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# ON MEAN FIELD GAMES MODELS FOR EXHAUSTIBLE COMMODITIES TRADE 

P. JAMESON GRABER AND CHARAFEDDINE MOUZOUNI


#### Abstract

We investigate a mean field game model for the production of exhaustible resources. In this model, firms produce comparable goods, strategically set their production rate in order to maximise profit, and leave the market as soon as they deplete their capacities. We examine the related Mean Field Game system and prove well-posedness for initial measure data by deriving suitable a priori estimates. Then, we show that feedback strategies which are computed from the Mean Field Game system provide $\varepsilon$-Nash equilibria to the corresponding N -Player Cournot game, for large values of N . This is done by showing tightness of the empirical process in the Skorokhod M1 topology, which is defined for distribution-valued processes.


## 1. Introduction

Since its introduction about ten years ago, the theory of the Mean Field Games has expanded tremendously, and has become an important tool in the study of dynamical and equilibrium behavior of large systems. The theory was introduced separately by a series of seminal papers by Lasry and Lions [27-29] and Caines et al. [4,5], and in lectures by Pierre-Louis Lions at the Collège de France, which were video-taped and made available on the internet [32]. The main idea is inspired from statistical physics literature, and consists in considering that a given player interacts with competitors through their statistical distribution in the space of possible states.

Mean Field Games (MFG) theory provides a methodology to produce approximate Nash equilibria for stochastic differential games with symmetric interactions and a large (but finite) number of players N . In these games, the exact equilibrium strategies could be determined by a system of coupled Hamilton-Jacobi-Bellman equations, derived from the dynamic programming principle. However, the dimension of the system in general increases in N , which makes this system extremely hard to solve either analytically or numerically, especially for large values of N. The Mean Field Game approach simplifies the modelling, and allows to compute an approximation of Nash equilibria by solving a system of two forward-backward coupled PDEs. This simplification justifies partly the interest in the MFG modelling for several applications.

In this paper we revisit a family of MFG models related to competing producers with exhaustible resources. The dynamic market evolution is driven by the use of certain existing reserves to produce and trade comparable goods. Producers disappear from market as soon as they exhaust their capacities, so that the fraction of remaining firms decreases

[^0]over time. This type of model was first introduced by Guéant, Lasry, and Lions [22], and addressed also by Chan and Sircar in [13], where it is referred to as "Bertrand \& Cournot Mean Field Games." In [14], the same authors use a similar MFG modelling approach, to discuss recent changes in global oil market. A more sophisticated model for the energy industry is proposed recently in [33], where producers have also the possibility to explore new resources to replenish their reserves.

From a mathematical standpoint, Bertrand \& Cournot MFG system consists in a system of a backward Hamilton-Jacobi-Bellman (HJB) equation to model a representative firm's value function, coupled with a forward Fokker-Planck equation to model the evolution of the distribution of the active firms' states. The exhaustibility condition gives rise to absorbing boundary conditions at zero. A rigorous analysis of this system was provided in [21], where authors show existence of smooth solutions to the system of equations, and uniqueness under a certain restriction. Unconditional uniqueness is proved in [20], in addition to the analysis of the case with Neumann boundary conditions.

Otherwise, very little is known so far on the rigorous link between the so called Bertrand \& Cournot MFG models, and the corresponding N-Player Bertrand-Cournot stochastic differential games. Indeed, the classical theory cannot be applied to this specific case for two main reasons: on the one hand, because of the absorbing boundary conditions; and on the other hand, because in our model players are coupled through their controls, and therefore belongs to the class of extended Mean Field Games (cf. [3,10,18,19]). This has motivated the present work, in which we analyse rigorously this question for Cournot competition.

We investigate the mean-field approximation for N-Player continuous-time Cournot game with linear price schedule, and exhaustible resources. In this context, the producers' state variable is the reserves level, and the strategic variable is the rate of production. Producers disappear from the market as soon as they deplete their reserves, and the remaining active producers are constrained to set a non-negative rate of production, in order to manage their remaining reserves and maximize sales profit. Due to this constraint on the production rate - which is natural from a modeling view point - we obtain a Hamiltonian function that is less regular in comparison to [13,20,21]. Market demand is assumed to be linear, so that the received market price is a non-increasing linear function of the total production across all producers. Further details and explanations about the model will be given in Section 1.2.

We start by studying the resulting system of coupled PDEs (the MFG system) by deriving suitable a priori estimates in Hölder spaces. We shall assume that the initial data is a probability measure that is supported on $(0, L]$, for some $L>0$, which entails that all producers start with positive reserves. Hence, our analysis completes that which is found in [20,21], by treating the case of a less regular Hamiltonian function and initial measure data. Next, we prove that the feedback control given by the solution of the Mean Field Game system provides an $\varepsilon$-Nash equilibrium (cf. Definition 1.3) to the corresponding N -Player Cournot game, where the error $\varepsilon$ is arbitrary small for large enough N . This result shows that the MFG model is indeed a good approximation to the game with finitely many players, and reinforces numerical methods based on the MFG approach. As in the classical theory, the key argument in the proof of this result is a suitable law
of large numbers. In our context, the main mathematical challenge comes from the fact that agents interact through the boundary behaviour, and are coupled by means of their chosen production strategies. To prove a tailor-made law of large numbers, we employ a compactness method borrowed from [23,30], by showing tightness of the empirical process in the space of distribution valued càdlàg processes, endowed with Skorokhod's M1 topology [30]. In contrast to the classical tools used so far, this method does not provide an exact quantification of the error $\varepsilon$, which is its main downside. Nevertheless, this approach has proven to be convenient for studying systems with absorbing boundary conditions. We also believe that it could be extended to the case of a systemic common noise, just as [30] contains an analysis of a stochastic McKean-Vlasov equation. However, we do not address this case here, finding the analysis of the stochastic HJB/FP-system somewhat out of reach under our assumptions on the data (cf. [9, Section 4] and the hypotheses found there).

For background on Skorokhod's topologies for real valued processes, we refer the reader to [38] and references therein. The M1 topology is extended to the space of tempered distributions, and to more general spaces in [30]. The fact that the feedback MFG control provides $\varepsilon$-Nash equilibria for the corresponding differential games with a large (but finite) number of players, was first noticed by Caines et al. [ 4,5 ] and further developed in several works (see e.g. [11,25] among many others). Cournot games with exhaustible resources and finite number of agents is investigated by Harris et al. in [24], and the corresponding MFG models were studied in $[13,14,22,33]$ with different variants, and numerical simulations. We refer the reader to $[3,12,22,29]$ for further background on Mean Field Game theory.

The paper is organized as follows. In the remainder of this section, we give some technical notations and preliminaries, introduce the mathematical description of the N -Player Cournot game with limited resources and the corresponding Mean Field Game, and state the main results of this paper. In Section 2 we prove existence and uniqueness of regular solutions to the MFG system by deriving suitable Hölder estimates. In Section 3 we show that the feedback control computed from the MFG system is an $\varepsilon$-Nash equilibrium to the N -Player game. For that purpose, we start by showing the weak convergence of the empirical process with respect to the M1 topology, then we deduce the main result by recalling the interpretation of the MFG system in terms of games with mean-field interactions.
1.1. Notations and preliminaries. Throughout this article we fix $\mathrm{L}>0$, define $\mathrm{Q}:=$ $(0, \mathrm{~L})$, and $\mathrm{Q}_{\mathrm{T}}:=(0, \mathrm{~T}) \times(0, \mathrm{~L})$. For any domain $\mathcal{D}$ in $\mathbb{R}$ or $\mathbb{R}^{2}$ we define $\overline{\mathcal{D}}$ to be the closure of $\mathcal{D}, \mathrm{L}^{s}(\mathcal{D}), 1 \leqslant s \leqslant \infty$ to be the Lebesgue space of $s$-integrable functions on $\mathcal{D}$; $L^{s}(\mathcal{D})_{+}$to be the set of elements $w \in L^{s}(\mathcal{D})$ such that $w(x) \geqslant 0$ for a.e. $x \in \mathcal{D} ; W_{s}^{k}(\mathcal{D})$, $k \in \mathbb{N}, 1 \leqslant s \leqslant \infty$, to be the Sobolev space of functions having a weak derivatives up to order $k$ which are $s$-summable on $\mathcal{D} ; \mathcal{C}(\mathcal{D})$ to be the space of all continuous functions on $\mathcal{D} ; \mathcal{C}_{0}(\mathcal{D})$ to be the space of all continuous functions on $\mathcal{D}$ that vanish at infinity $\left(\mathcal{C}_{0}(\mathcal{D})=\right.$ $\mathcal{C}(\mathcal{D})$ when $\mathcal{D}$ is compact); $\mathcal{C}^{\theta}(\mathcal{D})$ to be the space of all Hölder continuous functions with exponent $\theta$ on $\mathcal{D} ; \mathcal{C}_{c}^{\infty}(\mathcal{D})$ to be the set of smooth functions whose support is a compact
included in $\mathcal{D} ; \mathcal{S}_{\mathbb{R}}$ denotes the space of rapidly decreasing functions, and $\mathcal{S}_{\mathbb{R}}^{\prime}$ the space of tempered distributions.

For a subset $\mathcal{D} \subset \overline{Q_{T}}$, we also define $\mathcal{C}^{1,2}(\mathcal{D})$ to be the set of all functions on $\mathcal{D}$ which are locally continuously differentiable in $t$ and twice locally continuously differentiable in $x$, and by $W_{s}^{1,2}(\mathcal{D})$ the space of elements of $L^{s}(\mathcal{D})$ having weak derivatives of the form $\partial_{t}^{j} \partial_{x}^{k}$ with $2 j+k \leqslant 2$, endowed with the following norm:

$$
\|w\|_{\mathcal{W}_{\mathrm{s}}^{1,2}}:=\sum_{2 \mathfrak{j}+\mathrm{k} \leqslant 2}\left\|\partial_{\mathfrak{t}}^{\mathrm{j}} \partial_{\chi}^{\mathrm{k}} w\right\|_{\mathrm{L}^{\mathrm{s}}}
$$

The space of $\mathbb{R}$-valued Radon measures on $\mathcal{D}$ is denoted $\mathfrak{M}(\mathcal{D})$, which we identify with $\mathcal{C}_{0}(\mathcal{D})^{*}$ endowed with weak* topology, and $\mathcal{P}(\mathcal{D}), \tilde{\mathcal{P}}(\mathcal{D})$ are respectively the convex subset of probability measures on $\mathcal{D}$, and the convex subset of sub-probability measures: that is the set of positive radon measures $\mu$, s.t. $\mu(\mathcal{D}) \leqslant 1$. For any measure $\mu \in \mathfrak{M}(\mathcal{D})$, we denote by $\operatorname{supp}(\mu)$ the support of $\mu$.

Throughout the paper, we fix a complete filtered probability space $\left(\Omega, \mathcal{F}, \mathbb{F}=\left(\mathcal{F}_{t}\right)_{t \geqslant 0}, \mathbb{P}\right)$, and suppose that is rich enough to fulfill the assumptions that will be formulated in this article. We also fix constants $r, \sigma, T>0$, and denote by C a generic constant whose precise value may change from line to line. We also use the notation $C(\alpha, \beta, \gamma)$ and the like to point out the dependence of some constant on parameters $\alpha, \beta, \gamma$. Moreover, we use the notation $X \sim \mu$ to define a random variable $X$ with law $\mu$. For any $\mathbb{R}$-valued function $w$ we define the positive and negative parts of $w$, respectively:

$$
w^{+}:=\frac{1}{2}(|w|+w), \quad \text { and } \quad w^{-}:=\frac{1}{2}(|w|-w) ;
$$

and for any $x, y \in \mathbb{R}$ we use the following notation for the minimum and maximum, respectively:

$$
x \wedge y:=\frac{1}{2}(x+y-|x-y|) ; \quad \text { and } \quad x \vee y:=\frac{1}{2}(x+y+|x-y|) .
$$

Let us recall a few basic facts on stochastic differential equation with reflecting boundary in a half-line. Given a random variable $V$ that is supported on $(-\infty, L]$, we look for a pair of a.s. continuous and adapted processes $\left(X_{t}\right)_{t \geqslant 0}$ and $\left(\xi_{t}^{X}\right)_{t \geqslant 0}$ such that:

$$
\begin{align*}
& X_{t}=V+\int_{0}^{\mathrm{t}} \mathrm{~b}\left(\mathrm{~s}, \mathrm{X}_{\mathrm{s}}\right) \mathrm{ds}+\sigma \mathrm{W}_{\mathrm{t}}-\int_{0}^{\mathrm{t}} \mathbb{1}_{\left\{\mathrm{X}_{\mathrm{s}}=\mathrm{L}\right\}} \mathrm{d} \xi_{\mathrm{s}}^{\mathrm{X}} \in(-\infty, \mathrm{L}], \\
& \xi_{\mathrm{t}}^{\mathrm{X}}=\int_{0}^{\mathrm{t}} \mathbb{1}_{\left\{\mathrm{X}_{\mathrm{s}}=\mathrm{L}\right\}} \mathrm{d} \xi_{s}^{\mathrm{X}},  \tag{1.1a}\\
& \mathrm{X}_{0}=\mathrm{V}, \xi_{0}^{\mathrm{X}}=0, \quad \text { and } \xi^{\mathrm{X}} \text { is nondecreasing, }
\end{align*}
$$

where $\left(W_{t}\right)_{t \geqslant 0}$ is a $\mathbb{F}$-Wiener process that is independent of $V$. The random process $\left(\mathrm{X}_{\mathrm{t}}\right)_{\mathrm{t} \geqslant 0}$ is the reflected diffusion, $\left(\xi_{\mathrm{t}}^{\mathrm{X}}\right)_{\mathrm{t} \geqslant 0}$ is the local time, and the above set of equations is called the Skorokhod problem. Throughout the paper, we shall write problem (1.1a) in the following simple form:

$$
\mathrm{d} X_{\mathrm{t}}=\mathrm{b}\left(\mathrm{t}, X_{\mathrm{t}}\right) \mathrm{dt}+\sigma \mathrm{d} W_{\mathrm{t}}-\mathrm{d} \xi_{\mathrm{t}}^{\mathrm{X}}, \quad X_{0}=\mathrm{V}
$$

Suppose that the function $b$ is bounded, and satisfies for some $K>0$ the following condition:

$$
\begin{equation*}
|\mathbf{b}(\mathrm{t}, \mathrm{x})-\mathrm{b}(\mathrm{t}, \mathrm{y})| \leqslant \mathrm{K}|x-y| \tag{1.1b}
\end{equation*}
$$

for all $t \in[0, T]$, and $x, y \in(-\infty, L]$. Then, it is well-known (see e.g. $[1,15])$ that under these conditions, problem (1.1a) has a unique solution on $[0, \mathrm{~T}]$. Moreover, this solution is given explicitly by:

$$
\begin{equation*}
X_{t}:=\Gamma_{t}(Y), \quad \xi_{t}^{X}:=Y_{t}-\Gamma_{t}(Y) \tag{1.1c}
\end{equation*}
$$

where the process $\left(Y_{t}\right)_{t \in[0, T]}$ is the solution to

$$
\begin{equation*}
Y_{\mathrm{t}}=\mathrm{V}+\int_{0}^{\mathrm{t}} \mathrm{~b}\left(\mathrm{~s}, \Gamma_{\mathrm{s}}(\mathrm{Y})\right) \mathrm{d} s+\sigma W_{\mathrm{t}} \tag{1.1d}
\end{equation*}
$$

and where $\Gamma$ is the so called Skorokhod map, that is given by

$$
\Gamma_{t}(Y):=Y_{t}-\sup _{0 \leqslant s \leqslant t}\left(L-Y_{s}\right)^{-} .
$$

Furthermore, notice that

$$
\begin{equation*}
\xi_{t}^{X}-\xi_{t+h}^{X} \geqslant \inf _{v \in[0, h]}\left(Y_{t}-Y_{t+v}\right) \tag{1.1e}
\end{equation*}
$$

for any $t \in[0, T)$ and $h \in(0, T-t)$. In fact, when $\xi_{t}^{X}<\xi_{t+h}^{X}$, then

$$
0<\xi_{t+h}^{X}:=\sup _{0 \leqslant s \leqslant t+h}\left(L-Y_{s}\right)^{-}=\sup _{t \leqslant s \leqslant t+h}\left(L-Y_{s}\right)^{-}=\left(Y_{v_{0}}-L\right)
$$

for some $t \leqslant v_{0} \leqslant t+h$. Therefore

$$
\begin{aligned}
\xi_{t}^{X}-\xi_{t+h}^{X} & =\sup _{0 \leqslant s \leqslant t}\left(L-Y_{s}\right)^{-}-\sup _{0 \leqslant s \leqslant t+h}\left(L-Y_{s}\right)^{-} \\
& \geqslant\left(Y_{t}-L\right)-\left(Y_{v_{0}}-L\right) \geqslant \inf _{v \in[0, h]}\left(Y_{t}-Y_{t+v}\right) .
\end{aligned}
$$

This entails (1.1e) since the last inequality still holds when $\xi_{t}^{X}=\xi_{t+h}^{X}$.
Now we consider a boundary value problem for the Fokker-Planck equation. Let b in $\mathrm{L}^{2}\left(\mathrm{Q}_{\mathrm{T}}\right), \mathrm{m}_{0} \in \mathcal{P}(\overline{\mathrm{Q}})$, and consider the following Fokker-Planck equation

$$
\left\{\begin{array}{l}
\mathfrak{m}_{\mathrm{t}}-\frac{\sigma^{2}}{2} \mathfrak{m}_{x x}-(\mathrm{bm})_{x}=0 \text { in } Q_{T}  \tag{1.2a}\\
\mathfrak{m}(0)=\mathfrak{m}_{0} \text { in } \mathrm{Q}
\end{array}\right.
$$

complemented with the following mixed boundary conditions:

$$
\begin{equation*}
\mathfrak{m}(\mathrm{t}, 0)=0, \quad \text { and } \quad \frac{\sigma^{2}}{2} \mathfrak{m}_{\mathrm{x}}(\mathrm{t}, \mathrm{~L})+\mathrm{b}(\mathrm{t}, \mathrm{~L}) \mathfrak{m}(\mathrm{t}, \mathrm{~L})=0 \quad \text { on }(0, \mathrm{~T}) \tag{1.2b}
\end{equation*}
$$

Then we define a weak solution to (1.2a)-(1.2b) to be a function $m \in L^{1}\left(Q_{T}\right)_{+}$such that $\mathrm{m}|\mathrm{b}|^{2}$ in $\mathrm{L}^{1}\left(\mathrm{Q}_{\mathrm{T}}\right)$, and

$$
\begin{equation*}
\int_{0}^{\mathrm{T}} \int_{0}^{\mathrm{L}} \mathfrak{m}\left(-\phi_{\mathrm{t}}-\frac{\sigma^{2}}{2} \phi_{x x}+\mathrm{b} \phi_{x}\right) \mathrm{d} x \mathrm{dt}=\int_{0}^{\mathrm{L}} \phi(0, .) \mathrm{d} m_{0} \tag{1.2c}
\end{equation*}
$$

for every $\phi \in \mathcal{C}_{c}^{\infty}([0, T) \times \overline{\mathrm{Q}})$ satisfying

$$
\begin{equation*}
\phi(\mathrm{t}, 0)=\phi_{\chi}(\mathrm{t}, \mathrm{~L})=0, \quad \forall \mathrm{t} \in(0, \mathrm{~T}) . \tag{1.2d}
\end{equation*}
$$

This is the definition given by Porretta in [35]. The only difference is that here we consider mixed boundary conditions and measure initial data.

When $m_{0} \in \mathrm{~L}^{1}(\mathrm{Q})_{+}$, the problem (1.2a) endowed with periodic, Dirichlet or Neumann boundary conditions has several interesting features that were pointed out in [35, Section 3]. In particular, they are unique [35, Corollary 3.5] and enjoy some extra regularity [35, Proposition 3.10]. Note that these results still hold in the case of mixed boundary conditions (1.2b). Throughout the paper, we shall use the results of [35, Section 3] for (1.2a)-(1.2b).

In the case where $b$ is bounded, we shall use the fact that (1.2a)-(1.2b) admits a unique weak solution, for any $m_{0} \in \mathcal{P}(\overline{\mathrm{Q}})$. In fact, one can construct a solution by considering a suitable approximation of $m_{0}$, and then use the compactness results of [35, Proposition 3.10] in order to pass to the limit in $L^{1}\left(Q_{T}\right)$. The uniqueness is obtained by considering the dual equation, and using the same steps as for [35, Corollary 3.5] (cf. Proposition B.1).
1.2. Mathematical description of the model and main results. Let us now give a precise description of the problems considered in this paper. Consider a market with N producers of a given good whose strategic variable is the rate of production, and where raw materials are in limited supply. Concretely, one can think of energy producers that use exhaustible resources, such as oil, to produce and sell energy. Firms disappear from the market as soon as they deplete their reserves of raw materials.

Let us formalize this model in precise mathematical terms. Let $\left(W^{j}\right)_{1 \leqslant j \leqslant N}$ be a family of N independent $\mathbb{F}$-Wiener processes on $\mathbb{R}$, and consider the following system of Skorokhod problems:

$$
\left\{\begin{array}{l}
d X_{t}^{i}=-q_{t}^{i} d t+\sigma d W_{t}^{i}-d \xi_{t}^{x^{i}},  \tag{1.3}\\
X_{0}^{i}=V_{i}, \quad i=1, \ldots, N .
\end{array}\right.
$$

Here $\left(\mathrm{V}_{1}, \ldots, \mathrm{~V}_{\mathrm{N}}\right)$ is a vector of i.i.d and $\mathcal{F}_{0}$-measurable random variables, and we assume that $V_{1}, \ldots, V_{N}$ are independent of $W^{1}, \ldots, W^{N}$ respectively. Let us fix a common horizon $T>0$, and set

$$
\tau^{i}:=\inf \left\{t \geqslant 0: X_{t}^{i} \leqslant 0\right\} \wedge T .
$$

The stopped random process $\left(X_{t \wedge \tau^{i}}^{i}\right)_{t \in[0, T]}$ models the reserves level of the $i^{\text {th }}$ producer on the horizon T , which is gradually depleted according to a non-negative controlled rate of production $\left(q_{t}^{i}\right)_{t \in[0, T]}$. The stopping condition indicates that a firm can no longer replenish its reserves once they are exhausted, and the Wiener processes in (1.3) model the idiosyncratic fluctuations related to production. In addition, we consider $L$ to be an upper bound on the reserves level of any player. This latter assumption is also considered in $[20,21]$ and is taken into account by considering reflected dynamics in (1.3). Since the rate of production is always non-negative, note that reflection has practically no effect when $L$ is large compared to the initial reserves and $\sigma$.

Remark 1.1 (State constraints). Instead of reflecting boundary conditions, one could insist upon a hard state constraint of the form $X_{t}^{i} \leqslant L$. Some recent work on MFG with state constraints suggests this is possible [6-8], provided one correctly interprets the resulting
system of PDE (the Fokker-Planck equation presents a special challenge). In this work we take a more classical approach, for which probabilistic tools are more readily available.

The producers interact through the market. We assume that demand is linear, so that the price $p^{i}$ received by the firm $i$ reads:

$$
\begin{equation*}
p_{t}^{i}=1-\left(q_{t}^{i}+\kappa \bar{q}_{t}^{i}\right), \quad \text { where } \quad \bar{q}_{t}^{i}=\frac{1}{N-1} \sum_{j \neq i} q_{t}^{j} \mathbb{1}_{t<\tau^{j}}, \quad \text { for } 0 \leqslant t \leqslant T \text {. } \tag{1.4}
\end{equation*}
$$

Here $\kappa>0$ expresses the degree of market interaction, in proportion to which abundant total production will put downward pressure on all the prices. Note that only firms with nonempty reserves at $t \in[0, \mathrm{~T}]$ are taken into account in (1.4). The other firms are no longer present on the market. The producer $i$ chooses a non-negative production rate $q^{i}$ in order to maximize the following discounted profit functional:

$$
\mathcal{J}_{c}^{i, N}\left(q^{1}, \ldots, q^{N}\right):=\mathbb{E}\left\{\int_{0}^{T} e^{-r s}\left(1-\kappa \bar{q}_{s}^{i}-q_{s}^{i}\right) q_{s}^{i} \mathbb{1}_{s<\tau^{i}} d s+e^{-r T} u_{T}\left(X_{\tau^{i}}^{i}\right)\right\},
$$

where the terminal profit $u_{\top}$ is a smooth function satisfying $u_{\top}(0)=0$. Observe that firms can no longer earn revenue as soon as they deplete their reserves. We refer to $[14,24]$ for further explanations on the economic model and applications.

We denote by $\mathbb{A}_{c}$ the set of admissible controls for any player; that is the set of Markovian feedback controls, i.e. $q_{t}^{i}=q^{i}\left(t, X_{t}^{1}, \ldots, X_{t}^{N}\right)$; such that $\left(q_{t}^{i}\right)_{t \in[0, T]}$ is non-negative, satisfies

$$
\mathbb{E}\left[\int_{0}^{T}\left|q_{s}^{i}\right|^{2} \mathbb{1}_{s<\tau^{i}} \mathrm{ds}\right]<\infty,
$$

and the $i^{\text {th }}$ equation of (1.3) is well-posed in the classical sense. Restriction to Markovian controls rules out equilibria with undesirable properties such as non-credible threats (cf. [16, Chapter 13]).

