# Lyapunov criteria for uniform convergence of conditional distributions of absorbed Markov 

## processes

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#### Abstract

We study the quasi-stationary behavior of multidimensional processes absorbed when one of the coordinates vanishes. Our results cover competitive or weakly cooperative Lotka-Volterra birth and death processes and Feller diffusions with competitive Lotka-Volterra interaction. To this aim, we develop original non-linear Lyapunov criteria involving two Lyapunov functions, which apply to general Markov processes.


Keywords: stochastic Lotka-Volterra systems; multitype population dynamics; multidimensional birth and death process; multidimensional Feller diffusions; process absorbed on the boundary; quasi-stationary distribution; uniform exponential mixing property; Lyapunov function

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## 1 Introduction

We consider a Markov process $\left(X_{t}, t \geq 0\right)$ on a topological state space $E \cup \partial$, where $\partial \cap E=\emptyset$ and $\partial$ is absorbing, meaning that $X_{t} \in \partial$ for all $t \geq \tau_{\partial}$, where $\tau_{\partial}=\inf \left\{t \geq 0: X_{t} \in \partial\right\}$. We denote by $\mathbb{P}_{x}$ its law given $X_{0}=x$ and we assume that the process $X$ is absorbed in $\mathbb{P}_{x}$-a.s. finite time $\tau_{\partial}$ in $\partial$ for

[^0]all $x \in E$. A quasi-stationary distribution for $X$ is a probability measure $\nu_{Q S D}$ on $E$ such that
$$
\mathbb{P}_{\nu_{Q S D}}\left(X_{t} \in \cdot \mid t<\tau_{\partial}\right)=\nu_{Q S D}, \quad \forall t \geq 0
$$
where $\mathbb{P}_{\nu}=\int_{E} \mathbb{P}_{x} \nu(d x)$ for all probability measure $\nu$ on $E$.
Our goal is to prove the existence and uniqueness of a quasi-stationary distribution for two models which are very standard in ecology and evolution (cf. e.g. [5, 6, 3] and the references therein). In fact, we prove the uniform convergence in total variation of the law of $X_{t}$ given $X_{t} \notin \partial$ when $t \rightarrow+\infty$ to a unique quasi-stationary distribution. The first model is the multidimensional Lotka-Volterra birth and death process, and the second one the multidimensional competitive Lotka-Volterra (or logistic) Feller diffusion. These two models have received a lot of attention in the one-dimensional case, and a lot is known about their quasi-stationary behaviour [19, 2, 8, so we focus here on the multidimensional case, where the processes evolve on the state spaces $E \cup \partial=\mathbb{Z}_{+}^{d}$ for birth and death processes (with $\mathbb{Z}_{+}=\{0,1,2, \ldots\}$ ) and $E \cup \partial=\mathbb{R}_{+}^{d}$ for diffusion processes, where $d \geq 2$. The case where $\partial=\{(0, \ldots, 0)\}$ can be handled combining the results of the present paper and those known in the one-dimensional case [9, 12, 8] following the methods of [3, Thm. 1.1]. So we focus here on the case where absorption corresponds to the extinction of a single population, instead of the extinction of the whole population. This case corresponds to $\partial=\mathbb{Z}_{+}^{d} \backslash \mathbb{N}^{d}$ and $E=\mathbb{N}^{d}$ (where $\mathbb{N}=\{1,2, \ldots\}$ ) for multidimensional birth and death processes and $\partial=\mathbb{R}_{+}^{d} \backslash(0,+\infty)^{d}$ and $E=(0,+\infty)^{d}$ for multidimensional diffusions.

A general Lotka-Volterra birth and death process in dimension $d \geq 2$ is a Markov process $\left(X_{t}, t \geq 0\right)$ on $\mathbb{Z}_{+}^{d}$ with transition rates $q_{n, m}$ from $\bar{n}=$ $\left(n_{1}, \ldots, n_{d}\right) \in \mathbb{Z}_{+}^{d}$ to $m \neq n$ in $\mathbb{Z}_{+}^{d}$ given by

$$
q_{n, m}= \begin{cases}n_{i}\left(\lambda_{i}+\sum_{j=1}^{d} \gamma_{i j} n_{j}\right) & \text { if } m=n+e_{i}, \text { for some } i \in\{1, \ldots, d\} \\ n_{i}\left(\mu_{i}+\sum_{j=1}^{d} c_{i j} n_{j}\right) & \text { if } m=n-e_{i}, \text { for some } i \in\{1, \ldots, d\} \\ 0 & \text { otherwise }\end{cases}
$$

where $e_{i}=(0, \ldots, 0,1,0, \ldots, 0)$ where the 1 is at the $i$-th coordinate. Note that $q_{n, n-e_{i}}=0$ if $n_{i}=0$, so that the process remains in the state space $\mathbb{Z}_{+}^{d}$. Since in addition $q_{n, m}=0$ for all $n$ such that $n_{i}=0$ and $m$ such that $m_{i} \geq 1$, the set $\partial=\mathbb{Z}_{+}^{d} \backslash \mathbb{N}^{d}$ is absorbing for the process. We make the usual convention that

$$
q_{n, n}:=-q_{n}:=-\sum_{m \neq n} q_{n, m}
$$

From the biological point of view, the constant $\lambda_{i} \geq 0$ is the birth rate per individual of type $i \in\{1, \ldots, d\}$, the constant $\mu_{i} \geq 0$ is the death rate per individual of type $i, c_{i j} \geq 0$ is the rate of death of an individual of type $i$ from competition with an individual of type $j$, and $\gamma_{i j} \geq 0$ is the rate of birth of an individual of type $i$ from cooperation with (or predation of) an individual of type $j$. In general, a Lotka-Volterra process could be explosive if some of the $\gamma_{i j}$ are positive, but the assumptions of the next theorem ensure that it is not the case and that the process is almost surely absorbed in finite time.

Theorem 1.1. Consider a competitive Lotka-Volterra birth and death process $\left(X_{t}, t \geq 0\right)$ in $\mathbb{Z}_{+}^{d}$ as above. Assume that the matrix $\left(c_{i j}-\gamma_{i j}\right)_{1 \leq i, j \leq d}$ defines a positive operator on $\mathbb{R}_{+}^{d}$ in the sense that, for all $\left(x_{1}, \ldots, x_{d}\right) \in$ $\mathbb{R}_{+}^{d} \backslash\{0\}, \sum_{i j} x_{i}\left(c_{i j}-\gamma_{i j}\right) x_{j}>0$. Then the process has a unique quasistationary distribution $\nu_{Q S D}$ and there exist constants $C, \gamma>0$ such that, for all probability measures $\mu$ on $E=\mathbb{N}^{d}$,

$$
\begin{equation*}
\left\|\mathbb{P}_{\mu}\left(X_{t} \in \cdot \mid t<\tau_{\partial}\right)-\nu_{Q S D}\right\|_{T V} \leq C e^{-\lambda t}, \quad \forall t \geq 0 \tag{1.1}
\end{equation*}
$$

where $\|\cdot\|_{T V}$ is the usual total variation distance defined by $\|\mu\|_{T V}=$ $\sup _{f \in L^{\infty}(E),\|f\|_{\infty} \leq 1}|\mu(f)|$.

