tag{26}\\ \frac{\lambda}{c^{2} T} \mathbb{E}[W(X)] & \text { for the coefficients of } h\end{cases}
$$

where $W(X)$ is defined in Assumption 3. In the $\delta=0$ case, this gives the desired results.

- Bound on $\sum_{k=0}^{m-1} \Delta_{k}^{2}$. We use a projection argument again:

$$
\begin{aligned}
\sum_{k=0}^{m-1} \Delta_{k}^{2} & \leqslant \sum_{k=0}^{m-1} \frac{N_{T}}{c^{2} T^{2}} \sum_{i=1}^{N_{T}}\left(\int_{0}^{+\infty}\left[F\left(u, X_{i}, \hat{\rho}_{\delta}\right)-F\left(u, X_{i}, \rho_{\delta}\right)\right] \psi_{k}(u) \mathrm{d} u\right)^{2} \\
& \leqslant \frac{\hat{\lambda}}{c^{2} T} \sum_{i=1}^{N_{T}} \int_{0}^{+\infty}\left|F\left(u, X_{i}, \hat{\rho}_{\delta}\right)-F\left(u, X_{i}, \rho_{\delta}\right)\right|^{2} \mathrm{~d} u
\end{aligned}
$$

where $\hat{\lambda}:=\frac{N_{T}}{T}$. By Remark 1.2 we know that $\rho_{\delta} \in\left[\frac{\delta}{c}, \frac{\delta+\lambda}{c}\right]$ and $\hat{\rho}_{\delta} \in\left[\frac{\delta}{c}, \frac{\delta+\hat{\lambda}}{c}\right]$, so by the mean value theorem:

$$
\left|F\left(u, X_{i}, \hat{\rho}_{\delta}\right)-F\left(u, X_{i}, \rho_{\delta}\right)\right| \leqslant\left|\hat{\rho}_{\delta}-\rho_{\delta}\right| \sup _{\rho \geqslant \frac{\delta}{c}}\left|\frac{\partial F}{\partial \rho}\left(u, X_{i}, \rho\right)\right| .
$$

Since the function $t \mathrm{e}^{-\rho t} \mathbf{1}_{t>0}$ achieves its maximum at $t=\frac{1}{\rho}$, we see that:

$$
\sup _{\rho \geqslant \frac{\delta}{c}}\left|\frac{\partial F}{\partial \rho}\left(u, X_{i}, \rho\right)\right| \leqslant \begin{cases}\frac{c}{\mathrm{e} \delta} \mathbf{1}_{X_{i}>u} & \text { for the coefficients of } g \\ \frac{c}{\mathrm{e} \delta} \int_{u}^{X_{i}} w\left(x, X_{i}-x\right) \mathrm{d} x \mathbf{1}_{X_{i}>u} & \text { for the coefficients of } h\end{cases}
$$

Thus,

$$
\sum_{k=0}^{m-1} \Delta_{k}^{2} \leqslant \frac{\hat{\lambda}\left|\hat{\rho}_{\delta}-\rho_{\delta}\right|^{2}}{\mathrm{e}^{2} \delta^{2}} \times \begin{cases}\frac{1}{T} \sum_{i=1}^{N_{T}} X_{i} & \text { for the coefficients of } g  \tag{27}\\ \frac{1}{T} \sum_{i=1}^{N_{T}} W\left(X_{i}\right) & \text { for the coefficients of } h\end{cases}
$$

Using the decomposition (25) in 24 , we get $V_{m} \leqslant 2 \sum_{k=0}^{m-1} \mathbb{E}\left[Z_{k}^{2}\right]+2 \sum_{k=0}^{m-1} \mathbb{E}\left[\Delta_{k}^{2}\right]$. Combining (26) and (27) yields:

$$
\begin{gather*}
\mathbb{E}\left\|\hat{g}_{m_{2}}-g_{m_{2}}\right\|_{\mathrm{L}^{2}}^{2} \leqslant 2 \frac{\lambda}{c^{2} T} \mathbb{E}[X]+2 \mathbb{E}\left[\frac{\hat{\lambda}\left|\hat{\rho}_{\delta}-\rho_{\delta}\right|^{2}}{\mathrm{e}^{2} \delta^{2}} \frac{1}{T} \sum_{i=1}^{N_{T}} X_{i}\right],  \tag{28}\\
\mathbb{E}\left\|\hat{h}_{m_{3}}-h_{m_{3}}\right\|_{\mathrm{L}^{2}}^{2} \leqslant 2 \frac{\lambda}{c^{2} T} \mathbb{E}[W(X)]+2 \mathbb{E}\left[\frac{\hat{\lambda}\left|\hat{\rho}_{\delta}-\rho_{\delta}\right|^{2}}{\mathrm{e}^{2} \delta^{2}} \frac{1}{T} \sum_{i=1}^{N_{T}} W\left(X_{i}\right)\right] . \tag{29}
\end{gather*}
$$

We apply Hölder's inequality on the second term in 28 and we use Proposition 7.4 .

$$
\begin{aligned}
\mathbb{E}\left[\hat{\lambda}\left|\hat{\rho}_{\delta}-\rho_{\delta}\right|^{2} \frac{1}{T} \sum_{i=1}^{N_{T}} X_{i}\right] & \leqslant \mathbb{E}\left[\hat{\lambda}^{4}\right]^{1 / 4} \mathbb{E}\left[\left|\hat{\rho}_{\delta}-\rho_{\delta}\right|^{8}\right]^{1 / 4} \mathbb{E}\left[\left(\frac{1}{T} \sum_{i=1}^{N_{T}} X_{i}\right)^{2}\right]^{1 / 2} \\
& \left.\leqslant \frac{C(\lambda)}{c^{2}(1-\theta)^{2} T} \mathbb{E}\left[\left(\frac{1}{T} \sum_{i=1}^{N_{T}} X_{i}\right)\right)^{2}\right]^{1 / 2}
\end{aligned}
$$

with $C(\lambda)=\mathcal{O}\left(\lambda^{2}\right)$. We need to evaluate this last expectation:

$$
\mathbb{E}\left[\left(\frac{1}{T} \sum_{i=1}^{N_{T}} X_{i}\right)^{2}\right] \leqslant \mathbb{E}\left[\frac{N_{T}}{T^{2}} \sum_{i=1}^{N_{T}} X_{i}^{2}\right]=\mathbb{E}\left[\frac{N_{T}^{2}}{T^{2}}\right] \mathbb{E}\left[X^{2}\right]=\left(\frac{\lambda}{T}+\lambda^{2}\right) \mathbb{E}\left[X^{2}\right]
$$

Thus, we obtain:

$$
\mathbb{E}\left\|\hat{g}_{m}-g_{m}\right\|_{\mathrm{L}^{2}}^{2} \leqslant 2 \frac{\lambda}{c^{2} T} \mathbb{E}[X]+2 \frac{C(\lambda)}{c^{2} T(1-\theta)^{2} \delta^{2}} \mathbb{E}\left[X^{2}\right]^{1 / 2}
$$

with $C(\lambda)=\mathcal{O}\left(\lambda^{2}\right)$. We make the same reasoning for $h$, replacing $X_{i}$ by $W\left(X_{i}\right)$.

### 7.2 Proofs of Section 3

Let us recall some facts about Toeplitz matrices; the interested reader can find more details in the book of Böttcher and Grudsky 2000. Given $\left(\alpha_{n}\right)_{n \in \mathbb{Z}}$ a sequence of complex numbers, a Toeplitz matrix is an infinite matrix of the form:

$$
\left(\begin{array}{cccc}
\alpha_{0} & \alpha_{-1} & \alpha_{-2} & \cdots  \tag{30}\\
\alpha_{1} & \alpha_{0} & \alpha_{-1} & \cdots \\
\alpha_{2} & \alpha_{1} & \alpha_{0} & \cdots \\
\cdots & \cdots & \cdots & \cdots
\end{array}\right)
$$

The classical result from O . Toeplitz says that this matrix induces a bounded operator on $\ell^{2}(\mathbb{N})$ if and only if $\left(\alpha_{n}\right)_{n \in \mathbb{Z}}$ are the Fourier coefficients of some function $\alpha \in \mathrm{L}^{\infty}(\mathbb{T})$, where $\mathbb{T}$ denotes the complex unit circle. We denote both the matrix (30) and its induced operator on $\ell^{2}(\mathbb{N})$ by $\mathbf{T}(\alpha)$, the function $\alpha$ being called the symbol of the Toeplitz matrix. Finally, if $m \in \mathbb{N}^{*}$ and if $\mathbf{T}(\alpha)$ is a Toeplitz matrix, we denote by $\mathbf{T}_{m}(\alpha)$ the $m \times m$ matrix:

$$
\mathbf{T}_{m}(\alpha):=\left(\begin{array}{ccc}
\alpha_{0} & \cdots & \alpha_{-(m-1)}  \tag{31}\\
\vdots & \ddots & \vdots \\
\alpha_{m-1} & \cdots & \alpha_{0}
\end{array}\right)
$$

The operator norm of $\mathbf{T}(\alpha)$ depends on the properties of its symbol. In the case where $\alpha_{k}=0$ for all $k<0$, we have the following lemma.

Lemma 7.5. Let $\left(\alpha_{k}\right)_{k \geqslant 0} \in \ell^{1}(\mathbb{N})$ be a sequence of complex numbers. Then the Toeplitz matrix $\mathbf{T}(\alpha)$ is lower triangular and we have:

$$
\forall x \in \ell^{2}(\mathbb{N}), \quad \mathbf{T}(\alpha) \times x=\alpha * x
$$

In particular, $\|\mathbf{T}(\alpha)\|_{\mathrm{op}} \leqslant\|\alpha\|_{\ell^{1}}$.
Proof. The fact that $\mathbf{T}(\alpha)$ is lower triangular and that $\mathbf{T}(\alpha) \times x=\alpha * x$ is clear from the definition of a Toeplitz matrix. Then, Young's inequality for convolution yields $\|\alpha * x\|_{\ell^{2}} \leqslant\|\alpha\|_{\ell^{1}}\|x\|_{\ell^{2}}$.

Concerning the inverse of a Toeplitz matrix, its norm depends on the position of zero relatively to the range of the symbol. More precisely, we use the following result.

Lemma 7.6 (Lemma 3.8 in Böttcher and Grudsky 2000). Let $\alpha \in \mathrm{L}^{\infty}(\mathbb{T})$ and let $E(\alpha)$ be the convex hull of the essential range of $\alpha$. If $d:=\operatorname{dist}(0, E(\alpha))>0$, then $\mathbf{T}_{m}(\alpha)$ is invertible for all $m \geqslant 1$, and we have $\left\|\mathbf{T}_{m}^{-1}(\alpha)\right\|_{\mathrm{op}}<\frac{2}{d}$.
The matrix $\mathbf{A}_{m}$ defined by (14) is a Toeplitz matrix and its symbol is given by:

$$
\alpha(t):=\sum_{k=0}^{+\infty} \alpha_{k} t^{k}, \quad \text { with } \quad \alpha_{k}:= \begin{cases}1-\frac{b_{0}}{\sqrt{2}} & \text { if } k=0 \\ \frac{b_{k-1}-b_{k}}{\sqrt{2}} & \text { if } k \geqslant 1\end{cases}
$$

Let us notice that under Assumption 4, we have $\left(\alpha_{k}\right)_{k \geqslant 0} \in \ell^{1}(\mathbb{N})$ so the symbol $\alpha$ is continuous on $\mathbb{T}$, and thus $\alpha \in \mathrm{L}^{\infty}(\mathbb{T})$.