Now, we give a definition of Nash equilibria to this game:
Definition 1.2 (Nash equilibrium). A strategy profile $\left(q^{1, *}, \ldots, q^{N, *}\right)$ in $\prod_{i=1}^{N} \mathbb{A}_{c}$ is a Nash equilibrium of the $N$-Player Cournot game, if for any $i=1, \ldots, N$ and $q^{i} \in \mathbb{A}_{c}$

$$
\mathcal{J}_{\mathfrak{c}}^{i, N}\left(q^{i} ;\left(\mathfrak{q}^{\mathfrak{j}, *}\right)_{j \neq i}\right) \leqslant \mathcal{J}_{\mathfrak{c}}^{i, N}\left(q^{1, *}, \ldots, q^{N, *}\right) .
$$

In words, a Nash equilibrium is a set of admissible strategies such that each player has taken an optimal trajectory in view of the competitors' choices.

The existence of Nash equilibria for the N-Player Cournot game with exhaustible resources is addressed in [24]. In particular, the authors show the existence of a unique Nash equilibrium in the static (one period) case, and study numerically a specific duopoly example by using a convenient asymptotic expansion. In general, the analysis of equilibria for N -Player Cournot games is a challenging task both analytically and numerically, especially when N is large. In the case of exhaustible resources, the dynamic programming principle generates an even more complex PDE system because of the nonstandard boundary conditions which are obtained (cf. [24, Section 3.1]).

To remedy this problem several works have rather considered a Mean-Field model $[13,14,22,24,33]$ as an approximation to the initial $N$-Player game, when $N$ is large. More precisely, we introduce the following:

Definition 1.3 ( $\varepsilon$-Nash equilibrium). Let $\varepsilon>0$, and let $\left(\hat{q}^{1}, \ldots, \hat{q}^{N}\right)$ be an admissible strategy profile (i.e. an element of $\prod_{i=1}^{N} \mathbb{A}_{c}$ ). We say $\left(\hat{q}^{1}, \ldots, \hat{q}^{N}\right)$ provides an $\varepsilon$-Nash equilibrium to the game $\partial_{c}^{1, N}, \ldots, \partial_{c}^{N}, N$ provided that, for any $i=1, \ldots, N$ and $q^{i} \in \mathbb{A}_{c}$,

$$
\mathcal{J}_{c}^{i, N}\left(q^{i} ;\left(\hat{q}^{j}\right)_{j \neq i}\right) \leqslant \varepsilon+\mathcal{J}_{\mathcal{c}}^{i, N}\left(\hat{q}^{1}, \ldots, \hat{q}^{N}\right) .
$$

In words, an $\varepsilon$-Nash equilibrium is a set of admissible strategies such that each player has taken an almost optimal trajectory in view of the competitors' choices, where $\varepsilon$ measures the distance from optimality.

To construct $\varepsilon$-Nash equilibria, we turn to the corresponding mean field problem. Let us now consider a continuum of agents, producing and selling comparable goods. At time $t=0$, all the players have a positive capacity $x \in(0, L]$, and are distributed on $(0, L]$ according to $m_{0}$.

The remaining capacity (or reserves) of any atomic producer with a production rate $(\rho)_{t \geqslant 0}$ depletes according to

$$
\mathrm{d} X_{t}^{\rho}=-\rho_{\mathrm{t}} \mathbb{1}_{\mathrm{t}<\tau^{\rho}} \mathrm{dt}+\sigma \mathbb{1}_{\mathrm{t}<\tau^{\rho}} \mathrm{d} W_{\mathrm{t}}-\mathrm{d} \xi_{\mathrm{t}}^{\mathrm{X}^{\rho}}
$$

where

$$
\tau^{\rho}:=\inf \left\{t \geqslant 0: X_{t}^{\rho} \leqslant 0\right\} \wedge T,
$$

and $\left(W_{t}\right)_{t \in[0, T]}$ is a $\mathbb{F}$-Wiener process. A generic player which anticipates the total production $\bar{q}$ expects to receive the price

$$
p:=1-(\kappa \bar{q}+\rho)
$$

and solves the following optimization problem:

$$
\begin{equation*}
\max _{\rho \geqslant 0} \mathcal{J}_{c}(\rho):=\max _{\rho \geqslant 0} \mathbb{E}\left\{\int_{0}^{T} e^{-r s}\left(1-\kappa \bar{q}_{s}-\rho_{s}\right) \rho_{s} \mathbb{1}_{s<\tau^{\rho}} d s+e^{-r T} u_{T}\left(X_{T}^{\rho}\right)\right\} . \tag{1.5}
\end{equation*}
$$

The maximum in (1.5) is taken over all $\mathbb{F}$-adapted and non-negative processes $\left(\rho_{\mathbf{t}}\right)_{\mathbf{t} \in[0, \mathrm{~T}]}$ such that

$$
\mathbb{E}\left[\int_{0}^{T}\left|\rho_{s}\right|^{2} \mathbb{1}_{s<\tau^{\rho}} \mathrm{ds}\right]<\infty
$$

and $\left(X_{t}^{\rho}\right)_{t \in[0, T]}$ exists in the classical sense.

According to MFG theory, the equilibrium in this setting can be computed by solving the following coupled system of parabolic partial differential equations:

$$
\left\{\begin{array}{l}
\mathfrak{u}_{\mathrm{t}}+\frac{\sigma^{2}}{2}{u_{x x}-r u+q_{u}}_{2}=0 \quad \text { in } Q_{T},  \tag{1.6}\\
\mathfrak{m}_{\mathrm{t}}-\frac{\sigma^{2}}{2} \mathfrak{m}_{x x}-\left\{\mathfrak{q}_{u, m} \mathfrak{m}\right\}_{x}=0 \quad \text { in } Q_{T}, \\
\mathfrak{m}(t, 0)=0, \quad u(t, 0)=0, \quad u_{x}(t, L)=0 \quad \text { in }(0, T), \\
\mathfrak{m}(0)=m_{0}, \quad u(T, x)=u_{T}(x), \quad \text { in }[0, L] \\
\frac{\sigma^{2}}{2} m_{x}+q_{u, m} m=0 \quad \text { in }(0, T) \times\{L\},
\end{array}\right.
$$

where the function $q_{u, m}$ involved in the system is given by:

$$
\begin{equation*}
\mathrm{q}_{\mathfrak{u}, \mathfrak{m}}(\mathrm{t}, \mathrm{x}):=\frac{1}{2}\left(1-\kappa \overline{\mathfrak{q}}(\mathrm{t})-\mathbf{u}_{x}(\mathrm{t}, \mathrm{x})\right)^{+}, \quad \text { where } \quad \overline{\mathrm{q}}(\mathrm{t}):=\int_{0}^{\mathrm{L}} \mathrm{q}_{\mathfrak{u}, \mathfrak{m}}(\mathrm{t}, x) \mathfrak{m}(\mathrm{t}, \mathrm{x}) \mathrm{d} x, \tag{1.7}
\end{equation*}
$$

and $k>0$. Here $m$ is the density of a continuum of market actors, $q_{u, m}(t, x)$ is the optimal production rate of an atomic player with reserves $x$ at time $t$, and $u$ is the the game value function of an atomic player following the production policy $q_{u, m}$.

Let us assume that $u_{T}$ is a function in $\mathcal{C}^{2}(\bar{Q})$, such that the first derivative of $u_{T}$ denoted by $u_{T}^{\prime}$ fulfils

$$
\begin{equation*}
u_{\mathrm{T}}^{\prime} \geqslant 0 \quad \text { and } \quad u_{\mathrm{T}}(0)=u_{\mathrm{T}}^{\prime}(\mathrm{L})=0 \tag{F1}
\end{equation*}
$$

and that $m_{0}$ is a probability measure with support away from 0 , i.e.

$$
\begin{equation*}
\mathfrak{m}_{0} \in \mathcal{P}(\overline{\mathrm{Q}}), \text { and } \operatorname{supp}\left(\mathfrak{m}_{0}\right) \subset(0, \mathrm{~L}] . \tag{H2}
\end{equation*}
$$

We shall say that a pair $(u, m)$ is a solution to (1.6), if
(i) $u \in \mathcal{C}^{1,2}\left(Q_{T}\right), u, u_{x} \in \mathcal{C}\left(\overline{Q_{T}}\right)$;
(ii) $\mathfrak{m} \in \mathcal{C}([0, \mathrm{~T}] ; \mathfrak{M}(\overline{\mathrm{Q}})) \cap \mathrm{L}^{1}\left(\mathrm{Q}_{\mathrm{T}}\right)_{+}$, and $\|\mathfrak{m}(\mathrm{t})\|_{\mathrm{L}^{1}} \leqslant 1$ for every $\mathrm{t} \in(0, \mathrm{~T}]$;
(iii) the equation for $u$ holds in the classical sense, while the equation for $m$ holds in the weak sense (1.2c).
The following lemma establishes the connection between ( $u, m$ ) and problem (1.5):
Lemma 1.4. Let $(\mathfrak{u}, \mathfrak{m})$ be a solution to (1.6) and set $\rho_{\mathrm{t}}^{*}:=\mathrm{q}_{\mathrm{u}, \mathrm{m}}\left(\mathrm{t}, \mathrm{X}_{\mathrm{t}}^{\rho^{*}}\right)$. Then

$$
\begin{equation*}
\max _{\rho \geqslant 0} \mathcal{I}_{c}(\rho)=\mathcal{J}_{c}\left(\rho^{*}\right)=\int_{0}^{\mathrm{L}} u(0, .) d m_{0} . \tag{1.8}
\end{equation*}
$$

The proof of Lemma 1.4 is standard, and is given in Appendix A. We deduce that the MFG system (1.6) describes an equilibrium configuration for a Cournot game with exhaustible resources and a continuum of producers.

We are now in a position to state the main results of this paper.
Theorem 1.5 (Well-posedness). There exists a unique solution ( $\mathbf{u}, \mathrm{m}$ ) to system (1.6).

Theorem 1.6 (Existence of $\varepsilon$-Nash equilibria). Let $(\mathfrak{u}, \mathfrak{m})$ be the solution to the MFG system (1.6), and let $\mathrm{q}_{\mathrm{u}, \mathrm{m}}$ be given by (1.7). For any $\mathrm{N} \geqslant 1$ and $\mathrm{i} \in\{1, \ldots, \mathrm{~N}\}$ let

$$
\left\{\begin{array}{l}
d \hat{X}_{t}^{i}=-q_{u, m}\left(t, \hat{X}_{t}^{i}\right) d t+\sigma d W_{t}^{i}-d \xi_{t}^{\hat{X}^{i}}  \tag{1.9}\\
X_{0}^{i}=V_{i},
\end{array}\right.
$$

and set $\hat{\mathrm{q}}_{\mathrm{t}}^{\mathrm{i}}:=\mathrm{q}_{\mathrm{u}, \mathrm{m}}\left(\mathrm{t}, \hat{X}_{\mathrm{t}}^{\mathrm{i}}\right)$. Then for any $\varepsilon>0$, the strategy profile $\left(\hat{\mathrm{q}}^{1}, \ldots, \hat{\mathrm{q}}^{\mathrm{N}}\right)$ is admissible, i.e. belongs to $\prod_{i=1}^{\mathrm{N}} \mathbb{A}_{c}$, and provides an $\varepsilon$-Nash equilibrium to the game $\partial_{\mathfrak{c}}^{1, \mathrm{~N}}, \ldots, \partial_{c}^{\mathrm{N}, \mathrm{N}}$ for large N . Namely: $\forall \varepsilon>0, \exists \mathrm{~N}_{\varepsilon} \geqslant 1$ such that

$$
\begin{equation*}
\forall N \geqslant N_{\varepsilon}, \forall i=1, \ldots, N, \quad \mathcal{J}_{c}^{i, N}\left(q^{i} ;\left(\hat{q}^{j}\right)_{j \neq i}\right) \leqslant \varepsilon+\mathcal{J}_{c}^{i, N}\left(\hat{q}^{1}, \ldots, \hat{q}^{N}\right), \tag{1.10}
\end{equation*}
$$

for any admissible strategy $\mathrm{q}_{\mathrm{i}} \in \mathbb{A}_{\mathbf{c}}$.
The proof of Theorem 1.5 is given in Section 2, while Section 3 is devoted to the proof of Theorem 1.6.

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## 2. Analysis of the MFG system

This section is devoted to the analysis of (1.6). We begin by establishing some preliminary estimates having to do with (i) the existence and regularity of $q_{u, m}$ defined in (1.7), and (ii) some regularity properties of weak solutions to the Fokker-Planck equation. Then we prove some a priori bounds on solutions of the system (1.6), which gives way to our proof of Theorem 1.5 by means of classical fixed point theory.
2.1. Preliminary estimates. We start by giving an alternative convenient expression for the production rate function $q_{u, m}$. We aim to write $q_{u, m}$ as a functional of $u_{x}, m$ and the market price function $p_{u, m}$, that is defined by [13]:

$$
\begin{equation*}
p_{u, m}(t, x):=1-\left(q_{u}, \mathfrak{m}(t, x)+\kappa \bar{q}(t)\right) . \tag{2.1}
\end{equation*}
$$

The latter expression means that the price $p_{u, m}(t, x)$ received by an atomic player with reserves $x$ at time $t$, is a linear and nonincreasing function, of the player's production rate $q_{u, m}(t, x)$, and the aggregate production rate across all producers $\bar{q}(t)$. For any $\mu \in \mathfrak{M}(\overline{\mathrm{Q}})$, we define

$$
\begin{equation*}
a(\mu):=\frac{1}{1+\kappa \eta(\mu)} ; \quad c(\mu):=1-a(\mu) ; \quad \eta(\mu):=\int_{0}^{L} d|\mu| \tag{2.2}
\end{equation*}
$$

and set

$$
\begin{equation*}
\overline{\mathrm{p}}(\mathrm{t}):=\frac{1}{\eta(\mathrm{~m}(\mathrm{t}))} \int_{0}^{\mathrm{L}} p_{\mathrm{u}, \mathrm{~m}}(\mathrm{t}, \mathrm{x}) \mathrm{m}(\mathrm{t}, \mathrm{x}) \mathrm{d} x . \tag{2.3a}
\end{equation*}
$$

By integrating (2.1) with respect to $m$ and after a little algebra one recovers the following identity

$$
a(m(t))+c(m(t)) \bar{p}(t)=1-k \bar{q}(t)
$$

which entails

$$
\begin{equation*}
p_{u, m}(t, x)=a(m(t))+c(m(t)) \bar{p}(t)-q_{u, m}(t, x), \tag{2.3b}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathrm{q}_{\mathrm{u}, \mathrm{~m}}(\mathrm{t}, \mathrm{x})=\frac{1}{2}\left\{\mathrm{a}(\mathrm{~m}(\mathrm{t}))+\mathrm{c}(\mathrm{~m}(\mathrm{t})) \overline{\mathrm{p}}(\mathrm{t})-\mathrm{u}_{\mathrm{x}}(\mathrm{t}, \mathrm{x})\right\}^{+} . \tag{2.3c}
\end{equation*}
$$

This duality is also known as Bertrand and Cournot equivalence, and expresses the fact that the problem of controlling the rate of production by anticipating global production, is equivalent to the problem of controlling the selling price by anticipating the average price in the market and the rate of active producers. We omit the details and refer to [13, Section B.2]. For convenience, we shall often use (2.3c) as a definition for $q_{u, m}$.

In contrast to the systems studied in $[13,20,21], p_{u, m}$ has no explicit formula and is only defined as a fixed point through (2.3a)-(2.3c). The following Lemma makes that statement clear and point out a few facts on the market price function.

Lemma 2.1. Let $u \in \mathrm{~L}^{\infty}\left(0, \mathrm{~T} ; \mathrm{C}^{1}(\overline{\mathrm{Q}})\right)$, $\mathrm{m} \in \mathrm{L}^{1}\left(\mathrm{Q}_{\mathrm{T}}\right)_{+}$, and $\mathrm{k}>0$. Then the market price function $\mathrm{p}_{\mathrm{u}, \mathrm{m}}$ is well-defined through (2.3a)-(2.3c), belongs to $\mathrm{L}^{\infty}(0, \mathrm{~T} ; \mathrm{C}(\overline{\mathrm{Q}}))$, and satisfies

$$
\begin{equation*}
-\left\|u_{x}\right\|_{\infty} \leqslant p_{u, m} \leqslant 1 \tag{2.4}
\end{equation*}
$$

Moreover, if $\mathrm{u}_{\mathrm{x}}$ is non-negative, then $\mathrm{p}_{\mathrm{u}, \mathrm{m}}$ is non-negative as well.
Proof. Let $f: \mathbb{R}^{2} \rightarrow \mathbb{R}$ be given by $f(x, y)=x-\frac{1}{2}(x-y)^{+}$. Note that $f$ is 1 -Lipschitz in the first variable, and $\frac{1}{2}$-Lipschitz in the second. For any $p, w \in \mathbb{X}:=L^{\infty}(0, T ; \mathcal{C}(\bar{Q}))$, define

$$
\ell(\mathfrak{m}, \mathfrak{p})(\mathrm{t}):=\mathrm{a}(\mathfrak{m}(\mathrm{t}))+\mathrm{c}(\mathfrak{m}(\mathrm{t})) \overline{\mathrm{p}}(\mathrm{t}), \quad \text { where } \quad \overline{\mathfrak{p}}(\mathrm{t}):=\frac{1}{\eta(\mathfrak{m}(\mathrm{t}))} \int_{0}^{\mathrm{L}} \mathrm{p}(\mathrm{t}, \mathrm{x}) \mathfrak{m}(\mathrm{t}, \mathrm{x}) \mathrm{dx}
$$

and

$$
\Lambda(w, m, p)(t, x):=f(\ell(m, p)(t), w(t, x)) .
$$

We note the following inequalities for future reference:

$$
\begin{align*}
\left|\ell(m, p)(t)-\ell\left(m, p^{\prime}\right)(t)\right| & \leqslant \frac{k}{1+k}\left\|p(t, \cdot)-p^{\prime}(t, \cdot)\right\|_{\infty}  \tag{2.5a}\\
\left\|\Lambda(w, m, p)(t, \cdot)-\Lambda\left(w, m, p^{\prime}\right)(t, \cdot)\right\|_{\infty} & \leqslant \frac{\kappa}{1+\kappa}\left\|p(t, \cdot)-p^{\prime}(t, \cdot)\right\|_{\infty}  \tag{2.5b}\\
\left|\Lambda(w, m, p)(t, x)-\Lambda\left(w^{\prime}, m, p\right)(t, x)\right| & \leqslant \frac{1}{2}\left|w(t, x)-w^{\prime}(t, x)\right|  \tag{2.5c}\\
\left|\Lambda(w, m, p)(t, x)-\Lambda\left(w, m^{\prime}, p\right)(t, x)\right| & \leqslant\left|\ell(m, p)(t)-\ell\left(m^{\prime}, p\right)(t)\right| \tag{2.5d}
\end{align*}
$$

We aim to use Banach fixed point Theorem to show that

$$
\begin{equation*}
p=a(m)+c(m) \bar{p}-\frac{1}{2}\left\{a(m)+c(m) \bar{p}-u_{x}\right\}^{+} \tag{2.6}
\end{equation*}
$$

has a unique solution $p_{u, m} \in \mathbb{X}$, which satisfies (2.4). For any $p \in \mathbb{X}$, let us set

$$
\psi(\mathfrak{p}):=\Lambda\left(u_{x}, \mathfrak{m}, p\right)=a(m)+c(m) \bar{p}-\frac{1}{2}\left\{a(m)+c(m) \bar{p}-u_{x}\right\}^{+}
$$

Observe that $\psi(\mathbb{X}) \subset \mathbb{X}$, and $p \leqslant 1$ entails $\psi(p) \leqslant 1$. Moreover, if we suppose that $p \geqslant-\left\|u_{x}\right\|_{\infty}$, then it holds that

$$
\psi(p) \geqslant-c(m)\left\|u_{x}\right\|_{\infty}
$$

so that $\psi(p) \geqslant-\left\|u_{x}\right\|_{\infty}$, since $c(m)<1$. On the other hand, by appealing to (2.5b) we have

$$
\left\|\psi\left(p_{1}\right)-\psi\left(p_{2}\right)\right\|_{\mathbb{X}} \leqslant \frac{\kappa}{1+\kappa}\left\|p_{1}-p_{2}\right\|_{\mathbb{X}} \quad \forall p_{1}, p_{2} \in \mathbb{X} .
$$

Therefore by invoking Banach fixed point Theorem, and the estimates above we deduce the existence of a unique solution $p_{u, m} \in \mathbb{X}$ to problem (2.6) satisfying (2.4).

When $u_{x}$ is non-negative, note that $p \geqslant 0$ entails $\psi(p) \geqslant 0$, so that the same fixed point argument yields $p_{u, m} \geqslant 0$.

Next, we collect some facts related to the Fokker-Planck equation (1.2a)-(1.2b).
Lemma 2.2 (regularity of $\mathfrak{\eta}$ ). Let $m$ be a weak solution to (1.2a)-(1.2b), starting from some $\mathrm{m}_{0}$ satisfying ( $\mathcal{H} 2$ ). Suppose that b is bounded, and satisfies (1.1b). Then the map $\mathrm{t} \rightarrow \mathfrak{\eta}(\mathrm{t}):=$ $\eta(m(t))$ is continuous on $[0, \mathrm{~T}]$.

Moreover, if in addition $\mathrm{m}_{0}$ belongs to $\mathrm{L}^{1}(\mathrm{Q})$, then we have:
(i) the function $\mathrm{t} \rightarrow \boldsymbol{\eta}(\mathrm{t})$ is locally Hölder continuous on ( $0, \mathrm{~T}]$; namely, there exists $\gamma>0$ such that

$$
\begin{equation*}
\left|\boldsymbol{\eta}\left(\mathrm{t}_{1}\right)-\mathfrak{\eta}\left(\mathrm{t}_{2}\right)\right| \leqslant \mathrm{C}\left(\mathrm{t}_{0},\|\mathrm{~b}\|_{\infty}\right)\left|\mathrm{t}_{1}-\mathrm{t}_{2}\right|^{\gamma} \quad \forall \mathrm{t}_{1}, \mathrm{t}_{2} \in\left[\mathrm{t}_{0}, \mathrm{~T}\right] \tag{2.7a}
\end{equation*}
$$

for all $\mathrm{t}_{0} \in(0, \mathrm{~T})$;
(ii) for any $\alpha>0$ and $\phi \in \mathcal{C}^{\alpha}(\overline{\mathrm{Q}})$, there exists $\beta>0$ such that

$$
\begin{equation*}
\left|\int_{0}^{\mathrm{L}} \phi(x)\left(\mathfrak{m}\left(\mathrm{t}_{1}, x\right)-\mathfrak{m}\left(\mathrm{t}_{2}, x\right)\right) \mathrm{dx}\right| \leqslant \mathrm{C}\left(\mathrm{t}_{0},\|\mathrm{~b}\|_{\infty},\|\phi\|_{\mathrm{e}^{\alpha}}\right)\left|\mathrm{t}_{1}-\mathrm{t}_{2}\right|^{\beta} \quad \forall \mathrm{t}_{1}, \mathrm{t}_{2} \in\left[\mathrm{t}_{0}, \mathrm{~T}\right] \tag{2.7b}
\end{equation*}
$$

$$
\text { for all } \mathrm{t}_{0} \in(0, \mathrm{~T})
$$

Remark 2.3. This lemma shows that $\mathrm{t} \mapsto \mathrm{m}(\mathrm{t})$ is locally Hölder continuous in time in $(0, T]$ with respect to the $\left(\mathcal{C}^{\alpha}\right)^{*}$ topology; this is useful later to get equicontinuity for construction of a fixed point (cf. Section 2.3). Our method of proof does not allow us to show Hölder continuity on all of $[0, \mathrm{~T}]$, because it is based on heat kernel estimates, which degenerate as $t \rightarrow 0$ (cf. Equation (2.16)). However, we find it difficult to construct a counterexample.

Proof. The proof requires several steps and lies on the probabilistic interpretation of $m$ which we recall briefly here, and use in other parts of this paper.