This result was previously known only in the competitive case (i.e. when $\gamma_{i j}=0$ for all $\left.1 \leq i, j \leq d\right)$ and when the constants $c_{i j}$ satisfy some bounds (see [10]). As explained in the last reference, an important difficulty of this problem is the fact that the absorption rate (i.e. the rate of jump from a state in $E$ to a state in $\partial$ ) is not bounded with respect to the initial distribution. The existence of a quasi-stationary distribution for this kind of multi-dimensional birth and death processes can also be obtained using the theory of positive matrices, as exposed in [14], but the uniform exponential convergence (1.1) also implies uniqueness and several other properties, as recalled below.

A competitive Lotka-Volterra Feller diffusion process in dimension $d \geq 2$ is a Markov process $\left(X_{t}, t \geq 0\right)$ on $\mathbb{R}_{+}^{d}$, where $X_{t}=\left(X_{t}^{1}, \ldots, X_{t}^{d}\right)$, is a solution of the stochastic differential equation

$$
\begin{equation*}
d X_{t}^{i}=\sqrt{\gamma_{i} X_{t}^{i}} d B_{t}^{i}+X_{t}^{i}\left(r_{i}-\sum_{j=1}^{d} c_{i j} X_{t}^{j}\right) d t, \quad \forall i \in\{1, \ldots, d\}, \tag{1.2}
\end{equation*}
$$

where $\left(B_{t}^{1}, t \geq 0\right), \ldots,\left(B_{t}^{d}, t \geq 0\right)$ are independent standard Brownian motions. The Brownian terms and the linear drift terms correspond to classical

Feller diffusions, and the quadratic drift terms correspond to Lotka-Volterra interactions between coordinates of the process. The variances per individual $\gamma_{i}$ is a positive number, and the growth rates per individual $r_{i}$ can be any real number, for all $1 \leq i \leq d$. The competition parameters $c_{i j}$ are assumed nonnegative for all $1 \leq i, j \leq d$, which corresponds to competitive Lotka-Volterra interaction. It is well known that, in this case, there is global strong existence and pathwise uniqueness for the SDE $(1.2)$, and that it is almost surely absorbed in finite time in $\partial=\mathbb{R}_{+}^{d} \backslash(0,+\infty)^{d}$ if $c_{i i}>0$ for all $i \in\{1, \ldots, d\}$ (see [3] and Section 4).
Theorem 1.2. Consider a competitive Lotka-Volterra Feller diffusion $\left(X_{t}, t \geq\right.$ $0)$ in $\mathbb{R}_{+}^{d}$ as above. Assume that $c_{i i}>0$ for all $i \in\{1, \ldots, d\}$. Then the process has a unique quasi-stationary distribution $\nu_{Q S D}$ and there exist constants $C, \gamma>0$ such that, for all probability measures $\mu$ on $E=(0, \infty)^{d}$,

$$
\left\|\mathbb{P}_{\mu}\left(X_{t} \in \cdot \mid t<\tau_{\partial}\right)-\nu_{Q S D}\right\|_{T V} \leq C e^{-\lambda t}, \quad \forall t \geq 0
$$

This results was previously known in dimension 2 when the constants $c_{i j}$ and $\gamma_{i}$ satisfy $c_{12} \gamma_{1}=c_{21} \gamma_{2}$, which implies that the process (after some transformations) is a Kolmogorov diffusion (i.e. of the form $d Y_{t}=d W_{t}-$ $\nabla V\left(Y_{t}\right) d t$ for some Brownian motion $W$ and some $C^{2}$ function $V$, see [3]). Our result is valid in any dimension and has no restriction on the coefficients. One can also expect to extend our result to cooperative cases (e.g. with $c_{21}<0$ and $c_{12}<0$ as in [3]) by using our abstract Lyapunov criterion with Lyapunov functions combining those used to prove Theorems 1.1 and 1.2 .

To prove these two results, we make use of necessary and sufficient conditions for the convergence (1.1) obtained in 9]. These conditions are hard to check in practice (see e.g. [8, 4, 10]), so we first obtain in Section 2 several criteria based on generalized Lyapunov functions to check these conditions. These criteria apply to general Markov processes, not necessarily of the form of competitive Lotka-Volterra birth and death or diffusion processes. Our simplest Lyapunov criterion involves two bounded nonnegative functions $V$ and $\gamma$ such that $V(x) / \varphi(x) \rightarrow+\infty$ out of compact subsets of $E$ such that

$$
-L \varphi \leq C \mathbb{1}_{K}
$$

for some compact subset $K$ of $E$ and

$$
L V+a \frac{V^{1+\varepsilon}}{\varphi^{\varepsilon}} \leq b \varphi
$$

for some $\varepsilon>0$ and some constants $a>0$ and $b \geq 0$, where $L$ denotes (an extension of) the infinitesimal generator of the Markov process $X$. Together
with additional conditions, the existence of such Lyapunov functions $(V, \varphi)$ implies the convergence (1.1).

Many other properties can be deduced from (1.1). First, it implies the existence of a so-called mortality/extinction plateau: there exists $\lambda_{0}>0$ limit of the extinction rate of the population (see [20]). The constant $-\lambda_{0}$ is actually the largest non-trivial eigenvalue of the generator $L$ and satisfies

$$
\mathbb{P}_{\alpha}\left(t<\tau_{\partial}\right)=e^{-\lambda_{0} t}, \quad \forall t \geq 0 .
$$

In addition, (1.1) implies that $x \mapsto e^{\lambda_{0} t} \mathbb{P}_{x}\left(t<\tau_{\partial}\right)$ converges uniformly to a function $\eta: E \rightarrow(0,+\infty)$ when $t \rightarrow+\infty$ [11, Theorem 2.1]. Moreover, $\eta$ is the eigenfunction of $L$ corresponding to the eigenvalue $-\lambda_{0}$ [7, Prop. 2.3]. It also implies the existence and the exponential ergodicity of the associated $Q$-process, defined as the process $X$ conditionned to never be extinct (see [7, Thm. 3.1] for a precise definition). The convergence of the conditional laws of $X$ to the law of the $Q$-process holds also uniformly in total variation norm [11, Theorem 2.1], which entails conditional ergodic properties [11, Corollary 2.2].

Sections 3 and 4 are devoted to the study of (extensions of) competitive Lotka-Volterra birth and death processes and competitive Lotka-Volterra Feller diffusions and the proofs of Theorems 1.1 and 1.2 .

## 2 General Lyapunov citeria for exponential convergence of conditional distributions

In order to obtain results as general as possible, we consider a general framework inspired from [21]. We first define this framework in Subsection 2.1, and we then state and prove a first version of our Lyapunov criterion in Subsection 2.2 and a second (simpler) version in Subsection 2.3 .

### 2.1 Definitions and notations

We consider a right-continuous time-homogeneous Markov process $\left(\Omega, \mathcal{F},\left(\mathcal{F}_{t}\right)_{t \geq 0},\left(X_{t}, t \geq\right.\right.$ 0 ), $\left(\mathbb{P}_{x}\right)_{x \in E \cup\{\partial\}}$ ) (see [22]) with state space $E \cup \partial$ such that $E \cap \partial=\emptyset$ where $\partial$ and $E$ are topological spaces equiped with their Borel $\sigma$-fields. We assume that $\partial$ is an absorbing set for the process, which means that $X_{t} \in \partial$ for all $t \geq \tau_{\partial}$, where

$$
\tau_{\partial}=\inf \left\{t \geq 0: X_{t} \in \partial\right\}
$$

is such that, for all $x \in E$ and $t \geq 0, \mathbb{P}_{x}\left(\tau_{\partial}<+\infty\right)=1$ and $\mathbb{P}_{x}\left(t<\tau_{\partial}\right)>0$.