Proof of Lemma 3.1. We apply Lemma C. $1^{2}$ in Comte et al. 2017 to the coefficients of $g$ : the sequence $\left(\beta_{k}\right)_{k \geqslant 0}$, defined by $\beta_{0}:=\frac{b_{0}}{\sqrt{2}}$ and $\beta_{k}:=\frac{b_{k}-b_{k-1}}{\sqrt{2}}$ for $k \geqslant 1$, are the Fourier coefficients of the function $t \in \mathbb{T} \mapsto \mathcal{L} g\left(\frac{1+t}{1-t}\right) \in \mathbb{C}$. Thus, we have:

$$
\forall t \in \mathbb{T}, \quad \mathcal{L} g\left(\frac{1+t}{1-t}\right)=\sum_{k=0}^{+\infty} \beta_{k} t^{k}
$$

with the convention $\mathcal{L} g(\infty)=0$. Since $\alpha(t)=1-\sum_{k \geqslant 0} \beta_{k} t^{k}$, we get:

$$
\forall t \in \mathbb{T}, \quad \alpha(t)=1-\mathcal{L} g\left(\frac{1+t}{1-t}\right)
$$

We notice that if $t \in \mathbb{T} \backslash\{1\}$, then there exists $\omega \in \mathbb{R}$ such that $\frac{1+t}{1-t}=\mathrm{i} \omega$. Thus:

$$
\begin{aligned}
\forall t \in \mathbb{T} \backslash\{1\}, \quad \mathfrak{R e} \alpha(t) & =1-\mathfrak{R e}\left[\mathcal{L} g\left(\frac{1+t}{1-t}\right)\right] \\
& =1-\mathfrak{R e}[\mathcal{L} g(\mathrm{i} \omega)] \\
& =1-\mathfrak{R e}\left[\int_{0}^{+\infty} \mathrm{e}^{-\mathrm{i} \omega x} g(x) \mathrm{d} x\right] \\
& =1-\int_{0}^{+\infty} \cos (\omega x) g(x) \mathrm{d} x \\
& \geqslant 1-\int_{0}^{+\infty} g(x) \mathrm{d} x \geqslant 1-\theta
\end{aligned}
$$

This inequality holds for $t=1$ as well, hence $\alpha(\mathbb{T})$ is included in the half-plane $\{z \in \mathbb{C} \mid \mathfrak{R e}(z) \geqslant$ $1-\theta\}$, and so is its convex hull. By Lemma 7.6 .

$$
\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}} \leqslant \frac{2}{1-\theta}
$$

Remark 7.7. In their article, Zhang and Su 2018 show that $\inf _{|z|=1}|\alpha(z)| \geqslant 1-\theta>0$, that is $\operatorname{dist}(0, \alpha(\mathbb{T}))>0$, which is not sufficient to apply Lemma 7.6

[^2]
### 7.2.1 Proof of Theorem 3.2

Proposition 7.8. Under Assumption 4 if $\theta<\theta_{0}$ then it holds:

$$
\forall m \in \mathbb{N}^{*}, \quad \mathbb{E}\left\|\left(\widetilde{\mathbf{A}}_{m, 1}^{-1}-\mathbf{A}_{m}^{-1}\right) \mathbf{c}_{m}\right\|_{\ell^{2}}^{2} \leqslant \frac{C\left(\theta, \theta_{0}\right)}{\left(1-\theta_{0}\right)^{2}}\left\|\phi_{m}\right\|_{\mathrm{L}^{2}}^{2} \mathbb{E}\left\|\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right\|_{\mathrm{op}}^{2},
$$

where $C\left(\theta, \theta_{0}\right)$ is a constant satisfying $C\left(\theta, \theta_{0}\right) \lesssim \frac{\left(1-\theta_{0}\right)^{2}}{\left(\theta_{0}-\theta\right)^{2}}$.
Proof. We decompose the expectation according to the event $\Delta_{m}^{1}:=\left\{\left\|\hat{\mathbf{A}}_{m}^{-1}\right\|_{\mathrm{op}} \leqslant \frac{2}{1-\theta_{0}}\right\}$ :

$$
\begin{aligned}
\mathbb{E}\left\|\left(\widetilde{\mathbf{A}}_{m, 1}^{-1}-\mathbf{A}_{m}^{-1}\right) \mathbf{c}_{m}\right\|_{\ell^{2}}^{2} & =\left\|\mathbf{A}_{m}^{-1} \mathbf{c}_{m}\right\|_{\ell^{2}}^{2} \mathbb{P}\left[\Delta_{m}^{c}\right]+\mathbb{E}\left[\left\|\widetilde{\mathbf{A}}_{m, 1}^{-1}\left(\mathbf{A}_{m}-\hat{\mathbf{A}}_{m}\right) \mathbf{A}_{m}^{-1} \mathbf{c}_{m}\right\|_{\ell^{2}}^{2} \mathbf{1}_{\Delta_{m}}\right] \\
& \leqslant\left\|\mathbf{A}_{m}^{-1} \mathbf{c}_{m}\right\|_{\ell^{2}}^{2}\left(\mathbb{P}\left[\Delta_{m}^{c}\right]+\frac{4}{\left(1-\theta_{0}\right)^{2}} \mathbb{E}\left\|\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right\|_{\text {op }}^{2}\right) \\
& =\left\|\mathbf{a}_{m}\right\|_{\ell^{2}}^{2}\left(\mathbb{P}\left[\Delta_{m}^{c}\right]+\frac{4}{\left(1-\theta_{0}\right)^{2}} \mathbb{E}\left\|\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right\|_{\text {op }}^{2}\right)
\end{aligned}
$$

Since $\theta<\theta_{0}$ and $\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}} \leqslant \frac{2}{1-\theta}$ (see Lemma 3.1), we get:

$$
\begin{aligned}
\mathbb{P}\left[\Delta_{m}^{c}\right] \leqslant & \mathbb{P}\left[\left\|\hat{\mathbf{A}}_{m}^{-1}-\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}>\frac{2}{1-\theta_{0}}-\frac{2}{1-\theta}\right] \\
= & \mathbb{P}\left[\left\{\left\|\hat{\mathbf{A}}_{m}^{-1}-\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}>\frac{2}{1-\theta_{0}}-\frac{2}{1-\theta}\right\} \cap\left\{\left\|\mathbf{A}_{m}^{-1}\left(\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right)\right\|_{\mathrm{op}}<\frac{1}{2}\right\}\right] \\
& +\mathbb{P}\left[\left\{\left\|\hat{\mathbf{A}}_{m}^{-1}-\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}>\frac{2}{1-\theta_{0}}-\frac{2}{1-\theta}\right\} \cap\left\{\left\|\mathbf{A}_{m}^{-1}\left(\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right)\right\|_{\mathrm{op}} \geqslant \frac{1}{2}\right\}\right]
\end{aligned}
$$

- First term. We apply Theorem C.1 and we conclude using Markov's inequality:

$$
\begin{aligned}
& \mathbb{P}\left[\left\{\left\|\hat{\mathbf{A}}_{m}^{-1}-\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}>\frac{2}{1-\theta_{0}}-\frac{2}{1-\theta}\right\} \cap\left\{\left\|\mathbf{A}_{m}^{-1}\left(\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right)\right\|_{\mathrm{op}}<\frac{1}{2}\right\}\right] \\
& \leqslant \mathbb{P}\left[\left\{\frac{\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2}\left\|\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right\|_{\mathrm{op}}}{1-\left\|\mathbf{A}_{m}^{-1}\left(\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right)\right\|_{\mathrm{op}}}>\frac{2}{1-\theta_{0}}-\frac{2}{1-\theta}\right\} \cap\left\{\left\|\mathbf{A}_{m}^{-1}\left(\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right)\right\|_{\mathrm{op}}<\frac{1}{2}\right\}\right] \\
& \leqslant \mathbb{P}\left[\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2}\left\|\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right\|_{\mathrm{op}}>\frac{1}{1-\theta_{0}}-\frac{1}{1-\theta}\right] \\
& \leqslant \frac{(1-\theta)^{2}\left(1-\theta_{0}\right)^{2}}{\left(\theta_{0}-\theta\right)^{2}}\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{4} \mathbb{E}\left\|\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right\|_{\mathrm{op}}^{2} \leqslant \frac{16\left(1-\theta_{0}\right)^{2}}{\left(\theta_{0}-\theta\right)^{2}(1-\theta)^{2}} \mathbb{E}\left\|\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right\|_{\mathrm{op}}^{2}
\end{aligned}
$$

- Second term. We use Markov's inequality:

$$
\begin{aligned}
& \mathbb{P}\left[\left\{\left\|\hat{\mathbf{A}}_{m}^{-1}-\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}>\frac{2}{1-\theta_{0}}-\frac{2}{1-\theta}\right\} \cap\left\{\left\|\mathbf{A}_{m}^{-1}\left(\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right)\right\|_{\mathrm{op}} \geqslant \frac{1}{2}\right\}\right] \\
& \leqslant \mathbb{P}\left[\left\|\mathbf{A}_{m}^{-1}\left(\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right)\right\|_{\mathrm{op}} \geqslant \frac{1}{2}\right] \\
& \leqslant 4\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2} \mathbb{E}\left\|\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right\|_{\mathrm{op}}^{2} \leqslant \frac{16}{(1-\theta)^{2}} \mathbb{E}\left\|\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right\|_{\mathrm{op}}^{2} .
\end{aligned}
$$

Thus, we obtain:

$$
\mathbb{E}\left\|\left(\widetilde{\mathbf{A}}_{m, 1}^{-1}-\mathbf{A}_{m}^{-1}\right) \mathbf{c}_{m}\right\|_{\ell^{2}}^{2} \leqslant\left\|\phi_{m}\right\|_{\mathrm{L}^{2}}^{2}\left(\frac{16}{(1-\theta)^{2}}\left(1+\frac{\left(1-\theta_{0}\right)^{2}}{\left(\theta_{0}-\theta\right)^{2}}\right)+\frac{4}{\left(1-\theta_{0}\right)^{2}}\right) \mathbb{E}\left\|\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right\|_{\mathrm{op}}^{2}
$$

We can now prove Theorem 3.2
Proof of Theorem 3.2. By Pythagoras Theorem, $\left\|\phi-\hat{\phi}_{m}^{\mathrm{Lag}_{1}}\right\|_{\mathrm{L}^{2}}^{2}=\left\|\phi-\phi_{m}\right\|_{\mathrm{L}^{2}}^{2}+\left\|\phi_{m}-\hat{\phi}_{m}^{\mathrm{Lag}}\right\|_{\mathrm{L}^{2}}^{2}$.

$$
\begin{aligned}
\mathbb{E} \| \phi_{m}- & \hat{\phi}_{m}^{\mathrm{Lag}_{1}}\left\|_{\mathrm{L}^{2}}^{2}=\mathbb{E}\right\| \hat{\mathbf{a}}_{m}^{\mathrm{Lag}_{1}}-\mathbf{a}_{m}\left\|_{\ell^{2}}^{2}=\mathbb{E}\right\| \widetilde{\mathbf{A}}_{m, 1}^{-1} \hat{\mathbf{c}}_{m}-\mathbf{A}_{m}^{-1} \mathbf{c}_{m} \|_{\ell^{2}}^{2} \\
& \leqslant 3 \mathbb{E}\left\|\left(\widetilde{\mathbf{A}}_{m, 1}^{-1}-\mathbf{A}_{m}^{-1}\right) \mathbf{c}_{m}\right\|_{\ell^{2}}^{2}+3 \mathbb{E}\left\|\left(\mathbf{A}_{m}^{-1}-\widetilde{\mathbf{A}}_{m, 1}^{-1}\right)\left(\mathbf{c}_{m}-\hat{\mathbf{c}}_{m}\right)\right\|_{\ell^{2}}^{2}+3 \mathbb{E}\left\|\mathbf{A}_{m}^{-1}\left(\mathbf{c}_{m}-\hat{\mathbf{c}}_{m}\right)\right\|_{\ell^{2}}^{2}
\end{aligned}
$$

- First term. We apply Proposition 7.8 with Lemma 7.5

$$
\begin{aligned}
\mathbb{E}\left\|\left(\widetilde{\mathbf{A}}_{m, 1}^{-1}-\mathbf{A}_{m}^{-1}\right) \mathbf{c}_{m}\right\|_{\ell^{2}}^{2} & \leqslant C\left(\theta, \theta_{0}\right)\left\|\phi_{m}\right\|_{\mathrm{L}^{2}}^{2} \mathbb{E}\left\|\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right\|_{\mathrm{L}^{2}}^{2} \\
& \leqslant C\left(\theta, \theta_{0}\right)\left\|\phi_{m}\right\|_{\mathrm{L}^{2}}^{2} \mathbb{E}\left\|\hat{\mathbf{b}}_{m}-\mathbf{b}_{m}\right\|_{\ell^{1}}^{2} \\
& \leqslant C\left(\theta, \theta_{0}\right)\left\|\phi_{m}\right\|_{\mathrm{L}^{2}}^{2} m \mathbb{E}\left\|\hat{\mathbf{b}}_{m}-\mathbf{b}_{m}\right\|_{\ell^{2}}^{2} \\
& =C\left(\theta, \theta_{0}\right)\|\phi\|_{\mathrm{L}^{2}}^{2} m \mathbb{E}\left\|\hat{g}_{m}-g_{m}\right\|_{\mathrm{L}^{2}}^{2} .
\end{aligned}
$$

- Second term.

$$
\begin{aligned}
\mathbb{E}\left\|\left(\mathbf{A}_{m}^{-1}-\widetilde{\mathbf{A}}_{m, 1}^{-1}\right)\left(\mathbf{c}_{m}-\hat{\mathbf{c}}_{m}\right)\right\|_{\ell^{2}}^{2} & \leqslant \mathbb{E}\left[\left\|\mathbf{A}_{m}^{-1}-\widetilde{\mathbf{A}}_{m, 1}^{-1}\right\|_{\mathrm{op}}^{2}\left\|\mathbf{c}_{m}-\hat{\mathbf{c}}_{m}\right\|_{\ell^{2}}^{2}\right] \\
& \leqslant\left(\frac{8}{(1-\theta)^{2}}+\frac{8}{\left(1-\theta_{0}\right)^{2}}\right) \mathbb{E}\left\|\hat{h}_{m}-h_{m}\right\|_{\mathrm{L}^{2}}^{2} .
\end{aligned}
$$

- Third term.