Step 1 (probabilistic interpretation): Consider the reflected diffusion process governed by

$$
\begin{equation*}
\mathrm{d} X_{\mathrm{t}}=-\mathrm{b}\left(\mathrm{t}, X_{\mathrm{t}}\right) \mathrm{dt}+\sigma \mathrm{d} W_{\mathrm{t}}-\mathrm{d} \xi_{\mathrm{t}}^{\mathrm{x}}, \quad X_{0} \sim \mathrm{~m}_{0}, \tag{2.8a}
\end{equation*}
$$

where $X_{0}$ is $\mathcal{F}_{0}$-measurable, $\left(W_{t}\right)_{t \in[0, T]}$ is a $\mathbb{F}$-Wiener process that is independent of $X_{0}$, and set

$$
\begin{equation*}
\tau:=\inf \left\{t \geqslant 0: X_{t} \leqslant 0\right\} \wedge T . \tag{2.8b}
\end{equation*}
$$

By virtue of the regularity assumptions on $b$, equation (2.8a) is well-posed in the classical sense. Furthermore, since the process $\left(\xi_{t}^{X}\right)_{t \geqslant 0}$ is monotone, $\left(X_{t}\right)_{t \in[0, T]}$ is a continuous
semimartingale. Hence, by means of Itô's rule and the optional stopping theorem, we have for any test function $\phi \in \mathcal{C}_{c}^{\infty}([0, T) \times \overline{\mathrm{Q}})$ satisfying (1.2d):

$$
\mathbb{E}\left[\phi\left(0, X_{0}\right)\right]=\mathbb{E}\left[\int_{0}^{\tau}\left(-\phi_{\mathfrak{t}}\left(v, X_{v}\right)-\frac{\sigma^{2}}{2} \phi_{\chi x}\left(v, X_{v}\right)+\phi_{x}\left(v, X_{v}\right) \mathfrak{b}\left(v, X_{v}\right)\right) \mathrm{d} v\right],
$$

and thus the law of $X_{t}$ is a weak solution to the Fokker-Planck equation. The function $b$ being bounded, one sees that

$$
\mathbb{E}\left[\int_{0}^{T} \mathrm{~b}\left(s, X_{s}\right)^{2} \mathrm{ds}\right]<\infty .
$$

Therefore, by virtue of the uniqueness for (1.2a)-(1.2b) (cf. Proposition B.1), we obtain:

$$
\begin{equation*}
\int_{\mathcal{A}} m(t, x) d x=\mathbb{P}\left(t<\tau ; X_{t} \in A\right) \tag{2.8c}
\end{equation*}
$$

for every Borel set $A \in \bar{Q}$ and for a.e. $t \in(0, T)$.
Step 2: Now, let us show that $\mathrm{t} \rightarrow \mathbb{P}(\mathrm{t}<\tau)$ is right continuous on [0, T$]$. In fact, we have for any $\epsilon>0$ and $t \in[0, T]$

$$
\begin{align*}
\mathbb{P}(\mathrm{t}<\tau)-\mathbb{P}(\mathrm{t}+\mathrm{h}<\tau) & =\mathbb{P}(\mathrm{t}+\mathrm{h} \geqslant \tau ; \mathrm{t}<\tau)  \tag{2.9a}\\
& \leqslant \mathbb{P}\left(\mathrm{t}+\mathrm{h} \geqslant \tau ; X_{\mathrm{t}} \geqslant \epsilon\right)+\mathbb{P}\left(\mathrm{t}<\tau ; X_{\mathrm{t}}<\epsilon\right) .
\end{align*}
$$

On the one hand, for every $t \in[0, T]$

$$
\begin{equation*}
\lim _{\epsilon \rightarrow 0^{+}} \mathbb{P}\left(\mathrm{t}<\tau ; X_{\mathrm{t}}<\epsilon\right) \leqslant \lim _{\epsilon \rightarrow 0^{+}} \mathbb{P}\left(0<X_{\mathrm{t}}<\epsilon\right)=0 \tag{2.9b}
\end{equation*}
$$

thanks to the bounded convergence theorem. On the other hand

$$
\begin{aligned}
& \mathbb{P}\left(\mathrm{t}+\mathrm{h} \geqslant \tau ; X_{\mathrm{t}} \geqslant \epsilon\right) \leqslant \mathbb{P}\left(\inf _{v \in[t, t+h]} X_{v}-X_{\mathrm{t}} \leqslant-\epsilon\right) \\
& \leqslant \mathbb{P}\left(\inf _{v \in[0, h]} \sigma\left(W_{t+v}-W_{t}\right)+\left(\xi_{\mathrm{t}}^{\mathrm{X}}-\xi_{\mathrm{t}+\mathrm{h}}^{\mathrm{X}}\right) \leqslant-\epsilon+\mathrm{h}\|\mathrm{~b}\|_{\infty}\right),
\end{aligned}
$$

where we have used the fact that the local time is nondecreasing and $b$ is bounded. Furthermore, by using (1.1e), it holds that

$$
\xi_{t}^{X}-\xi_{t+h}^{X} \geqslant \inf _{v \in[0, h]}\left(Y_{t}-Y_{t+v}\right) \geqslant \sigma \inf _{v \in[0, h]}\left(W_{t}-W_{t+v}\right)-h\|b\|_{\infty} .
$$

Therefore

$$
\mathbb{P}\left(t+h \geqslant \tau ; X_{t} \geqslant \epsilon\right) \leqslant \mathbb{P}\left(\sup _{v \in[0, h]} B_{v}-\inf _{v \in[0, h]} B_{v} \geqslant \frac{\epsilon-2 h\|b\|_{\infty}}{\sigma}\right),
$$

where $\left(B_{t}\right)_{t \geqslant 0}$ is a Wiener process.
Now, choose $\epsilon=\epsilon(\mathrm{h}):=\mathrm{h}^{1 / 2} \log (1 / \mathrm{h})$. We have $\epsilon(\mathrm{h}) \rightarrow 0$ as $\mathrm{h} \rightarrow 0^{+}$, and by using Markov's inequality and the distribution of the maximum of Brownian motion we get:

$$
\begin{equation*}
\mathbb{P}\left(t+h \geqslant \tau ; X_{t} \geqslant \epsilon\right) \leqslant \frac{2 \sigma}{\epsilon(h)-2 h\|b\|_{\infty}} \mathbb{E}\left|B_{h}\right| \leqslant \frac{2 \sigma}{\log (1 / h)-2\|b\|_{\infty} h^{1 / 2}} . \tag{2.9c}
\end{equation*}
$$

Thus $0 \leqslant \mathbb{P}(\mathrm{t}<\tau)-\mathbb{P}(\mathrm{t}+\mathrm{h}<\tau) \rightarrow 0$ as $\mathrm{h} \rightarrow 0^{+}$.

Step 3 (Hölder estimates): Now, we prove (2.7a)-(2.7b). At first, note that (2.8c) entails

$$
\begin{equation*}
\int_{0}^{\mathrm{L}} \phi(x) \mathfrak{m}(\mathrm{t}, \mathrm{x}) \mathrm{d} x=\mathbb{E}\left[\phi\left(\mathrm{X}_{\mathrm{t}}\right) \mathbb{1}_{\mathrm{t}<\tau}\right] \tag{2.10}
\end{equation*}
$$

for a.e. $t \in(0, T)$ and for any $\phi \in \mathcal{C}(\bar{Q})$. Actually (2.10) holds for every $t \in[0, T]$, since the RHS and LHS of (2.10) are both right continuous on $[0, \mathrm{~T}]$, and $m_{0}$ is supported on ( $0, \mathrm{~L}$ ]. Indeed, on the one hand $\mathrm{t} \rightarrow \int_{0}^{\mathrm{L}} \phi(\mathrm{x}) \mathfrak{m}(\mathrm{t}, \mathrm{x}) \mathrm{dx}$ is continuous on $[0, \mathrm{~T}]$ for any continuous function $\phi$ on $\bar{Q}$, since $m \in \mathcal{C}\left([0, T] ; L^{1}(Q)\right)$ (cf. [35, Theorem 3.6]). On the other hand, for any $\phi \in \mathcal{C}(\bar{Q})$

$$
\begin{align*}
\mathbb{E} \mid \phi\left(\mathrm{X}_{\mathrm{t}+\mathrm{h}}\right) \mathbb{1}_{\mathrm{t}+\mathrm{h}<\tau}- & \phi\left(\mathrm{X}_{\mathrm{t}}\right) \mathbb{1}_{\mathrm{t}<\tau} \mid  \tag{2.11}\\
& \leqslant\|\phi\|_{\infty}(\mathbb{P}(\mathrm{t}<\tau)-\mathbb{P}(\mathrm{t}+\mathrm{h}<\tau))+\mathbb{E}\left|\phi\left(\mathrm{X}_{\mathrm{t}+\mathrm{h}}\right)-\phi\left(\mathrm{X}_{\mathrm{t}}\right)\right|,
\end{align*}
$$

so that

$$
\lim _{h \rightarrow 0^{+}} \mathbb{E}\left|\phi\left(X_{t+h}\right) \mathbb{1}_{t+h<\tau}-\phi_{\epsilon}\left(X_{t}\right) \mathbb{1}_{t<\tau}\right|=0
$$

thanks to (2.9a)-(2.9c), and the bounded convergence theorem.
Now, let us fix $\epsilon>0$ and define $\phi_{\epsilon}=\phi_{\epsilon}(x)$ to be a smooth cut-off function on $[0, \mathrm{~L}]$, which satisfies the following conditions:

$$
\begin{equation*}
0 \leqslant \phi_{\epsilon} \leqslant 1 ; 0 \leqslant \phi_{\epsilon}^{\prime} \leqslant 2 / \epsilon ; \quad \phi_{\epsilon} \mathbb{1}_{[0, \epsilon]}=0 ; \quad \phi_{\epsilon} \mathbb{1}_{[2 \epsilon, \mathrm{~L}]}=1 . \tag{2.12}
\end{equation*}
$$

As a first step, we aim to derive an estimation of the concentration of mass at the origine. Namely, we want to show that for an arbitrary $k>1$,

$$
\begin{equation*}
\int_{0}^{\mathrm{L}}\left(1-\phi_{\epsilon}(x)\right) \mathrm{m}(\mathrm{t}, \mathrm{x}) \mathrm{d} x \leqslant \mathrm{C}\left(\mathrm{k},\|\mathrm{~b}\|_{\infty}\right)\left(1-e^{-\pi^{2} \mathrm{t} / 4 \mathrm{~L}^{2}}\right)^{-1 / 2 \mathrm{k}} \epsilon^{1 / 2 \mathrm{k}} \quad \forall \mathrm{t} \in(0, \mathrm{~T}] . \tag{2.13}
\end{equation*}
$$

Given (2.10), this is equivalent to showing that

$$
\begin{equation*}
\mathbb{E}\left[\left(1-\phi_{\epsilon}\left(X_{t}\right)\right) \mathbb{1}_{\mathrm{t}<\tau}\right] \leqslant C\left(k,\|\mathrm{~b}\|_{\infty}\right)\left(1-e^{-\pi^{2} \mathrm{t} / 4 \mathrm{~L}^{2}}\right)^{-1 / 2 \mathrm{k}} \epsilon^{1 / 2 \mathrm{k}} \quad \forall \mathrm{t} \in(0, \mathrm{~T}] \tag{2.14}
\end{equation*}
$$

holds for any $k>1$. Apply Girsanov's Theorem with the following change of measure:

$$
\left.\frac{\mathrm{d} \mathbb{Q}}{\mathrm{~d} \mathbb{P}}\right|_{\mathcal{F}_{\mathrm{t}}}=\exp \left\{-\sigma^{-1} \int_{0}^{\mathrm{t}} \mathrm{~b}\left(\mathrm{~s}, \mathrm{X}_{\mathrm{s}}\right) \mathrm{d} W_{s}-\frac{\sigma^{-2}}{2} \int_{0}^{\mathrm{t}} \mathrm{~b}\left(\mathrm{~s}, \mathrm{X}_{\mathrm{s}}\right)^{2} \mathrm{~d} s\right\}=: \Psi_{\mathrm{t}} .
$$

Under $\mathbb{Q}$, the process $\left(X_{t}\right)_{t \in[0, T]}$ is a reflected Brownian motion at $L$, with initial condition $X_{0}$, thanks to (1.1c). Moreover, by virtue of Hölder inequality, we have for every $k>1$ :

$$
\begin{aligned}
& \mathbb{E}_{\mathbb{P}}\left[\left(1-\phi_{\epsilon}\left(X_{t}\right)\right) \mathbb{1}_{\mathrm{t}<\tau}\right] \\
& =\mathbb{E}_{\mathbb{Q}}\left[\Psi_{\mathrm{t}}^{-1}\left(1-\phi_{\epsilon}\left(X_{\mathrm{t}}\right)\right) \mathbb{1}_{\mathrm{t}<\tau}\right] \leqslant \mathbb{E}_{\mathbb{Q}}\left[\Psi_{\mathrm{t}}^{-\mathrm{k}^{\prime}}\right]^{1 / k^{\prime}} \mathbb{E}_{\mathbb{Q}}\left[\left(1-\phi_{\epsilon}\left(X_{\mathrm{t}}\right)\right)^{\mathrm{k}} \mathbb{1}_{\mathrm{t}<\tau}\right]^{1 / \mathrm{k}} \\
& \quad \leqslant \mathbb{E}_{\mathbb{P}}\left[\Psi_{\mathrm{t}}^{\left.1-\mathrm{k}^{\prime}\right]^{1 / k^{\prime}} \mathbb{E}_{\mathbb{Q}}\left[\left(1-\phi_{\epsilon}\left(X_{\mathrm{t}}\right)\right)^{\mathrm{k}} \mathbb{1}_{\mathrm{t}<\tau}\right]^{1 / k}}\right. \\
& \quad \leqslant \mathbb{E}_{\mathbb{P}}\left[\exp \left\{\mathrm{C}(\mathrm{k}, \sigma) \int_{0}^{\mathrm{t}} \mathrm{~b}\left(\mathrm{~s}, X_{s}\right)^{2} \mathrm{ds}\right\}\right]^{1 / 2 \mathrm{k}^{\prime}} \mathbb{E}_{\mathbb{Q}}\left[\left(1-\phi_{\epsilon}\left(X_{\mathrm{t}}\right)\right)^{\mathrm{k}} \mathbb{1}_{\mathrm{t}<\tau}\right]^{1 / \mathrm{k}} .
\end{aligned}
$$

Indeed, one checks that

$$
\mathbb{E}_{\mathbb{P}}\left[\Psi_{\mathrm{t}}^{1-\mathrm{k}^{\prime}}\right]^{1 / \mathrm{k}^{\prime}} \leqslant \mathbb{E}_{\mathbb{P}}\left[Z_{\mathrm{t}}\right]^{1 / 2 \mathrm{k}^{\prime}} \mathbb{E}_{\mathbb{P}}\left[\exp \left\{2\left(\frac{1-\mathrm{k}^{\prime}}{\sigma}\right)^{2} \int_{0}^{\mathrm{t}} \mathrm{~b}\left(\mathrm{~s}, \mathrm{X}_{\mathrm{s}}\right)^{2} \mathrm{ds}\right\}\right]^{1 / 2 \mathrm{k}^{\prime}}
$$

where $\left(Z_{t}\right)_{t \geqslant 0}$ is a super-martingale. Hence, using the fact that $b$ is bounded, we obtain

$$
\begin{equation*}
\mathbb{E}_{\mathbb{P}}\left[\left(1-\phi_{\epsilon}\left(X_{t}\right)\right) \mathbb{1}_{t<\tau}\right] \leqslant C\left(k,\|\mathfrak{b}\|_{\infty}\right) \mathbb{E}_{\mathbb{Q}}\left[\left(1-\phi_{\epsilon}\left(X_{t}\right)\right)^{k} \mathbb{1}_{t<\tau}\right]^{1 / k} \tag{2.15}
\end{equation*}
$$

Now

$$
\mathbb{E}_{\mathbb{Q}}\left[\left(1-\phi_{\epsilon}\left(X_{t}\right)\right)^{\mathrm{k}} \mathbb{1}_{\mathrm{t}<\tau}\right]=\int_{0}^{\mathrm{L}}\left(1-\phi_{\epsilon}(x)\right)^{\mathrm{k}} w(\mathrm{t}, \mathrm{x}) \mathrm{d} x
$$

where $w$ solves

$$
w_{\mathrm{t}}=\frac{1}{2} w_{x x}, w(\mathrm{t}, 0)=0, w_{x}(\mathrm{t}, \mathrm{~L})=0,\left.w\right|_{\mathrm{t}=0}=\mathrm{m}_{0}
$$

We can compute $w$ via Fourier series, namely

$$
w(t, x)=\sum_{n \geqslant 1} A_{n} e^{-\lambda_{n}^{2} t / 2} \sin \left(\lambda_{n} x\right), \quad A_{n}:=\frac{2}{L} \int_{0}^{L} \sin \left(\lambda_{n} y\right) d m_{0}(y), \quad \lambda_{n}:=\frac{(2 n-1) \pi}{2 L} .
$$

Note that

$$
\begin{aligned}
\int_{0}^{\mathrm{L}}\left(1-\phi_{\epsilon}(x)\right)^{\mathrm{k}} w(\mathrm{t}, \mathrm{x}) \mathrm{d} x & \leqslant(2 \epsilon)^{1 / 2}\|w(\mathrm{t}, \cdot)\|_{\mathrm{L}^{2}} \\
& \leqslant \epsilon^{1 / 2}\left(\sum_{n \geqslant 1} \mathrm{~L}\left|A_{n}\right|^{2} e^{-\lambda_{n}^{2} \mathrm{t}}\right)^{1 / 2} \text { (Parseval) } \\
& \leqslant \epsilon^{1 / 2}\left(\frac{4}{\mathrm{~L}\left(1-e^{-\pi^{2} \mathrm{t} / 4 \mathrm{~L}^{2}}\right)}\right)^{1 / 2}
\end{aligned}
$$

So (2.15) now yields

$$
\begin{equation*}
\mathbb{E}_{\mathbb{P}}\left[\left(1-\phi_{\epsilon}\left(X_{t}\right)\right) \mathbb{1}_{\mathrm{t}<\tau}\right] \leqslant C\left(k,\|b\|_{\infty}\right)\left(\frac{4}{\mathrm{~L}\left(1-e^{-\pi^{2} \mathrm{t} / 4 \mathrm{~L}^{2}}\right)}\right)^{1 / 2 \mathrm{k}} \epsilon^{1 / 2 \mathrm{k}} \tag{2.16}
\end{equation*}
$$

which is (2.14). This in turn implies (2.13).
Furthermore, note that for any $1<s<3 / 2$,

$$
\begin{equation*}
\left\|\mathfrak{m}\left(\mathrm{t}_{1}\right)-\mathfrak{m}\left(\mathrm{t}_{2}\right)\right\|_{W_{s}^{-1}} \leqslant \mathrm{C}\left(\left\|\mathrm{~m}_{0}\right\|_{\mathrm{L}^{1}},\left\|\mathfrak{m}|\mathbf{b}|^{2}\right\|_{\mathrm{L}^{1}}\right)\left|\mathrm{t}_{1}-\mathrm{t}_{2}\right|^{1-1 / s} \quad \forall \mathrm{t}_{1}, \mathrm{t}_{2} \in[0, \mathrm{~T}], \tag{2.17}
\end{equation*}
$$

where $W_{s}^{-1}(\mathrm{Q})$ is the dual space of $\mathrm{W}_{0}^{1, s^{\prime}}(\mathrm{Q}):=\left\{v \in \mathrm{~W}_{s^{\prime}}^{1}(\mathrm{Q}): v(0)=0\right\}$. This claim follows from [35, Proposition 3.10(iii)], where we obtain the estimate

$$
\begin{align*}
&\|\mathfrak{m}\|_{L^{\infty}\left(0, T ; \mathrm{L}^{1}(\mathrm{Q})\right)}+\|\nabla \mathfrak{m}\|_{\mathrm{L}^{s}\left(\mathrm{Q}_{\mathrm{T}}\right)}+\|\mathfrak{m}\|_{\mathrm{L}^{\nu}\left(\mathrm{Q}_{\mathrm{T}}\right)}+\left\|\mathfrak{m}_{\mathrm{t}}\right\|_{\mathrm{L}^{s}\left(0, \mathrm{~T} ; W_{s}^{-1}(\mathrm{Q})\right)}  \tag{2.18}\\
& \leqslant \mathrm{C}\left(\left\|\mathrm{~m}_{0}\right\|_{\mathrm{L}^{1}},\left\|\mathfrak{m}|\mathrm{~b}|^{2}\right\|_{\mathrm{L}^{1}}\right)
\end{align*}
$$

for any s up to $3 / 2$ and $v$ up to 3 . In particular, (2.17) follows from the estimate on $\left\|m_{\mathfrak{t}}\right\|_{\mathrm{L}^{s}\left(0, \mathrm{~T} ; W_{s}^{-1}(\mathrm{Q})\right)}$. Now, fix $0<\mathrm{t}_{1} \leqslant \mathrm{t}_{2} \leqslant \mathrm{~T}$, and let $\phi_{\epsilon}$ be the cut-off function that is defined in (2.12). Based on the specifications of (2.12), observe that

$$
\left\|\phi_{\epsilon}\right\|_{W_{s^{\prime}}^{1}} \leqslant \mathrm{C}^{-1 / \mathrm{s}} .
$$

Since $\phi_{\epsilon}$ satisfies Neumann boundary conditions at $x=L$ and Dirichlet at $x=0$, it is a valid test function and we can appeal to the estimates above to obtain for any $k>1$,

$$
\begin{align*}
& \left|\mathfrak{\eta}\left(\mathrm{t}_{1}\right)-\mathfrak{\eta}\left(\mathrm{t}_{2}\right)\right|=\left|\int_{0}^{\mathrm{L}}\left\{\left(1-\phi_{\epsilon}(x)\right)+\phi_{\epsilon}(x)\right\}\left(\mathfrak{m}\left(\mathrm{t}_{1}, x\right)-\mathfrak{m}\left(\mathrm{t}_{2}, x\right)\right) \mathrm{dx}\right|  \tag{2.19}\\
& \leqslant \int_{0}^{\mathrm{L}}\left|1-\phi_{\epsilon}(x) \| \mathfrak{m}\left(\mathrm{t}_{1}, x\right)-\mathfrak{m}\left(\mathrm{t}_{2}, x\right)\right| \mathrm{d} x+\left|\int_{0}^{\mathrm{L}} \phi_{\epsilon}(x)\left(\mathfrak{m}\left(\mathrm{t}_{1}, x\right)-\mathfrak{m}\left(\mathrm{t}_{2}, x\right)\right) \mathrm{dx}\right| \\
& \leqslant \int_{0}^{\mathrm{L}}\left(1-\phi_{\epsilon}(x)\right)\left(\mathfrak{m}\left(\mathrm{t}_{1}, x\right)+\mathfrak{m}\left(\mathrm{t}_{2}, x\right)\right) \mathrm{d} x+\left|\int_{0}^{\mathrm{L}} \phi_{\epsilon}(x)\left(\mathfrak{m}\left(\mathrm{t}_{1}, x\right)-\mathfrak{m}\left(\mathrm{t}_{2}, x\right)\right) \mathrm{dx}\right| \\
& \leqslant C\left(k,\|b\|_{\infty}\right)\left(1-e^{-\pi^{2} t_{1} / 4 L^{2}}\right)^{-1 / 2 k} \epsilon^{1 / 2 k}+\left\|\phi_{\epsilon}\right\|_{W_{s^{\prime}}^{1}}\left\|m\left(t_{1}\right)-\mathfrak{m}\left(t_{2}\right)\right\|_{W_{s}^{-1}} \\
& \leqslant C\left(k,\|b\|_{\infty}\right)\left(1-e^{-\pi^{2} t_{1} / 4 L^{2}}\right)^{-1 / 2 k} \epsilon^{1 / 2 k}+C \epsilon^{-1 / s}\left|t_{1}-t_{2}\right|^{1-1 / s},
\end{align*}
$$

where we have used (2.13) in the penultimate line and (2.17) in the ultimate line. Given $0<\gamma<(s-1) /(s+2)$, we take $\epsilon=\left|\mathrm{t}_{1}-\mathrm{t}_{2}\right|^{s(1-\gamma)-1}$ and then set $k=\frac{s(1-\gamma)-1}{2 \gamma}>1$ to obtain (2.7a).

Finally, let $\phi \in \mathcal{C}^{\alpha}(\bar{Q})$ for some $\alpha>0$, an let $t_{0} \in(0, T)$. In view of (2.10), we have for every $\mathrm{t}_{1}, \mathrm{t}_{2} \in\left[\mathrm{t}_{0}, \mathrm{~T}\right]$,

$$
\begin{aligned}
&\left|\int_{0}^{\mathrm{L}} \phi(x)\left(m\left(\mathrm{t}_{1}, x\right)-m\left(\mathrm{t}_{2}, x\right)\right) \mathrm{d} x\right| \leqslant \mathbb{E}\left|\phi\left(X_{\mathrm{t}_{1}}\right) \mathbb{1}_{\mathrm{t}_{1}<\tau}-\phi\left(\mathrm{X}_{\mathrm{t}_{2}}\right) \mathbb{1}_{\mathrm{t}_{2}<\tau}\right| \\
& \leqslant\|\phi\|_{\mathrm{e}^{\alpha}}\left(\left|\mathfrak{\eta}\left(\mathrm{t}_{1}\right)-\mathfrak{\eta}\left(\mathrm{t}_{2}\right)\right|+\mathbb{E}\left|X_{\mathrm{t}_{1}}-\mathrm{X}_{\mathrm{t}_{2}}\right|^{\alpha}\right) .
\end{aligned}
$$

Hence, by using (2.7a) and the Burkholder-Davis-Gundy inequality [37, Thm IV.42.1], we deduce the desired result:

$$
\left|\int_{0}^{\mathrm{L}} \phi(x)\left(\mathfrak{m}\left(\mathrm{t}_{1}, x\right)-\mathfrak{m}\left(\mathrm{t}_{2}, x\right)\right) \mathrm{dx}\right| \leqslant C\left(\mathrm{t}_{0},\|\mathrm{~b}\|_{\infty}\right)\|\phi\| \mathrm{e}^{\alpha}\left|\mathrm{t}_{1}-\mathrm{t}_{2}\right|^{\beta}
$$

for some $\beta>0$.
Step 4 (general data): Now, we suppose that $\mathrm{m}_{0}$ is a probability measure satisfying ( $\left.\mathcal{H} 2\right)$, and not necessarily an element of $L^{1}(Q)$. Let us choose a sequence $\left(m_{0}^{n}\right) \subset L^{1}(Q)_{+}$, which converges weakly (in the sense of measures) to $m_{0}$, such that

$$
\begin{equation*}
\left\|m_{0}^{n}\right\|_{L^{1}} \leqslant \int_{0}^{\mathrm{L}} \mathrm{dm}_{0} \leqslant 1 \tag{2.20}
\end{equation*}
$$

and let $\mathfrak{m}^{n}$ to be the weak solution to (1.2a)-(1.2b) starting from $m_{0}^{n}$. The function $b$ being bounded, we can use [35, Proposition 3.10] to extract a subsequence of $\left(m^{n}\right)$, which converges to $m$ in $L^{1}\left(Q_{T}\right)$. Owing to (2.7a), the sequence $\eta^{n}:=\eta\left(m^{n}\right)$ is equicontinuous. Hence, one can extract further a subsequence to deduce that $\eta$ is continuous on $(0, T]$. Combining this conclusion with the fact that $t \rightarrow \mathbb{P}(t<\tau)$ is right continuous on $[0, \mathrm{~T}]$ and (2.8c), we deduce in particular that

$$
\begin{equation*}
\eta(\mathrm{t})=\mathbb{P}(\mathrm{t}<\tau), \quad \forall \mathrm{t} \in(0, \mathrm{~T}] . \tag{2.21}
\end{equation*}
$$

Now, since $m_{0}$ is supported on $(0, L]$ one has $\eta\left(m_{0}\right)=\eta(0)=\mathbb{P}(0<\tau)=1$, which in turn entails that $\eta$ is continuous on $[0, \mathrm{~T}]$ thanks to (2.9a)-(2.9c) and (2.21). The proof is complete.