Throughout the paper, we let $\left\{O_{n}, n \in \mathbb{N}\right\}$ and $\left\{U_{n}, n \in \mathbb{N}\right\}$ denote two fixed increasing families of subsets of $E$ such that

$$
E=\bigcup_{n \in \mathbb{N}} O_{n} \quad \text { and } \quad \operatorname{cl}\left(O_{n}\right) \subset \operatorname{int}\left(U_{n+1}\right) \subset \operatorname{cl}\left(U_{n+1}\right) \subset \operatorname{int}\left(O_{n+1}\right), \forall n \in \mathbb{N}
$$

where $\operatorname{cl}(A)$ and $\operatorname{int}(A)$ denote respectively the closure and the interior in $E$ of $A \subset E$. We assume that, for all $n \in \mathbb{N}$,

$$
T_{n}:=\inf \left\{t \geq 0: X_{t} \notin O_{n}\right\}
$$

is a $\left(\mathcal{F}_{t}\right)_{t \geq 0}$-stopping time and

$$
\tau_{U_{n}}:=\inf \left\{t \geq 0: X_{t} \in U_{n}\right\}
$$

is a $\left(\mathcal{F}_{t}\right)_{t \geq 0}$-stopping time to which the strong Markov property applies for $\left(X_{t}\right)_{t \geq 0}$. For example, the fact that $T_{n}$ and $\tau_{U_{n}}$ are stopping times holds true when $X$ is continuous, $\left(\mathcal{F}_{t}\right)_{t \geq 0}$ is the natural filtration, $O_{n}$ are open sets and $U_{n}$ are closed sets. This is also true if $X$ is only right-continuous, $\left(\mathcal{F}_{t}\right)_{t \geq 0}$ is right-continuous, $O_{n}$ are closed sets and $U_{n}$ are open subsets. However, in the last case, the strong Markov property at time $\tau_{U_{n}}$ is harder to check. This is true for Feller processes, but for general stochastic differential equations, the well-posedness of the martingale problem only implies the strong Markov property on the natural filtration (see e.g. the books of Le Gall [17] or Rogers and Williams [22, 23]).

We assume that the process $\left(X_{t}, t \geq 0\right)$ is regularly absorbed in the sense that

$$
\begin{equation*}
\tau_{\partial}=\lim _{n \rightarrow+\infty} T_{n}, \quad \text { a.s. under } \mathbb{P}_{x}, \forall x \in E \tag{2.1}
\end{equation*}
$$

Thus, $\tau_{\partial}$ is a $\left(\mathcal{F}_{t}\right)_{t \geq 0}$-stopping time.
We also assume that the paths of $\left(X_{t}, t \geq 0\right)$ are càdlàg (right-continuous and admitting left limits) up to time $T_{n}$ for all $n \in \mathbb{N}$. Note that the process $X$ needs not be càdlàg since it may not admit a left limit at time $\tau_{\partial}=\lim _{n} T_{n}$.

We shall make use of the following weakened notion of generator for $X$.
Definition 2.1. We say that a measurable function $V: E \cup \partial \rightarrow \mathbb{R}$ belongs to the domain $\mathcal{D}(L)$ of the weakened generator $L$ of $X$ if there exists a measurable function $W: E \rightarrow \mathbb{R}$ such that, for all $n \in \mathbb{N}, t \geq 0$ and $x \in E$,
$\mathbb{E}_{x} V\left(X_{t \wedge T_{n}}\right)=V(x)+\mathbb{E}_{x}\left[\int_{0}^{t \wedge T_{n}} W\left(X_{s}\right) d s\right]$ and $\mathbb{E}_{x}\left|\int_{0}^{t \wedge T_{n}} W\left(X_{s}\right) d s\right|<\infty$,
and we define $L V=W$ on $E$. We also define $L V(x)=0$ for all $x \in \partial$.

So far, the operator $L$ may be multi-valued, but it is enough for us to define $L V$ as any $W$ satisfying the above property. Of course, this definition extends the usual definition of the extended generator of $X$, and therefore of the usual weak and infinitesimal generators (cf. [13, 21]).

We introduce the following notation.
Definition 2.2. We say that a function $f: E \rightarrow \mathbb{R}$ converges to a limit $\ell$ out of $\left(O_{n}\right)_{n \geq 1}$ and we write

$$
\lim _{x \notin O_{n}, n \rightarrow \infty} f(x)=\ell
$$

if, for all sequence $\left(x_{k}\right)_{k \geq 1}$ in $E$ such that $\left\{k \in \mathbb{N}: x_{k} \in O_{n}\right\}$ is finite for all $n \geq 1$,

$$
\lim _{k \rightarrow \infty} f\left(x_{k}\right)=\ell
$$

Equivalently, this means that

$$
\lim _{n \rightarrow+\infty} \sup _{x \in E \backslash O_{n}}|f(x)-\ell|=0 .
$$

We also define the set of admissible functions to which our Lyapunov criteria apply.

Definition 2.3. We say that a couple $(V, \varphi)$ of functions $V$ and $\varphi$ measurable from $E \cup \partial$ to $\mathbb{R}$ is an admissible couple of functions if
(i) $V$ and $\varphi$ are bounded nonnegative on $E \cup \partial$ and positive on $E, V(x)=$ $\varphi(x)=0$ for all $x \in \partial$ and

$$
\begin{equation*}
\inf _{x \in E} \frac{V(x)}{\varphi(x)}>0 \tag{2.3}
\end{equation*}
$$

(ii) We have the convergences

$$
\begin{equation*}
\lim _{x \notin O_{n}, n \rightarrow \infty} \frac{V(x)}{\varphi(x)}=+\infty \tag{2.4}
\end{equation*}
$$

and

$$
\begin{equation*}
\lim _{n \rightarrow+\infty} V\left(X_{T_{n}}\right)=0 \quad \text { a.s. } \tag{2.5}
\end{equation*}
$$

(iii) $V$ and $\varphi$ belong to the domain of the weakened generator $L$ of $X, L V$ is bounded from above and $L \varphi$ is bounded from below.

Note that since $V\left(X_{t}\right)=V\left(X_{t}\right) \mathbb{1}_{t<\tau_{2}},(2.5)$ is actually equivalent to

$$
\begin{equation*}
\lim _{n \rightarrow+\infty} V\left(X_{t \wedge T_{n}}\right)=V\left(X_{t}\right), \quad \forall t \geq 0, \text { a.s. } \tag{2.6}
\end{equation*}
$$

Moreover, since $V$ is bounded, (2.4) implies that

$$
\lim _{x \notin O_{n}, n \rightarrow \infty} \varphi(x)=0 .
$$

Therefore, using the fact that $X_{T_{n+1}} \notin O_{n}$ for all $n \geq 1$ since $X$ is right continuous, an admissible couple of functions $(V, \varphi)$ also satisfies

$$
\begin{equation*}
\lim _{n \rightarrow+\infty} \varphi\left(X_{t \wedge T_{n}}\right)=\varphi\left(X_{t}\right), \quad \forall t \geq 0, \text { a.s. } \tag{2.7}
\end{equation*}
$$

### 2.2 A non-linear Lyapunov criterion

Given a couple of admissible functions $(V, \varphi)$, our first Lyapunov criterion is based on the following assumption.