$$
\begin{aligned}
& \mathbb{E}\left\|\mathbf{A}_{m}^{-1}\left(\mathbf{c}_{m}-\hat{\mathbf{c}}_{m}\right)\right\|_{\ell^{2}}^{2} \leqslant\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2} \mathbb{E}\left\|\mathbf{c}_{m}-\hat{\mathbf{c}}_{m}\right\|_{\ell^{2}}^{2} \\
& \leqslant \frac{4}{(1-\theta)^{2}} \mathbb{E}\left\|\hat{h}_{m}-h_{m}\right\|_{\mathrm{L}^{2}}^{2} \\
& \mathbb{E}\left\|\phi_{m}-\hat{\phi}_{m}^{\mathrm{Lag}_{1}}\right\|_{\mathrm{L}^{2}}^{2} \leqslant 3 \times \frac{C\left(\theta, \theta_{0}\right)}{\left(1-\theta_{0}\right)^{2}}\|\phi\|_{\mathrm{L}^{2}}^{2} m \mathbb{E}\left\|\hat{g}_{m}-g_{m}\right\|_{\mathrm{L}^{2}}^{2}+\frac{60}{\left(1-\theta_{0}\right)^{2}} \mathbb{E}\left\|\hat{h}_{m}-h_{m}\right\|_{\mathrm{L}^{2}}^{2},
\end{aligned}
$$

with $C\left(\theta, \theta_{0}\right) \lesssim \frac{\left(1-\theta_{0}\right)^{2}}{(1-\theta)^{2}}$. To conclude, we use the upper bounds established in the proof of Theorem 2.5. If $\delta=0$, we have:

$$
\mathbb{E}\left\|\hat{g}_{m}-g_{m}\right\|_{\mathrm{L}^{2}}^{2} \leqslant \frac{\lambda}{c^{2} T} \mathbb{E}[X], \quad \mathbb{E}\left\|\hat{h}_{m}-h_{m}\right\|_{\mathrm{L}^{2}}^{2} \leqslant \frac{\lambda}{c^{2} T} \mathbb{E}[W(X)],
$$

and if $\delta>0$, we have:

$$
\begin{gathered}
\mathbb{E}\left\|\hat{g}_{m}-g_{m}\right\|_{\mathrm{L}^{2}}^{2} \leqslant \frac{C(\lambda)}{c^{2} T}\left(\mathbb{E}[X]+\frac{\mathbb{E}\left[X^{2}\right]^{\frac{1}{2}}}{(1-\theta)^{2} \delta^{2}}\right), \\
\mathbb{E}\left\|\hat{h}_{m}-h_{m}\right\|_{\mathrm{L}^{2}}^{2} \leqslant \frac{C(\lambda)}{c^{2} T}\left(\mathbb{E}[W(X)]+\frac{\mathbb{E}\left[W(X)^{2}\right]^{\frac{1}{2}}}{(1-\theta)^{2} \delta^{2}}\right),
\end{gathered}
$$

with $C(\lambda)=\mathcal{O}\left(\lambda^{2}\right)$.

### 7.2.2 Proof of Proposition 3.4

Let us introduce the sequence of functions $\left(D_{k}\right)_{k \geqslant 0}$ as:

$$
D_{k}(x):= \begin{cases}\frac{\Psi_{0}(x)}{\sqrt{2}} & \text { if } k=0 \\ \frac{\Psi_{k}(x)-\Psi_{k-1}(x)}{\sqrt{2}} & \text { if } k \geqslant 1 .\end{cases}
$$

so we can rewrite $\mathbf{A}_{m}=\mathbf{I}_{m}-\frac{\lambda}{c} \mathbf{T}_{m}(\mathbb{E}[D(X)])$ and $\hat{\mathbf{A}}_{m}=\mathbf{I}_{m}-\frac{1}{c T} \sum_{i=1}^{N_{T}} \mathbf{T}_{m}\left(D\left(X_{i}\right)\right)$, with $\mathbf{T}_{m}(\bullet)$ defined by (31). Now, the difference between $\hat{\mathbf{A}}_{m}$ and $\mathbf{A}_{m}$ can be decomposed as:

$$
\begin{equation*}
\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}=\frac{1}{c T} \sum_{i=1}^{N_{T}}\left\{\mathbf{T}_{m}\left(D\left(X_{i}\right)\right)-\mathbb{E}\left[\mathbf{T}_{m}\left(D\left(X_{i}\right)\right)\right]\right\}+\frac{N_{T}}{c T} \mathbf{T}_{m}(\mathbb{E}[D(X)])-\frac{\lambda}{c} \mathbf{T}_{m}(\mathbb{E}[D(X)]) \tag{32}
\end{equation*}
$$

The next lemma gives a control on the first term in the decomposition 32).

Lemma 7.9. Let $\mathbf{S}_{n}:=\sum_{i=1}^{n} \mathbf{Z}_{i}$, with $\mathbf{Z}_{i}:=\mathbf{T}_{m}\left(D\left(X_{i}\right)\right)-\mathbb{E}\left[\mathbf{T}_{m}\left(D\left(X_{i}\right)\right)\right]$. Then for $p \geqslant 1$ and $\log m \geqslant p$, we have:

$$
\mathbb{E}\left\|\mathbf{S}_{n}\right\|_{\mathrm{op}}^{2 p} \leqslant C(p)\left[(n \mu m \log m)^{p}+(m \log m)^{2 p}\right]
$$

with $C(p)$ a constant depending on $p$.
Proof. We want to apply Theorem C. 2 First, we need upper bounds on $\left\|\mathbf{Z}_{i}\right\|_{\text {op }}$ and $\lambda_{\max }\left(\mathbb{E}\left[\mathbf{S}_{n}^{\top} \mathbf{S}_{n}\right]\right)$.

- Bound on $\left\|\mathbf{Z}_{i}\right\|_{\mathrm{op}}$ :

$$
\begin{aligned}
\left\|\mathbf{Z}_{i}\right\|_{\mathrm{op}} & =\sup _{\|x\|_{\ell^{2}} \leqslant 1}\left\|\left(\mathbf{T}_{m}\left(D\left(X_{i}\right)\right)-\mathbb{E}\left[\mathbf{T}_{m}\left(D\left(X_{i}\right)\right)\right]\right) x\right\|_{\ell^{2}} \\
& =\sup _{\|x\|_{\ell^{2}} \leqslant 1}\left\|\left(D\left(X_{i}\right)-\mathbb{E}\left[D\left(X_{i}\right)\right]\right) * x\right\|_{\ell^{2}} \\
& \leqslant\left\|D\left(X_{i}\right)-\mathbb{E}\left[D\left(X_{i}\right)\right]\right\|_{\ell^{1}} \\
& \leqslant \sqrt{2} \sum_{k=0}^{m-1}\left|\Psi_{k}\left(X_{i}\right)-\mathbb{E}\left[\Psi_{k}\left(X_{i}\right)\right]\right| \\
& \leqslant 2 \sqrt{2} \sum_{k=0}^{m-1}\left\|\Psi_{k}\right\|_{\infty}
\end{aligned}
$$

By Lemma 3.3 there exists an absolute constant $C>0$ such that $\left\|\Psi_{k}\right\|_{\infty} \leqslant C$, hence $\left\|\mathbf{Z}_{i}\right\|_{\mathrm{op}} \leqslant C 2 \sqrt{2} m$.

- Bound on $\lambda_{\max }\left(\mathbb{E}\left[\mathbf{S}_{n}^{\top} \mathbf{S}_{n}\right]\right)$ :

$$
\begin{aligned}
\lambda_{\max }\left(\mathbb{E}\left[\mathbf{S}_{n}^{\top} \mathbf{S}_{n}\right]\right) & =\sup _{\|x\|_{\ell^{2}=1}} x^{\top} \mathbb{E}\left[\mathbf{S}_{n}^{\top} \mathbf{S}_{n}\right] x \\
& =n \sup _{\|x\|_{\ell^{2}}=1} x^{\top} \mathbb{E}\left[\mathbf{Z}_{1}^{\top} \mathbf{Z}_{1}\right] x \\
& =n \sup _{\|x\|_{\ell^{2}}=1} \mathbb{E}\left[\left\|\mathbf{Z}_{1} x\right\|_{\ell^{2}}^{2}\right] \\
& =n \sup _{\|x\|_{\ell^{2}=1}} \mathbb{E}\left[\left\|\left(D\left(X_{1}\right)-\mathbb{E}\left[D\left(X_{1}\right)\right]\right) * x\right\|_{\ell^{2}}^{2}\right] .
\end{aligned}
$$

If $x \in \mathbb{R}^{m}$, we have:

$$
\begin{aligned}
\mathbb{E}\left[\left\|\left(D\left(X_{1}\right)-\mathbb{E}\left[D\left(X_{1}\right)\right]\right) * x\right\|_{\ell^{2}}^{2}\right] & =\sum_{j=0}^{m-1} \mathbb{E}\left[\left(\left\{D\left(X_{1}\right)-\mathbb{E}\left[D\left(X_{1}\right)\right]\right\} * x\right)_{j}^{2}\right] \\
& =\sum_{j=0}^{m-1} \operatorname{Var}\left[\left(D\left(X_{1}\right) * x\right)_{j}\right] \leqslant \sum_{j=0}^{m-1} \mathbb{E}\left[\left(D\left(X_{1}\right) * x\right)_{j}^{2}\right]
\end{aligned}
$$

and by Cauchy-Schwarz inequality:

$$
\begin{aligned}
\left(D\left(X_{1}\right) * x\right)_{j}^{2} & \leqslant\left(\sum_{k=0}^{j} D_{k}\left(X_{1}\right)^{2}\right)\left(\sum_{k=0}^{j} x_{k}^{2}\right) \\
& \leqslant\|x\|_{\ell^{2}}^{2} \sum_{k=0}^{j} \Psi_{k}\left(X_{1}\right)^{2} \leqslant\|x\|_{\ell^{2}}^{2}\left\|\mathbf{1}_{X_{1}>\bullet}\right\|_{\mathrm{L}^{2}}^{2}=\|x\|_{\ell^{2}}^{2} X_{1}
\end{aligned}
$$

because $\Psi_{k}\left(X_{1}\right)=\left\langle\mathbf{1}_{X_{1}>\bullet}, \psi_{k}\right\rangle$ and $\left(\psi_{k}\right)$ is an orthonormal basis of $\mathrm{L}^{2}\left(\mathbb{R}_{+}\right)$. Hence, we obtain $\lambda_{\text {max }}\left(\mathbb{E}\left[\mathbf{S}^{\top} \mathbf{S}\right]\right) \leqslant n m \mu$.
We want apply Theorem C. 2 to our matrix $\mathbf{S}_{n}$, which is not Hermitian. We use the following trick, called the Paulsen dialtation. For $\mathbf{M}$ a rectangular matrix, we define:

$$
\mathbf{M} \mapsto \mathcal{H}(\mathbf{M})=\left(\begin{array}{cc}
0 & \mathbf{M} \\
\mathbf{M}^{\dagger} & 0
\end{array}\right)
$$

where $\mathbf{M}^{\dagger}$ denotes the conjugate transpose of $\mathbf{M}$. Now, $\mathcal{H}(\mathbf{M})$ is an Hermitian matrix, and:

$$
\mathcal{H}(\mathbf{M})^{2}=\left(\begin{array}{cc}
\mathbf{M M}^{\dagger} & 0 \\
0 & \mathbf{M}^{\dagger} \mathbf{M}
\end{array}\right)
$$

hence $\lambda_{\max }\left(\mathcal{H}(\mathbf{M})^{2}\right)=\|\mathbf{M}\|_{\text {op }}^{2}$ and $\lambda_{\max }(\mathcal{H}(\mathbf{M}))=\|\mathbf{M}\|_{\text {op }}$. We can now apply Theorem C. 2 for $\mathbf{M}=\mathbf{S}_{n}$, we have:

$$
\mathcal{H}\left(\mathbf{S}_{n}\right)=\left(\begin{array}{cc}
0 & \sum_{i} \mathbf{Z}_{i} \\
\sum_{i} \mathbf{Z}_{i}^{\top} & 0
\end{array}\right)=\sum_{i}\left(\begin{array}{cc}
0 & \mathbf{Z}_{i} \\
\mathbf{Z}_{i}^{\top} & 0
\end{array}\right)=\sum_{i} \mathcal{H}\left(\mathbf{Z}_{i}\right)
$$

thus for $p \geqslant 1$ and $r \geqslant \max (2 p, 2 \log m)$, we get that:

$$
\begin{aligned}
\left(\mathbb{E}\left\|\mathbf{S}_{n}\right\|_{\text {op }}^{2 p}\right)^{1 / 2 p} & =\left(\mathbb{E} \lambda_{\max }\left(\mathcal{H}\left(\sum_{i} \mathbf{Z}_{i}\right)\right)^{2 p}\right)^{1 / 2 p} \\
& \leqslant \sqrt{\mathrm{e} r} \lambda_{\max }^{1 / 2}\left(\sum_{i} \mathbb{E} \mathcal{H}\left(\mathbf{Z}_{i}\right)^{2}\right)+2 \mathrm{e} r\left(\mathbb{E} \max _{i} \lambda_{\max }\left(\mathcal{H}\left(\mathbf{Z}_{i}\right)\right)^{2 p}\right)^{1 / 2 p} \\
& \leqslant \sqrt{\mathrm{e} r \lambda_{\max }\left(\mathbb{E} \mathbf{S}_{n}^{\top} \mathbf{S}_{n}\right)}+2 \mathrm{e} r\left(\mathbb{E} \max _{i}\left\|\mathbf{Z}_{i}\right\|_{\text {op }}^{2 p}\right)^{1 / 2 p} \\
& \leqslant \sqrt{\mathrm{ernm} \mathrm{\mu}}+C 4 \sqrt{2} \mathrm{erm}
\end{aligned}
$$