Remark 2.4. When $m_{0}$ satisfies ( $\left.\mathcal{H} 2\right)$ and does not necessarily belong to $L^{1}(Q)$, the probabilistic characterisation (2.10) still holds for every $t \in[0, T]$. In fact, using the same approximation techniques as in Lemma 2.2- Step 4, and appealing to (2.7b) and (2.8c), it holds that

$$
\int_{0}^{\mathrm{L}} \phi(x) \mathfrak{m}(\mathrm{t}, \mathrm{x}) \mathrm{d} x=\mathbb{E}\left[\phi\left(\mathrm{X}_{\mathrm{t}}\right) \mathbb{1}_{\mathrm{t}<\tau}\right]
$$

for every $t \in[0, T], \alpha>0$ and $\phi \in \mathcal{C}^{\alpha}(\overline{\mathrm{Q}})$. Thus, (2.10) ensues by using density arguments.
2.2. A priori estimates. Now, we collect several a priori estimates for system (1.6).

Lemma 2.5. Suppose that $(u, m)$ satisfies the system (1.6) such that $m \in L^{1}\left(Q_{T}\right)_{+}$, and $u$ belongs to $\mathrm{W}_{\mathrm{s}}^{1,2}\left(\mathrm{Q}_{\mathrm{T}}\right)$ for large enough $\mathrm{s}>1$. Then, we have:
(i) the maps $\mathfrak{u}$ and $\mathfrak{u}_{x}$ are non-negative; in particular

$$
\begin{equation*}
0 \leqslant \mathrm{q}_{\mathfrak{u}, \mathrm{m}} \leqslant 1 / 2 ; \tag{2.22}
\end{equation*}
$$

(ii) there exists $\theta>0$ and a constant $\mathrm{c}_{0}>0$ such that

$$
\begin{equation*}
\|u\|_{C^{\theta}\left(\overline{Q_{T}}\right)},\left\|u_{x}\right\|_{C^{\theta}\left(\overline{Q_{T}}\right)} \leqslant c_{0} \tag{2.23}
\end{equation*}
$$

where $\mathrm{c}_{0}$ depends only on T and data. In addition, we have

$$
\left\|u_{x x}\right\|_{C^{\theta}\left(Q^{\prime}\right)}, \leqslant c_{1}\left(Q^{\prime}, \theta\right) \forall Q^{\prime} \subset \subset(0, T) \times(0, L] ;
$$

If in addition $\mathrm{m}_{0}$ belongs to $\mathrm{L}^{1}(\mathrm{Q})$, then there exists a Hölder exponent $\theta>0$ such that

$$
\left\|p_{u, m}\right\|_{\mathrm{C}^{\theta}\left(\left[\mathrm{t}_{0}, \mathrm{~T}\right] \times[0, \mathrm{~L}]\right)} \leqslant \mathrm{c}_{2}\left(\mathrm{t}_{0}, \theta\right), \quad \forall \mathrm{t}_{0} \in(0, \mathrm{~T}),
$$

and

$$
\left\|u_{t}\right\|_{C^{\theta}\left(Q^{\prime}\right)} \leqslant c_{2}\left(Q^{\prime}, \theta\right) \forall Q^{\prime} \subset \subset(0, T) \times(0, L] .
$$

Proof. For large enough $s>1$, we know that $u, u_{x} \in \mathcal{C}\left(\overline{Q_{T}}\right)$ thanks to Sobolev-Hölder embeddings. In view of

$$
-u_{t}-\frac{\sigma^{2}}{2} u_{x x}+r u \geqslant 0
$$

one easily deduces that $u \geqslant e^{-r \top} \min _{x} u_{\top}$, which entails in particular that $u \geqslant 0$ thanks to $(\mathcal{H} 1)$. Thus, the minimum is attained at $u(t, 0)=0$, so that $u_{x}(t, 0) \geqslant 0$ for all $t \in[0, T]$. Differentiating the first equation in (1.6) we have that $\mathfrak{u}_{\mathrm{x}}$ is a generalised solution (cf. [26, Chapter III]) of the following parabolic equation:

$$
\mathbf{u}_{x t}+\frac{\sigma^{2}}{2} \mathbf{u}_{\mathrm{xxx}}-\mathbf{r} \mathbf{u}_{x}-\mathbf{q}_{\mathbf{u}, \mathfrak{m}} \mathbf{u}_{\mathrm{xx}}=0
$$

By virtue of the maximum principle [26, Theorem III.7.1] we infer that $u_{x} \geqslant 0$, since $u_{x}(t, 0), u_{x}(t, L)$ and $u_{\top}^{\prime}$ are all non-negative functions. Therefore (2.22) follows straightforwardly from (2.3c) thanks to Lemma 2.1.

Note that $u$ solves a parabolic equation with bounded coefficients. Since compatibility conditions of order zero are fulfilled thanks to ( $\mathcal{H} 1)$, then from [26, Theorem IV.9.1] we have an estimate on $u$ in $W_{k}^{1,2}\left(Q_{T}\right)$ for arbitrary $k>1$, namely

$$
\begin{equation*}
\|\mathfrak{u}\|_{W_{k}^{1,2}\left(Q_{T}\right)} \leqslant C\left(\left\|q_{u, m}\right\|_{L^{k}\left(Q_{T}\right)}+\left\|u_{T}\right\|_{W_{k}^{2-\frac{2}{k}}\left(Q_{T}\right)}\right) \leqslant C\left(\left\|q_{u, m}\right\|_{L^{\infty}\left(Q_{T}\right)}+\left\|u_{T}\right\|_{W_{k}^{2-\frac{2}{k}}\left(Q_{T}\right)}\right) . \tag{2.24}
\end{equation*}
$$

This estimate depends only on $\mathrm{T}, \mathrm{k}$ and data, thanks to (2.22). We deduce (2.23) thanks to Sobolev-Hölder embeddings.

Now, let $\phi \in \mathcal{C}_{c}^{\infty}((0, T) \times(0,+\infty))$. Observe that $w=\phi u_{x}$ satisfies

$$
w_{\mathrm{t}}+\frac{\sigma^{2}}{2} w_{x x}-\mathrm{rw}-\mathrm{q}_{\mathfrak{u}, \mathrm{m}} w_{x}=\phi_{\mathrm{t}} \mathrm{u}_{\mathrm{x}}+\sigma^{2} \phi_{\chi} \mathrm{u}_{x x}+\frac{\sigma^{2}}{2} \phi_{\chi x} \mathrm{u}_{x}-\mathrm{q}_{\mathrm{u}, \mathrm{~m}} \phi_{\chi} \mathrm{u}_{x}
$$

For any $k>1$, the right-hand side is bounded in $L^{k}\left(Q_{T}\right)$ with a constant that depends only on $\phi$, and previous estimates. Since $w$ has homogeneous boundary conditions, we deduce from [26, Theorem IV.9.1] that $\left\|w_{x}\right\|_{C^{\theta}\left(\overline{Q_{\bar{T}}}\right)}$ is bounded by a constant depending only on the norm of $\phi$ and previous estimates. The local Hölder estimate on $\mathfrak{u}_{x x}$ then follows.

Let $p(t, x)=p_{u, m}(t, x)$. Recall that $p(t, x)=f\left(\ell(m, p)(t), u_{x}(t, x)\right)$ where $f(x, y):=$ $x-\frac{1}{2}(x-y)^{+}$(cf. Lemma 2.1). Since $f$ is 1 -Lipschitz in the first variable and $\frac{1}{2}$-Lipschitz in the second, we deduce that

$$
\begin{equation*}
\left|p\left(t_{1}, x_{1}\right)-p\left(t_{2}, x_{2}\right)\right| \leqslant\left|\ell(m, p)\left(t_{1}\right)-\ell(m, p)\left(t_{2}\right)\right|+\frac{1}{2}\left|u_{x}\left(t_{1}, x_{1}\right)-u_{x}\left(t_{2}, x_{2}\right)\right| . \tag{2.25}
\end{equation*}
$$

In particular, for each $t$,

$$
\begin{equation*}
\left|\mathfrak{p}\left(\mathrm{t}, \mathrm{x}_{1}\right)-\mathfrak{p}\left(\mathrm{t}, \mathrm{x}_{2}\right)\right| \leqslant \frac{1}{2}\left|\mathfrak{u}_{\mathrm{x}}\left(\mathrm{t}, \mathrm{x}_{1}\right)-\mathfrak{u}_{\mathrm{x}}\left(\mathrm{t}, \mathrm{x}_{2}\right)\right| \tag{2.26}
\end{equation*}
$$

which, by (2.23), implies that $p(t, \cdot)$ is Hölder continuous for every $t$.
Now, we further assume that $m_{0} \in L^{1}(Q)_{+}$to use (2.7a)-(2.7b). We shall use the following function which is introduced in Lemma 2.1:

$$
\ell(\mathfrak{m}, \mathfrak{p})(\mathrm{t})=\mathrm{a}(\mathfrak{m}(\mathrm{t}))+\mathrm{c}(\mathfrak{m}(\mathrm{t})) \overline{\mathrm{p}}(\mathrm{t}), \quad \text { where } \quad \overline{\mathfrak{p}}(\mathrm{t})=\frac{1}{\eta(\mathfrak{m}(\mathrm{t}))} \int_{0}^{\mathrm{L}} p(\mathrm{t}, \mathrm{x}) \mathfrak{m}(\mathrm{t}, \mathrm{x}) \mathrm{dx}
$$

Fix $t_{0} \in(0, T)$ and for $t_{1}, t_{2}$ in $\left[t_{0}, T\right]$ write
(2.27) $\quad \ell(m, p)\left(t_{1}\right)-\ell(m, p)\left(t_{2}\right)=a\left(m\left(t_{1}\right)\right)-a\left(m\left(t_{2}\right)\right)$

$$
\begin{aligned}
& +\kappa\left(a\left(\mathfrak{m}\left(t_{1}\right)\right)-a\left(\mathfrak{m}\left(t_{2}\right)\right)\right) \int_{0}^{L} p\left(t_{1}, .\right) d m\left(t_{1}\right) \\
& +\kappa \mathfrak{k a}\left(\mathfrak{m}\left(t_{2}\right)\right) \int_{0}^{L} p\left(t_{1}, .\right) d\left(\mathfrak{m}\left(t_{1}\right)-\mathfrak{m}\left(t_{2}\right)\right) \\
& +
\end{aligned}
$$

where we have used the fact that $c(m)=\kappa a(m) \eta(m)$. Observe that $\eta \rightarrow \frac{1}{1+\kappa \eta}$ is $\kappa-$ Lipschitz in the $\eta$ variable, and recall that $p\left(t_{1}, \cdot\right)$ is Hölder continuous. Moreover, by
virtue of (2.23) we know that $\mathrm{q}_{\mathrm{u}, \mathrm{m}}$ satisfies (1.1b). Therefore, using the upper bound on $a(m), c(m)$ and (2.7a)-(2.7b) we infer that

$$
\begin{equation*}
\left|\ell(m, p)\left(t_{1}\right)-\ell(m, p)\left(t_{2}\right)\right| \leqslant C\left|t_{1}-t_{2}\right|^{\beta}+\frac{\kappa}{1+\kappa}\left\|p\left(t_{1}, \cdot\right)-p\left(t_{2}, \cdot\right)\right\|_{\infty} \tag{2.28}
\end{equation*}
$$

Note that the constant in (2.28) depend only on $c_{0}$ and $\kappa$ thanks to (2.22), (2.23) and Lemma 2.1. Using now (2.28) in (2.25), and choosing $\theta$ small enough, we deduce

$$
\begin{equation*}
\frac{1}{1+\kappa}\left\|p\left(t_{1}, \cdot\right)-p\left(t_{2}, \cdot\right)\right\|_{\infty} \leqslant C\left|t_{1}-t_{2}\right|^{\beta}+\frac{1}{2}\left\|u_{x}\left(t_{1}, \cdot\right)-u_{x}\left(t_{2}, \cdot\right)\right\|_{\infty} \leqslant C\left|t_{1}-t_{2}\right|^{\theta} \tag{2.29}
\end{equation*}
$$

Putting together (2.26) and (2.29) we infer that $p$ has a Hölder estimate, whereupon by (2.28) so does $\ell(m, p)$. Thus $q_{u, m}$ also has a Hölder estimate, and so does $u_{t}$ by the HJB equation satisfied by $u$.
2.3. Well-posedness. We are now in position to prove the main result of this section.

Proof of Theorem 1.5. The proof requires several steps, the key arguments being precisely the estimates collected in Lemmas 2.1-2.5.

Step 1 (data in $\mathrm{L}^{1}$ ): We suppose that $\mathrm{m}_{0}$ is an element of $\mathrm{L}^{1}(\mathrm{Q})$ satisfying ( $\left.\mathcal{H} 2\right)$. Define $\mathbb{X}$ to be the space of couples $(v, v)$, such that $v$ and $v_{x}$ are globally continuous on $\overline{Q_{T}}$, and $v$ belongs to $\mathrm{L}^{1}\left(\mathrm{Q}_{\mathrm{T}}\right)_{+}$. The functional space $\mathbb{X}$ endowed with the norm:

$$
\|(v, v)\|_{\mathbb{X}}:=\|v\|_{\infty}+\left\|v_{\chi}\right\|_{\infty}+\|v\|_{\mathrm{L}^{1}}
$$

is a Banach space. Consider the map $\mathbb{T}:(\nu, \nu, \lambda) \in \mathbb{X} \times[0,1] \rightarrow(w, \mu)$ where $(w, \mu)$ are given by the following parametrized system of coupled partial differential equations:

$$
\left\{\begin{array}{l}
\text { (i) } w_{\mathrm{t}}+\frac{\sigma^{2}}{2} w_{x x}-\mathrm{rw}+\lambda^{2} \mathrm{q}_{v, v}^{2}=0 \quad \text { in } \mathrm{Q}_{\mathrm{T}}, \\
\text { (ii) } \quad \mu_{\mathrm{t}}-\frac{\sigma^{2}}{2} \mu_{x x}-\left\{\lambda \mathrm{q}_{v, v} \mu\right\}_{x}=0 \quad \text { in } \mathrm{Q}_{\mathrm{T}},  \tag{2.30}\\
\text { (iii) } \quad \mu(\mathrm{t}, 0)=0, \quad w(\mathrm{t}, 0)=0, \quad w_{x}(\mathrm{t}, \mathrm{~L})=0 \quad \text { in }[0, \mathrm{~T}], \\
\text { (iv) } \quad \mu(0)=\lambda m_{0}, \quad w(\mathrm{~T}, x)=\lambda u_{\mathrm{T}}(x) \quad \text { in }[0, \mathrm{~L}], \\
(v) \\
\frac{\sigma^{2}}{2} \mu_{x}+\lambda \mathrm{q}_{v, v} \mu=0 \quad \text { in }[0, \mathrm{~T}] \times\{\mathrm{L}\} .
\end{array}\right.
$$

By virtue of Lemma 2.1, the map $\mathrm{q}_{v, v}$ is well-defined for any $(v, v) \in \mathbb{X}$, and satisfies

$$
\begin{equation*}
\left|q_{v, v}\right| \leqslant C\left(1+\left\|v_{x}\right\|_{\infty}\right) . \tag{2.31}
\end{equation*}
$$

In view of [26, Theorem IV.9.1], the function $w$ exists and is bounded in $W_{s}^{1,2}\left(\mathrm{Q}_{\mathrm{T}}\right)$ for any $s>1$, by a constant which depends on $\left\|v_{x}\right\|_{\infty}$ and data. (Note that the required compatibility conditions hold owing to ( $\mathcal{H} 1$ ). Although [26, Theorem IV.9.1] is stated for Dirichlet boundary conditions, its proof is readily adapted to Neumann or mixed boundary conditions as in the present context; cf. the discussion in the first paragraph of [26, Section IV.9]). We deduce that

$$
\|w\|_{\mathrm{e}^{\alpha}}+\left\|w_{x}\right\|_{\mathrm{e}^{\alpha}} \leqslant \mathrm{C}\left(\mathrm{~T}, \mathrm{~L}, \mathbf{u}_{\mathrm{T}},\left\|v_{x}\right\|_{\infty}\right)
$$

for some $\alpha>0$. On the other hand, it is well known (see e.g. [26, Chapter III]) that for any $(v, v) \in \mathbb{X}$, equation (2.30)(ii) has a unique weak solution $\mu$. Therefore, $\mathbb{T}$ is well-defined.

Let us now prove that $\mathbb{T}$ is continuous and compact. Suppose $\left(v_{n}, v_{n}, \lambda_{n}\right)$ is a a bounded sequence in $\mathbb{X} \times[0,1]$ and let $\left(w_{n}, \mu_{n}\right)=\mathbb{T}\left(v_{n}, v_{n}, \lambda_{n}\right)$. To prove compactness, we show that, up to a subsequence, $\left(w_{n}, \mu_{n}\right)$ converges to some $(w, \mu)$ in $\mathbb{X}$. Since $\left(v_{n}\right)_{\chi}$ is uniformly bounded, by virtue of [35, Proposition 3.10], the sequence $\mu_{n}$ is relatively compact in $\mathrm{L}^{1}\left(\mathrm{Q}_{\mathrm{T}}\right)_{+}$, thanks to (2.31) (cf. (2.32) below where more details are given). Since $w_{n}$ and $\left(w_{n}\right)_{x}$ are uniformly bounded in $\mathcal{C}^{\alpha}\left(\bar{Q}_{T}\right)$, by the Ascoli-Arzelà Theorem and uniform convergence of the derivative there exists some $w$ such that $w, w_{x}$ are continuous in $\overline{\mathrm{Q}_{T}}$ and, passing to a subsequence, $w_{n} \rightarrow w$ and $\left(w_{n}\right)_{x} \rightarrow w_{x}$ uniformly, where in fact $w_{n} \rightharpoonup w$ weakly in $W_{s}^{1,2}\left(\mathrm{Q}_{\mathrm{T}}\right)$ for any $s>1$. This is what we wanted to show.
To prove continuity, we assume $\left(v_{n}, v_{n}, \lambda_{n}\right) \rightarrow(v, v, \lambda)$ in $\mathbb{X} \times[0,1]$. It is enough to show that, after passing to a subsequence, $\mathbb{T}\left(v_{n}, v_{n}, \lambda_{n}\right) \rightarrow \mathbb{T}(v, v, \lambda)$. By the preceding argument, we can assume $\mathbb{T}\left(v_{n}, v_{n}, \lambda_{n}\right) \rightarrow(w, \mu)$. We can also use estimates (2.5b)-(2.5d) to deduce that $q_{v_{n}, v_{n}} \rightarrow q_{v, v}$ a.e. (cf. the proof of Equation (2.35) below), and since $q_{v_{n}, v_{n}}$ is uniformly bounded we can also assert $q_{v_{n}, v_{n}} \rightarrow q_{v, v}$ in $L^{s}$ for any $s \geqslant 1$. Then we deduce that $(w, \mu)$ is a solution of (2.30) for the given $(v, v, \lambda)$. Therefore, $(w, \mu)=\mathbb{T}(v, v, \lambda)$, as desired.

Now, let $(\mathfrak{u}, \mathfrak{m}) \in \mathbb{X}$ and $\lambda \in[0,1]$ so that $(u, m)=\mathbb{T}(u, m, \lambda)$. Then $(u, m)$ satisfies assumptions of Lemma 2.5 with $\mathfrak{m}_{0}, u_{\top}, q_{u, m}$ replaced by $\lambda m_{0}, \lambda u_{\top}$ and $\lambda q_{\mathfrak{u}, \mathfrak{m}}$, respectively. Since the bounds of Lemma 2.5 carry through uniformly in $\lambda \in[0,1]$ we infer that

$$
\|(u, m)\|_{\mathbb{X}} \leqslant 1 \vee c_{0},
$$

where $\mathfrak{c}_{0}>0$ is the constant of Lemma 2.5. In addition, for $\lambda=0$ we have $\mathbb{T}(u, m, 0)=$ $(0,0)$. Therefore, by virtue of Leray-Schauder fixed point Theorem (see e.g. [17, Theorem 11.6]), we deduce the existence of a solution ( $u, m$ ) in $\mathbb{X}$ to system (1.6).

Step 2 (measure data): We deal now with general $\mathrm{m}_{0}$, i.e. a probability measure that is supported on $(0, L]$. Let $\left(m_{0}^{n}\right) \subset L^{1}(Q)_{+}$be a sequence of functions, which converges weakly (in the sense of measures) to $m_{0}$, and such that

$$
\left\|\mathfrak{m}_{0}^{\mathfrak{n}}\right\|_{\mathrm{L}^{1}} \leqslant \int_{0}^{\mathrm{L}} \mathrm{~d} \mathfrak{m}_{0} \leqslant 1, \text { and } \operatorname{supp}\left(\mathfrak{m}_{0}^{\mathfrak{n}}\right) \subset(0, \mathrm{~L}]
$$

For any $n \geqslant 1$, define $\left(u^{n}, m^{n}\right)$ to be a solution in $\mathbb{X}$ to system (1.6) starting from $m_{0}^{n}$.
In view of [35, Proposition 3.10 (iii)] and (2.22), the corresponding solutions $\mathrm{m}^{n}$ to the non-local Fokker-Planck equation lie in a relatively compact set of $L^{1}\left(Q_{T}\right)$. Moreover, it holds that

$$
\begin{equation*}
m^{n} \geqslant 0 \text { and } \sup _{0 \leqslant t \leqslant T}\left\|m^{n}(t)\right\|_{L^{1}} \leqslant \int_{0}^{L} d m_{0} \tag{2.32}
\end{equation*}
$$

Passing to a subsequence we have $m^{n} \rightarrow m$ in $L^{1}\left(Q_{T}\right), m^{n}(t) \rightarrow m(t)$ in $L^{1}(Q)$ for a.e. $t$ in $(0, T)$, and $m^{n} \rightarrow m$ for a.e. $(t, x)$ in $Q_{T}$. It follows that $m \in L^{1}\left(Q_{T}\right)_{+}$and

$$
\begin{equation*}
\|m(t)\|_{L^{1}} \leqslant 1 \quad \text { for a.e. } t \in(0, T) \tag{2.33}
\end{equation*}
$$

In addition, we know that $q_{u, m}$ fulfils the assumptions of Lemma 2.2. Thus $t \rightarrow\|m(t)\|_{L^{1}}$ is continuous on $(0, T]$, so that (2.33) holds for avery $t \in(0, T]$. Furthermore, we can
appeal to the probabilistic characterisation (2.10), thanks to Remark 2.4, to get

$$
\begin{aligned}
\left|\int_{0}^{\mathrm{L}} \phi(x)(\mathfrak{m}(\mathrm{t}+\mathrm{h}, \mathrm{x})-\mathfrak{m}(\mathrm{t})) \mathrm{dx}\right| \leqslant \mathbb{E} \mid & \phi\left(X_{\mathrm{t}+\mathrm{h}}\right) \mathbb{1}_{\mathrm{t}+\mathrm{h}<\tau}-\phi\left(\mathrm{X}_{\mathrm{t}}\right) \mathbb{1}_{\mathrm{t}<\tau} \mid \\
& \leqslant\|\phi\|_{\infty}|\mathfrak{\eta}(\mathrm{t})-\mathfrak{\eta}(\mathrm{t}+\mathrm{h})|+\mathbb{E}\left|\phi\left(\mathrm{X}_{\mathrm{t}+\mathrm{h}}\right)-\phi\left(\mathrm{X}_{\mathrm{t}}\right)\right|
\end{aligned}
$$

for every $\phi \in \mathcal{C}(\bar{Q})$, and $t \in[0, T]$. Now owing to Lemma $2.2, \eta$ is continuous on $[0, T]$. Hence, by taking the limit in the last estimation we infer that

$$
\lim _{h \rightarrow 0} \int_{0}^{\mathrm{L}} \phi(x)(\mathfrak{m}(\mathrm{t}+\mathrm{h}, \mathrm{x})-\mathfrak{m}(\mathrm{t})) \mathrm{d} x=0
$$

thanks to the bounded convergence theorem. Consequently the map $t \rightarrow m(t)$ is continuous on $[0, \mathrm{~T}]$ with respect to the strong topology of $\mathfrak{M}(\overline{\mathrm{Q}})$.