Assumption 1. The couple of admissible functions $(V, \varphi)$ satisfies that there exist constants $\varepsilon, A, B>0$ such that, for all probability measure $\mu$ on $E$,

$$
\begin{equation*}
\mu(L V)-\mu(V) \frac{\mu(L \varphi)}{\mu(\varphi)} \leq A \mu(\varphi)-B \frac{\mu(V)^{1+\varepsilon}}{\mu(\varphi)^{\varepsilon}} \tag{2.8}
\end{equation*}
$$

and there exist constants $r_{0}, p_{0}>0$ such that, for $n$ large enough,

$$
\begin{equation*}
\mathbb{P}_{x}\left(r_{0}<\tau_{\partial}\right) \leq p_{0} V(x), \quad \forall x \in E \backslash O_{n} \tag{2.9}
\end{equation*}
$$

Assumption (2.8) can be seen as a nonlinear Lyapunov criterion. We also need to assume a local Doeblin property.

Assumption 2. There exist $k_{0} \in \mathbb{N}, \theta_{0}, \theta_{1}, a_{1}>0$ and a probability measure $\nu$ on $E$ such that, for all $x \in O_{k_{0}}$ and all $s \in\left[\theta_{1}, \theta_{1}+\theta_{0}\right]$,

$$
\mathbb{P}_{x}\left(X_{s} \in \cdot\right) \geq a_{1} \nu
$$

In addition, we assume the following weak irreducibility: for all $n \geq k_{0}$, there exists $s_{n} \geq 0$ such that

$$
D_{n}:=\inf _{x \in O_{n}} \mathbb{P}_{x}\left(X_{s_{n}} \in O_{k_{0}}\right)>0
$$

The following assumption is about the uniform boundedness of exponential moments of the hitting time of the sets $U_{n} \cup \partial$.

Assumption 3. For all $\lambda>0$, there exists $n \geq 1$ such that

$$
\begin{equation*}
\sup _{x \in E} \mathbb{E}_{x}\left(e^{\lambda\left(\tau_{U_{n}} \wedge \tau_{\partial}\right)}\right)<\infty \tag{2.10}
\end{equation*}
$$

Our last assumption is a form of local Harnack inequality, uniform in time.

Assumption 4. For all $n \geq 0$, there exists a constant $C_{n}$ such that, for all $t \geq 0$,

$$
\begin{equation*}
\sup _{x \in O_{n}} \mathbb{P}_{x}\left(t<\tau_{\partial}\right) \leq C_{n} \inf _{x \in O_{n}} \mathbb{P}_{x}\left(t<\tau_{\partial}\right) \tag{2.11}
\end{equation*}
$$

We are now able to state our first result, which will be applied in the subsequent sections.

Theorem 2.4. Assume that the process $\left(X_{t}, t \geq 0\right)$ is regularly absorbed and that there exists a couple of admissible functions $(V, \varphi)$ satisfying Assumption 1. Assume also that Assumptions 2, 3 and 4 are satisfied. Then the process $X$ admits a unique quasi-stationary distribution $\nu_{Q S D}$ and there exist constants $C, \gamma>0$ such that for all probability measure $\mu$ on $E$,

$$
\begin{equation*}
\left\|\mathbb{P}_{\mu}\left(X_{t} \in \cdot \mid t<\tau_{\partial}\right)-\nu_{Q S D}\right\|_{T V} \leq C e^{-\gamma t}, \quad \forall t \geq 0 \tag{2.12}
\end{equation*}
$$

The proof of Theorem 2.4 is based on the next result.
Proposition 2.5. Assume that the process $\left(X_{t}, t \geq 0\right)$ is regularly absorbed and that there exists a couple of admissible functions $(V, \varphi)$ satisfying Assumption 1. Then, for all $x \in E$ and $t \geq 0$,

$$
\begin{equation*}
\frac{\mathbb{E}_{x}\left[V\left(X_{t}\right)\right]}{\mathbb{E}_{x}\left[\varphi\left(X_{t}\right)\right]}=\frac{V(x)}{\varphi(x)}+\int_{0}^{t}\left\{\frac{\mathbb{E}_{x}\left[L V\left(X_{s}\right)\right]}{\mathbb{E}_{x}\left[\varphi\left(X_{s}\right)\right]}-\frac{\mathbb{E}_{x}\left[V\left(X_{s}\right)\right]}{\mathbb{E}_{x}\left[\varphi\left(X_{s}\right)\right]} \frac{\mathbb{E}_{x}\left[L \varphi\left(X_{s}\right)\right]}{\mathbb{E}_{x}\left[\varphi\left(X_{s}\right)\right]}\right\} d s \tag{2.13}
\end{equation*}
$$

where the value of the integral in the r.h.s. is well-defined since, for all $t \geq 0$,

$$
\frac{\mathbb{E}_{x}\left[L V\left(X_{s}\right)\right]}{\mathbb{E}_{x}\left[\varphi\left(X_{s}\right)\right]}-\frac{\mathbb{E}_{x}\left[V\left(X_{s}\right)\right]}{\mathbb{E}_{x}\left[\varphi\left(X_{s}\right)\right]} \frac{\mathbb{E}_{x}\left[L \varphi\left(X_{s}\right)\right]}{\mathbb{E}_{x}\left[\varphi\left(X_{s}\right)\right]} \in L^{1}([0, t])
$$

Proof. Using the Definition 2.1 of the weakened infinitesimal generator, we have for all $n \geq 1$
$\mathbb{E}_{x} V\left(X_{t \wedge T_{n}}\right)=V(x)+\mathbb{E}_{x} \int_{0}^{t \wedge T_{n}} L V\left(X_{s}\right) d s=V(x)+\mathbb{E}_{x} \int_{0}^{t} L V\left(X_{s}\right) \mathbb{1}_{s<T_{n}} d s$.

Applying (2.6) and Lebesgue's theorem to the left-hand side and Fatou's lemma to the right-hand side (using that $L V$ is bounded from above), we obtain

$$
\mathbb{E}_{x} V\left(X_{t}\right) \leq V(x)+\mathbb{E}_{x}\left[\int_{0}^{t \wedge \tau_{\partial}} L V\left(X_{s}\right)\right] d s=V(x)+\int_{0}^{t} \mathbb{E}_{x}\left[L V\left(X_{s}\right)\right] d s
$$

Since $V \geq 0$ and $L V$ is bounded from above, we deduce that $\mathbb{E}_{x} L V\left(X_{s}\right) \in$ $L^{1}([0, t])$ and $L V\left(X_{s}\right) \in L^{1}(\Omega \times[0, t])$. Therefore, we can actually apply Lebesgue's Theorem to the right-hand side of (2.14) and hence

$$
\mathbb{E}_{x} V\left(X_{t}\right)=V(x)+\int_{0}^{t} \mathbb{E}_{x}\left[L V\left(X_{s}\right)\right] d s
$$

The same argument applies to $-\varphi$ using (2.7):

$$
\begin{equation*}
\mathbb{E}_{x} \varphi\left(X_{t}\right)=\varphi(x)+\int_{0}^{t} \mathbb{E}_{x}\left[L \varphi\left(X_{s}\right)\right] d s \tag{2.15}
\end{equation*}
$$

Therefore (cf. e.g. [1, Thm. VIII. 2 and Lem. VIII.2]), for all $T>0, t \mapsto$ $\left(\mathbb{E}_{x} V\left(X_{t}\right), \mathbb{E}_{x} \varphi\left(X_{t}\right)\right)$ belongs to the Sobolev space $W^{1,1}\left([0, T], \mathbb{R}^{2}\right)$ (the set of functions from $[0, T]$ to $\mathbb{R}^{2}$ in $L^{1}$ admitting a derivative in the sense of distributions in $L^{1}$ ).