If $\log m \geqslant p$, then $r=2 \log m$ and we get $\mathbb{E}\left\|\mathbf{S}_{n}\right\|_{\text {op }}^{2 p} \lesssim 2^{2 p-1}(n \mu m \log m)^{p}+2^{6 p-1}(m \log m)^{2 p}$.
Now we can prove Proposition 3.4
Proof of Proposition 3.4 From the decomposition (32), we get:

$$
\mathbb{E}\left\|\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right\|_{\mathrm{op}}^{2 p} \leqslant 2^{2 p-1} \frac{1}{(c T)^{2 p}} \mathbb{E}\left\|\mathbf{S}_{N_{T}}\right\|_{\mathrm{op}}^{2 p}+2^{2 p-1} \frac{\mathbb{E}\left|N_{T}-\lambda T\right|^{2 p}}{(c T)^{2 p}}\left\|\mathbf{T}_{m}(\mathbb{E}[D(X)])\right\|_{\mathrm{op}}^{2 p}
$$

For the first term, we apply Lemma 7.9 conditional on $N_{T}$ :

$$
\begin{aligned}
\frac{1}{(c T)^{2 p}} \mathbb{E}\left\|\mathbf{S}_{N_{T}}\right\|_{\mathrm{op}}^{2 p} & \leqslant C(p)\left[\frac{\mathbb{E}\left[N_{T}^{p}\right] \mu^{p}(m \log m)^{p}}{(c T)^{2 p}}+\left(\frac{m \log m}{c T}\right)^{2 p}\right] \\
& =C(p)\left[\mu^{p} \mathbb{E}\left[\left(\frac{N_{T}}{c T}\right)^{p}\right]\left(\frac{m \log m}{c T}\right)^{p}+\left(\frac{m \log m}{c T}\right)^{2 p}\right]
\end{aligned}
$$

with $\mathbb{E}\left[\left(\frac{N_{T}}{c T}\right)^{p}\right]=\mathcal{O}\left(\lambda^{p}\right)$. For the second term, we know from Corollary 7.3 that $\mathbb{E}\left[\left(N_{T}-\lambda T\right)^{2 p}\right]=$ $\mathcal{O}\left(\lambda^{p} T^{p}\right)$, and:

$$
\begin{aligned}
\left\|\mathbf{T}_{m}(\mathbb{E}[D(X)])\right\|_{\mathrm{op}} \leqslant \sum_{k=0}^{m-1}\left|\mathbb{E}\left[D_{k}(X)\right]\right| & \leqslant \sqrt{2} \sum_{k=0}^{m-1}\left|\mathbb{E}\left[\Psi_{k}(X)\right]\right| \\
& \leqslant \sqrt{2 m}\left(\sum_{k=0}^{m-1} \mathbb{E}\left[\Psi_{k}(X)^{2}\right]\right)^{1 / 2} \\
& =\sqrt{2 m}\left(\mathbb{E}\left[\sum_{k=0}^{m-1}\left\langle\mathbf{1}_{X>\bullet}, \psi_{k}\right\rangle^{2}\right]\right)^{1 / 2} \leqslant \sqrt{2 m \mu}
\end{aligned}
$$

thus:

$$
\frac{\mathbb{E}\left|N_{T}-\lambda T\right|^{2 p}}{(c T)^{2 p}}\left\|\mathbf{T}_{m}(\mathbb{E}[D(X)])\right\|_{\mathrm{op}}^{2 p} \leqslant \frac{\mathcal{O}\left(\lambda^{p}\right)}{T^{p}} \mu^{p} m^{p} .
$$

### 7.2.3 Proof of Theorem 3.5

The following results are based on the proofs of Lemma 3.1 and Corollary 3.2 in Comte and Mabon 2017.

Proposition 7.10. If $m \log m \leqslant c T$, then it holds:

$$
\mathbb{E}\left\|\widetilde{\mathbf{A}}_{m, 2}^{-1}-\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p} \leqslant C(p, \lambda)\left(\mu^{p}\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p}\right) \wedge\left(\left(\mu^{p}+\mu^{2 p}\right)\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{4 p}\left(\frac{m \log m}{c T}\right)^{p}\right)
$$

with $C(p, \lambda)=\mathcal{O}\left(\lambda^{p} \vee \lambda^{2 p}\right)$.
Proof. We decompose the expectation according to the event $\Delta_{m}^{2}:=\left\{\left\|\hat{\mathbf{A}}_{m}^{-1}\right\|_{\mathrm{op}}^{1} \leqslant \frac{c T}{m \log m}\right\}$ :

$$
\begin{align*}
\mathbb{E}\left[\left\|\mathbf{A}_{m}^{-1}-\widetilde{\mathbf{A}}_{m, 2}^{-1}\right\|_{\mathrm{op}}^{2 p}\right] & =\mathbb{E}\left[\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p} \mathbf{1}_{\Delta_{m}^{c}}+\left\|\hat{\mathbf{A}}_{m}^{-1}\left(\mathbf{A}_{m}-\hat{\mathbf{A}}_{m}\right) \mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p} \mathbf{1}_{\Delta_{m}}\right] \\
& =\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p} \mathbb{P}\left[\Delta_{m}^{c}\right]+\mathbb{E}\left[\left\|\hat{\mathbf{A}}_{m}^{-1}\left(\mathbf{A}_{m}-\hat{\mathbf{A}}_{m}\right) \mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p} \mathbf{1}_{\Delta_{m}}\right] \tag{33}
\end{align*}
$$

We now give two bounds on 33, depending on the value of $\left\|\mathbf{A}_{m}^{-1}\right\|_{\text {op }}$.

- First case: $\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}>\frac{1}{2} \sqrt{\frac{c T}{m \log m}}$.

Starting from Equation (33) and using the set $\Delta_{m}^{2}$, we have that:

$$
\begin{aligned}
\mathbb{E}\left[\left\|\mathbf{A}_{m}^{-1}-\widetilde{\mathbf{A}}_{m, 2}^{-1}\right\|_{\mathrm{op}}^{2 p}\right] & \leqslant\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p}+\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p} \mathbb{E}\left[\left\|\hat{\mathbf{A}}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p}\left\|\mathbf{A}_{m}-\hat{\mathbf{A}}_{m}\right\|_{\mathrm{op}}^{2 p} \mathbf{1}_{\Delta_{m}}\right] \\
& \leqslant\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p}+\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p}\left(\frac{c T}{m \log m}\right)^{p} \mathbb{E}\left[\left\|\mathbf{A}_{m}-\hat{\mathbf{A}}_{m}\right\|_{\mathrm{op}}^{2 p}\right]
\end{aligned}
$$

We apply Proposition 3.4 and get:

$$
\begin{aligned}
\mathbb{E}\left[\left\|\mathbf{A}_{m}^{-1}-\widetilde{\mathbf{A}}_{m, 2}^{-1}\right\|_{\mathrm{op}}^{2 p}\right] & \leqslant\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p}+\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p}\left(\frac{c T}{m \log m}\right)^{p} C(p, \lambda) \mu^{p}\left(\frac{m \log m}{c T}\right)^{p} \\
& \leqslant\left(1+C(p, \lambda) \mu^{p}\right)\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p}
\end{aligned}
$$

with $C(p, \lambda)=\mathcal{O}\left(\lambda^{p}\right)$.

- Second case: $\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}<\frac{1}{2} \sqrt{\frac{c T}{m \log m}}$.

Starting from (33) again, we get:

$$
\mathbb{E}\left[\left\|\mathbf{A}_{m}^{-1}-\widetilde{\mathbf{A}}_{m, 2}^{-1}\right\|_{\mathrm{op}}^{2 p}\right] \leqslant\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p} \mathbb{P}\left[\Delta_{m}^{c}\right]+\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p} \mathbb{E}\left[\left\|\mathbf{A}_{m}-\hat{\mathbf{A}}_{m}\right\|_{\mathrm{op}}^{2 p}\left\|\hat{\mathbf{A}}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p} \mathbf{1}_{\Delta_{m}}\right] .
$$

1. Upper bound on $\mathbb{E}\left[\left\|\mathbf{A}_{m}-\hat{\mathbf{A}}_{m}\right\|_{\mathrm{op}}^{2 p}\left\|\hat{\mathbf{A}}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p} \mathbf{1}_{\Delta_{m}}\right]$.

First let us notice that

$$
\left\|\hat{\mathbf{A}}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p} \leqslant 2^{2 p-1}\left\|\hat{\mathbf{A}}_{m}^{-1}-\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p}+2^{2 p-1}\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p}
$$

Applying Proposition 3.4, we get:

$$
\begin{align*}
& \mathbb{E}\left[\left\|\mathbf{A}_{m}-\hat{\mathbf{A}}_{m}\right\|_{\mathrm{op}}^{2 p}\left\|\hat{\mathbf{A}}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p} \mathbf{1}_{\Delta_{m}}\right] \\
& \leqslant 2^{2 p-1}\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p} \mathbb{E}\left[\left\|\mathbf{A}_{m}-\hat{\mathbf{A}}_{m}\right\|_{\mathrm{op}}^{2 p} \mathbf{1}_{\Delta_{m}}\right]+2^{2 p-1} \mathbb{E}\left[\left\|\mathbf{A}_{m}-\hat{\mathbf{A}}_{m}\right\|_{\mathrm{op}}^{2 p}\left\|\hat{\mathbf{A}}_{m}^{-1}-\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p} \mathbf{1}_{\Delta_{m}}\right] \\
& \leqslant 2^{2 p-1}\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p} \mathbb{E}\left[\left\|\mathbf{A}_{m}-\hat{\mathbf{A}}_{m}\right\|_{\mathrm{op}}^{2 p} \mathbf{1}_{\Delta_{m}}\right]+2^{2 p-1}\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p} \mathbb{E}\left[\left\|\mathbf{A}_{m}-\hat{\mathbf{A}}_{m}\right\|_{\mathrm{op}}^{4 p}\left\|\hat{\mathbf{A}}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p} \mathbf{1}_{\Delta_{m}}\right] \\
& \leqslant C(p, \lambda) \mu^{p}\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p}\left(\frac{m \log m}{c T}\right)^{p}+\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p}\left(\frac{c T}{m \log m}\right)^{p} C(2 p, \lambda) \mu^{2 p}\left(\frac{m \log m}{c T}\right)^{2 p} \\
& \leqslant C^{\prime}(p, \lambda)\left(\mu^{p}+\mu^{2 p}\right)\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p}\left(\frac{m \log m}{c T}\right)^{p}, \tag{34}
\end{align*}
$$

with $C^{\prime}(p, \lambda)=\mathcal{O}\left(\lambda^{p} \vee \lambda^{2 p}\right)$.
2. Upper bound on $\mathbb{P}\left[\Delta_{m}^{\mathrm{c}}\right]=\mathbb{P}\left[\left\|\hat{\mathbf{A}}_{m}^{-1}\right\|_{\mathrm{op}}>\sqrt{\frac{c T}{m \log m}}\right]$.

From the triangular inequality:

$$
\left\|\hat{\mathbf{A}}_{m}^{-1}\right\|_{\mathrm{op}} \leqslant\left\|\hat{\mathbf{A}}_{m}^{-1}-\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}+\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}
$$

we obtain:

$$
\mathbb{P}\left[\left\|\hat{\mathbf{A}}_{m}^{-1}\right\|_{\mathrm{op}}>\sqrt{\frac{c T}{m \log m}}\right] \leqslant \mathbb{P}\left[\left\|\hat{\mathbf{A}}_{m}^{-1}-\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}>\sqrt{\frac{c T}{m \log m}}-\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}\right]
$$

Moreover we have assumed that $\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}<\frac{1}{2} \sqrt{\frac{c T}{m \log m}}$, so:

$$
\mathbb{P}\left[\left\|\hat{\mathbf{A}}_{m}^{-1}\right\|_{\mathrm{op}}>\sqrt{\frac{c T}{m \log m}}\right] \leqslant \mathbb{P}\left[\left\|\hat{\mathbf{A}}_{m}^{-1}-\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}>\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}\right]
$$

Now let us rewrite this probability, as:

$$
\begin{align*}
\mathbb{P} & {\left[\left\|\hat{\mathbf{A}}_{m}^{-1}-\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}>\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}\right] } \\
= & \mathbb{P}\left[\left\{\left\|\hat{\mathbf{A}}_{m}^{-1}-\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}>\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}\right\} \cap\left\{\left\|\mathbf{A}_{m}^{-1}\left(\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right)\right\|_{\mathrm{op}}<\frac{1}{2}\right\}\right] \\
& +\mathbb{P}\left[\left\{\left\|\hat{\mathbf{A}}_{m}^{-1}-\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}>\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}\right\} \cap\left\{\left\|\mathbf{A}_{m}^{-1}\left(\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right)\right\|_{\mathrm{op}} \geqslant \frac{1}{2}\right\}\right] \\
\leqslant & \mathbb{P}\left[\left\{\left\|\hat{\mathbf{A}}_{m}^{-1}-\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}>\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}\right\} \cap\left\{\left\|\mathbf{A}_{m}^{-1}\left(\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right)\right\|_{\mathrm{op}}<\frac{1}{2}\right\}\right] \\
& +\mathbb{P}\left[\left\|\mathbf{A}_{m}^{-1}\left(\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right)\right\|_{\mathrm{op}} \geqslant \frac{1}{2}\right] . \tag{35}
\end{align*}
$$