On the other hand, by Lemma 2.5 we have that $u^{n}, u_{x}^{n}$ are uniformly bounded in $\mathcal{C}^{\theta}\left(\bar{Q}_{T}\right)$, and $u_{t}^{n}, u_{x x}^{n}$ are uniformly bounded in $\mathcal{C}^{\Theta}\left(Q^{\prime}\right)$ for each $Q^{\prime} \subset \subset(0, T) \times(0, L]$. Thus, up to a subsequence we obtain that $u, u_{x} \in \mathcal{C}\left(\overline{Q_{T}}\right)$, and

$$
\begin{equation*}
\mathbf{u}^{\mathfrak{n}} \rightarrow \mathbf{u} \in \mathfrak{C}^{1,2}((0, \mathrm{~T}) \times(0, \mathrm{~L}]) \tag{2.34}
\end{equation*}
$$

where the convergence is in the $\mathcal{C}^{1,2}$ norm on arbitrary compact subsets of $(0, T) \times(0, L]$.
To show that the Hamilton-Jacobi equation holds in a classical sense and the FokkerPlanck equation holds in the sense of distributions, it remains to show that

$$
\begin{equation*}
\mathrm{q}_{\mathfrak{u}^{n}, \mathfrak{m}^{n}} \rightarrow \mathrm{q}_{\mathrm{u}, \mathrm{~m}} \text { a.e. } \tag{2.35}
\end{equation*}
$$

at least on a subsequence. Set $p^{n}=p_{u^{n}, \mathfrak{m}^{n}}=\Lambda\left(u_{x}^{n}, m^{n}, p^{n}\right)$ and $p=p_{u, m}=\Lambda\left(u_{x}, m, p\right)$, with $\wedge$ defined in Lemma 2.1. Using (2.5b)-(2.5d) we get

$$
\begin{equation*}
\left\|p^{n}(t, \cdot)-p(t, \cdot)\right\|_{\infty} \leqslant\left\|\Lambda\left(u_{x}^{n}, m^{n}, p^{n}\right)(t, \cdot)-\Lambda\left(u_{x}, m^{n}, p^{n}\right)(t, \cdot)\right\|_{\infty} \tag{2.36}
\end{equation*}
$$

$$
+\left\|\Lambda\left(u_{x}, m^{n}, p^{n}\right)(t, \cdot)-\Lambda\left(u_{x}, m^{n}, p\right)(t, \cdot)\right\|_{\infty}+\left\|\Lambda\left(u_{x}, m^{n}, p\right)(t, \cdot)-\Lambda\left(u_{x}, m, p\right)(t, \cdot)\right\|_{\infty}
$$

$$
\leqslant \frac{1}{2}\left\|u_{x}^{n}-u_{x}\right\|_{\infty}+\frac{\kappa}{1+\kappa}\left\|p^{n}(t, \cdot)-p(t, \cdot)\right\|_{\infty}+\left|\ell\left(m^{n}, p\right)(t)-\ell(m, p)(t)\right|
$$

which means

$$
\begin{equation*}
\left\|p^{n}(t, \cdot)-p(t, \cdot)\right\|_{\infty} \leqslant \frac{1+k}{2}\left\|u_{x}^{n}-u_{x}\right\|_{\infty}+(1+\kappa)\left|\ell\left(m^{n}, p\right)(t)-\ell(m, p)(t)\right| . \tag{2.37}
\end{equation*}
$$

Noting that (up to a subsequence) $m^{n}(t) \rightarrow m(t)$ in $L^{1}(Q)$ a.e., we use the fact that $a(m), c(m), \eta(m)$ are all continuous with respect to this metric to deduce that

$$
\begin{equation*}
\left|\ell\left(m^{n}, p\right)(t)-\ell(m, p)(t)\right| \rightarrow 0 \text { a.e. } t \in(0, T) \tag{2.38}
\end{equation*}
$$

from which we conclude that

$$
\begin{equation*}
\left\|p^{n}(t, \cdot)-p(t, \cdot)\right\|_{\infty} \rightarrow 0 \text { a.e. } t \in(0, T) . \tag{2.39}
\end{equation*}
$$

Now from (2.39) and (2.5a) we have

$$
\begin{equation*}
\left|\ell\left(m, p^{n}\right)(t)-\ell(m, p)(t)\right| \rightarrow 0 \text { a.e. } t \in(0, T) . \tag{2.40}
\end{equation*}
$$

Combining (2.38) and (2.40) we see that $\ell\left(\mathfrak{m}^{n}, p^{n}\right) \rightarrow \ell(m, p)$ a.e. We deduce (2.35) from the definition (2.3c). Therefore $\left(u^{n}, m^{\mathfrak{n}}\right)$ converges to some $(u, m)$ which is a solution to (1.6) with initial data $\mathrm{m}_{0}$.

Step 3 (uniqueness): Let $\left(\mathfrak{u}_{\mathfrak{i}}, \mathfrak{m}_{\mathfrak{i}}\right), \mathfrak{i}=1,2$ be two solutions of (1.6). We set

$$
\mathrm{G}_{\mathrm{i}}:=\mathrm{q}_{\mathfrak{u}_{\mathrm{i}}, \mathfrak{m}_{\mathfrak{i}}} \text { and } \overline{\mathrm{G}}_{i}:=\int_{0}^{\mathrm{L}} \mathrm{q}_{\mathfrak{u}_{\mathfrak{i}}, \mathfrak{m}_{\mathfrak{i}}}(\mathrm{t}, \mathrm{y}) \mathrm{dm}(\mathrm{~m})
$$

From (1.7), we know that

$$
\begin{equation*}
\mathrm{G}_{i}=\frac{1}{2}\left(1-\kappa \overline{\mathrm{G}}_{i}-\mathbf{u}_{\mathrm{i}, \mathrm{x}}\right)^{+} . \tag{2.41}
\end{equation*}
$$

Let $\mathfrak{u}=\mathfrak{u}_{1}-\mathfrak{u}_{2}, \mathrm{~m}=\mathrm{m}_{1}-\mathrm{m}_{2}, \mathrm{G}=\mathrm{G}_{1}-\mathrm{G}_{2}, \overline{\mathrm{G}}=\overline{\mathrm{G}}_{1}-\overline{\mathrm{G}}_{2}$. Using $(\mathrm{t}, \mathrm{x}) \rightarrow \mathrm{e}^{-r t} \mathfrak{u}(\mathrm{t}, \mathrm{x})$ as a test function in the equations satisfied by $m_{1}, m_{2}$, with some algebra yields
(2.42) $0=\int_{0}^{T} e^{-r t} \int_{0}^{\mathrm{L}}\left(\mathrm{G}_{2}^{2}-\mathrm{G}_{1}^{2}-\mathrm{G}_{1} \mathrm{u}_{\mathrm{x}}\right) \mathrm{m}_{1}+\left(\mathrm{G}_{1}^{2}-\mathrm{G}_{2}^{2}+\mathrm{G}_{2} \mathrm{u}_{\mathrm{x}}\right) \mathrm{m}_{2} \mathrm{dx} \mathrm{dt}$

$$
=\int_{0}^{\mathrm{T}} e^{-r t} \int_{0}^{\mathrm{L}}\left(\mathrm{G}_{1}-\mathrm{G}_{2}\right)^{2}\left(\mathfrak{m}_{1}+\mathrm{m}_{2}\right) \mathrm{dxdt}+\int_{0}^{\mathrm{T}} e^{-r t} \int_{0}^{\mathrm{L}}\left(2 \mathrm{G}+\mathrm{u}_{\mathrm{x}}\right)\left(\mathrm{G}_{2} \mathrm{~m}_{2}-\mathrm{G}_{1} \mathrm{~m}_{1}\right) \mathrm{dxdt} .
$$

Now since $\mathrm{G}_{2}=0$ on the set where $1-\kappa \overline{\mathrm{G}}_{2}(\mathrm{t})-\mathrm{u}_{2, \mathrm{x}}<0$, we can write

$$
\begin{aligned}
\left(2 \mathrm{G}+\mathrm{u}_{\mathrm{x}}\right) \mathrm{G}_{2} & =\left(\left(1-\kappa \overline{\mathrm{G}}_{1}-\mathfrak{u}_{1, \mathrm{x}}\right)^{+}-\left(1-\kappa \overline{\mathrm{G}}_{2}(\mathrm{t})-\mathfrak{u}_{2, \mathrm{x}}\right)+\mathfrak{u}_{1, \mathrm{x}}-\mathfrak{u}_{2, \mathrm{x}}\right) \mathrm{G}_{2} \\
& =\left(-\kappa \overline{\mathrm{G}}+\left(1-\kappa \overline{\mathrm{G}}_{1}-\mathfrak{u}_{1, \mathrm{x}}\right)^{-}\right) \mathrm{G}_{2} .
\end{aligned}
$$

Similarly we can write

$$
\begin{aligned}
\left(2 \mathrm{G}+\mathrm{u}_{\mathrm{x}}\right) \mathrm{G}_{1} & =\left(\left(1-\kappa \overline{\mathrm{G}}_{1}-\mathfrak{u}_{1, x}\right)-\left(1-\kappa \overline{\mathrm{G}}_{2}(\mathrm{t})-\mathfrak{u}_{2, \mathrm{x}}\right)^{+}+\mathfrak{u}_{1, x}-\mathfrak{u}_{2, \mathrm{x}}\right) \mathrm{G}_{1} \\
& =\left(-\kappa \overline{\mathrm{G}}-\left(1-\kappa \overline{\mathrm{G}}_{2}-\mathfrak{u}_{2, \mathrm{x}}\right)^{-}\right) \mathrm{G}_{1} .
\end{aligned}
$$

Thus we compute

$$
\begin{aligned}
\int_{0}^{\mathrm{L}}\left(2 \mathrm{G}+\mathrm{u}_{x}\right)\left(\mathrm{G}_{2} \mathrm{~m}_{2}-\mathrm{G}_{1} \mathrm{~m}_{1}\right) \mathrm{d} x \mathrm{dt}=\kappa \overline{\mathrm{G}}^{2} & +\int_{0}^{\mathrm{L}}\left(1-\kappa \overline{\mathrm{G}}_{1}-\mathrm{u}_{1, x}\right)^{-} \mathrm{G}_{2} m_{2} \mathrm{~d} x \mathrm{dt} \\
& +\int_{0}^{\mathrm{L}}\left(1-\kappa \overline{\mathrm{G}}_{2}-\mathfrak{u}_{2, x}\right)^{-} \mathrm{G}_{1} m_{1} d x d t \geqslant \kappa \overline{\mathrm{G}}^{2} .
\end{aligned}
$$

So from (2.42) we conclude

$$
\begin{equation*}
\int_{0}^{T} e^{-r t} \int_{0}^{\mathrm{L}}\left(\mathrm{G}_{1}-\mathrm{G}_{2}\right)^{2}\left(\mathrm{~m}_{1}+\mathrm{m}_{2}\right) \mathrm{d} x \mathrm{dt}+\mathrm{k} \int_{0}^{\mathrm{T}} \mathrm{e}^{-\mathrm{rt}}\left(\overline{\mathrm{G}}_{1}-\overline{\mathrm{G}}_{2}\right)^{2} \mathrm{dt}=0 . \tag{2.43}
\end{equation*}
$$

In particular, $\overline{\mathrm{G}}_{1} \equiv \overline{\mathrm{G}}_{2}$. We can then appeal to uniqueness for the Hamilton-Jacobi equation to get $\mathfrak{u}_{1} \equiv \mathfrak{u}_{2}$ (cf. [26, Chapter V]). By (2.41), this entails that $\mathrm{G}_{1} \equiv \mathrm{G}_{2}$, and so $\mathrm{m}_{1} \equiv \mathrm{~m}_{2}$ by uniqueness for the Fokker-Planck equation.

## 3. Application of the MFG approach

This section is devoted to the proof of Theorem 1.6. Namely, we show that the optimal feedback strategy, computed from the MFG system (1.6), provides an $\varepsilon$-Nash equilibria
for the $N$-Player Cournot game, where the error $\varepsilon$ is arbitrarily small as $N \rightarrow \infty$. Throughout this section $(u, m)$ is the solution to (1.6) starting from some probability measure $m_{0}$ satisfying $(\mathcal{H} 2)$, and the function $\mathrm{q}_{\mathrm{u}, \mathrm{m}}$ is given by (1.7). Moreover, we define

$$
\left\{\begin{array}{l}
\mathrm{d} \hat{X}_{\mathrm{t}}^{\mathrm{i}}=-\mathrm{q}_{\mathrm{u}, \mathrm{~m}}\left(\mathrm{t}, \hat{X}_{\mathrm{t}}^{\mathrm{i}}\right) \mathrm{dt}+\sigma \mathrm{d} W_{\mathrm{t}}^{i}-\mathrm{d} \xi_{\mathrm{t}}^{\hat{X}^{i}}  \tag{3.1}\\
X_{0}^{i}=V_{i},
\end{array}\right.
$$

and set $\hat{q}_{t}^{i}:=q_{u, m}\left(t, \hat{X}_{t}^{i}\right)$. We recall that the objective functional is defined as

$$
\mathfrak{d}_{c}^{i, N}\left(q^{1}, \ldots, q^{N}\right):=\mathbb{E}\left\{\int_{0}^{T} e^{-r s}\left(1-\kappa \bar{q}_{s}^{i}-q_{s}^{i}\right) q_{s}^{i} \mathbb{1}_{s<\tau^{i}} d s+e^{-r T} u_{T}\left(X_{\tau^{i}}^{i}\right)\right\} .
$$

Our goal is to prove that

$$
\mathcal{J}_{\mathfrak{c}}^{i, N}\left(q^{i} ;\left(\hat{q}^{j}\right)_{j \neq i}\right) \leqslant \varepsilon+\mathcal{I}_{\mathfrak{c}}^{i, N}\left(\hat{q}^{1}, \ldots, \hat{q}^{N}\right) \quad \forall q^{i} \in \mathcal{A}_{c}, \forall i=1, \ldots, N
$$

as long as N is large enough.
Let us set

$$
\hat{\tau}^{i}:=\inf \left\{t \geqslant 0: \hat{X}_{t}^{i} \leqslant 0\right\} \wedge T
$$

and define the following process:

$$
\begin{equation*}
\hat{v}_{t}^{N}:=\frac{1}{\mathrm{~N}} \sum_{\mathrm{k}=1}^{\mathrm{N}} \delta_{\hat{\mathrm{x}}_{\mathrm{t}}^{k}} \mathbb{1}_{\mathrm{t}<\hat{\tau}^{k}}, \quad \forall \mathrm{t} \in[0, \mathrm{~T}], \tag{3.2}
\end{equation*}
$$

where $\delta_{x}$ denotes the Dirac delta measure of the point $x \in \mathbb{R}$. Observe that the above definition makes sense because the stochastic dynamics ( $\hat{X}^{1}, \ldots, \hat{X}^{N}$ ) exists in the strong sense owing to Lemma 2.5. In particular, the strategy profile $\left(\hat{q}^{1}, \ldots, \hat{q}^{N}\right)$ defined in Theorem 1.6 belongs to $\prod_{i=1}^{N} \mathbb{A}_{\mathcal{c}}$. Moreover, by using the probabilistic characterization (2.8c), note that for any measurable and bounded function $\phi$ on $\overline{\mathrm{Q}}$ we have

$$
\begin{equation*}
\mathbb{E}\left[\int_{0}^{\mathrm{L}} \phi \mathrm{~d} \hat{v}_{\mathrm{t}}^{\mathrm{N}}\right]=\int_{0}^{\mathrm{L}} \phi \mathrm{dm}(\mathrm{t}), \quad \text { for a.e. } \mathrm{t} \in(0, \mathrm{~T}) . \tag{3.3}
\end{equation*}
$$

The above identity is not strong enough to show Theorem 1.6 and we need a stronger condition (cf. (3.15)). Therefore, we need to work harder in order to get more information on the asymptotic behavior of the empirical process (3.2) when $\mathrm{N} \rightarrow \infty$.

We aim to prove that the empirical process $\left(\hat{v}^{\mathrm{N}}\right)_{\mathrm{N} \geqslant 1}$ converges in law to the deterministic measure $m$ in a suitable function space, by using arguments borrowed from [23,30]. For this, we start by showing the existence of sub-sequences ( $\hat{v}^{\mathrm{N}^{\prime}}$ ) that converges in law to some limiting process $v^{*}$. Then, we show that $v^{*}$ belongs to $\tilde{\mathcal{P}}(\overline{\mathrm{Q}})$ and satisfies the same equation as $m$. Finally, we invoke the uniqueness of weak solutions to the Fokker-Planck equation to deduce full weak convergence toward $m$.

The crucial step consists in showing that the sequence of the laws of $\left(\hat{v}^{N}\right)_{N \geqslant 1}$ is relatively compact on a suitable topological space. This is where the machinery of [30] is convenient. In order to use the analytical tools of that paper, we view the empirical process as a random variable on the space of càdlàg (right continuous and has left-hand limits) functions, mapping $[0, \mathrm{~T}]$ into the space of tempered distributions. This function space is denoted $D_{\mathbb{S}_{\mathbb{R}}^{\prime}}$ and is endowed with the so called Skorokhod's M1 topology. Note
that there are no measurability issues owing to [30, Proposition 2.7]. Moreover, by virtue of [34], the process $\left(\hat{v}_{t}^{N}\right)_{t \in[0, T]}$ has a version that is càdlàg in the strong topology of $\mathcal{S}_{\mathbb{R}}^{\prime}$ for every $N \geqslant 1$, since $\hat{v}_{t}^{N}(\phi):=\int_{\mathbb{R}} \phi \mathrm{d} \hat{v}_{t}^{N}$ is a real-valued càdlàg process, for every $\phi \in \mathcal{S}_{\mathbb{R}}$ and $N \geqslant 1$. We refer the reader to [30] for the construction of ( $D_{\mathcal{S}_{\mathbb{R}}^{\prime}}, M 1$ ), and to [38] for general background on Skorokhod's topologies. We shall denote by ( $D_{\mathbb{R}}$, M1) the space of $\mathbb{R}$-valued càdlàg functions mapping $[0, T]$ to $\mathbb{R}$, endowed with Skorokhod's M1 topology.

The main strengths of working with the M1 topology in our context, are based on the following facts:

- tightness on $\left(\mathrm{D}_{\mathcal{S}_{\mathfrak{R}}}, \mathrm{M} 1\right)$ implies the relative compactness on ( $\mathrm{D}_{\mathcal{S}_{\mathbb{R}}^{\prime}}, \mathrm{M} 1$ ) thanks to [30, Theorem 3.2]);
- the proof of tightness on $\left(\mathrm{D}_{\mathcal{S}_{\mathbb{R}}}, M 1\right)$ is reduced through the canonical projection to the study of tightness in ( $D_{\mathbb{R}}, M 1$ ), for which we have suitable characterizations [30,38];
- bounded monotone real-valued processes are automatically tight on ( $\mathrm{D}_{\mathbb{R}}, \mathrm{M} 1$ ); this is an important feature, that enables to prove tightness of the sequence of empirical process laws, by using a suitable decomposition.
It is also important to note that this approach could be generalized to deal with the case of a systemic noise, by using a martingale approach as in [23, Lemma 5.9]. We do not deal with that case in this paper.

More generally, one can replace $\mathcal{S}_{\mathbb{R}}^{\prime}$ by any dual space of a countably Hilbertian nuclear space (cf. [30] and references therein). Although the class $\mathcal{S}_{\mathbb{R}}^{\prime}$ seems to be excessively large for our purposes, we recover measure-valued processes by means of Riesz representation theorem (cf. [23, Proposition 5.3] for an example in the same context).

Throughout this part, we shall use the symbol $\Rightarrow$ to denote convergence in law. The key technical lemma of this section is the following:

Lemma 3.1. As $\mathrm{N} \rightarrow \infty$, we have $\hat{\gamma}^{\mathrm{N}} \Rightarrow \mathrm{m}$ on $\left(\mathrm{D}_{\mathcal{S}_{\mathfrak{R}}^{\prime}}\right.$, M 1 ), i.e. for every continuous bounded real-valued function $\Psi$ on $\left(\mathrm{D}_{\mathcal{S}_{\mathbb{R}}^{\prime}}, \mathrm{M} 1\right)$, it holds that

$$
\lim _{N} \mathbb{E}\left[\Psi\left(\hat{v}^{N}\right)\right]=\Psi(\mathfrak{m}) .
$$

The bulk of this section is devoted to the proof of Lemma 3.1. The proof of Theorem 1.6 is completed in Section 3.2.

### 3.1. Proof of Lemma 3.1.

3.1.1. Tightness. At first, we aim to prove the tightness of $\left(\hat{v}^{N}\right)_{N \geqslant 1}$ on the space ( $\left.\mathrm{D}_{\boldsymbol{S}_{\mathbb{R}}^{\prime}}, \mathrm{M} 1\right)$; that is, for every $\phi \in \mathcal{S}_{\mathbb{R}}$ and for all $\varepsilon>0$, there exists a compact subset $K$ of ( $D_{\mathbb{R}}, M 1$ ) such that:

$$
\mathbb{P}\left(\hat{\gamma}^{\mathrm{N}}(\phi) \in K\right)>1-\varepsilon \quad \text { for all } \mathrm{N} \geqslant 1 .
$$

For that purpose, we shall use a convenient characterization of tightness in ( $\mathrm{D}_{\mathbb{R}}$, M1) (cf. [38, Theorem 12.12.3]).