In particular, $t \mapsto \mathbb{E}_{x} \varphi\left(X_{t}\right)$ is continuous on $[0, T]$. Since $\mathbb{P}_{x}\left(t<\tau_{\partial}\right)>0$ and hence $\mathbb{E}_{x} \varphi\left(X_{t}\right)>0$ for all $t \in[0, T]$, we deduce that $\inf _{t \in[0, T]} \mathbb{E}_{x} \varphi\left(X_{t}\right)>$ 0 . Therefore, we deduce from standard properties of $W^{1,1}$ functions [1, Cor. VIII. 9 and Cor. VIII.10] that $t \mapsto \mathbb{E}_{x} V\left(X_{t}\right) / \mathbb{E}_{x} \varphi\left(X_{t}\right)$ also belongs to $W^{1,1}$ and admits as derivative

$$
t \mapsto \frac{\mathbb{E}_{x} L V\left(X_{t}\right)}{\mathbb{E}_{x} \varphi\left(X_{t}\right)}-\frac{\mathbb{E}_{x} V\left(X_{t}\right) \mathbb{E}_{x} L \varphi\left(X_{t}\right)}{\mathbb{E}_{x}\left[\varphi\left(X_{t}\right)\right]^{2}} \in L^{1}([0, T])
$$

Hence we have proved (2.13).
Before proving Theorem 2.4, let us prove that Assumption2is equivalent to a seemingly stronger assumption.

Lemma 2.6. Assumption ${ }^{2}$ implies that, for all $n_{0} \geq k_{0}$ and for all $t_{0}>0$, there exist $t_{1}, c_{1}>0$ such that, for all $s \in\left[t_{1}, t_{1}+t_{0}\right]$ and $x \in O_{n_{0}}$,

$$
\begin{equation*}
\mathbb{P}_{x}\left(X_{s} \in \cdot\right) \geq c_{1} \nu \tag{2.16}
\end{equation*}
$$

In addition, for all $n \geq n_{0} \geq k_{0}$,

$$
\begin{equation*}
\inf _{x \in O_{n}} \mathbb{P}_{x}\left(X_{s_{n}} \in O_{n_{0}}\right) \geq D_{n}>0 \tag{2.17}
\end{equation*}
$$

Proof. Equation (2.17) is an immediate consequence of the definition of $D_{n}$ and the fact that $O_{k_{0}} \subset O_{n_{0}}$.

To prove (2.16), let $m_{0} \geq k_{0}$ be such that $\nu\left(O_{m_{0}}\right)>0$. Then, for all $x \in O_{m_{0}}$ and all $s \in\left[s_{m_{0}}+\theta_{1}, s_{m_{0}}+\theta_{1}+\theta_{0}\right]$,

$$
\mathbb{P}_{x}\left(X_{s} \in \cdot\right) \geq D_{m_{0}} a_{1} \nu(\cdot)
$$

Hence, for all $\ell \in \mathbb{N}, x \in O_{m_{0}}$ and $s \in\left[\ell\left(s_{m_{0}}+\theta_{1}\right), \ell\left(s_{m_{0}}+\theta_{1}+\theta_{0}\right)\right]$,

$$
\mathbb{P}_{x}\left(X_{s} \in \cdot\right) \geq\left(D_{m_{0}} a_{1}\right)^{\ell} \nu\left(O_{m_{0}}\right)^{\ell-1} \nu(\cdot)
$$

Therefore, for all $n_{0} \geq k_{0}, x \in O_{n_{0}}$ and $s \in\left[s_{n_{0}}+\ell\left(s_{m_{0}}+\theta_{1}\right), s_{n_{0}}+\ell\left(s_{m_{0}}+\right.\right.$ $\left.\left.\theta_{1}+\theta_{0}\right)\right]$,

$$
\mathbb{P}_{x}\left(X_{s} \in \cdot\right) \geq D_{n_{0}}\left(D_{m_{0}} a_{1}\right)^{\ell} \nu\left(O_{m_{0}}\right)^{\ell-1} \nu(\cdot)
$$

This ends the proof of Lemma 2.6.
Proof of Theorem 2.4. The convergence (2.12) can be proved using the necessary and sufficient criterion obtained in [7] for general Markov processes. This criterion is given by the two conditions (A1) and (A2) below.

There exists a probability measure $\nu$ on $E$ such that
(A1) there exist $t_{0}^{\prime}, c_{1}^{\prime}>0$ such that for all $x \in E$,

$$
\mathbb{P}_{x}\left(X_{t_{0}^{\prime}} \in \cdot \mid t_{0}^{\prime}<\tau_{\partial}\right) \geq c_{1}^{\prime} \nu(\cdot)
$$

(A2) there exists $c_{2}^{\prime}>0$ such that for all $x \in E$ and $t \geq 0$,

$$
\mathbb{P}_{\nu}\left(t<\tau_{\partial}\right) \geq c_{2}^{\prime} \mathbb{P}_{x}\left(t<\tau_{\partial}\right) .
$$

The proof is divided in two steps, where (A1) and (A2) are proved respectively.

Step 1: Proof of (A1).
Assumption 2.8) and Proposition 2.5 imply that

$$
\frac{\mathbb{E}_{x} V\left(X_{t}\right)}{\mathbb{E}_{x} \varphi\left(X_{t}\right)} \leq \frac{V(x)}{\varphi(x)}+A t-B \int_{0}^{t}\left(\frac{\mathbb{E}_{x} V\left(X_{s}\right)}{\mathbb{E}_{x} \varphi\left(X_{s}\right)}\right)^{1+\varepsilon} d s
$$

Fix $t>0$ and assume that for all $s \in[0, t], \mathbb{E}_{x} V\left(X_{s}\right) / \mathbb{E}_{x} \varphi\left(X_{s}\right) \geq a$, where $a=(2 A / B)^{1 /(1+\varepsilon)}$. Then

$$
\frac{\mathbb{E}_{x} V\left(X_{t}\right)}{\mathbb{E}_{x} \varphi\left(X_{t}\right)} \leq \frac{V(x)}{\varphi(x)}-\frac{B}{2} \int_{0}^{t}\left(\frac{\mathbb{E}_{x} V\left(X_{s}\right)}{\mathbb{E}_{x} \varphi\left(X_{s}\right)}\right)^{1+\varepsilon} d s
$$

Integrating this differential inequality entails

$$
\frac{\mathbb{E}_{x} V\left(X_{t}\right)}{\mathbb{E}_{x} \varphi\left(X_{t}\right)} \leq\left[\left(\frac{\varphi(x)}{V(x)}\right)^{\varepsilon}+\frac{\varepsilon B t}{2}\right]^{-1 / \varepsilon} \leq\left(\frac{2}{\varepsilon B t}\right)^{1 / \varepsilon}
$$