To control the second term, we apply Markov inequality and Proposition 3.4

$$
\begin{align*}
\mathbb{P}\left[\left\|\mathbf{A}_{m}^{-1}\left(\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right)\right\|_{\mathrm{op}} \geqslant \frac{1}{2}\right] & \leqslant \mathbb{P}\left[\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}\left\|\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right\|_{\mathrm{op}} \geqslant \frac{1}{2}\right] \\
& \leqslant C(p, \lambda) \mu^{p}\left(\frac{m \log m}{c T}\right)^{p}\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p} \tag{36}
\end{align*}
$$

Next to control the first term on the right hand side of Equation (35), we apply Theorem C.1:

$$
\begin{align*}
& \mathbb{P}\left[\left\{\left\|\hat{\mathbf{A}}_{m}^{-1}-\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}>\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}\right\} \cap\left\{\left\|\mathbf{A}_{m}^{-1}\left(\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right)\right\|_{\mathrm{op}}<\frac{1}{2}\right\}\right] \\
& \quad \leqslant \mathbb{P}\left[\left\{\frac{\left\|\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right\|_{\mathrm{op}}\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2}}{1-\left\|\mathbf{A}_{m}^{-1}\left(\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right)\right\|_{\mathrm{op}}}>\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}\right\} \cap\left\{\left\|\mathbf{A}_{m}^{-1}\left(\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right)\right\|_{\mathrm{op}}<\frac{1}{2}\right\}\right] \\
& \quad \leqslant \mathbb{P}\left[\left\|\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right\|_{\mathrm{op}}>\frac{1}{2}\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{-1}\right] . \tag{37}
\end{align*}
$$

We apply Markov inequality again, along with Proposition 3.4.

$$
\begin{aligned}
& \mathbb{P}\left[\left\{\left\|\hat{\mathbf{A}}_{m}^{-1}-\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}>\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}\right\} \cap\left\{\left\|\mathbf{A}_{m}^{-1}\left(\hat{\mathbf{A}}_{m}-\mathbf{A}_{m}\right)\right\|_{\mathrm{op}}<\frac{1}{2}\right\}\right] \\
& \\
& \leqslant C(p, \lambda) \mu^{p}\left(\frac{m \log m}{c T}\right)^{p}\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p} .
\end{aligned}
$$

So starting from Equation (35) and gathering Equations (36) and (37) gives:

$$
\begin{equation*}
\mathbb{P}\left[\left\|\hat{\mathbf{A}}_{m}^{-1}\right\|_{\mathrm{op}}>\sqrt{\frac{c T}{m \log m}}\right] \leqslant C(p, \lambda) \mu^{p}\left(\frac{m \log m}{c T}\right)^{p}\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2 p} \tag{38}
\end{equation*}
$$

with $C(p, \lambda)=\mathcal{O}\left(\lambda^{p}\right)$.

Finally gathering Equations (34) with (38), we get that

$$
\mathbb{E}\left[\left\|\mathbf{A}_{m}^{-1}-\widetilde{\mathbf{A}}_{m, 2}^{-1}\right\|_{\mathrm{op}}^{2 p}\right] \leqslant C(p, \lambda)\left(\mu^{p}+\mu^{2 p}\right)\left(\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{4} \frac{m \log m}{c T}\right)^{p} .
$$

with $C(p, \lambda)=\mathcal{O}\left(\lambda^{p} \vee \lambda^{2 p}\right)$.
The next proposition is a variant of the last one. It gives a better bound than applying directly Proposition 7.10 to $\mathbb{E}\left[\left\|\widetilde{\mathbf{A}}_{m, 2}^{-1}-\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2}\left\|\mathbf{c}_{m}\right\|_{\ell^{2}}^{2}\right]$.
Proposition 7.11. If $m \log m \leqslant c T$, then it holds:

$$
\mathbb{E}\left\|\left(\widetilde{\mathbf{A}}_{m, 2}^{-1}-\mathbf{A}_{m}^{-1}\right) \mathbf{c}_{m}\right\|_{\ell^{2}}^{2} \leqslant C(\lambda)\left\|\phi_{m}\right\|_{\mathrm{L}^{2}}^{2}\left(\mu \wedge\left\{\left(\mu+\mu^{2}\right)\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2} \frac{m \log m}{c T}\right\}\right)
$$

with $C(\lambda)=\mathcal{O}\left(\lambda \vee \lambda^{2}\right)$.
Proof. The proof follows the lines of the proof of of Proposition 7.10 but starting from the following decomposition:

$$
\begin{aligned}
\mathbb{E}\left\|\left(\mathbf{A}_{m}^{-1}-\widetilde{\mathbf{A}}_{m, 2}^{-1}\right) \mathbf{c}_{m}\right\|_{\ell^{2}}^{2} & =\left\|\mathbf{A}_{m}^{-1} \mathbf{c}_{m}\right\|_{\ell^{2}}^{2} \mathbb{P}\left[\Delta_{m}^{c}\right]+\mathbb{E}\left[\left\|\widehat{\mathbf{A}}_{m}^{-1}\left(\mathbf{A}_{m}-\widehat{\mathbf{A}}_{m}\right) \mathbf{A}_{m}^{-1} \mathbf{c}_{m}\right\|_{\ell^{2}}^{2} \mathbf{1}_{\Delta_{m}}\right] \\
& =\left\|\mathbf{a}_{m}\right\|_{\ell^{2}}^{2} \mathbb{P}\left[\Delta_{m}^{c}\right]+\mathbb{E}\left[\left\|\widehat{\mathbf{A}}_{m}^{-1}\left(\mathbf{A}_{m}-\widehat{\mathbf{A}}_{m}\right) \mathbf{A}_{m}^{-1} \mathbf{c}_{m}\right\|_{\ell^{2}}^{2} \mathbf{1}_{\Delta_{m}}\right],
\end{aligned}
$$

It yields the following upper bound

$$
\mathbb{E}\left[\left\|\left(\mathbf{A}_{m}^{-1}-\widetilde{\mathbf{A}}_{m, 2}^{-1}\right) \mathbf{c}_{m}\right\|_{\ell^{2}}^{2}\right] \leqslant\left\|\mathbf{a}_{m}\right\|_{\ell^{2}}^{2} \mathbb{P}\left[\Delta_{m}^{c}\right]+\left\|\mathbf{a}_{m}\right\|_{\ell^{2}}^{2} \mathbb{E}\left[\left\|\widehat{\mathbf{A}}_{m}^{-1}\right\|_{\mathrm{op}}^{2}\left\|\mathbf{A}_{m}-\widehat{\mathbf{A}}_{m}\right\|_{\mathrm{op}}^{2} \mathbf{1}_{\Delta_{m}}\right]
$$

Following the proof of Proposition 7.10, we get:

$$
\mathbb{E}\left[\left\|\left(\mathbf{A}_{m}^{-1}-\widetilde{\mathbf{A}}_{m, 2}^{-1}\right) \mathbf{c}_{m}\right\|_{\ell^{2}}^{2}\right] \leqslant C(\lambda)\left\|\mathbf{a}_{m}\right\|_{\mathrm{L}^{2}}^{2}\left(\mu \wedge\left\{\left(\mu+\mu^{2}\right)\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2} \frac{m \log m}{c T}\right\}\right)
$$

with $C(\lambda)=\mathcal{O}\left(\lambda \vee \lambda^{2}\right)$.
Now we can prove Theorem 3.5.
Proof of Theorem 3.5. By Pythagoras Theorem, $\left\|\phi-\hat{\phi}_{m}^{\mathrm{Lag}_{2}}\right\|_{\mathrm{L}^{2}}^{2}=\left\|\phi-\phi_{m}\right\|_{\mathrm{L}^{2}}^{2}+\left\|\phi_{m}-\hat{\phi}_{m}^{\mathrm{Lag}_{2}}\right\|_{\mathrm{L}^{2}}^{2}$. In the proof of Theorem 2.5. we saw that:

$$
\mathbb{E}\left\|\hat{\mathbf{c}}_{m}-\mathbf{c}_{m}\right\|_{\ell^{2}}^{2}=\mathbb{E}\left\|\hat{h}_{m}-h_{m}\right\|_{\mathrm{L}^{2}}^{2} \leqslant \frac{\lambda}{c^{2} T} \mathbb{E}[W(X)]
$$

We decompose the variance term in three terms:

$$
\begin{aligned}
& \mathbb{E}\left\|\phi_{m}-\hat{\phi}_{m}^{\mathrm{Lag}_{2}}\right\|_{\mathrm{L}^{2}}^{2}=\mathbb{E}\left\|\hat{\mathbf{a}}_{m}^{\mathrm{Lag}_{2}}-\mathbf{a}_{m}\right\|_{\ell^{2}}^{2}=\mathbb{E}\left\|\widetilde{\mathbf{A}}_{m, 2}^{-1} \hat{\mathbf{c}}_{m}-\mathbf{A}_{m}^{-1} \mathbf{c}_{m}\right\|_{\ell^{2}}^{2} \\
& \quad \leqslant 3 \mathbb{E}\left\|\left(\widetilde{\mathbf{A}}_{m, 2}^{-1}-\mathbf{A}_{m}^{-1}\right) \mathbf{c} \mathbf{c}_{m}\right\|_{\ell^{2}}^{2}+3 \mathbb{E}\left\|\left(\mathbf{A}_{m}^{-1}-\widetilde{\mathbf{A}}_{m, 2}^{-1}\right)\left(\mathbf{c}_{m}-\hat{\mathbf{c}}_{m}\right)\right\|_{\ell^{2}}^{2}+3 \mathbb{E}\left\|\mathbf{A}_{m}^{-1}\left(\mathbf{c}_{m}-\hat{\mathbf{c}}_{m}\right)\right\|_{\ell^{2}}^{2}
\end{aligned}
$$

For the first term, we apply Proposition 7.10 For the second term, we use the fact that $\widetilde{\mathbf{A}}_{m, 2}^{-1}$ and $\hat{\mathbf{c}}_{m}$ are independent, and we apply Proposition 7.11.

$$
\mathbb{E}\left\|\left(\mathbf{A}_{m}^{-1}-\widetilde{\mathbf{A}}_{m, 2}^{-1}\right)\left(\mathbf{c}_{m}-\hat{\mathbf{c}}_{m}\right)\right\|_{\ell^{2}}^{2} \leqslant \mathbb{E}\left\|\mathbf{A}_{m}^{-1}-\widetilde{\mathbf{A}}_{m, 2}^{-1}\right\|_{\mathrm{op}}^{2} \times \mathbb{E}\left\|\mathbf{c}_{m}-\hat{\mathbf{c}}_{m}\right\|_{\ell^{2}}^{2}=\mathcal{O}\left(\frac{1}{T^{2}}\right) .
$$

For the third term:

$$
\mathbb{E}\left\|\mathbf{A}_{m}^{-1}\left(\mathbf{c}_{m}-\hat{\mathbf{c}}_{m}\right)\right\|_{\ell^{2}}^{2} \leqslant\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2} \mathbb{E}\left\|\hat{\mathbf{c}}_{m}-\mathbf{c}_{m}\right\|_{\ell^{2}}^{2} \leqslant\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2} \frac{\lambda}{c^{2} T} \mathbb{E}[W(X)]
$$

We apply Lemma 3.1 and we obtain the following bound, with $C(\lambda)=\mathcal{O}\left(\lambda \vee \lambda^{2}\right)$ :

$$
\begin{aligned}
\mathbb{E}\left\|\phi_{m}-\hat{\phi}_{m}^{\mathrm{Lag}_{2}}\right\|_{\mathrm{L}^{2}}^{2} & \leqslant 3\left\|\mathbf{A}_{m}^{-1}\right\|_{\mathrm{op}}^{2} \frac{C(\lambda)}{c T}\left(\left\|\phi_{m}\right\|_{\mathrm{L}^{2}}^{2}\left(\mu+\mu^{2}\right) m \log (m)+\frac{\mathbb{E}[W(X)]}{c}\right)+\mathcal{O}\left(\frac{1}{T^{2}}\right) \\
& \leqslant 12 \frac{C(\lambda)}{c T(1-\theta)^{2}}\left(\left\|\phi_{m}\right\|_{\mathrm{L}^{2}}^{2}\left(\mu+\mu^{2}\right) m \log (m)+\frac{\mathbb{E}[W(X)]}{c}\right)+\mathcal{O}\left(\frac{1}{T^{2}}\right)
\end{aligned}
$$