We start by controlling the concentration of mass at the origin:

Lemma 3.2. For every $t \in[0, \mathrm{~T}]$, we have

$$
\sup _{\mathrm{N} \geqslant 1} \mathbb{E} \hat{v}_{\mathrm{t}}^{\mathrm{N}}(0, \varepsilon) \rightarrow 0, \quad \text { as } \varepsilon \rightarrow 0
$$

Proof. Let us fix $\varepsilon>0$. Note that, for every $\mathrm{t} \in[0, \mathrm{~T}]$

$$
\mathbb{E} \hat{v}_{t}^{N}(0, \varepsilon)=\frac{1}{N} \sum_{i=1}^{N} \mathbb{P}\left(\hat{X}_{t}^{i} \in(0, \varepsilon) ; t<\hat{\tau}^{i}\right)
$$

Thus, on the one hand

$$
\sup _{\mathrm{N} \geqslant 1} \mathbb{E} \hat{\boldsymbol{\gamma}}_{0}^{\mathrm{N}}(0, \varepsilon)=\int_{0}^{\varepsilon} \mathrm{dm}_{0} \rightarrow 0, \quad \text { as } \varepsilon \rightarrow 0 .
$$

On the other hand, we have for any $t \in(0, T]$

$$
\begin{equation*}
\sup _{N \geqslant 1} \mathbb{E} \hat{\gamma}_{t}^{N}(0, \varepsilon) \leqslant \sup _{N \geqslant 1} N^{-1} \sum_{i=1}^{N} \mathbb{E}\left[\left(1-\phi_{\varepsilon}\left(\hat{X}_{t}^{i}\right)\right) \mathbb{1}_{t<\hat{\tau}^{i}}\right] \tag{3.4}
\end{equation*}
$$

where $\phi_{\varepsilon}$ is the cut-off function defined in (2.12). Thus, by virtue of (2.16) we obtain

$$
\sup _{N \geqslant 1} \mathbb{E} \hat{v}_{t}^{N}(0, \varepsilon) \leqslant C\left(L, t,\left\|q_{u, m}\right\|_{\infty}\right) \varepsilon^{1 / 4}
$$

which entails the desired result.
The second ingredient is the control of the mass loss increment:
Lemma 3.3. For every $t \in[0, T]$ and $\lambda>0$

$$
\lim _{h \rightarrow 0} \lim \sup _{N} \mathbb{P}\left(\left|\eta\left(\hat{v}_{t}^{N}\right)-\eta\left(\hat{v}_{t+h}^{N}\right)\right| \geqslant \lambda\right)=0,
$$

where the map $\mu \rightarrow \eta(\mu)$ is defined in (2.2).
Proof. The proof is inspired by [23, Proposition 4.7]. Let $\varepsilon, h>0$ and $t \in[0, T]$, we have

$$
\begin{align*}
\mathbb{P}\left(\eta\left(\hat{v}_{t}^{N}\right)-\eta\right. & \left.\left(\hat{v}_{t+h}^{N}\right) \geqslant \lambda\right)  \tag{3.5}\\
& \leqslant \mathbb{P}\left(\hat{v}_{t}^{N}(0, \varepsilon) \geqslant \lambda / 2\right)+\mathbb{P}\left(\eta\left(\hat{v}_{t}^{N}\right)-\eta\left(\hat{v}_{t+h}^{N}\right) \geqslant \lambda ; \hat{v}_{t}^{N}(0, \varepsilon)<\lambda / 2\right) .
\end{align*}
$$

The reason why we use the latter decomposition will be clear in (3.6). Owing to Markov's inequality and Lemma 3.2, one has

$$
\lim \sup _{N} \mathbb{P}\left(\hat{\gamma}_{t}^{N}(0, \varepsilon) \geqslant \lambda / 2\right) \leqslant 2 \lambda^{-1} \sup _{N} \mathbb{E} \hat{v}_{t}^{N}(0, \varepsilon) \rightarrow 0, \quad \text { as } \varepsilon \rightarrow 0
$$

Now we deal with the second part in estimate (3.5). Define $\mathcal{J}_{\mathfrak{t}}$ to be the following random set of indices:

$$
\mathcal{J}_{t}:=\left\{1 \leqslant \mathfrak{i} \leqslant N: \hat{X}_{t}^{i} \geqslant \varepsilon\right\} ;
$$

then, we have

$$
\begin{aligned}
& \mathbb{P}\left(\eta\left(\hat{v}_{t}^{N}\right)-\eta\left(\hat{v}_{t+h}^{N}\right) \geqslant \lambda ; \hat{v}_{t}^{N}(0, \varepsilon)<\lambda / 2\right) \\
& \leqslant \sum_{\# \mathcal{J} \geqslant N(1-\lambda / 2)} \mathbb{P}\left(\eta\left(\hat{v}_{t}^{N}\right)-\eta\left(\hat{v}_{t+h}^{N}\right) \geqslant \lambda \mid \mathcal{J}_{t}=\mathcal{J}\right) \mathbb{P}\left(\mathcal{J}_{t}=\mathcal{J}\right),
\end{aligned}
$$

where \#J denotes the number of elements of $\mathcal{J} \subseteq\{1,2, \ldots, N\}$. Thus, we reduce the problem to the estimation of the dynamics increments; using the same steps as for (2.9c) we have

$$
\begin{align*}
& \mathbb{P}\left(\eta\left(\hat{v}_{t}^{N}\right)-\eta\left(\hat{v}_{t+h}^{N}\right) \geqslant \lambda \mid \mathcal{J}_{t}=\mathcal{J}\right)  \tag{3.6}\\
& \leqslant \mathbb{P}\left(\#\left\{i \in \mathcal{J}: \inf _{s \in[t, t+h]} \hat{X}_{s}^{i}-\hat{X}_{t}^{i} \leqslant-\varepsilon\right\} \geqslant \lambda N / 2 \mid \mathcal{J}_{\mathfrak{t}}=\mathcal{J}\right) \\
& \leqslant \mathbb{P}\left(\#\left\{i \in \mathcal{J}: \sup _{s \in[0, h]} B_{s}^{i}-\inf _{s \in[0, h]} B_{s}^{i} \geqslant \frac{\varepsilon-h}{\sigma}\right\} \geqslant \lambda N / 2\right),
\end{align*}
$$

where we have used the uniform bound on $q_{u, m}$ of Lemma 2.5, and where $\left(B^{i}\right)_{1 \leqslant i \leqslant N}$ is a family of independent Wiener processes. By symmetry, this final probability depends only on \#J, so that the right hand side above is maximized when $\mathcal{J}=\{1, \ldots, N\}$. We infer that
$\mathbb{P}\left(\eta\left(\hat{v}_{t}^{N}\right)-\eta\left(\hat{v}_{t+h}^{N}\right) \geqslant \lambda ; \hat{v}_{t}^{N}(0, \varepsilon)<\lambda / 2\right) \leqslant \mathbb{P}\left(\frac{1}{N} \sum_{i=1}^{N} \mathbb{1}_{\left\{\sup _{s \in[0, h]} B_{s}^{i}-\inf _{s \in[0, h]} B_{s}^{i} \geqslant \frac{\varepsilon-h}{\sigma}\right\}} \geqslant \lambda / 2\right)$.
In the same way as for $(2.9 \mathrm{c})$, we choose $\varepsilon(\mathrm{h})=\mathrm{h}^{1 / 2} \log (1 / \mathrm{h})$ so that $\lim _{h \rightarrow 0^{+}} \varepsilon(\mathrm{h})=0$, and use Markov's inequality to get

$$
\mathbb{P}\left(\eta\left(\hat{v}_{t}^{N}\right)-\eta\left(\hat{v}_{t+h}^{N}\right) \geqslant \lambda ; \hat{v}_{t}^{N}(0, \varepsilon)<\lambda / 2\right) \leqslant \frac{4 \sigma}{\lambda\left(\log (1 / h)-h^{1 / 2}\right)} .
$$

This entails the desired result by taking the limit $h \rightarrow 0^{+}$.
Now we deal with the case of a left hand limit. Let $t \in(0, T]$ and $h \mapsto \varepsilon(h)$ as defined above. Using a similar decomposition as before, we have for small enough $h>0$

$$
\begin{aligned}
& \mathbb{P}\left(\eta\left(\hat{v}_{t-h}^{N}\right)-\eta\left(\hat{v}_{t}^{N}\right) \geqslant \lambda\right) \\
& \leqslant \mathbb{P}\left(\hat{v}_{t-h}^{N}(0, \varepsilon) \geqslant \lambda / 2\right)+\mathbb{P}\left(\eta\left(\hat{v}_{t-h}^{N}\right)-\eta\left(\hat{v}_{t}^{N}\right) \geqslant \lambda ; \hat{v}_{t-h}^{N}(0, \varepsilon)<\lambda / 2\right) .
\end{aligned}
$$

Appealing to Markov's inequality, estimate (3.4), and estimate (2.16) of Section 2, we have for small enough $h>0$

$$
\mathbb{P}\left(\hat{v}_{t-h}^{N}(0, \varepsilon) \geqslant \lambda / 2\right) \leqslant 2 \lambda^{-1} \mathbb{E} \hat{v}_{t-h}^{N}(0, \varepsilon) \leqslant 2 C \lambda^{-1}\left(1-e^{-\pi^{2} t / 8 L^{2}}\right)^{-1 / 4} \varepsilon^{1 / 4}
$$

whence

$$
\lim _{h \rightarrow 0^{+}} \lim \sup _{N} \mathbb{P}\left(\hat{\gamma}_{t-h}^{N}(0, \varepsilon(h)) \geqslant \lambda / 2\right)=0 .
$$

On the other hand, we show by using the same steps as in (3.6) that

$$
\mathbb{P}\left(\eta\left(\hat{v}_{t-h}^{N}\right)-\eta\left(\hat{v}_{t}^{N}\right) \geqslant \lambda ; \hat{v}_{t-h}^{N}(0, \varepsilon)<\lambda / 2\right) \leqslant \frac{4 \sigma}{\lambda\left(\log (1 / h)-h^{1 / 2}\right)} .
$$

This entails the desired result by taking the limit $h \rightarrow 0^{+}$.
We are now in position to show tightness on ( $\mathrm{D}_{\mathcal{S}_{\mathfrak{R}}^{\prime}}, \mathrm{M} 1$ ).
Proposition 3.4 (Tightness). The sequence of the laws of $\left(\hat{v}^{\mathrm{N}}\right)_{\mathrm{N} \geqslant 1}$ is tight on the space $\left(\mathrm{D}_{\mathcal{S}_{\mathbb{R}}}, \mathrm{M} 1\right)$.

Proof. We present a brief sketch to explain the main arguments, and refer to [23, Proposition 5.1] for a similar proof.

Thanks to [30, Theorem 3.2], it is enough to show that the sequence of the laws of $\left(\hat{v}^{\mathrm{N}}(\phi)\right)_{\mathrm{N} \geqslant 1}$ is tight on $\left(\mathrm{D}_{\mathbb{R}}, \mathrm{M} 1\right)$ for any $\phi \in \mathcal{S}_{\mathbb{R}}$. To prove this, one can use the conditions of [38, Theorem 12.12.3], which can be rewritten in a convenient form by virtue of [2]. From [30, Proposition 4.1], we are done if we achieve the two following steps:
(1) find $\alpha, \beta, c>0$, such that

$$
\mathbb{P}\left(\mathrm{H}_{\mathbb{R}}\left(\hat{v}_{\mathrm{t}_{1}}^{\mathrm{N}}(\phi), \hat{v}_{\mathrm{t}_{2}}^{\mathrm{N}}(\phi), \hat{v}_{\mathrm{t}_{3}}^{\mathrm{N}}(\phi)\right) \geqslant \lambda\right) \leqslant \mathrm{c} \lambda^{-\alpha}\left|\mathrm{t}_{3}-\mathrm{t}_{1}\right|^{1+\beta}
$$

for any $\mathrm{N} \geqslant 1, \lambda>0$ and $0 \leqslant \mathrm{t}_{1}<\mathrm{t}_{2}<\mathrm{t}_{3} \leqslant \mathrm{~T}$, where

$$
\mathrm{H}_{\mathbb{R}}\left(\mathrm{x}_{1}, \mathrm{x}_{2}, \mathrm{x}_{3}\right):=\inf _{0 \leqslant \gamma \leqslant 1}\left|x_{2}-(1-\gamma) \mathrm{x}_{1}-\gamma x_{3}\right| \quad \text { for } x_{1}, x_{2}, x_{3} \in \mathbb{R} ;
$$

(2) show that

$$
\lim _{h \rightarrow 0^{+}} \lim _{N} \mathbb{P}\left(\sup _{t \in(0, h)}\left|\hat{v}_{t}^{N}(\phi)-\hat{v}_{0}^{N}(\phi)\right|+\sup _{t \in(\mathrm{~T}-\mathrm{h}, \mathrm{~T})}\left|\hat{v}_{\mathrm{T}}^{\mathrm{N}}(\phi)-\hat{v}_{t}^{N}(\phi)\right| \geqslant \lambda\right)=0 .
$$

The key step is to consider the following decomposition [30, Proposition 4.2]:

$$
\begin{equation*}
\bar{v}_{t}^{N}(\phi):=\frac{1}{N} \sum_{k=1}^{N} \phi\left(\hat{X}_{t}^{k} \wedge \hat{\tau}^{k}\right)=\hat{v}_{t}^{N}(\phi)+\phi(0) \varepsilon_{t}^{N} \tag{3.7}
\end{equation*}
$$

where

$$
\varepsilon_{t}^{N}:=1-\eta\left(\hat{v}_{t}^{N}\right)
$$

is the exit rate process, which quantifies the fraction of firms out of market. Since $\left(\varepsilon_{t}^{N}\right)_{t \in[0, T]}$ is monotone increasing we have

$$
\inf _{0 \leqslant \gamma \leqslant 1}\left|\varepsilon_{t_{2}}^{N}-(1-\gamma) \varepsilon_{t_{1}}^{N}-\gamma \varepsilon_{t_{3}}^{N}\right|=0
$$

so that

$$
\mathrm{H}_{\mathbb{R}}\left(\hat{v}_{\mathrm{t}_{1}}^{N}(\phi), \hat{v}_{\mathrm{t}_{2}}^{N}(\phi), \hat{v}_{\mathrm{t}_{3}}^{N}(\phi)\right) \leqslant\left|\bar{v}_{\mathrm{t}_{1}}^{\mathrm{N}}(\phi)-\bar{v}_{\mathrm{t}_{2}}^{N}(\phi)\right|+\left|\bar{v}_{\mathrm{t}_{2}}^{\mathrm{N}}(\phi)-\bar{v}_{\mathrm{t}_{3}}^{\mathrm{N}}(\phi)\right| .
$$

Thus, by virtue of Markov's inequality

$$
\begin{aligned}
\mathbb{P}\left(\mathrm{H}_{\mathbb{R}}\left(\hat{v}_{\mathrm{t}_{1}}^{\mathrm{N}}(\phi), \hat{v}_{\mathrm{t}_{2}}^{N}(\phi), \hat{v}_{\mathrm{t}_{3}}^{N}(\phi)\right)\right. & \geqslant \lambda) \\
& \leqslant 8 \lambda^{-4}\left(\mathbb{E}\left|\bar{v}_{\mathrm{t}_{1}}^{N}(\phi)-\bar{v}_{\mathrm{t}_{2}}^{N}(\phi)\right|^{4}+\mathbb{E}\left|\bar{v}_{\mathrm{t}_{2}}^{\mathrm{N}}(\phi)-\bar{v}_{\mathrm{t}_{3}}^{\mathrm{N}}(\phi)\right|^{4}\right) .
\end{aligned}
$$

Therefore, we deduce requirement (1) from the following estimate:
(3.8) $\forall \mathrm{s}, \mathrm{t} \in[0, \mathrm{~T}]$,

$$
\mathbb{E}\left|\bar{v}_{t}^{N}(\phi)-\bar{v}_{s}^{N}(\phi)\right|^{4} \leqslant\left\|\phi_{x}\right\|_{\infty}^{4} \frac{1}{N} \sum_{k=1}^{N} \mathbb{E}\left|\hat{X}_{t \wedge \hat{\tau}^{k}}^{k}-\hat{X}_{s \wedge \hat{\tau}^{k}}^{k}\right|^{4} \leqslant C\left\|\phi_{x}\right\|_{\infty}^{4}|t-s|^{2} ;
$$

where we have used Hölder's inequality and the Burkholder-Davis-Gundy inequality [37, Thm IV.42.1].

The second requirement is also obtained by using the latter estimate, decomposition (3.7), and Lemma 3.3. In fact, we have

$$
\begin{aligned}
& \mathbb{P}\left(\sup _{t \in(0, h)}\left|\hat{v}_{t}^{N}(\phi)-\hat{v}_{0}^{N}(\phi)\right| \geqslant \lambda\right) \\
& \leqslant \mathbb{P}\left(\sup _{t \in(0, h)}\left|\bar{v}_{t}^{N}(\phi)-\bar{v}_{0}^{N}(\phi)\right| \geqslant \lambda / 2\right)+\mathbb{P}\left(|\phi(0)| \mathcal{E}_{h}^{N} \geqslant \lambda / 2\right),
\end{aligned}
$$

so that the desired result follows thanks to (3.8), and Lemma 3.3. By the same way, we deal with the second term $\mathbb{P}\left(\sup _{t \in(T-h, T)}\left|\hat{v}_{T}^{N}(\phi)-\hat{v}_{t}^{N}(\phi)\right| \geqslant \lambda\right)$.
3.1.2. Full convergence. We arrive now at the final ingredient for the proof of Lemma 3.1. Let us set

$$
\mathcal{C}^{\text {test }}:=\left\{\phi \in \mathcal{C}_{c}^{\infty}([0, \mathrm{~T}) \times \overline{\mathrm{Q}}) \mid \phi(\mathrm{t}, 0)=\phi_{\chi}(\mathrm{t}, \mathrm{~L})=0, \quad \forall \mathrm{t} \in(0, \mathrm{~T})\right\} .
$$

We start by deriving an equation for $\left(\hat{v}_{t}^{N}\right)_{t \in[0, T]}$.
Proposition 3.5. For every $\mathrm{N} \geqslant 1$ and $\phi \in \mathcal{C}^{\text {test }}$, it holds that

$$
\int_{0}^{\mathrm{L}} \phi(0, .) \mathrm{d} \hat{v}_{0}^{\mathrm{N}}=\int_{0}^{\mathrm{T}} \int_{0}^{\mathrm{L}}\left(-\phi_{\mathrm{t}}-\frac{\sigma^{2}}{2} \phi_{x x}+\mathrm{q}_{\mathrm{u}, \mathrm{~m}} \phi_{x}\right) \mathrm{d} \hat{v}_{s}^{\mathrm{N}} \mathrm{~d} s+\mathrm{I}_{\mathrm{N}}(\phi) \quad \text { a.s. }
$$

where

$$
\mathrm{I}_{\mathrm{N}}(\phi):=-\frac{\sigma}{\mathrm{N}} \sum_{\mathrm{k}=1}^{\mathrm{N}} \int_{0}^{\mathrm{T}} \phi_{x}\left(\mathrm{~s}, \hat{X}_{\mathrm{s}}^{\mathrm{k}}\right) \mathbb{1}_{\mathrm{s}<\hat{\tau}^{k}} \mathrm{~d} W_{\mathrm{s}}^{\mathrm{k}}
$$

Proof. Let us consider $\phi \in \mathcal{C}^{\text {test }}$. First observe that for any $k \in\{1, \ldots, N\}$, and $t \in[0, T]$

$$
\hat{X}_{t}^{k} \wedge \hat{\tau}^{k}=V_{k}-\int_{0}^{t} \hat{q}_{s}^{k} \mathbb{1}_{s<\tau \hat{q}^{k}} d s+\sigma W_{t \wedge \tau^{\hat{q}^{k}}}^{k}-\xi_{t}^{\hat{X}^{k}}
$$

Hence, for any $k \in\{1, \ldots, N\}$, the random process $\left(\hat{X}_{t \wedge \hat{\tau}^{k}}^{k}\right)_{t \in[0, T]}$ is a continuous semimartingale, and by applying Itô's rule we have:

$$
\begin{aligned}
& \phi\left(T, \hat{X}_{\hat{\tau}^{k}}^{k}\right)-\phi\left(0, V_{k}\right)+\int_{0}^{T} \phi_{x}\left(s, \hat{X}_{s \wedge \hat{\tau}^{k}}^{k}\right) \mathrm{d} \xi_{s}^{\hat{X}^{k}} \\
& =\int_{0}^{T}\left\{\frac{\sigma^{2}}{2} \phi_{x x}\left(s, \hat{X}_{s}^{k}\right)-q_{u, m}\left(s, \hat{X}_{s}^{k}\right) \phi_{x}\left(s, \hat{X}_{s}^{k}\right)\right\} \mathbb{1}_{s<\hat{\tau}^{k}} \mathrm{~d} s \\
& \quad+\int_{0}^{T} \phi_{t}\left(s, \hat{X}_{s \wedge \hat{\tau}^{k}}^{k}\right) \mathrm{d} s+\sigma \int_{0}^{T} \phi_{x}\left(s, \hat{X}_{s}^{k}\right) \mathbb{1}_{s<\hat{\tau}^{k}} d W_{s}^{k} .
\end{aligned}
$$

By using the boundary conditions satisfied by $\phi$, and noting that $\phi_{\mathfrak{t}}(\mathrm{t}, 0)=0$ for any $\mathrm{t} \in(0, \mathrm{~T})$, we deduce that

$$
\begin{aligned}
-\phi\left(0, \mathrm{~V}_{\mathrm{k}}\right)- & \sigma \int_{0}^{T} \phi_{x}\left(\mathrm{~s}, \hat{X}_{s}^{k}\right) \mathbb{1}_{s<\hat{\tau}^{k}} \mathrm{~d} W_{s}^{k} \\
& =\int_{0}^{T}\left\{\phi_{\mathrm{t}}\left(s, \hat{X}_{s}^{k}\right)+\frac{\sigma^{2}}{2} \phi_{x x}\left(s, \hat{X}_{s}^{k}\right)-q_{u, m}\left(s, \hat{X}_{s}^{k}\right) \phi_{x}\left(s, \hat{X}_{s}^{k}\right)\right\} \mathbb{1}_{s<\hat{\tau}^{k}} d s
\end{aligned}
$$

The desired result follows by summing over $k \in\{1, \ldots, N\}$, and multiplying by $N^{-1}$.

By virtue of [30, Theorem 3.2], the tightness of the sequence of laws of $\left(\hat{v}^{\mathrm{N}}\right)_{\mathrm{N} \geqslant 1}$ ensures that this sequence is relatively compact on ( $\mathrm{D}_{\mathcal{S}_{\mathbb{R}}^{\prime}}, \mathrm{M} 1$ ). Consequently, Proposition 3.4 entails the existence of a subsequence (still denoted $\left.\left(\hat{v}^{\mathrm{N}}\right)_{\mathrm{N} \geqslant 1}\right)$ such that

$$
\hat{v}^{\mathrm{N}} \Rightarrow \hat{v}^{*}, \quad \text { on }\left(\mathrm{D}_{\mathcal{S}_{\mathbb{R}}^{\prime}}, \mathrm{M} 1\right)
$$

Thanks to [30, Proposition 2.7 (i)],

$$
\forall \phi \in \mathcal{S}_{\mathbb{R}}, \quad \hat{v}^{\mathrm{N}}(\phi) \Rightarrow \hat{v}^{*}(\phi), \quad \text { as } \mathrm{N} \rightarrow \infty, \quad \text { on }\left(\mathrm{D}_{\mathbb{R}}, \mathrm{M} 1\right) .
$$

To avoid possible confusion about multiple distinct limit points, we will denote $\hat{v}^{*}$ any limiting processes that realizes one of these limiting laws. First, we note that $\hat{v}^{*}$ is a $\tilde{\mathcal{P}}(\overline{\mathrm{Q}})$-valued process:

Proposition 3.6. For every $\mathrm{t} \in[0, \mathrm{~T}], \hat{v}_{\mathrm{t}}^{*}$ is almost surely supported on $\overline{\mathrm{Q}}$ and belongs to $\tilde{\mathcal{P}}(\overline{\mathrm{Q}})$.
Proof. This follows from the "Portmanteau theorem" and the Riesz representation theorem. We omit the details and refer to [23, Proposition 5.3].

Next, we recover the partial differential equation satisfied by the process $\left(\hat{v}_{t}^{*}\right)_{t \in[0, \mathrm{~T}]}$.
Lemma 3.7. For every $\phi \in \mathcal{C}^{\text {test }}$, it holds that

$$
\int_{0}^{\mathrm{L}} \phi(0, .) \mathrm{d} m_{0}+\int_{0}^{\mathrm{T}} \int_{0}^{\mathrm{L}}\left(\phi_{\mathrm{t}}+\frac{\sigma^{2}}{2} \phi_{x x}-\mathrm{q}_{\mathrm{u}, \mathrm{~m}} \phi_{x}\right) \mathrm{d} \hat{v}_{s}^{*} \mathrm{ds}=0 \quad \text { a.s. }
$$

Proof. Let us consider $\phi \in \mathcal{C}^{\text {test }}$ and set:

$$
\mu(\phi):=\int_{0}^{\mathrm{L}} \phi(0, .) \mathrm{d} m_{0}+\int_{0}^{\mathrm{T}} \int_{0}^{\mathrm{L}}\left(\phi_{\mathrm{t}}+\frac{\sigma^{2}}{2} \phi_{x x}-\mathrm{q}_{\mathfrak{u}, \mathfrak{m}} \phi_{x}\right) \mathrm{d} \hat{v}_{s}^{*} \mathrm{ds} ;
$$

and

$$
\mu_{\mathrm{N}}(\phi):=\int_{0}^{\mathrm{L}} \phi(0, .) \mathrm{d} m_{0}+\int_{0}^{\mathrm{T}} \int_{0}^{\mathrm{L}}\left(\phi_{\mathrm{t}}+\frac{\sigma^{2}}{2} \phi_{x x}-\mathrm{q}_{\mathrm{u}, \mathrm{~m}} \phi_{x}\right) \mathrm{d} \hat{\gamma}_{\mathrm{s}}^{\mathrm{N}} \mathrm{ds} .
$$

Owing to Proposition 3.5 we have

$$
\mu_{\mathrm{N}}(\phi)=\mathrm{I}_{\mathrm{N}}(\phi)+\int_{0}^{\mathrm{L}} \phi(0, .) \mathrm{d}\left(\mathrm{~m}_{0}-\hat{v}_{0}^{\mathrm{N}}\right) .
$$

Note that

$$
\mathbb{E} \mathrm{I}_{\mathrm{N}}(\phi)^{2} \leqslant \mathrm{C}\left\|\phi_{\chi}\right\|_{\infty}^{2} \mathrm{~N}^{-1} .
$$

Hence, by appealing to Horowitz-Karandikar inequality (see e.g. [36, Theorem 10.2.1]) we deduce that

$$
\mathbb{E} \mu_{\mathrm{N}}^{2}(\phi) \leqslant \mathrm{C}\left\|\phi_{\chi}\right\|_{\infty}^{2} \mathrm{~N}^{-2 / 5} .
$$

Consequently, to conclude the proof it is enough to show that

$$
\mu_{N}(\phi) \Rightarrow \mu(\phi) \quad \text { as } N \rightarrow \infty .
$$

Let $\mathcal{A}$ be the set of elements in $D_{\mathcal{S}_{\mathbb{R}}^{\prime}}$ that take values in $\tilde{\mathcal{P}}(\overline{\mathrm{Q}})$, and consider a sequence $\left(\psi^{\mathrm{N}}\right) \subset \mathcal{A}$ which converges to some $\psi$ in $\mathcal{A}$ with respect to the M1 topology. Let $\mathbf{q}_{\mathrm{u}, \mathrm{m}}$ be a continuous function on $[0, \mathrm{~T}] \times \mathbb{R}$, which satisfies the following conditions:

$$
\begin{equation*}
\mathbf{q}_{\mathbf{u},\left.\mathfrak{m}\right|_{\bar{Q}_{\mathrm{T}}}} \equiv \mathrm{q}_{\mathbf{u}, \mathrm{m}} ; \quad\left\|\mathbf{q}_{\mathbf{u}, \mathrm{m}}\right\|_{\infty}=\left\|\mathbf{q}_{\mathbf{u}, \mathrm{m}}\right\|_{\infty} ; \quad \forall \mathrm{t} \in[0, \mathrm{~T}], \operatorname{supp} \mathbf{q}_{\mathbf{u}, \mathrm{m}}(\mathrm{t}, .) \subset(-\mathrm{L}, 2 \mathrm{~L}) . \tag{3.9a}
\end{equation*}
$$

We also define the sequence

$$
\begin{equation*}
\mathbf{q}_{\mathbf{u}, \mathfrak{m}}^{\mathrm{n}}(\mathrm{t}, \mathrm{x}):=\left(\mathbf{q}_{\mathrm{u}, \mathrm{~m}}(\mathrm{t}, .) * \xi_{\mathrm{n}}\right)(x), \quad n \geqslant 1, \tag{3.9b}
\end{equation*}
$$

where $\xi_{n}(x):=\mathfrak{n} \xi(n x)$ is a compactly supported mollifier on $\mathbb{R}$.
We have

$$
\begin{aligned}
& J:=\left|\int_{0}^{T} \int_{0}^{L} q_{u, m} \phi_{\chi} d \psi_{s}^{N} d s-\int_{0}^{T} \int_{0}^{L} q_{u, m} \phi_{\chi} d \psi_{s} d s\right| \\
& =\left|\int_{0}^{T} \int_{\mathbb{R}} \mathbf{q}_{u, m} \phi_{\chi} d \psi_{s}^{N} d s-\int_{0}^{T} \int_{\mathbb{R}} \boldsymbol{q}_{u, m} \phi_{\chi} d \psi_{s} d s\right| \\
& \leqslant 2\left\|\phi_{x}\right\|_{\infty}\left\|\mathbf{q}_{\mathbf{u}, \mathrm{m}}^{\mathrm{n}}-\mathbf{q}_{\mathrm{u}, \mathrm{~m}}\right\|_{\infty} \\
& +\left|\int_{0}^{T} \int_{\mathbb{R}} \mathbf{q}_{u, m}^{n} \phi_{x} \mathrm{~d}\left(\psi_{\mathrm{s}}^{\mathrm{N}}-\psi_{\mathrm{s}}\right) \mathrm{ds}\right|=: \mathrm{J}_{1}+\mathrm{J}_{2} .
\end{aligned}
$$

Since $\mathbf{q}_{\mathfrak{u}, \mathfrak{m}}^{\mathfrak{n}}(s,.) \phi_{\chi}(s,.) \in \mathcal{S}_{\mathbb{R}}$ for any $s \in[0, T]$, then $J_{2}$ vanishes as $\psi^{N} \rightarrow \psi$. On the other hand, note that $J_{1}$ also vanishes as $n \rightarrow+\infty$ so that we obtain $\lim _{N} J=0$. Moreover, one easily checks that

$$
\int_{0}^{T} \int_{0}^{\mathrm{L}} \mathrm{Fd} \psi_{s}^{\mathrm{N}} \mathrm{~d} s \rightarrow \int_{0}^{\mathrm{T}} \int_{0}^{\mathrm{L}} \mathrm{Fd} \psi_{s} \mathrm{~d} s, \quad \mathrm{~F} \equiv \phi_{\mathrm{t}}, \phi_{x x} \text { as } \quad \mathrm{N} \rightarrow+\infty .
$$

Therefore, by virtue of the continuous mapping theorem, we obtain that $\mu_{N}(\phi) \Rightarrow \mu(\phi)$, which concludes the proof.