Choosing $t=t_{0}$ such that $\left(\frac{2}{\varepsilon B t_{0}}\right)^{1 / \varepsilon}<a$ gives a contradiction, hence there exists $u_{x} \in\left[0, t_{0}\right]$ such that $\mathbb{E}_{x} V\left(X_{u_{x}}\right) / \mathbb{E}_{x} \varphi\left(X_{u_{x}}\right)<a$. Note that the value of $u_{x}$ may depend on $x \in E$ but the value of $t_{0}$ is independent of $x \in E$. We also define $n_{0}$ as the first integer $n$ such that

$$
\inf _{y \in E \backslash O_{n}} \frac{V(y)}{\varphi(y)} \geq 2 a
$$

whose existence is ensured by Assumption (2.4).
Now,

$$
\begin{aligned}
a & \geq \frac{\mathbb{E}_{x} V\left(X_{u_{x}}\right)}{\mathbb{E}_{x} \varphi\left(X_{u_{x}}\right)} \\
& \geq \frac{\mathbb{E}_{x}\left[V\left(X_{u_{x}}\right) \mathbb{1}_{X_{u_{x}} \in E \backslash O_{n_{0}}}\right]+\mathbb{E}_{x}\left[V\left(X_{u_{x}}\right) \mathbb{1}_{X_{u_{x}} \in O_{n_{0}}}\right]}{\sup _{y \in E \backslash O_{n_{0}}} \frac{\varphi(y)}{V(y)} \mathbb{E}_{x}\left[V\left(X_{u_{x}}\right) \mathbb{1}_{X_{u_{x}} \in E \backslash O_{n_{0}}}\right]+\sup _{y \in E} \frac{\varphi(y)}{V(y)} \mathbb{E}_{x}\left[V\left(X_{u_{x}}\right) \mathbb{1}_{X_{u_{x}} \in O_{n_{0}}}\right]} \\
& \geq \frac{\mathbb{E}_{x}\left[V\left(X_{u_{x}}\right) \mathbb{1}_{X_{u_{x}} \in E \backslash O_{n_{0}}}\right]+\mathbb{E}_{x}\left[V\left(X_{u_{x}}\right) \mathbb{1}_{\left.X_{u_{x} \in O_{n_{0}}}\right]}^{\frac{1}{2 a} \mathbb{E}_{x}\left[V\left(X_{u_{x}}\right) \mathbb{1}_{X_{u_{x}} \in E \backslash O_{n_{0}}}\right]+\sup _{y \in E} \frac{\varphi(y)}{V(y)} \mathbb{E}_{x}\left[V\left(X_{u_{x}}\right) \mathbb{1}_{X_{u_{x}} \in O_{n_{0}}}\right]} .\right.}{}
\end{aligned}
$$

Since $a>\mathbb{E}_{x} V\left(X_{u_{x}}\right) / \mathbb{E}_{x} \varphi\left(X_{u_{x}}\right) \geq 1 / \sup _{y \in E}(\varphi(y) / V(y))$, we deduce that

$$
\mathbb{E}_{x}\left[V\left(X_{u_{x}}\right) \mathbb{1}_{X_{u_{x}} \in O_{n_{0}}}\right] \geq \frac{1}{2\left(a \sup _{y \in E} \frac{\varphi(y)}{V(y)}-1\right)} \mathbb{E}_{x}\left[V\left(X_{u_{x}}\right) \mathbb{1}_{X_{u_{x}} \in E \backslash O_{n_{0}}}\right]
$$

On the one hand, this last inequality, Assumption (2.9) and Markov's property entails (where $s_{n_{0}}$ is taken from (2.17) of Lemma 2.6)

$$
\begin{aligned}
\mathbb{P}_{x}\left(u_{x}\right. & \left.+s_{n_{0}}\left[\frac{r_{0}}{s_{n_{0}}}\right]<\tau_{\partial}\right) \leq \mathbb{P}_{x}\left(u_{x}+r_{0}<\tau_{\partial}\right) \\
& \leq \mathbb{P}_{x}\left(X_{u_{x}} \in O_{n_{0}}\right)+\mathbb{E}_{x}\left[\mathbb{P}_{X_{u_{x}}}\left(r_{0}<\tau_{\partial}\right) \mathbb{1}_{X_{u_{x}} \in E \backslash O_{n_{0}}}\right] \\
& \leq \mathbb{P}_{x}\left(X_{u_{x}} \in O_{n_{0}}\right)+p_{0} \mathbb{E}_{x}\left[V\left(X_{u_{x}}\right) \mathbb{1}_{X_{u_{x}} \in E \backslash O_{n_{0}}}\right] \\
& \leq \mathbb{P}_{x}\left(X_{u_{x}} \in O_{n_{0}}\right)+2 p_{0}\left(a \sup _{y \in E} \frac{\varphi(y)}{V(y)}-1\right) \mathbb{E}_{x}\left[V\left(X_{u_{x}}\right) \mathbb{1}_{X_{u_{x}} \in O_{n_{0}}}\right] \\
& \leq \mathbb{P}_{x}\left(X_{u_{x}} \in O_{n_{0}}\right)\left[1+2 p_{0}\|V\|_{\infty}\left(a \sup _{y \in E} \frac{\varphi(y)}{V(y)}-1\right)\right]
\end{aligned}
$$

On the other hand, using (2.17) and Markov's property, for all $x \in E$, and all $k \geq 1$,

$$
\begin{aligned}
\mathbb{P}_{x}\left(X_{u_{x}+k s_{n_{0}}} \in O_{n_{0}}\right) & \geq \mathbb{P}_{x}\left(X_{u_{x}} \in O_{n_{0}}\right) \inf _{y \in O_{n_{0}}} \mathbb{P}_{y}\left(X_{k s_{n_{0}}} \in O_{n_{0}}\right) \\
& \geq D_{n_{0}}^{k} \mathbb{P}_{x}\left(X_{u_{x}} \in O_{n_{0}}\right)
\end{aligned}
$$

Therefore, choosing $k=\left\lceil\frac{r_{0}}{s_{n}}\right\rceil$, we have proved that

$$
\begin{aligned}
\mathbb{P}_{x}\left(X_{u_{x}+k s_{n_{0}}} \in O_{n_{0}} \mid u_{x}\right. & \left.+k s_{n_{0}}<\tau_{\partial}\right) \\
& \geq \frac{D_{n_{0}}^{k}}{1+2 p_{0}\|V\|_{\infty}\left(a \sup _{y \in E} \varphi(y) / V(y)-1\right)}=: q
\end{aligned}
$$

Now, since $u_{x} \in\left[0, t_{0}\right]$, Equation 2.16$)$ of Lemma 2.6 entails $\mathbb{P}_{y}\left(X_{t_{1}+t_{0}-u_{x}} \in\right.$ $\Gamma) \geq c_{1} \nu(\Gamma)$ for all $y \in O_{n_{0}}$ and all $\Gamma \subset E$ measurable. Hence, denoting $v_{x}=u_{x}+k s_{n_{0}}$,

$$
\begin{aligned}
& \mathbb{P}_{x}\left(X_{t_{1}+t_{0}+k s_{n_{0}}} \in \Gamma \mid t_{1}+t_{0}+k s_{n_{0}}<\tau_{\partial}\right) \\
&=\frac{\mathbb{E}_{x}\left[\mathbb{1}_{v_{x}<\tau_{\partial}} \mathbb{P}_{X_{v_{x}}}\left(X_{t_{1}+t_{0}-u_{x}} \in \Gamma\right)\right]}{\mathbb{P}_{x}\left(t_{1}+t_{0}+k s_{n_{0}}<\tau_{\partial}\right)} \\
& \geq c_{1} \nu(\Gamma) \frac{\mathbb{P}_{x}\left(X_{v_{x}} \in O_{n_{0}}\right)}{\mathbb{P}_{x}\left(t_{1}+t_{0}+k s_{n_{0}}<\tau_{\partial}\right)} \\
& \geq c_{1} \nu(\Gamma) \mathbb{P}_{x}\left(X_{v_{x}} \in O_{n_{0}} \mid v_{x}<\tau_{\partial}\right) \\
& \geq c_{1} q \nu(\Gamma) .
\end{aligned}
$$

This concludes the proof of (A1), with $c_{1}^{\prime}=c_{1} q$ and $t_{1}^{\prime}=t_{1}+t_{0}+k s_{n_{0}}$.
Step 2: Proof of (A2).
Since $\cup_{n} O_{n}=E$, we can fix $n_{1} \geq n_{0}$ such that $\nu\left(O_{n_{1}}\right) \geq 1 / 2$. Let $m \geq n_{0}$ and $0 \leq s \leq t$, and define $\ell$ the smallest integer such that $s \leq s_{m}+\ell\left(t_{1}^{\prime}+s_{n_{1}}\right)$.