### 7.3 Proofs of Section 4

Proof of Lemma 4.6. It follows from direct calculation using (6):

$$
\begin{aligned}
\left\langle F, \psi_{k}\right\rangle & =\int_{0}^{+\infty} C \exp (-\gamma x) \psi_{k}(x) \mathrm{d} x \\
& =C \sqrt{2} \sum_{j=0}^{k}\binom{k}{j} \frac{(-2)^{j}}{j!} \int_{0}^{+\infty} x^{j} \mathrm{e}^{-(1+\gamma) x} \mathrm{~d} x \\
& =C \sqrt{2} \sum_{j=0}^{k}\binom{k}{j} \frac{(-2)^{j}}{(1+\gamma)^{j+1}} \\
& =\frac{C \sqrt{2}}{\gamma+1}\left(1-\frac{2}{\gamma+1}\right)^{k}=\frac{C \sqrt{2}}{\gamma+1}\left(\frac{\gamma-1}{\gamma+1}\right)^{k}
\end{aligned}
$$

Since $\gamma$ is positive, we have $\left|\frac{\gamma-1}{\gamma+1}\right|<1$ and we can compute the geometric series:

$$
\sum_{k=m}^{+\infty}\left\langle F, \psi_{k}\right\rangle^{2}=\frac{2 C}{(\gamma+1)^{2}} \frac{\left(\frac{\gamma-1}{\gamma+1}\right)^{2 m}}{1-\left(\frac{\gamma-1}{\gamma+1}\right)^{2}}=\frac{C^{2}}{2 \gamma}\left(\frac{\gamma-1}{\gamma+1}\right)^{2 m}
$$

Proof of Proposition 4.7. For the ruin probability and the Laplace transform of the ruin time, we start from (17) and we apply Lemma 4.6 For the expected jump size causing the ruin, we also start from (17) and we write $\phi=F_{1}+F_{2}$ with:

$$
F_{1}(u):=\mu(1+2 \theta) \mathrm{e}^{-\frac{1-\theta}{\mu} u} \mathbf{1}_{u>0}, \quad F_{2}(u):=\mu \mathrm{e}^{-u / \mu} \mathbf{1}_{u>0}
$$

Hence, $\left\|\phi-\phi_{m}\right\|_{\mathrm{L}^{2}}^{2} \leqslant 2 \sum_{k=m}^{+\infty}\left\langle F_{1}, \psi_{k}\right\rangle^{2}+2 \sum_{k=m}^{+\infty}\left\langle F_{2}, \psi_{k}\right\rangle^{2}$. We apply Lemma 4.6.

$$
\begin{aligned}
\left\|\phi-\phi_{m_{1}}\right\|_{L^{2}}^{2} & \leqslant \frac{\mu^{2}(1+2 \theta)^{2}}{\frac{1-\theta}{\mu}}\left(\frac{\frac{1-\theta}{\mu}-1}{\frac{1-\theta}{\mu}+1}\right)^{2 m}+\mu^{3}\left(\frac{\frac{1}{\mu}-1}{\frac{1}{\mu}+1}\right)^{2 m} \\
& =\frac{\mu^{3}(1+2 \theta)^{2}}{1-\theta}\left(\frac{1-\theta-\mu}{1-\theta+\mu}\right)^{2 m}+\mu^{3}\left(\frac{1-\mu}{1+\mu}\right)^{2 m} \\
& \leqslant \frac{\mu^{3}(1+2 \theta)^{2}}{1-\theta}\left(\left|\frac{1-\theta-\mu}{1-\theta+\mu}\right| \vee\left|\frac{1-\mu}{1+\mu}\right|\right)^{2 m}
\end{aligned}
$$

## Acknowledgments

I want to thank Fabienne Comte and Céline Duval for their helpful advice and their support of my work. I also want to thank fedja for their help with Lemma 3.3 regarding the primitives of the Laguerre functions, and Emmanuel Rio for his hints for Lemma 7.1

## Funding

This work was supported by a grant from Région Île-de-France.

## References

Milton Abramowitz and Irene A. Stegun. Handbook of Mathematical Functions: with Formulas, Graphs, and Mathematical Tables. Number 55 in Applied Mathematics Series. National Bureau of Standards, New York, NY, 10th edition, 1972.

Søren Asmussen and Hansjörg Albrecher. Ruin probabilities, volume 14 of Advanced series on statistical science and applied probability. World Scientific, Singapore ; New Jersey, 2nd edition, 2010.

Bruno Bongioanni and José L. Torrea. What is a Sobolev space for the Laguerre function systems? Studia Mathematica, 192(2):147-172, 2009.

Albrecht Böttcher and Sergei M. Grudsky. Toeplitz Matrices, Asymptotic Linear Algebra, and Functional Analysis. Birkhäuser Basel, Basel, 2000.
R. Y. Chen, A. Gittens, and J. A. Tropp. The masked sample covariance estimator: an analysis using matrix concentration inequalities. Information and Inference, 1(1):2-20, December 2012.

Fabienne Comte and Valentine Genon-Catalot. Adaptive Laguerre density estimation for mixed Poisson models. Electronic Journal of Statistics, 9(1):1113-1149, 2015.

Fabienne Comte and Gwennaëlle Mabon. Laguerre deconvolution with unknown matrix operator. Mathematical Methods of Statistics, 26(4):237-266, October 2017.

Fabienne Comte, Charles-A. Cuenod, Marianna Pensky, and Yves Rozenholc. Laplace deconvolution on the basis of time domain data and its application to dynamic contrast-enhanced imaging. Journal of the Royal Statistical Society: Series B (Statistical Methodology), 79(1):69-94, January 2017.

Kristof Croux and Noël Veraverbeke. Nonparametric estimators for the probability of ruin. Insurance: Mathematics and Economics, 9(2-3):127-130, September 1990.
fedja. Proving that the primitives of the laguerre functions are uniformly bounded. MathOverflow, February 2021. URL https://mathoverflow.net/q/383022.

Edward W. Frees. Nonparametric Estimation of the Probability of Ruin. ASTIN Bulletin, 16(S1): S81-S90, April 1986.

Hans U. Gerber and Elias S.W. Shiu. On the Time Value of Ruin. North American Actuarial Journal, 2(1):48-72, January 1998.
Gwennaëlle Mabon. Adaptive Deconvolution on the Non-negative Real Line: Adaptive deconvolution on R+. Scandinavian Journal of Statistics, 44(3):707-740, September 2017.

Esterina Masiello. On semiparametric estimation of ruin probabilities in the classical risk model. Scandinavian Actuarial Journal, 2014(4):283-308, June 2014.

Pascal Massart. The Tight Constant in the Dvoretzky-Kiefer-Wolfowitz Inequality. The Annals of Probability, 18(3):1269-1283, July 1990.
R. Mnatsakanov, L. L. Ruymgaart, and F. H. Ruymgaart. Nonparametric estimation of ruin probabilities given a random sample of claims. Mathematical Methods of Statistics, 17(1):35-43, March 2008.
S. M. Pitts. Nonparametric estimation of compound distributions with applications in insurance. Annals of the Institute of Statistical Mathematics, 46(3):537-555, September 1994.

Konstadinos Politis. Semiparametric Estimation for Non-Ruin Probabilities. Scandinavian Actuarial Journal, 2003(1):75-96, January 2003.

Yasutaka Shimizu. Estimation of the expected discounted penalty function for Lévy insurance risks. Mathematical Methods of Statistics, 20(2):125-149, June 2011.

Yasutaka Shimizu. Non-parametric estimation of the Gerber-Shiu function for the Wiener-Poisson risk model. Scandinavian Actuarial Journal, 2012(1):56-69, March 2012.

Wen Su and Wenguang Yu. Asymptotically Normal Estimators of the Gerber-Shiu Function in Classical Insurance Risk Model. Mathematics, 8(10):1638, September 2020.

Wen Su, Benxuan Shi, and Yunyun Wang. Estimating the Gerber-Shiu function under a risk model with stochastic income by Laguerre series expansion. Communications in Statistics - Theory and Methods, pages 1-23, May 2019a.

Wen Su, Yaodi Yong, and Zhimin Zhang. Estimating the Gerber-Shiu function in the perturbed compound Poisson model by Laguerre series expansion. Journal of Mathematical Analysis and Applications, 469(2):705-729, January 2019b.

Zhimin Zhang and Wen Su. A new efficient method for estimating the Gerber-Shiu function in the classical risk model. Scandinavian Actuarial Journal, 2018(5):426-449, May 2018.

Zhimin Zhang and Wen Su. Estimating the Gerber-Shiu function in a Lévy risk model by Laguerre series expansion. Journal of Computational and Applied Mathematics, 346:133-149, January 2019.

## A Model reduction procedure

We propose a model reduction procedure to choose the dimensions $m_{2}$ and $m_{3}$, defined by 12 . We explain the method for the choice of $m_{2}$ in the case $\delta=0$.

Let us assume we have estimated the $M$ first coefficients of $g$, for a large $M$. By Remark 2.6, we know that the best estimator is $\hat{g}_{M}$. Our goal is to choose $\hat{m}_{2}$ smaller than $M$ that achieves a similar MISE. This provides a parsimonious version of the estimator without degrading its MISE. By Theorem 2.5 the MISE of $\hat{g}_{m}$ is given by:

$$
\mathbb{E}\left\|g-\hat{g}_{m}\right\|_{\mathrm{L}^{2}}^{2} \leqslant\left\|g-g_{m}\right\|_{\mathrm{L}^{2}}^{2}+\frac{\lambda}{c^{2} T} \mathbb{E}[X] .
$$

Ideally, we would like to choose the first $m$ such that the bias term $\left\|g-g_{m}\right\|_{\mathrm{L}^{2}}^{2}$ is smaller than the variance term $\frac{\lambda}{c^{2} T} \mathbb{E}[X]$. Since these terms are unknown, we estimate them by $\sum_{k=m}^{M-1} \hat{b}_{k}^{2}$ and $\frac{1}{(c T)^{2}} \sum_{i=1}^{N_{T}} X_{i}$ respectively. We choose $\hat{m}_{2}$ as:

$$
\begin{equation*}
\hat{m}_{2}=\min \left\{1 \leqslant m \leqslant M-1 \left\lvert\, \sum_{k=m}^{M-1} \hat{b}_{k}^{2} \leqslant \frac{\kappa_{2}}{(c T)^{2}} \sum_{i=1}^{N_{T}} X_{i}\right.\right\}, \tag{39}
\end{equation*}
$$

with $\kappa_{2}$ an adjustment constant. The next proposition shows that the MISE of $\hat{g}_{\hat{m}_{2}}$ does not exceed the MISE of $\hat{g}_{M}$ by more than $\kappa_{2} \times$ (variance term).

Proposition A.1. Let $\kappa_{2}>0$, if $\hat{m}_{2}$ is chosen as then the MISE of $\hat{g}_{\hat{m}_{2}}$ is:

$$
\mathbb{E}\left\|g-\hat{g}_{\hat{m}_{2}}\right\|_{\mathrm{L}^{2}}^{2} \leqslant\left\|g-g_{M}\right\|_{\mathrm{L}^{2}}^{2}+\left(1+\kappa_{2}\right) \frac{\lambda}{c^{2} T} \mathbb{E}[X] .
$$

Proof. By Pythagoras Theorem:

$$
\begin{aligned}
\left\|g-\hat{g}_{\hat{m}_{2}}\right\|_{\mathrm{L}^{2}}^{2} & =\left\|g-\hat{g}_{M}\right\|_{\mathrm{L}^{2}}^{2}+\left\|\hat{g}_{M}-\hat{g}_{\hat{m}_{2}}\right\|_{\mathrm{L}^{2}}^{2} \\
& =\left\|g-\hat{g}_{M}\right\|_{\mathrm{L}^{2}}^{2}+\sum_{k=\hat{m}_{2}}^{M-1} \hat{b}_{k}^{2} \\
& \leqslant\left\|g-\hat{g}_{M}\right\|_{\mathrm{L}^{2}}^{2}+\frac{\kappa_{2}}{(c T)^{2}} \sum_{i=1}^{N_{T}} X_{i} .
\end{aligned}
$$

We take the expectation, and we apply Theorem 2.5

$$
\mathbb{E}\left\|g-\hat{g}_{\hat{m}_{2}}\right\|_{\mathrm{L}^{2}}^{2} \leqslant\left\|g-g_{M}\right\|_{\mathrm{L}^{2}}^{2}+\left(1+\kappa_{2}\right) \frac{\lambda}{c^{2} T} \mathbb{E}[X] .
$$

|  | With Model Reduction | Without Model Reduction |
| :---: | :---: | :---: |
| Ruin Probability | 1.06 | 1.06 |
|  | $[0.86,1.26]$ | $[0.86,1.25]$ |
|  | $\overline{m_{2}}=4.2, \overline{\hat{m}_{3}}=4.1$ | $m_{2}=m_{3}=N_{T}$ |
| Jump size causing the ruin | 40.1 | 41.6 |
|  | $[33.7,46.9]$ | $[34.7,48.4]$ |
| Laplace transform, $\delta=0.1$ | $\hat{m}_{2}=4.1, \overline{\hat{m}_{3}}=4.3$ | $m_{2}=m_{3}=N_{T}$ |
|  | 0.098 |  |
|  | $[0.075,0.110]$ | $[0.079,0.117]$ |
| $\hat{m}_{2}=3.9, \overline{\hat{m}_{3}}=4.0$ | $m_{2}=m_{3}=N_{T}$ |  |

Table 4: Comparison between the MISE of the Laguerre-Fourier estimator with and without model reduction. In each case, we chose the following parameters: $X \sim \operatorname{Exp}(1 / 2), \lambda=1.25, c=3$, $T=80$. With this set of parameters, $\mathbb{E}\left[N_{T}\right]=100$. Each cell displays an estimation of the MISE over 200 samples with a $95 \%$ confidence interval, and the mean selected models $\hat{m}_{2}$ and $\hat{m}_{3}$. In every case, $m_{1}$ is equal to $N_{T}$.