We are now in position to prove Lemma 3.1.
Proof of Lemma 3.1. From Lemma 3.7, we know that $\mathrm{d}^{\circ}{ }^{*}=\mathrm{d} \hat{\mathrm{v}}_{\mathrm{t}}^{*} \mathrm{dt}$ and $\mathrm{dm}=\mathrm{dm}(\mathrm{t}) \mathrm{dt}$ both satisfy (almost surely) the same Fokker-Planck equation in the sense of measures (cf. Appendix B). By invoking the uniqueness of solutions to that equation (cf. Proposition B.1), we deduce that $\hat{v}^{*} \equiv \mathrm{~m}$ almost surely. Since all converging sub-sequences converge weakly toward $m$, we infer that $\hat{v}^{N} \Rightarrow m$, on ( $D_{\mathcal{S}_{\mathbb{R}}^{\prime}}, M 1$ ).
3.2. Proof of Theorem 1.6. We start by collecting the following technical result whose proof is given in Appendix A.
Lemma 3.8. Fix $n \geqslant 1$, define $\mathcal{A}$ to be all elements in $\mathrm{D}_{\mathcal{S}_{\mathbb{R}}^{\prime}}$ that take values in $\tilde{\mathcal{P}}(\overline{\mathrm{Q}})$, and let $\Psi_{m}$ (resp. $\Psi_{\mathbf{q}}^{\mathrm{n}}$ ) be the map defined from $\mathrm{D}_{\mathcal{S}_{\mathbb{R}}^{\prime}}$ into $\mathrm{D}_{\mathcal{S}_{\mathbb{R}}}$ (resp. from $\mathcal{A}$ into $\mathrm{D}_{\mathbb{R}}$ ) such that

$$
\Psi_{\mathfrak{m}}(v)(\mathrm{t}):=v(\mathrm{t})-\mathfrak{m}(\mathrm{t}) \quad \text { and } \quad \Psi_{\mathbf{q}}^{\mathfrak{n}}(v)(\mathrm{t}):=\left|\int_{\mathbb{R}} \mathbf{q}_{\mathfrak{u}, \mathfrak{m}}^{\mathfrak{n}}(\mathrm{t}, .) \mathrm{d} v(\mathrm{t})\right| .
$$

Then $\Psi_{\mathbf{m}}, \Psi_{\mathbf{q}}^{\mathrm{q}}$ are continuous with respect to the M1 topology.
Let us now explain the proof of Theorem 1.6. We shall proceed by contradiction, assuming that (1.10) does not hold. Then there exists $\varepsilon_{0}>0$, a sequence of integers $N_{k}$ such that $\lim _{k} N_{k}=+\infty$, and sequences $\left(i_{k}\right) \subset\left\{1, \ldots, N_{k}\right\},\left(q^{i_{k}}\right) \subset \mathbb{A}_{c}$, such that

$$
\begin{equation*}
\mathcal{J}_{c}^{i_{k}, N_{k}}\left(q^{i_{k}} ;\left(\hat{\mathbf{q}}^{j}\right)_{j \neq i_{k}}\right)>\varepsilon_{0}+\mathcal{J}_{c}^{i_{k}}, N_{k}\left(\hat{\mathbf{q}}^{1}, \ldots, \hat{\mathbf{q}}^{N}\right), \quad \forall k \geqslant 0 . \tag{3.10}
\end{equation*}
$$

We derive a contradiction by estimating the difference between $\mathcal{\partial}_{c}^{i_{k}, N_{k}}$ and the mean field objective $\mathcal{J}_{c}$, which we recall from Section 1.2:

$$
\begin{equation*}
\mathcal{J}_{c}(\rho):=\mathbb{E}\left\{\int_{0}^{T} e^{-r s}\left(1-\kappa \bar{q}_{s}-\rho_{s}\right) \rho_{s} \mathbb{1}_{s<\tau^{\rho}} d s+e^{-r T} u_{T}\left(X_{T}^{\rho}\right)\right\}, \bar{q}=\int_{0}^{L} q_{u, m} d m \tag{3.11}
\end{equation*}
$$

where

$$
\mathrm{d} X_{\mathrm{t}}^{\rho}=-\rho_{\mathrm{t}} \mathbb{1}_{\mathrm{t}<\tau^{\rho}} \mathrm{dt}+\sigma \mathbb{1}_{\mathrm{t}<\tau^{\rho}} \mathrm{d} W_{\mathrm{t}}-\mathrm{d} \xi_{\mathrm{t}}^{\mathrm{X}^{\rho}}
$$

Using Lemma 3.1, we will show that this difference goes to zero as $\mathrm{N}_{\mathrm{k}} \rightarrow+\infty$.
Let us set for any $k \geqslant 0$,

$$
\left\{\begin{array}{l}
d X_{t}^{i_{k}}:=-q_{t}^{i_{k}} d t+\sigma d W_{t}^{i_{k}}-d \xi_{t}^{x_{k}^{i_{k}}}, X_{0}^{i_{k}}=V_{i_{k}} \\
\tau^{i_{k}}:=\inf \left\{t \geqslant 0: X_{t}^{i_{k}} \leqslant 0\right\} \wedge T
\end{array}\right.
$$

and define

$$
z_{1, \mathrm{~T}}^{k}:=\int_{0}^{\mathrm{T}} \mathfrak{q}_{s}^{i_{k}} \mathbb{1}_{s<\tau^{i_{k}}} d s, \quad \text { and } \quad z_{2, \mathrm{~T}}^{\mathrm{k}}:=\int_{0}^{\mathrm{T}}\left|q_{s}^{i_{k}}\right|^{2} \mathbb{1}_{s<\tau^{i_{k}}} d s .
$$

Recall that all elements of $\mathbb{A}_{c}$ are non-negative, so that $z_{1, T}^{k} \geqslant 0$ for any $k \geqslant 0$. We start by collecting estimates on $\left(z_{1, T}^{i_{k}}\right)_{k \geqslant 0}$ and $\left(z_{2, T}^{i_{k}}\right)_{k \geqslant 0}$. Observe that for any $t \in[0, T]$,

$$
X_{t \wedge \tau^{i_{k}}}^{i_{k}}=V_{k}-\int_{0}^{t} q_{s}^{i_{k}} \mathbb{1}_{s<\tau^{i_{k}}} d s+\sigma W_{t \wedge \tau^{i_{k}}}^{i_{k}}-\xi_{t}^{X^{i_{k}}}, \quad \forall k \geqslant 0
$$

Since the local time is nondecreasing, we infer that

$$
0 \leqslant z_{1, T}^{k} \leqslant V_{i_{k}}-X_{\tau^{i_{k}}}^{i_{k}}+\sigma W_{\tau^{i_{k}}}^{i_{k}}, \quad \forall k \geqslant 0
$$

holds almost surely. By means of the optional stopping theorem, we deduce that

$$
\begin{equation*}
\sup _{k \geqslant 0} \mathbb{E}\left[\mathcal{Z}_{1, \mathrm{~T}}^{\mathrm{k}}\right] \leqslant \mathrm{L} \tag{3.12}
\end{equation*}
$$

Moreover, recall that

$$
\mathcal{J}_{c}^{i_{k}, N_{k}}\left(q^{i_{k}} ;\left(\hat{\mathbf{q}}^{j}\right)_{j \neq i_{k}}\right)=\mathbb{E}\left\{\int_{0}^{T} e^{-r s}\left(1-\kappa \overline{\hat{q}}_{s}^{i_{k}}-q_{s}^{i_{k}}\right) q_{s}^{i_{k}} \mathbb{1}_{s<\tau^{i_{k}}} d s+e^{-r T} u_{T}\left(X_{\tau^{i_{k}}}^{i_{k}}\right)\right\},
$$

where for any $k \geqslant 0$

$$
\overline{\hat{\mathbf{q}}}_{s}^{\mathfrak{i}_{\mathrm{k}}}=\frac{1}{\mathrm{~N}_{\mathrm{k}}-1} \sum_{\mathfrak{j} \neq \mathfrak{i}_{k}} \mathrm{q}_{\mathrm{u}, \mathrm{~m}}\left(\mathrm{~s}, \hat{X}_{s}^{\mathrm{j}}\right) \mathbb{1}_{\mathrm{s}<\hat{\tau}^{j}} .
$$

Thus, for any $k \geqslant 0$

$$
e^{-r T} \mathbb{E}\left[z_{2, \mathrm{~T}}^{k}\right] \leqslant\left\|u_{T}\right\|_{\infty}+\mathbb{E}\left\{\int_{0}^{T} e^{-r s}\left|1-\kappa \overline{\hat{q}}_{s} \mathfrak{i}_{k}\right| q_{s}^{i_{k}} \mathbb{1}_{s<\tau_{k}} d s\right\}-\partial_{c}^{i_{k}, N_{k}}\left(q^{i_{k}} ;\left(\hat{q}^{j}\right)_{j \neq i_{k}}\right) .
$$

By virtue of (3.10) and the uniform bound on $\mathrm{q}_{\mathrm{u}, \mathrm{m}}$ that is given in (2.22), we deduce that

$$
e^{-r \mathrm{~T}} \mathbb{E}\left[\mathcal{z}_{2, \mathrm{~T}}^{\mathrm{k}}\right] \leqslant 2\left\|u_{\mathrm{T}}\right\|_{\infty}+(\kappa+1) \sup _{\mathrm{k} \geqslant 0} \mathbb{E}\left[\mathcal{Z}_{1, \mathrm{~T}}^{\mathrm{k}}\right]+\mathrm{C}(\kappa, \mathrm{~T}),
$$

so that

$$
\begin{equation*}
\sup _{k \geqslant 0} \mathbb{E}\left[\mathcal{Z}_{2, \mathrm{~T}}^{\mathrm{k}}\right] \leqslant \mathrm{C}\left(\mathrm{~T}, \kappa,\left\|u_{\mathrm{T}}\right\|_{\infty}, \mathrm{L}\right) \tag{3.13}
\end{equation*}
$$

On the other hand, we have for any $k \geqslant 0$,

$$
\begin{aligned}
& \mathcal{J}_{c}^{\mathfrak{i}_{k}, N_{k}}\left(\mathfrak{q}^{i_{k}} ;\left(\hat{q}^{j}\right)_{\mathfrak{j} \neq i_{k}}\right) \\
& \leqslant \mathbb{E}\left\{\int_{0}^{T} e^{-r s}\left(1-k \int_{0}^{L} q_{u, m}(s, .) d \hat{v}_{s}^{N_{k}}-q_{s}^{i_{k}}\right) q_{s}^{i_{k}} \mathbb{1}_{s<\tau^{i_{k}}} d s+e^{-r T} u_{T}\left(X_{\tau^{i_{k}}}^{i_{k}}\right)\right\} \\
& +
\end{aligned}
$$

Thus, for any $k \geqslant 0$

$$
\begin{aligned}
& \mathcal{J}_{c}^{i_{k}, N_{k}}\left(q^{i_{k}} ;\left(\hat{q}^{j}\right)_{j \neq i_{k}}\right)-\mathcal{J}_{c}\left(q^{i_{k}}\right)-C N_{k}^{-1} \\
& \leqslant \kappa \mathbb{E}\left[\int_{0}^{T} e^{-r s} q_{s}^{i_{k}} \mathbb{1}_{s<\tau^{i^{i}}}\left|\int_{\mathbb{R}} \mathbf{q}_{u, m}(s, .) \mathrm{d}\left(\mathfrak{m}(s)-\hat{v}_{s}^{N_{k}}\right)\right| \mathrm{ds}\right] \\
& \leqslant \kappa \mathbb{E}\left[\int_{0}^{T} e^{-r s} q_{s}^{i_{k}} \mathbb{1}_{s<\tau^{q^{i}}}\left|\int_{\mathbb{R}} \mathbf{q}_{u, m}^{n}(s, .) d\left(m(s)-\hat{v}_{s}^{N_{k}}\right)\right| d s\right] \\
& +\kappa \mathbb{E}\left[\int_{0}^{T} e^{-r s} q_{s}^{i_{k}} \mathbb{1}_{s<\tau^{q^{i}}} d s\right]\left\|\mathbf{q}_{\mathbf{u}, \mathfrak{m}}^{n}-\mathbf{q}_{\mathbf{u}, m}\right\|_{\infty},
\end{aligned}
$$

where $\mathcal{J}_{\mathfrak{c}}$ is given by (3.11) and $\mathbf{q}_{u, m}, \mathbf{q}_{\mathfrak{u}, m}^{n}$ are given by (3.9a)-(3.9b).
Let us fix $\varepsilon>0$. Since $\left(\mathbf{q}_{u, m}^{n}\right)_{n \geqslant 1}$ converges uniformly toward $\mathbf{q}_{u, m}$ on $[0, T] \times \mathbb{R}$, we can choose $n$ large enough and independently of $k \geqslant 0$ so that

$$
\begin{align*}
& \mathcal{J}_{c}^{i_{k}, N_{k}}\left(q^{i_{k}} ;\left(\hat{q}^{j}\right)_{j \neq i_{k}}\right)-\mathcal{J}_{c}\left(q^{i_{k}}\right)  \tag{3.14}\\
\leqslant & \kappa \mathbb{E}\left[z_{2, T}^{k}\right]^{1 / 2} \mathbb{E}\left[\int_{0}^{T}\left|\int_{\mathbb{R}} q_{u, m}^{n}(s, .) d\left(\hat{v}_{s}^{N_{k}}-m(s)\right)\right|^{2} d s\right]^{1 / 2}+\kappa \varepsilon \mathbb{E}\left[z_{1, T}^{k}\right]+C N_{k}^{-1} .
\end{align*}
$$

Appealing to Lemma 3.1, Lemma 3.8 and the continuous mapping theorem we have

$$
\begin{equation*}
\lim _{N} \mathbb{E}\left[\int_{0}^{\mathrm{T}}\left|\int_{\mathbb{R}} \mathbf{q}_{\mathfrak{u}, \mathfrak{m}}^{\mathrm{n}}(\mathrm{~s}, .) \mathrm{d}\left(\hat{v}_{\mathrm{s}}^{\mathrm{N}_{k}}-\mathfrak{m}(\mathrm{s})\right)\right|^{2} \mathrm{ds}\right]=0 \tag{3.15}
\end{equation*}
$$

Thus, by combining (3.12), (3.13), and (3.14):

$$
\mathcal{J}_{c}^{i_{k}, N_{k}}\left(q^{i_{k}} ;\left(\hat{q}^{j}\right)_{j \neq i_{k}}\right)-\mathcal{J}_{\mathcal{c}}\left(\mathfrak{q}^{i_{k}}\right) \leqslant C\left(T, \kappa,\left\|u_{T}\right\|_{\infty}, L\right) \varepsilon
$$

for big enough $k \geqslant 0$. Whence, by means of Lemma 1.4:

$$
\mathcal{J}_{c}^{i_{k}, N_{k}}\left(q^{i_{k}} ;\left(\hat{q}^{j}\right)_{j \neq i_{k}}\right) \leqslant C\left(T, \kappa,\left\|u_{T}\right\|_{\infty}, L\right) \varepsilon+\mathcal{J}_{c}\left(\rho^{*}\right)
$$

for big enough $k \geqslant 0$. In the same manner, one can show that

$$
\mathcal{J}_{c}\left(\rho^{*}\right) \leqslant C \varepsilon+\mathcal{I}_{c}^{i_{k}, N_{k}}\left(\hat{q}^{1}, \ldots, \hat{q}^{N}\right)
$$

holds for big enough $k \geqslant 0$. Hence, going back to (3.10) and using the above estimates, we obtain

$$
\varepsilon_{0}<\mathrm{C}\left(\mathrm{~T}, \kappa,\left\|u_{\mathrm{T}}\right\|_{\infty}, \mathrm{L}\right) \varepsilon .
$$

We deduce the desired contradiction by choosing $\varepsilon$ suitably small.

## Appendix A. Proofs of some elementary or technical results

We start by giving a proof to Lemma 1.4.
Proof of Lemma 1.4. This kind of verification results is standard: one checks that the candidate optimal control is indeed the maximum using the equation satisfied by $\mathfrak{u}$; which is the value function. Let $\rho$ be an admissible control ( $\mathbb{F}$-adapted and satisfying the constraints). Since the local time is monotone, then $X^{\rho}$ is a semimartingale and with the use of Itô's rule we obtain

$$
\begin{aligned}
& \mathbb{E}\left[e^{-r T} u_{T}\left(X_{\tau^{\rho}}^{\rho}\right)\right]= \\
& \mathbb{E}\left[u\left(0, X_{0}^{\rho}\right)+\int_{0}^{\tau^{\rho}} e^{-r s}\{ \right.
\end{aligned} \begin{aligned}
& \left.\left.u_{t}\left(s, X_{s}^{\rho}\right)-r u\left(s, X_{s}^{\rho}\right)-\rho_{s} u_{x}\left(s, X_{s}^{\rho}\right)+\frac{\sigma^{2}}{2} u_{x x}\left(s, X_{s}^{\rho}\right)\right\} d s\right] \\
& \\
& =\mathbb{E}\left[u\left(0, X_{0}^{\rho}\right)-\int_{0}^{\tau^{\rho}} e^{-r s}\left\{q_{u, m}^{2}\left(s, X_{s}^{\rho}\right)+\rho_{s} u_{x}\left(s, X_{s}^{\rho}\right)\right\} d s\right],
\end{aligned}
$$

where we have used the boundary value problem satisfied by $u$ and the fact that $u_{t}, u_{x}, u_{x x}$ are continuous on $(0, \mathrm{~T}) \times(0, \mathrm{~L}]$ (cf. (2.34)).

By using definition (1.7), note that

$$
\mathbf{q}_{\mathbf{u}, \mathfrak{m}}^{2}=\frac{1}{4}\left|\left(1-\kappa \bar{q}-u_{x}\right) \vee 0\right|^{2}=\sup _{\rho \geqslant 0} \rho\left(1-\kappa \bar{q}-\rho-u_{x}\right)=q_{u, m}\left(1-\kappa \bar{q}-q_{u}, \mathfrak{m}-u_{x}\right) .
$$

Therefore

$$
\mathbb{E}\left[e^{-r \mathrm{r}} \mathfrak{u}_{\mathrm{T}}\left(X_{\tau^{\rho}}^{\rho}\right)\right] \leqslant \mathbb{E}\left[\mathfrak{u}\left(0, X_{0}^{\rho}\right)-\int_{0}^{\tau^{\rho}} e^{-\mathrm{rs}} \rho_{s}\left(1-\kappa \bar{q}-\rho_{s}\right) \mathrm{ds}\right],
$$

so that

$$
\int_{0}^{\mathrm{L}} u(0, .) d m_{0}=\mathbb{E}\left[u\left(0, X_{0}^{\rho}\right)\right] \geqslant \mathbb{E}\left[\int_{0}^{\tau^{\rho}} e^{-r s}\left(1-\kappa \bar{q}-\rho_{s}\right) \rho_{s} d s+e^{-r T} u_{\mathrm{T}}\left(X_{\tau^{\rho}}^{\rho}\right)\right]
$$

By virtue of Lemma 2.5, we know that the process $\left(X_{t}^{\rho^{*}}\right)_{t \in[0, T]}$ exists in the strong sense. Replacing $\rho$ by $\rho^{*}$ in the above computations, inequalities become equalities and we easily infer that

$$
\mathcal{J}_{\mathfrak{c}}\left(\rho^{*}\right)=\int_{0}^{\mathrm{L}} \mathfrak{u}(0, .) \mathrm{dm}_{0} .
$$

Thus (1.8) is proved.
Next, we give a proof to Lemma 3.8.

Proof of Lemma 3.8. Throughout the proof, we shall use notations of $[30,38]$.
Step 1 (continuity in $\mathcal{S}_{\mathbb{R}}^{\prime}$ ): By virtue of Theorem 1.5 , we know that $t \rightarrow \mathfrak{m}(t)$ is continuous on $[0, \mathrm{~T}]$ with respect to the strong topology of $\mathcal{S}_{\mathbb{R}}^{\prime}$. Let $\phi \in \mathcal{S}_{\mathbb{R}}^{\prime}$, we aim to compute the modulus of continuity of $\mathrm{t} \rightarrow \int_{\mathbb{R}} \phi \mathrm{dm}(\mathrm{t})$. For this, we shall appeal to the probabilistic characterization (2.10), thanks to Remark 2.4. We have for any $h>0$

$$
\begin{align*}
\left|\int_{\mathbb{R}} \phi \mathrm{d}(\mathfrak{m}(\mathrm{t}+\mathrm{h})-\mathfrak{m}(\mathrm{t}))\right| & \leqslant \mathbb{E}\left|\phi\left(\mathrm{X}_{\mathrm{t}+\mathrm{h}}\right) \mathbb{1}_{\mathrm{t}+\mathrm{h}<\tau}-\phi\left(\mathrm{X}_{\mathrm{t}}\right) \mathbb{1}_{\mathrm{t}<\tau}\right|  \tag{A.1}\\
& \leqslant \mathrm{C}\|\phi\|_{\mathrm{e}^{1}}\left(\mathbb{P}(\mathrm{t}<\tau)-\mathbb{P}(\mathrm{t}+\mathrm{h}<\tau)+\mathbb{E}\left|X_{\mathrm{t}+\mathrm{h}}-\mathrm{X}_{\mathrm{t}}\right|\right)
\end{align*}
$$

Following the same steps as for (2.9a)-(2.9c), and using Burkholder-Davis-Gundy inequality, we obtain for small enough $h>0$

$$
\left|\int_{\mathbb{R}} \phi \mathrm{d}(\mathfrak{m}(\mathrm{t}+\mathrm{h})-\mathfrak{m}(\mathrm{t}))\right| \leqslant \mathrm{C}\|\phi\|_{\mathcal{C}^{1}} \omega_{\mathfrak{m}}(\mathrm{h})
$$

where

$$
\omega_{\mathfrak{m}}(h):=h^{1 / 2}+\left(\log (1 / h)-h^{1 / 2}\right)^{-1}+\sup _{s \in[0, \mathrm{~T}]} \int_{0}^{\mathrm{L}}\left(1-\phi_{h^{1 / 2}} \log (1 / \mathrm{h})(x)\right) \mathfrak{m}(s, x) d x
$$

and $\phi_{\epsilon}$ is the cut-off function defined in (2.12). In order to get $\lim _{h \rightarrow 0^{+}} \omega_{m}(h)=0$, we need to prove that

$$
\lim _{h \rightarrow 0^{+}} \sup _{s \in[0, \mathrm{~T}]} \int_{0}^{\mathrm{L}}\left(1-\phi_{\mathrm{h}^{1 / 2}} \log (1 / \mathrm{h})(x)\right) \mathfrak{m}(s, x) \mathrm{dx}=0 .
$$

This ensues easily from Dini's Lemma, by choosing the sequence $\left(\phi_{\epsilon}\right)_{\epsilon>0}$ to be monotonically increasing.