Using (2.17), we obtain, for all $x \in O_{m}$,

$$
\begin{aligned}
& \mathbb{P}_{x}\left(t<\tau_{\partial}\right) \\
& \geq \mathbb{P}_{x}\left(X_{s_{m}} \in O_{n_{0}}, X_{s_{m}+t_{1}^{\prime}} \in O_{n_{1}}, X_{s_{m}+t_{1}^{\prime}+s_{n_{1}}} \in O_{n_{0}}, X_{s_{m}+2 t_{1}^{\prime}+s_{n_{1}}} \in O_{n_{1}},\right. \\
& \left.\quad \ldots, X_{s_{m}+\ell\left(t_{1}^{\prime}+s_{n_{1}}\right)} \in O_{n_{0}}, t-s+s_{m}+\ell\left(t_{1}^{\prime}+s_{n_{1}}\right)<\tau_{\partial}\right) \\
& \geq D_{m}\left(\frac{c_{1}^{\prime}}{2} D_{n_{1}}\right)^{\ell} \inf _{x \in O_{0}} \mathbb{P}_{x}\left(t-s<\tau_{\partial}\right) \\
& \geq D_{m}\left(\frac{c_{1}^{\prime}}{2} D_{n_{1}}\right)^{1+s /\left(t_{1}^{\prime}+s_{n_{1}}\right)} \inf _{x \in O_{m}} \mathbb{P}_{x}\left(t-s<\tau_{\partial}\right) \\
& \geq \frac{c_{1}^{\prime} D_{m} D_{n_{1}}}{2} \inf _{x \in O_{m}} \mathbb{P}_{x}\left(t-s<\tau_{\partial}\right) \exp (-\lambda s),
\end{aligned}
$$

where

$$
\lambda:=\frac{1}{t_{1}^{\prime}+s_{n_{1}}} \log \left(\frac{2}{c_{1}^{\prime} D_{n_{1}}}\right)>0
$$

is independent of $m \geq n_{0}$ and $x \in O_{m}$. Using Assumption 4, we obtain for all $0 \leq s \leq t$

$$
\begin{equation*}
\inf _{x \in O_{m}} \mathbb{P}_{x}\left(t<\tau_{\partial}\right) \geq \frac{c_{1} D_{m} D_{n_{1}}}{2 C_{m}} \sup _{x \in O_{m}} \mathbb{P}_{x}\left(t-s<\tau_{\partial}\right) \exp (-\lambda s) \tag{2.18}
\end{equation*}
$$

Now, we apply Assumption 3 for $\lambda$ as defined above: there exists $m \geq$ $n_{0} \vee n_{1}$ such that

$$
\begin{equation*}
M:=\sup _{x \in E} \mathbb{E}_{x}\left[\exp \left(\lambda\left(\tau_{U_{m}} \wedge \tau_{\partial}\right)\right)\right]<\infty \tag{2.19}
\end{equation*}
$$

Recall that we assumed that the strong Markov property applies at the stopping time $\tau_{U_{m}}$. Note also that $X_{\tau_{U_{m}}} \in O_{m}$ on the event $\left\{\tau_{U_{m}}<\infty\right\}$. Hence, using (2.19) and the strong Markov property,

$$
\begin{aligned}
\mathbb{P}_{x}\left(t<\tau_{\partial}\right) & =\mathbb{P}_{x}\left(t<\tau_{U_{m}} \wedge \tau_{\partial}\right)+\mathbb{P}_{x}\left(\tau_{U_{m}} \leq t<\tau_{\partial}\right) \\
& \leq M e^{-\lambda t}+\int_{0}^{t} \sup _{y \in O_{m}} \mathbb{P}_{y}\left(t-s<\tau_{\partial}\right) \mathbb{P}_{y}\left(\tau_{U_{m}} \in d s\right)
\end{aligned}
$$

Now, for all $s \leq t, \mathbb{P}_{y}\left(\tau_{U_{m}} \in d s\right) \leq \mathbb{P}_{y}\left(\tau_{U_{m}} \wedge \tau_{\partial} \in d s\right)$. Hence, using (2.18)
twice and the fact that $\nu\left(O_{n_{1}}\right) \geq 1 / 2$, we have for all $x \in E$

$$
\begin{aligned}
\mathbb{P}_{x}\left(t<\tau_{\partial}\right) \leq & M e^{-\lambda t}+\int_{0}^{t} \sup _{y \in O_{m}} \mathbb{P}_{y}\left(t-s<\tau_{\partial}\right) \mathbb{P}_{y}\left(\tau_{U_{m}} \wedge \tau_{\partial} \in d s\right) \\
\leq & \frac{2 C_{m}}{c_{1} D_{m} D_{n_{1}}}\left[M \inf _{y \in O_{m}} \mathbb{P}_{y}\left(t<\tau_{\partial}\right)\right. \\
& \left.\quad+\inf _{y \in O_{m}} \mathbb{P}_{y}\left(t<\tau_{\partial}\right) \int_{0}^{t} e^{\lambda s} \mathbb{P}_{x}\left(\tau_{U_{m}} \wedge \tau_{\partial} \in d s\right)\right] \\
\leq & \frac{4 M C_{m}}{c_{1} D_{m} D_{n_{1}}} \inf _{y \in O_{n_{1}}} \mathbb{P}_{y}\left(t<\tau_{\partial}\right) \\
\leq & \frac{8 M C_{m}}{c_{1} D_{m} D_{n_{1}}} \mathbb{P}_{\nu}\left(t<\tau_{\partial}\right) .
\end{aligned}
$$

This is (A2).

### 2.3 A simpler Lyapunov criterion

Condition 2.8 may be hard to check in practice, so we give a stronger condition, easier to check.

Proposition 2.7. Assume that $(V, \varphi)$ is a couple of admissible functions such that $\inf _{x \in O_{m}} \varphi(x)>0$ for all integer $m$. Then Condition (2.8) is implied by the following two conditions:
(a) There exist an integer $n$ and a constant $C \geq 0$ such that

$$
\begin{equation*}
-L \varphi \leq C \mathbb{1}_{O_{n}} \tag{2.20}
\end{equation*}
$$

(b) There exist constants $C^{\prime}>0$ and $C^{\prime \prime} \geq 0$ such that

$$
\begin{equation*}
L V+C^{\prime} \frac{V^{1+\varepsilon}}{\varphi^{\varepsilon}} \leq C^{\prime \prime} \varphi \tag{2.21}
\end{equation*}
$$

Note that Conditions (a) and (b) imply that $L V$ is bounded from above and $L \varphi$ is bounded from below (which are parts of the conditions for a couple of functions to be admissible). This result immediately entails the following corollary, which will be used in the applications provided in Sections 3 and 4 .