The same goes for $\hat{m}_{3}$ : we estimate the bias term by $\sum_{k=m}^{M-1} \hat{c}_{k}^{2}$ and the variance term by $\frac{1}{(c T)^{2}} \sum_{i=1}^{N_{T}} W\left(X_{i}\right)$; we choose $\hat{m}_{3}$ as:

$$
\hat{m}_{3}=\min \left\{1 \leqslant m \leqslant M-1 \left\lvert\, \sum_{k=m}^{M-1} \hat{c}_{k}^{2} \leqslant \frac{\kappa_{3}}{(c T)^{2}} \sum_{i=1}^{N_{T}} W\left(X_{i}\right)\right.\right\} .
$$

By the same arguments, the MISE of $\hat{h}_{m_{3}}$ is given by:

$$
\mathbb{E}\left\|h-\hat{h}_{\hat{m}_{3}}\right\|_{\mathrm{L}^{2}}^{2} \leqslant\left\|h-h_{M}\right\|_{\mathrm{L}^{2}}^{2}+\left(1+\kappa_{3}\right) \frac{\lambda}{c^{2} T} \mathbb{E}[W(X)] .
$$

In the case $\delta>0$, we choose the same $\hat{m}_{2}$ and $\hat{m}_{3}$ as in the case $\delta=0$. By the same arguments, we obtain:

$$
\begin{aligned}
\mathbb{E}\left\|g-\hat{g}_{\hat{m}_{2}}\right\|_{\mathrm{L}^{2}}^{2} & \leqslant\left\|g-g_{M}\right\|_{\mathrm{L}^{2}}^{2}+\frac{C(\lambda)}{c^{2} T}\left(\mathbb{E}[X]+\frac{\mathbb{E}\left[X^{2}\right]^{\frac{1}{2}}}{(1-\theta)^{2} \delta^{2}}\right)+\kappa_{2} \frac{\lambda}{c^{2} T} \mathbb{E}[X] \\
& \leqslant\left\|g-g_{M}\right\|_{\mathrm{L}^{2}}^{2}+\left(1+\kappa_{2}\right) \frac{C(\lambda)}{c^{2} T}\left(\mathbb{E}[X]+\frac{\mathbb{E}\left[X^{2}\right]^{\frac{1}{2}}}{(1-\theta)^{2} \delta^{2}}\right) \\
\mathbb{E}\left\|h-\hat{h}_{\hat{m}_{3}}\right\|_{\mathrm{L}^{2}}^{2} & \leqslant\left\|h-h_{M}\right\|_{\mathrm{L}^{2}}^{2}+\frac{C(\lambda)}{c^{2} T}\left(\mathbb{E}[W(X)]+\frac{\mathbb{E}\left[W(X)^{2}\right]^{\frac{1}{2}}}{(1-\theta)^{2} \delta^{2}}\right)+\kappa_{3} \frac{\lambda}{c^{2} T} \mathbb{E}[W(X)] \\
& \leqslant\left\|h-h_{M}\right\|_{\mathrm{L}^{2}}^{2}+\left(1+\kappa_{3}\right) \frac{C(\lambda)}{c^{2} T}\left(\mathbb{E}[W(X)]+\frac{\mathbb{E}\left[W(X)^{2}\right]^{\frac{1}{2}}}{(1-\theta)^{2} \delta^{2}}\right)
\end{aligned}
$$

Numerically, we compared the MISE's of the Laguerre-Fourier estimator with and without the model reduction procedure for $\hat{m}_{2}$ and $\hat{m}_{3}$, with the choice $\kappa_{2}=\kappa_{3}=0.3$. We show the results in Table 4 We see that the model reduction procedure does not affect the MISE of the estimator and we emphasize that the selected dimensions are far lower than the maximum dimension ( $\hat{m}$ 's are less than 10 whereas the maximum dimension is 100 ).

## B Uniform bound on the primitives of the Laguerre functions

In this section, we prove that the primitives of the Laguerre function are uniformly bounded. The sketch of the proof comes from fedja 2021.

Proof of Lemma 3.3. Let $u_{k}(x):=L_{k}(x) \mathrm{e}^{-x / 2}$. We first notice that the complete integral of $u_{k}$ is uniformly bounded:

$$
\begin{aligned}
\int_{0}^{+\infty} u_{k}(x) \mathrm{d} x & =2 \int_{0}^{+\infty} L_{k}(2 x) \mathrm{e}^{-x} \mathrm{~d} x \\
& =2 \sum_{j=0}^{k}\binom{k}{j} \frac{(-2)^{j}}{j!} \int_{0}^{+\infty} x^{j} \mathrm{e}^{-x} \mathrm{~d} x \\
& =2 \sum_{j=0}^{k}\binom{k}{j}(-2)^{j}=2(-1)^{k}
\end{aligned}
$$

We will show that $\left|\int_{0}^{x} u_{k}\right| \leqslant C\left|\int_{0}^{\infty} u_{k}\right|=2 C$ for an absolute constant $C>0$. The $k$-th Laguerre polynomial satisfies the ODE:

$$
x L_{k}^{\prime \prime}+(1-x) L_{k}^{\prime}+k L_{k}=0
$$

thus, the function $u_{k}$ satisfies:

$$
\begin{equation*}
x u_{k}^{\prime \prime}+u_{k}^{\prime}+\left(k+\frac{1}{2}-\frac{1}{4} x\right) u_{k}=0 \tag{40}
\end{equation*}
$$

To kill the first derivative, we consider $v_{k}(x):=u_{k}\left(x^{3 / 2}\right) x^{1 / 2}$. The functions $u_{k}$ and $v_{k}$ have the same partial integrals (up to the constant $2 / 3$ ):

$$
\forall x \geqslant 0, \quad \int_{0}^{x} v_{k}(t) \mathrm{d} t=\frac{2}{3} \int_{0}^{x^{2 / 3}} u_{k}(t) \mathrm{d} t .
$$

The first two derivatives of $v_{k}$ are:

$$
\begin{gathered}
v_{k}^{\prime}(x)=\frac{3}{2} x u_{k}^{\prime}\left(x^{3 / 2}\right)+\frac{1}{2} x^{-1 / 2} u\left(x^{3 / 2}\right) \\
v_{k}^{\prime \prime}(x)=\frac{9}{4} x^{3 / 2} u^{\prime \prime}\left(x^{3 / 2}\right)+\frac{9}{4} u^{\prime}\left(x^{3 / 2}\right)-\frac{1}{4} x^{-3 / 2} u\left(x^{3 / 2}\right)
\end{gathered}
$$

so using the ODE 40) for $u$ yields the following ODE for $v$ :

$$
v_{k}^{\prime \prime}+\Phi_{k} v_{k}=0
$$

where $\Phi_{k}$ is given by:

$$
\Phi_{k}(x):=\frac{9}{4}\left(\frac{k+\frac{1}{2}}{x^{1 / 2}}-\frac{1}{4} x\right)+\frac{1}{4 x^{2}} .
$$

The important properties of this function are that it is convex and decreasing.
Since the Laguerre polynomials have simple zeros, the function $v_{k}$ has $(k+1)$ simple zeros (the zeros of the $k$-th Laguerre polynomial, and 0 ) so the integral of $v_{k}$ can be decomposed as an alternating sum $\int_{0}^{\infty} v_{k}=: I_{k}=A_{0}-A_{1}+A_{2}-\cdots+(-1)^{k} A_{k}$, where the $A_{i}$ 's are the unsigned areas of the excursions of $v_{k}$ (see Figure 1). Based on the following lemma (proven later), we claim that $A_{0}<A_{1}<\cdots<A_{k}$.

Lemma B.1. Let $F_{1}, F_{2}$ be two $\mathcal{C}^{1}$ functions defined on an open interval I containing 0, and let $y_{1}, y_{2}$ be the solutions of the following $O D E$ :

$$
\left\{\begin{array}{l}
y_{1}^{\prime \prime}+F_{1}(x) y_{1}=0 \\
y_{2}^{\prime \prime}+F_{2}(x) y_{2}=0 \\
y_{1}(0)=y_{2}(0)=0 \\
y_{2}^{\prime}(0) \geqslant y_{1}^{\prime}(0)>0
\end{array}\right.
$$

Let $M>0$ such that $y_{1}$ and $y_{2}$ are positive on $J:=(0, M) \subseteq I$. If $F_{1}>F_{2}$ on $J$, then we have $y_{1} \leqslant y_{2}$ on $J$.


Figure 1: Graph of $v_{4}$. We see that the area of the excursions is increasing.

Indeed, let $z>0$ be a zero of $v_{k}$ and assume w.l.o.g. that $v_{k}$ is positive after $z$ and negative before. Let $y_{1}(x):=-v_{k}(z-x)$ and $y_{2}(x):=v_{k}(z+x)$ (i.e. $y_{1}$ is the central inversion of $v_{k}$ with respect to $z$ ). Thus, $y_{1}$ and $y_{2}$ satisfy the ODE:

$$
\left\{\begin{array}{l}
y_{1}^{\prime \prime}+\Phi_{k}(z-x) y_{1}=0 \\
y_{2}^{\prime \prime}+\Phi_{k}(z+x) y_{2}=0
\end{array}\right.
$$

with the initial conditions $y_{1}(0)=y_{2}(0)=v_{k}(z)=0$ and $y_{1}^{\prime}(0)=y_{2}^{\prime}(0)=v_{k}^{\prime}(z)>0$. Since $\Phi_{k}$ is decreasing, we have $\Phi_{k}(z-x)>\Phi_{k}(z+x)$, so Lemma B.1 yields that $y_{1}(x) \leqslant y_{2}(x)$ for $x>0$ as long as $y_{1}(x)$ and $y_{2}(x)$ are positive. Hence, the area of the excursion preceding $z$ is smaller than the area of the excursion following $z$.

Now let $z$ be the last zero of $v_{k}$, and let us assume w.l.o.g. that $v_{k}$ is positive after $z$ (otherwise we consider $-v_{k}$, it satisfies the same ODE than $v_{k}$ ). In this case, $I_{k}=A_{0}-A_{1}+\ldots-A_{k-1}+A_{k}$ with $A_{0}<\ldots<A_{k-1}<A_{k}$. Thus the maximum value of $\left|\int_{0}^{x} v_{k}\right|$ is attained either by taking the complete integral (it is the maximum of $\int_{0}^{x} v_{k}$ ), either by leaving out the last excursion (it is the minimum of $\int_{0}^{x} v_{k}$ ). We show that the second option is dominated by the first one: there exists an absolute constant $C>0$ such that $A_{k}-I_{k} \leqslant C I_{k}$. To do that it suffices to show that $A_{k-1} \leqslant c A_{k}$ for absolute constant $c \in(0,1)$, hence $A_{k}-I_{k} \leqslant \frac{c}{1-c} I_{k}$.

The strategy is to compare the function $v_{k}$ to the Airy function of the first kind. This function is solution of the ODE:

$$
\left\{\begin{array}{l}
y^{\prime \prime}-x y=0  \tag{41}\\
\lim _{x \rightarrow+\infty} y(x)=0 .
\end{array}\right.
$$

Let $z_{\ell}^{*}$ and $z_{p}^{*}$ be respectively the last and the penultimate zeros of Ai. We recall that $z_{\ell}^{*}$ is negative (all the zeros of Ai are negative), and that Ai is negative on $\left(z_{p}^{*}, z_{\ell}^{*}\right)$ and positive on $\left(z_{\ell}^{*},+\infty\right)$.
Lemma B.2. Let Ai be the Airy function of the first kind, and let $z$ be the last zero of $v_{k}$. There exists a decreasing linear function $\Upsilon_{k}$ satisfying $\Upsilon_{k}(z)=\Phi_{k}(z)$, and there exists $a>0$ and $b \in \mathbb{R}$, such that the function $\operatorname{Ai}(a x+b)$ satisfies:

$$
\left\{\begin{array}{l}
y^{\prime \prime}+\Upsilon_{k}(x) y=0 \\
y(z)=0
\end{array}\right.
$$

and such that it stays positive on $(z,+\infty)$ and tends to 0 at $+\infty$.
Let $w_{k}(x):=\varepsilon \operatorname{Ai}(a x+b)$ where $a$ and $b$ are given by Lemma B.2 and $\varepsilon>0$ is small enough such that $0<w_{k}^{\prime}(z)<v_{k}^{\prime}(z)$. Let $z_{\ell}$ and $z_{p}$ be respectively the last and the penultimate zeros of $w_{k}$. By definition of $w_{k}$, its last zero is the same than $v_{k}$ 's, that is $z_{\ell}=z$. Moreover, the zeros of $w_{k}$ and Ai are linked by the linear transformation $x \mapsto a x+b$ : we have $z_{\ell}^{*}=a z_{\ell}+b$ and $z_{p}^{*}=a z_{p}+b$.