Step 2 (continuity of $\Psi_{m}$ ): Let $\epsilon>0, x, y \in D_{\mathcal{S}_{\mathbb{R}^{\prime}}}$, be any bounded subset of $\mathcal{S}_{\mathbb{R}}$, and $\lambda_{x}:=\left(z_{x}, t_{x}\right), \lambda_{y}:=\left(z_{y}, t_{y}\right)$ be a parametric representations of the graphs of $x$ and $y$ respectively, such that

$$
g_{B}\left(\lambda_{x}, \lambda_{y}\right):=\sup _{s \in[0,1]} p_{B}\left(z_{x}(s)-z_{y}(s)\right) \vee\left|t_{x}(s)-t_{y}(s)\right| \leqslant \epsilon,
$$

where $p_{B}(v):=\sup _{x \in B}|v(x)|$. Note that $\lambda_{x}, \lambda_{y}$ depend on $\epsilon$, but we do not use the subscript $\epsilon$ in order to simplify the notation. We have

$$
\begin{aligned}
g_{B}\left(\lambda_{x}, \lambda_{y}\right) \geqslant \sup _{s \in[0,1]} p_{B}\left(z_{x}(s)\right. & \left.-m\left(t_{x}(s)\right)-z_{y}(s)+\mathfrak{m}\left(t_{y}(s)\right)\right) \vee\left|t_{x}(s)-t_{y}(s)\right| \\
& -\sup _{s \in[0,1]} \max p_{B}\left(\mathfrak{m}\left(t_{x}(s)\right)-\mathfrak{m}\left(t_{y}(s)\right)\right) \vee\left|t_{x}(s)-t_{y}(s)\right| .
\end{aligned}
$$

Since the map $t \rightarrow \mathfrak{m}(t) \in \mathcal{S}_{\mathbb{R}}^{\prime}$ is continuous, observe that

$$
\lambda_{v}^{\prime}: s \rightarrow\left(z_{v}(s)-\mathfrak{m}\left(\mathrm{t}_{v}(\mathrm{~s})\right), \mathrm{t}_{v}(\mathrm{~s})\right), \quad v \equiv x, y
$$

is a parametric representation of the graph

$$
\gamma_{v}^{\prime}:=\left\{(w, \mathrm{t}) \in \mathcal{S}_{\mathbb{R}}^{\prime} \times[0, \mathrm{~T}]: w \in\left[v\left(\mathrm{t}^{-}\right)-\mathfrak{m}(\mathrm{t}), v(\mathrm{t})-\mathfrak{m}(\mathrm{t})\right]\right\}, \quad v \equiv \mathrm{x}, \mathrm{y} .
$$

## Consequently

(A.2)

$$
\begin{aligned}
d_{B, M 1}\left(\Psi_{m}(x), \Psi_{m}(y)\right) \leqslant g_{B}\left(\lambda_{x}, \lambda_{y}\right)+\sup _{s \in[0,1]} & p_{B}\left(m\left(t_{x}(s)\right)-m\left(t_{y}(s)\right)\right) \vee\left|t_{x}(s)-t_{y}(s)\right| \\
& \leqslant 2 \epsilon+\sup _{s \in[0,1]} p_{B}\left(m\left(t_{x}(s)\right)-m\left(t_{y}(s)\right)\right) .
\end{aligned}
$$

Hence, by using the estimation of Step 1, we infer that:

$$
\begin{equation*}
\mathbf{d}_{B, \mathrm{M} 1}\left(\Psi_{m}(x), \Psi_{m}(y)\right) \leqslant C(B) \omega_{m}(\epsilon), \tag{A.3}
\end{equation*}
$$

which in turn implies that $\Psi_{\mathrm{m}}$ is continuous.
Step 3 (continuity of $\Psi_{\mathbf{q}}^{\mathfrak{n}}$ ): Let us fix $n \geqslant 1$. Note that $\mathbf{q}_{\mathfrak{u}, m}^{n}$ maps $[0, T]$ into $\mathcal{S}_{\mathbb{R}}$, and the following holds:

$$
\begin{equation*}
\sup _{t \in[0, T]} \sup _{x \in \mathbb{R}}\left|x^{\alpha} \partial_{\chi}^{\beta} \mathbf{q}_{u, m}^{n}(t, x)\right| \leqslant C(L, \alpha) n^{\beta} \int_{\mathbb{R}}\left|\partial_{\chi}^{\beta} \xi(y)\right| d y, \quad \forall \alpha, \beta \in \mathbb{N} . \tag{A.4}
\end{equation*}
$$

Owing to (A.4), we have $\mathbf{q}_{\mathbf{u}, \mathrm{m}}^{n}([0, T]) \subset B_{n}$, where $B_{n}$ is a bounded subset of $\mathcal{S}_{\mathbb{R}}$. Let $\epsilon>0, x, y \in \mathcal{A}$, and $\lambda_{x}:=\left(z_{x}, t_{x}\right), \lambda_{y}:=\left(z_{y}, t_{y}\right)$ be a parametric representations of the graphs of $x$ and $y$ respectively such that

$$
\mathrm{g}_{\mathrm{B}_{\mathrm{n}}}\left(\lambda_{x}, \lambda_{y}\right) \leqslant \epsilon .
$$

We have

$$
\begin{aligned}
& g_{B_{n}}\left(\lambda_{x}, \lambda_{y}\right) \geqslant \sup _{s \in[0,1]}\left|\int_{0}^{L} \mathbf{q}_{u, m}^{n}\left(t_{x}(s), .\right) d\left(z_{x}(s)-z_{y}(s)\right)\right| V\left|t_{x}(s)-t_{y}(s)\right| \\
& \geqslant \sup _{s \in[0,1]}\left|\int_{0}^{L} \mathbf{q}_{u, m}^{n}\left(t_{x}(s), .\right) d z_{x}(s)-\int_{0}^{L} \mathbf{q}_{u, m}^{n}\left(t_{y}(s), .\right) d z_{y}(s)\right| V\left|t_{x}(s)-t_{y}(s)\right| \\
& -\sup _{s \in[0,1]}\left|\int_{0}^{\mathrm{L}}\left(\mathbf{q}_{\mathfrak{u}, \mathfrak{m}}^{n}\left(\mathrm{t}_{\mathrm{x}}(\mathrm{~s}), .\right)-\mathbf{q}_{\mathbf{u}, \mathfrak{m}}^{\mathrm{n}}\left(\mathrm{t}_{\mathrm{y}}(\mathrm{~s}), .\right)\right) \mathrm{d} z_{\mathrm{y}}(\mathrm{~s})\right| \vee\left|\mathrm{t}_{\mathrm{x}}(\mathrm{~s})-\mathrm{t}_{\mathrm{y}}(\mathrm{~s})\right| .
\end{aligned}
$$

Thus, it holds that

$$
\begin{aligned}
& \sup _{s \in[0,1]}\left|\int_{0}^{L} \mathbf{q}_{\mathfrak{u}, \mathfrak{m}}^{n}\left(t_{x}(s), .\right) d z_{x}(s)-\int_{0}^{L} \mathbf{q}_{u, m}^{n}\left(t_{y}(s), .\right) d z_{y}(s)\right| V\left|t_{x}(s)-t_{y}(s)\right| \\
& \leqslant 2 \epsilon+\sup _{s \in[0,1]}\left|\int_{0}^{\mathrm{L}}\left(\mathbf{q}_{\mathbf{u}, \mathfrak{m}}^{\mathrm{n}}\left(\mathrm{t}_{\mathrm{x}}(\mathrm{~s}), .\right)-\mathbf{q}_{\mathbf{u}, \mathfrak{m}}^{\mathrm{n}}\left(\mathrm{t}_{\mathrm{y}}(\mathrm{~s}), .\right)\right) \mathrm{d} z_{\mathrm{y}}(\mathrm{~s})\right| \leqslant 2 \epsilon+\omega_{2}^{\mathrm{n}}(\epsilon) .
\end{aligned}
$$

where $\omega_{2}^{n}$ is the continuity modulus of $\mathbf{q}_{\mathbf{u}, \mathrm{m}}^{n}$. By noting that

$$
\lambda_{v}^{\prime \prime}: s \rightarrow\left(\int_{0}^{\mathrm{L}} \mathbf{q}_{u, m}^{n}\left(\mathrm{t}_{v}(\mathrm{~s}), .\right) \mathrm{d} z_{v}(\mathrm{~s}), \mathrm{t}_{v}(\mathrm{~s})\right), \quad v \equiv x, y
$$

is a parametric representation of the graph

$$
\gamma_{v}^{\prime \prime}:=\left\{(w, \mathrm{t}) \in \mathcal{S}_{\mathbb{R}}^{\prime} \times[0, \mathrm{~T}]: w \in\left[\int_{0}^{\mathrm{L}} \mathbf{q}_{\mathfrak{u}, \mathfrak{m}}^{\mathrm{n}}(\mathrm{t}, .) \mathrm{d} v\left(\mathrm{t}^{-}\right), \int_{0}^{\mathrm{L}} \mathbf{q}_{\mathfrak{u}, \mathfrak{m}}^{\mathrm{n}}(\mathrm{t}, .) \mathrm{d} v(\mathrm{t})\right]\right\}, \quad v \equiv \mathrm{x}, \mathrm{y},
$$

we deduce that

$$
\mathbf{d}_{\mathrm{M} 1}\left(\Psi_{\mathbf{q}}^{\mathfrak{n}}(x), \Psi_{\mathbf{q}}^{\mathfrak{n}}(y)\right) \leqslant 2 \epsilon+\omega_{2}^{\mathfrak{n}}(\epsilon) .
$$

The proof is complete.

## Appendix B. On uniqueness for solutions of Fokker-Planck equations

In this part, we show that problem (1.2a)-(1.2b) admits at most one weak solution in a wide class of positive Radon measures. We believe that this result is well-known, and we explain the proof for lack of precise reference.

Let us start by generalizing the notion of weak solution that is given in (1.2c). For any $\mathrm{m}_{0} \in \mathcal{P}(\overline{\mathrm{Q}})$, we define a measure-valued weak solution to (1.2a)-(1.2b) to be a measure $\mathbf{m}$ on $\overline{\mathrm{Q}}_{\mathrm{T}}$ of the type

$$
\mathrm{d} \mathbf{m}=\mathrm{dm}(\mathrm{t}) \mathrm{dt},
$$

with $\mathfrak{m}(t) \in \tilde{\mathcal{P}}(\bar{Q})$ for all $t \in[0, T]$, and $t \rightarrow \mathfrak{m}(t, \mathcal{A})$ measurable on $[0, T]$ for any Borel set $A \subset \bar{Q}$; such that

$$
\|\mathrm{b}\|_{\mathrm{L}_{\mathrm{m}}^{2}}^{2}:=\int_{0}^{\mathrm{T}} \int_{0}^{\mathrm{L}}|\mathrm{~b}|^{2} \mathrm{~d} \mathbf{m}<\infty
$$

and

$$
\begin{equation*}
\int_{0}^{\mathrm{T}} \int_{0}^{\mathrm{L}}\left(-\phi_{\mathrm{t}}-\frac{\sigma^{2}}{2} \phi_{x x}+\mathrm{b} \phi_{x}\right) \mathrm{d} \mathbf{m}=\int_{0}^{\mathrm{L}} \phi(0, .) \mathrm{d} m_{0} \tag{B.1}
\end{equation*}
$$

for every $\phi \in \mathcal{C}^{\text {test }}$. We claim that such a solution is unique:
Proposition B.1. There is at most one measure-valued weak solution to (1.2a)-(1.2b).
Proof. Our approach is similar to [35, Section 3.1]. Let $\mathbf{m}$ be a measure-valued weak solution to (1.2a)-(1.2b), and consider the following dual problem:

$$
\left\{\begin{array}{l}
-w_{\mathrm{t}}-\frac{\sigma^{2}}{2} w_{x x}+\mathbf{b} w_{x}=\psi \quad \text { in } \mathrm{Q}_{\mathrm{T}}  \tag{B.2}\\
w(\mathrm{t}, 0)=w_{x}(\mathrm{t}, \mathrm{~L})=0 \quad \text { in }(0, \mathrm{~T}) \\
w(\mathrm{~T}, x)=0 \quad \text { in } \mathrm{Q}
\end{array}\right.
$$

where $\psi, \mathbf{b} \in \mathcal{C}^{\infty}\left(\overline{\mathrm{Q}}_{\mathrm{T}}\right)$. Let $w$ be a smooth solution to (B.2). Since $w^{2}$ is smooth, we have:

$$
\int_{0}^{\mathrm{T}} \int_{0}^{\mathrm{L}}\left\{-\left(w^{2}\right)_{\mathrm{t}}-\frac{\sigma^{2}}{2}\left(w^{2}\right)_{x x}+\mathrm{b}\left(w^{2}\right)_{x}\right\} \mathrm{d} \mathbf{m}=\int_{0}^{\mathrm{L}} w^{2}(0, .) \mathrm{d} \mathrm{~m}_{0} .
$$

By (B.2) we thus have

$$
\int_{0}^{\mathrm{T}} \int_{0}^{\mathrm{L}} w\left(\psi-\mathbf{b} w_{x}\right) \mathrm{d} \mathbf{m}-\frac{\sigma^{2}}{2} \int_{0}^{\mathrm{T}} \int_{0}^{\mathrm{L}}\left|w_{x}\right|^{2} \mathrm{~d} \mathbf{m}+\frac{\sigma^{2}}{2} \int_{0}^{\mathrm{T}} \int_{0}^{\mathrm{L}} \mathrm{~b} w w_{x} \mathrm{~d} \mathbf{m}=\int_{0}^{\mathrm{L}} w^{2}(0, .) \mathrm{d} m_{0},
$$

so that

$$
\frac{\sigma^{2}}{4} \int_{0}^{\mathrm{T}} \int_{0}^{\mathrm{L}}\left|w_{x}\right|^{2} \mathrm{~d} \mathbf{m} \leqslant \mathrm{C}\left(\|w\|_{\infty}^{2} \int_{0}^{\mathrm{T}} \int_{0}^{\mathrm{L}}|\mathrm{~b}-\mathbf{b}|^{2} \mathrm{~d} \mathbf{m}+\|\psi\|_{\infty}\|w\|_{\infty}\right)
$$

Hence, from the maximum principle:

$$
\begin{equation*}
\int_{0}^{\mathrm{T}} \int_{0}^{\mathrm{L}}\left|w_{x}\right|^{2} \mathrm{~d} \mathbf{m} \leqslant \mathrm{C}\|\psi\|_{\infty}^{2}\left(1+\|\mathrm{b}-\mathbf{b}\|_{\mathrm{L}_{\mathrm{m}}^{2}}^{2}\right) \tag{B.3}
\end{equation*}
$$

Now, let $\mathbf{m}_{1}, \mathbf{m}_{2}$ be two measure-valued weak solutions to (1.2a)-(1.2b). We know that

$$
\mathrm{b} \in \mathrm{~L}_{\mathbf{m}_{1}}^{2}\left(\mathrm{Q}_{\mathrm{T}}\right) \cap \mathrm{L}_{\mathbf{m}_{2}}^{2}\left(\mathrm{Q}_{\mathrm{T}}\right) .
$$

Thus, $\mathrm{b} \in \mathrm{L}_{\mathbf{m}}^{2}\left(\mathrm{Q}_{\mathrm{T}}\right)$, where $\mathbf{m}=\mathbf{m}_{1}+\mathbf{m}_{2}$. Let $\mathrm{b}^{\epsilon}$ be a sequence of smooth functions converging to $b$ in $L_{\mathbf{m}}^{2}\left(Q_{T}\right)$. Since $m$ is regular, note that such a sequence exists by density of smooth functions in $L_{\mathbf{m}}^{2}\left(Q_{T}\right)$. The measures $\mathbf{m}_{1}, \mathbf{m}_{2}$ being positive, $\mathrm{b}^{\epsilon}$ converges toward $b$ in $L_{\mathbf{m}_{1}}^{2}\left(Q_{T}\right) \cap L_{\mathbf{m}_{2}}^{2}\left(Q_{T}\right)$ as well. Now, let us consider $w^{\epsilon}$ to be a solution to the dual problem that is obtained by replacing $\mathbf{b}$ by $\mathrm{b}^{\epsilon}$ in (B.2). By using $w^{\epsilon}$ as a test function, we obtain

$$
\begin{equation*}
\int_{0}^{\mathrm{T}} \int_{0}^{\mathrm{L}} \psi \mathrm{~d}\left(\mathbf{m}_{1}-\mathbf{m}_{2}\right)=\int_{0}^{\mathrm{T}} \int_{0}^{\mathrm{L}}\left(\mathrm{~b}-\mathrm{b}^{\epsilon}\right) w_{x}^{\epsilon} \mathrm{d} \mathbf{m}_{2}-\int_{0}^{\mathrm{T}} \int_{0}^{\mathrm{L}}\left(\mathrm{~b}-\mathrm{b}^{\epsilon}\right) w_{x}^{\epsilon} \mathrm{d} \mathbf{m}_{1}=: \mathrm{I}_{2}^{\epsilon}-\mathrm{I}_{1}^{\epsilon} . \tag{B.4}
\end{equation*}
$$

By virtue of (B.3), we have for $\mathbf{j}=1,2$ :

$$
\left\|w_{x}^{\epsilon}\right\|_{\mathrm{L}_{\mathrm{m}_{\mathrm{j}}}^{2}} \leqslant \mathrm{C}\|\psi\|_{\infty}\left(1+\left\|\mathrm{b}-\mathrm{b}^{\epsilon}\right\|_{\mathrm{L}_{\mathrm{m}_{\mathrm{j}}}^{2}}\right) \leqslant \mathrm{C}
$$

so that

$$
\left|\mathrm{I}_{\mathrm{j}}^{\epsilon}\right| \leqslant\left\|w_{x}^{\epsilon}\right\|_{L_{\mathbf{m}_{\mathrm{j}}^{2}}^{2}}\left\|\mathrm{~b}-\mathrm{b}^{\epsilon}\right\|_{\mathrm{L}_{\mathrm{m}_{\mathrm{j}}^{2}}^{2}} \leqslant \mathrm{C}\left\|\mathrm{~b}-\mathrm{b}^{\epsilon}\right\|_{\mathrm{L}_{\mathrm{m}_{\mathrm{j}}}^{2}} \rightarrow 0, \quad \text { as } \epsilon \rightarrow 0 .
$$

Consequently, for any smooth function $\psi$

$$
\int_{0}^{\mathrm{T}} \int_{0}^{\mathrm{L}} \psi \mathrm{~d}\left(\mathbf{m}_{1}-\mathbf{m}_{2}\right)=0,
$$

which entails $\mathbf{m}_{1} \equiv \mathbf{m}_{2}$ and concludes the proof.
[1] R. F. Anderson and S. Orey, Small random perturbation of dynamical systems with reflecting boundary, Nagoya Mathematical Journal 60 (1976), 189-216.
[2] F. Avram and M. S. Taqqu, Probability bounds for M-Skorohod oscillations, Stochastic Processes and their Applications 33 (1989), no. 1.
[3] A. Bensoussan, J. Frehse, and P. Yam, Mean field games and mean field type control theory, Vol. 101, Springer, 2013.
[4] P.-E. Caines, M. Huang, and P. Malhamé, Large-population cost-coupled LQG problems with nonuniform agents: individual-mass behavior and decentralized $\epsilon$-Nash equilibria., IEEE Trans. Automat. Control 52(9) (2007), 1560-1571.
[5] ___, Large population stochastic dynamic games: closed-loop McKean-Vlasov systems and the Nash certainty equivalence principle., Commun. Inf. Syst. 6(3) (2006), 221-251.
[6] P. Cannarsa and R. Capuani, Existence and uniqueness for Mean Field Games with state constraints, PDE Models for Multi-Agent Phenomena, 2018, pp. 49-71.
[7] P. Cannarsa, R. Capuani, and P. Cardaliaguet, $\mathrm{C}^{1,1}$-smoothness of constrained solutions in the calculus of variations with application to mean field games, arXiv preprint arXiv:1806.08966 (2018).
[8] P. Cannarsa, R. Capuani, P. Cardaliaguet, and Piermarco and Capuani Cannarsa Rossana and Cardaliaguet, Mean Field Games with state constraints: from mild to pointwise solutions of the PDE system, arXiv preprint arXiv:1812.11374 (2018).
[9] P. Cardaliaguet, F. Delarue, J.-M. Lasry, and P.-L. Lions, The master equation and the convergence problem in mean field games, arXiv:1509.02505v1 (2015).
[10] P. Cardaliaguet and C.-A. Lehalle, Mean field game of controls and an application to trade crowding, Mathematics and Financial Economics 12 (2018), no. 3, 335-363.
[11] R. Carmona and F. Delarue, Probabilistic analysis of mean-field games, SIAM Journal on Control and Optimization 51 (2013), no. 4, 2705-2734.
[12] $\qquad$ Probabilistic theory of mean-field games with applications I-II, Vol. 83-84, Springer, 2018.
[13] P. Chan and R. Sircar, Bertrand and Cournot mean field games, Applied Mathematics \& Optimization 71 (2015), no. 3, 533-569.
[14]_, Fracking, renewables \& mean field games, SIAM Review 59 (2017), no. 3, 588-615.
[15] M. Freidlin, Functional integration and partial differential equations, Annals of Mathematics Studies, Princeton University Press, 1985.
[16] D. Fudenberg and J. Tirole, Game Theory, MIT Press, 1991.
[17] D. Gilbarg and N. S. Trudinger, Elliptic partial differential equations of second order, Classics in Mathematics, Springer-Verlag, Berlin, 2001.
[18] D. Gomes, S. Patrizi, and V. Voskanyan, On the existence of classical solutions for stationary extended mean field games, Nonlinear Analysis: Theory, Methods \& Applications 99 (2014), 49-79.
[19] D. Gomes and V. Voskanyan, Extended deterministic mean-field games, SIAM Journal on Control and Optimization 54 (2016), no. 2, 1030-1055.
[20] P. J. Graber and C. Mouzouni, Variational mean field games for market competition, arXiv preprint arXiv:1707.07853 (2017).
[21] P. J. Graber and A. Bensoussan, Existence and uniqueness of solutions for Bertrand and Cournot mean field games, Applied Mathematics \& Optimization 77 (2018), no. 1, 47-71.
[22] O. Guéant, J.-M. Lasry, and P.-L. Lions, Mean field games and applications, Paris-princeton lectures on mathematical finance 2010, 2011, pp. 205-266.
[23] B. Hambly and S. Ledger, A stochastic McKean-Vlasov equation for absorbing diffusions on the half-line, The Annals of Applied Probability 27 (2017), no. 5, 2698-2752.
[24] C. Harris, S. Howison, and R. Sircar, Games with exhaustible resources, SIAM Journal on Applied Mathematics 70 (2010), no. 7, 2556-2581.
[25] V. N. Kolokoltsov, J. Li, and W. Yang, Mean field games and nonlinear Markov processes, arXiv preprint arXiv:1112.3744 (2011).
[26] O.A. Ladyzenskaja, V.A. Solonnikov, and N.N. Ural'ceva, Linear and quasilinear equations of parabolic type, Translations of Mathematical Monographs, vol. 23, American Mathematical Society, Providence, R.I., 1967.
[27] J.-M. Lasry and P.-L. Lions, Jeux à champ moyen. I-Le cas stationnaire, Comptes Rendus Mathématique 343 (2006), no. 9, 619-625.
[28] , Jeux à champ moyen. II-Horizon fini et contrôle optimal, Comptes Rendus Mathématique 343 (2006), no. 10, 679-684.
[29] _, Mean field games, Japanese Journal of Mathematics 2 (2007), no. 1, 229-260.
[30] S. Ledger, Skorokhod's M1 topology for distribution-valued processes, Electronic Communications in Probability 21 (2016).
[31] A. Ledvina and R. Sircar, Dynamic Bertrand Oligopoly, Applied Mathematics \& Optimization 63 (2011), no. 1, 11-44.
[32] P.-L. Lions, Cours au Collège de France, www.college-de-france.fr.
[33] M. Ludkovski and X. Yang, Mean Field Game Approach to Production and Exploration of Exhaustible Commodities, arXiv preprint arXiv:1710.05131 (2017).
[34] I. Mitoma, On the sample continuity of $\mathcal{S}^{\prime}$-processes, Journal of Mathematical Society of Japan 35 (1983), 629-636.
[35] A. Porretta, Weak solutions to Fokker-Planck equations and mean field games, Archive for Rational Mechanics and Analysis 216 (2015), no. 1, 1-62.
[36] T. Rachev and L. Rüschendorf, Mass Transportation Problems. Vol I: Theory; Vol II: Applications, Probability and Its Applications, Springer-Verlag New York, 1998.
[37] L. C. G. Rogers and D. Williams, Diffusions, Markov processes, and martingales. Vol 2, Cambridge Mathematical Library. Cambridge University Press, 2000.
[38] W. Whitt, Stochastic-process limits: an introduction to stochastic-process limits and their application to queues, Springer Science \& Business Media, 2002.
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