Let us consider:

$$
\mathcal{W}:=\operatorname{det}\left(\begin{array}{cc}
v_{k} & w_{k} \\
v_{k}^{\prime} & w_{k}^{\prime}
\end{array}\right)=v_{k} w_{k}^{\prime}-v_{k}^{\prime} w_{k}
$$

the Wronskian of $v_{k}$ and $w_{k}$. Since $\mathcal{W}$ vanishes at $z$ and $+\infty$, we have:

$$
0=\int_{z}^{+\infty} \mathcal{W}^{\prime}(x) \mathrm{d} x=\int_{z}^{+\infty}\left(\Phi_{k}(x)-\Upsilon_{k}(x)\right) \underbrace{v_{k}(x) w_{k}(x)}_{>0} \mathrm{~d} x
$$

so the sign of $\Phi_{k}-\Upsilon_{k}$ must change on $(z,+\infty)$. Since $\Phi_{k}-\Upsilon_{k}$ is convex and $\left(\Phi_{k}-\Upsilon_{k}\right)(z)=0$, this is possible only if it is negative and then positive. Hence, $\mathcal{W}$ is first decreasing and then increasing on $(z,+\infty)$. Since it starts from zero at $z$ and tends to zero at $+\infty$, we conclude that $\mathcal{W}$ is negative on $(z,+\infty)$.

It follows that $w_{k}<v_{k}$ on $(z,+\infty)$. Indeed, by contradiction if $x_{0}:=\inf \left\{x \in(z,+\infty) \mid w_{k}(x) \geqslant\right.$ $\left.v_{k}(x)\right\}$ existed and was finite, we would have $v_{k}\left(x_{0}\right)=w_{k}\left(x_{0}\right)>0$ and $v_{k}^{\prime}\left(x_{0}\right) \leqslant w_{k}^{\prime}\left(x_{0}\right)$. Thus, we would have $\mathcal{W}\left(x_{0}\right)=v_{k}\left(x_{0}\right) w_{k}^{\prime}\left(x_{0}\right)-v_{k}^{\prime}\left(x_{0}\right) w_{k}\left(x_{0}\right) \geqslant 0$; contradiction. We conclude that the area of the last excursion of $w_{k}$ is less than $v_{k}$ 's: $\int_{z_{\ell}}^{\infty} w_{k} \leqslant A_{k}$.

We have seen that $\Phi_{k}-\Upsilon_{k}$ was negative then positive on $(z,+\infty)$. By convexity, it has to be positive on $(0, z)$, i.e. $\Phi_{k}>\Upsilon_{k}$ on the left of $z$. We apply Lemma B. 1 to $-v_{k}(z-x)$ and $-w_{k}(z-x)$, and we conclude that the area of the penultimate excursion of $v_{k}$ is less than $w_{k}$ 's: $A_{k-1} \leqslant\left|\int_{z_{p}}^{z_{\ell}} w_{k}\right|$.

We conclude that $A_{k-1} / A_{k}$ is bounded above by the ratio of the areas of the penultimate excursion to the last excursion of $w_{k}$. By a linear change of variable, this ratio is equal to the ratio of the areas of the penultimate excursion to the last excursion of the Airy function:

$$
\frac{A_{k-1}}{A_{k}} \leqslant \frac{\left|\int_{z_{p}}^{z_{\ell}} w_{k}(x) \mathrm{d} x\right|}{\int_{z_{\ell}}^{+\infty} w_{k}(x) \mathrm{d} x}=\frac{\left|\int_{z_{p}^{*}}^{z_{\ell}^{*}} \operatorname{Ai}(x) \mathrm{d} x\right|}{\int_{z_{\ell}^{*}}^{+\infty} \operatorname{Ai}(x) \mathrm{d} x}=: c
$$

This is an absolute constant, we just need to prove it is smaller than 1 to end the proof. The function Ai satisfies an ODE of the form $y^{\prime \prime}+F(x) y=0$ with $F(x)=-x$ a decreasing function, thus by considering the functions $y_{1}(x)=-\operatorname{Ai}\left(z_{\ell}^{*}-x\right)$ and $y_{2}:=\operatorname{Ai}\left(z_{\ell}^{*}+x\right)$, we can apply again Lemma B. 1 (as we did with $v_{k}$ ), to conclude that $y_{1}(x) \leqslant y_{2}(x)$ for $x>0$ as long as $y_{1}(x)$ is positive ( $y_{2}$ being positive for every $x>0$ ). This proves that $\left|\int_{z_{p}^{*}}^{z_{\ell}^{*}} \mathrm{Ai}\right|<\int_{z_{\ell}^{*}}^{\infty} \mathrm{Ai}$, that is $c<1$.
Proof of Lemma B.1. First, let us consider the case $y_{2}^{\prime}(0)>y_{1}^{\prime}(0)$. Let $\mathcal{W}:=y_{1} y_{2}^{\prime}-y_{1}^{\prime} y_{2}$ be the Wronskian of $y_{1}$ and $y_{2}$. Then $\mathcal{W}^{\prime}=\left(F_{1}-F_{2}\right) y_{1} y_{2}$ is positive on $J$ and $\mathcal{W}(0)=0$, so $\mathcal{W}>0$ on $J$.

By contradiction, suppose there exists $x \in J$ such that $y_{1}(x) \geqslant y_{2}(x)>0$ and consider $x_{0}:=\inf \left\{x \in J \mid y_{1}(x) \geqslant y_{2}(x)>0\right\}$. Since $y_{2}^{\prime}(0)>y_{1}^{\prime}(0)$, we know that $y_{1}$ is below $y_{2}$ on a right neighborhood of 0 , so $x_{0}>0$. By continuity, we have $y_{1}\left(x_{0}\right)=y_{2}\left(x_{0}\right)>0$, so we must have $y_{1}^{\prime}\left(x_{0}\right) \geqslant y_{2}^{\prime}\left(x_{0}\right)$; otherwise, we would have $y_{1}(x) \geqslant y_{2}(x)>0$ on a left neighborhood of $x_{0}$, which is impossible if $x_{0}>0$. Thus $\mathcal{W}\left(x_{0}\right)=y_{1}\left(x_{0}\right)\left[y_{2}^{\prime}\left(x_{0}\right)-y_{1}^{\prime}\left(x_{0}\right)\right] \leqslant 0$ with $x_{0} \in J$; contradiction.

Now consider the case $y_{2}^{\prime}(0)=y_{1}^{\prime}(0)$. For $\varepsilon>0$ small enough, we consider $y_{\varepsilon}$ the solution of:

$$
\left\{\begin{array}{l}
y_{\varepsilon}^{\prime \prime}+F_{1}(x) y_{\varepsilon}=0 \\
y_{\varepsilon}(0)=0 \\
y_{\varepsilon}^{\prime}(0)=y_{2}^{\prime}(0)-\varepsilon>0
\end{array}\right.
$$

Applying the first case, we have $y_{\varepsilon} \leqslant y_{2}$ while $y_{\varepsilon}$ and $y_{2}$ are positive. Taking $\varepsilon \rightarrow 0$, we obtain $y_{1} \leqslant y_{2}$.

Proof of Lemma B.2. The Airy function is solution of 41, so $\operatorname{Ai}(a x+b)($ with $a>0)$ satisfies:

$$
\left\{\begin{array}{l}
y^{\prime \prime}-a^{2}(a x+b) y=0 \\
\lim _{x \rightarrow+\infty} y(x)=0
\end{array}\right.
$$

We need to determine $(a, b)$ such that the two following conditions hold:

1. Let $\Upsilon_{k}(x):=-a^{2}(a x+b)$, we need to choose $a, b$ such that $\Upsilon_{k}(z)=\Phi_{k}(z)$.
2. Let $z_{\ell}^{*}$ be the last zero of Ai. We know that $z_{\ell}^{*}<0$ and that Ai is positive on $\left(z_{\ell}^{*},+\infty\right)$, thus we need to choose $a, b$ such that $a z+b=z_{\ell}^{*}$ so $\operatorname{Ai}(a x+b)$ stays positive on $(z,+\infty)$ and vanishes at $z$.

Thus, $(a, b) \in \mathbb{R}_{+}^{*} \times \mathbb{R}$ must be solution of:

$$
\left\{\begin{array} { l } 
{ - a ^ { 2 } ( a z + b ) = \Phi _ { k } ( z ) }  \tag{42}\\
{ a z + b = z _ { \ell } ^ { * } }
\end{array} \Longleftrightarrow \left\{\begin{array}{l}
a^{2}=-\frac{\Phi_{k}(z)}{z_{\ell}^{*}} \\
a z+b=z_{\ell}^{*}
\end{array}\right.\right.
$$

Since $z_{\ell}^{*}<0$, this system has a solution iff $\Phi_{k}(z)>0$. By contradiction, if we had $\Phi_{k}(z) \leqslant 0$, then $\Phi_{k}$ would be negative on $(z,+\infty)$. Since the function $v_{k}$ is positive after $z$ and satisfies $v_{k}^{\prime \prime}(x)=-\Phi_{k}(x) v_{k}(x)$, then $v_{k}$ would be strictly convex on $(z,+\infty)$. But the function $v_{k}$ starts from zero at $z$ with a positive derivative, stays positive, and tends to 0 at $+\infty$, so it cannot be strictly convex on $(z,+\infty)$; contradiction. Thus, the system (42) has a solution.

## C Miscellaneous results

Theorem C.1. Let $\mathbf{A}, \mathbf{B}$ be $(m \times m)$ matrices. If $\mathbf{A}$ is invertible and $\left\|\mathbf{A}^{-1} \mathbf{B}\right\|_{\mathrm{op}}<1$, then $\mathbf{A}+\mathbf{B}$ is invertible and it holds:

$$
\left\|(\mathbf{A}+\mathbf{B})^{-1}-\mathbf{A}^{-1}\right\|_{\mathrm{op}} \leqslant \frac{\left\|\mathbf{A}^{-1}\right\|_{\mathrm{op}}^{2}\|\mathbf{B}\|_{\mathrm{op}}}{1-\left\|\mathbf{A}^{-1} \mathbf{B}\right\|_{\mathrm{op}}}
$$

Proof. Since $\left\|\mathbf{A}^{-1} \mathbf{B}\right\|_{\mathrm{op}}<1$, its Neumann series is normally convergent and we have:

$$
\sum_{k=0}^{+\infty}(-1)^{k}\left(\mathbf{A}^{-1} \mathbf{B}\right)^{k}=\left(\mathbf{I}_{m}+\mathbf{A}^{-1} \mathbf{B}\right)^{-1}
$$

Hence:

$$
\begin{aligned}
\left\|(\mathbf{A}+\mathbf{B})^{-1}-\mathbf{A}^{-1}\right\|_{\mathrm{op}} & \leqslant\left\|\mathbf{A}^{-1}\right\|_{\mathrm{op}}\left\|\left(\mathbf{I}_{m}+\mathbf{A}^{-1} \mathbf{B}\right)^{-1}-\mathbf{I}_{m}\right\|_{\mathrm{op}} \\
& \leqslant\left\|\mathbf{A}^{-1}\right\|_{\mathrm{op}} \sum_{k=1}^{+\infty}\left\|\mathbf{A}^{-1} \mathbf{B}\right\|_{\mathrm{op}}^{k} \\
& =\left\|\mathbf{A}^{-1}\right\|_{\mathrm{op}} \frac{\left\|\mathbf{A}^{-1} \mathbf{B}\right\|_{\mathrm{op}}}{1-\left\|\mathbf{A}^{-1} \mathbf{B}\right\|_{\mathrm{op}}} \leqslant \frac{\left\|\mathbf{A}^{-1}\right\|_{\mathrm{op}}^{2}\|\mathbf{B}\|_{\mathrm{op}}}{1-\left\|\mathbf{A}^{-1} \mathbf{B}\right\|_{\mathrm{op}}} .
\end{aligned}
$$

Theorem C. 2 (Theorem A. 1 in Chen et al. 2012]). Suppose that $q \geqslant 2$ and fix $r \geqslant \max (q, 2 \log p)$. Consider a finite sequence $\left\{\mathbf{Y}_{i}\right\}$ of independent, symmetric, random, self-adjoint matrices with dimension $p \times p$. Then:

$$
\left[\mathbb{E} \lambda_{\max }\left(\sum_{i} \mathbf{Y}_{i}\right)^{q}\right]^{1 / q} \leqslant \sqrt{\operatorname{er} \lambda_{\max }\left(\sum_{i} \mathbb{E} \mathbf{Y}_{i}^{2}\right)}+2 e r\left[\mathbb{E} \max _{i} \lambda_{\max }^{q}\left(\mathbf{Y}_{i}\right)\right]^{1 / q}
$$


[^0]:    *florian.dussap@u-paris.fr

[^1]:    ${ }^{1}$ alternatively, we could split the data $\left\{X_{1}, \ldots, X_{N_{T}}\right\}$ in two parts: we use half of the data to estimate $\widetilde{\mathbf{A}}_{m, 2}^{-1}$, and the other half to estimate $\hat{\mathbf{c}}_{m}$.

[^2]:    ${ }^{2}$ this lemma is stated for the generalized Laguerre basis, which depends on a parameter $a$. This parameter is equal to 1 in our case.