Corollary 2.8. Assume that the process $\left(X_{t}, t \geq 0\right)$ is regularly absorbed and that there exists a couple of admissible functions $(V, \varphi)$ satisfying (2.9) and Conditions (a) and (b) of Proposition 2.7 and such that $\inf _{x \in O_{m}} \varphi(x)>0$
for all integer m. Assume also that Assumptions 2, 3 and 4 are satisfied. Then the process $X$ admits a unique quasi-stationary distribution $\nu_{Q S D}$ and there exist constants $C, \gamma>0$ such that, for all probability measure $\mu$ on $E$,

$$
\left\|\mathbb{P}_{\mu}\left(X_{t} \in \cdot \mid t<\tau_{\partial}\right)-\nu_{Q S D}\right\|_{T V} \leq C e^{-\gamma t}, \quad \forall t \geq 0
$$

Proof of Proposition 2.7. Let $\mu$ be any probability measure on $E$. On the one hand, it follows from (2.20) that

$$
-\frac{\mu(L \varphi)}{\mu(\varphi)} \leq \frac{C}{\inf _{x \in O_{n}} \varphi(x)}
$$

On the other hand, Hölder's inequality with $p=(1+\varepsilon) / \varepsilon$ and $q=1+\varepsilon$ yields

$$
\begin{aligned}
\mu(V)^{1+\varepsilon} & \leq \mu\left(\frac{V^{q}}{\varphi^{\varepsilon q /(1+\varepsilon)}}\right)^{(1+\varepsilon) / q} \mu\left(\varphi^{\varepsilon p /(1+\varepsilon)}\right)^{(1+\varepsilon) / p} \\
& =\mu\left(\frac{V^{1+\varepsilon}}{\varphi^{\varepsilon}}\right) \mu(\varphi)^{\varepsilon} .
\end{aligned}
$$

Hence we deduce from (2.21) that

$$
\begin{aligned}
\mu(L V)- & \frac{\mu(L \varphi)}{\mu(\varphi)} \mu(V) \\
& \leq C^{\prime \prime} \mu(\varphi)-C^{\prime} \mu\left(\frac{V^{1+\varepsilon}}{\varphi^{\varepsilon}}\right)+\frac{C}{\inf _{x \in O_{n}} \varphi(x)} \mu(V) \\
& \leq C^{\prime \prime} \mu(\varphi)+\mu\left[V\left(\frac{C}{\inf _{x \in O_{n}} \varphi(x)}-\frac{C^{\prime} V^{\varepsilon}}{2 \varphi^{\varepsilon}}\right)\right]-\frac{C^{\prime}}{2} \frac{\mu(V)^{1+\varepsilon}}{\mu(\varphi)^{\varepsilon}}
\end{aligned}
$$

Now, because of Assumption (2.4), there exists $m$ large enough such that, for all $x \in E \backslash O_{m}$,

$$
\frac{C^{\prime} V^{\varepsilon}(x)}{2 \varphi^{\varepsilon}(x)} \geq \frac{C}{\inf _{y \in O_{n}} \varphi(y)}
$$

Therefore,

$$
\begin{aligned}
& \mu(L V)-\frac{\mu(L \varphi)}{\mu(\varphi)} \mu(V) \\
& \leq C^{\prime \prime} \mu(\varphi)+\frac{C}{\inf _{x \in O_{n}} \varphi(x)} \mu\left(V \mathbb{1}_{O_{m}}\right)-\frac{C^{\prime}}{2} \frac{\mu(V)^{1+\varepsilon}}{2 \mu(\varphi)^{\varepsilon}} \\
& \leq\left(C^{\prime \prime}+\frac{C \sup _{x \in E} V(x)}{\inf _{x \in O_{n}} \varphi(x) \inf _{x \in O_{m}} \varphi(x)}\right) \mu(\varphi)-\frac{C^{\prime}}{2} \frac{\mu(V)^{1+\varepsilon}}{2 \mu(\varphi)^{\varepsilon}}
\end{aligned}
$$

This is (2.8).

## 3 Application to multidimensional birth and death processes absorbed when one of the coordinates hits 0

We consider general multitype birth and death processes in continuous time, taking values in $\mathbb{Z}_{+}^{d}$ for some $d \geq 2$. Let $\left(X_{t}, t \geq 0\right)$ be a Markov process on $\mathbb{Z}_{+}^{d}$ with transition rates

$$
\text { from } n=\left(n_{1}, \ldots, n_{d}\right) \text { to } \begin{cases}n+e_{j} & \text { with rate } n_{j} b_{j}(n), \\ n-e_{j} & \text { with rate } n_{j} d_{j}(n)\end{cases}
$$

for all $1 \leq j \leq d$, with $e_{j}=(0, \ldots, 0,1,0, \ldots, 0)$, where the nonzero coordinate is the $j$-th one, $b(n)=\left(b_{1}(n), \ldots, b_{r}(n)\right)$ and $d(n)=\left(d_{1}(n), \ldots, d_{r}(n)\right)$ are functions from $\mathbb{Z}_{+}^{d}$ to $\mathbb{R}_{+}^{d}$. This model represents a density-dependent population dynamics with $d$ types of individuals (say $d$ species), where $b_{i}(n)$ (resp. $\left.d_{i}(n)\right)$ represents the reproduction rate (resp. death rate) per individuals of species $i$ when the population is in state $n$.

Note that the forms of the birth and death rates imply that, once a coordinate $X_{t}^{j}$ of the process hits 0 , it remains equal to 0 . This corresponds to the extinction of the population of type $j$. Hence, the set $\partial:=\mathbb{Z}_{+}^{r} \backslash \mathbb{N}^{r}$ is absorbing for the process $X$ and we denote as usual by $\tau_{\partial}$ its absorption time.

We define for all $1 \leq k \leq d$

$$
\begin{aligned}
\bar{d}(k) & =\sup _{n \in \mathbb{N}^{d},|n|=k}|n| \sum_{i=1}^{d} \mathbb{1}_{n_{i}=1} d_{i}(n), \\
\text { and } \quad \underline{d}(k) & =\inf _{n \in \mathbb{N}^{d},|n|=k} \sum_{i=1}^{d} n_{i}\left[\mathbb{1}_{n_{i} \neq 1} d_{i}(n)-b_{i}(n)\right],
\end{aligned}
$$

where $|n|:=n_{1}+\ldots+n_{d}$. We shall assume
Assumption 5. There exists $\eta>0$ small enough such that, for all $k \in \mathbb{N}$ large enough,

$$
\begin{equation*}
\underline{d}(k) \geq \eta \bar{d}(k) \tag{3.1}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{\underline{d}(k)}{k^{1+\eta}} \xrightarrow[k \rightarrow+\infty]{ }+\infty \tag{3.2}
\end{equation*}
$$

Note that, assuming that the set $O_{n}$ is finite for all $n$, any function $f: \mathbb{Z}_{+}^{d} \rightarrow \mathbb{R}$ is in the domain of the weakened infinitesimal generator of $X$ and, for all $n \in \mathbb{N}_{+}^{d}$,

$$
L f(n)=\sum_{j=1}^{d}\left[f\left(n+e_{j}\right)-f(n)\right] n_{j} b_{j}(n)+\sum_{j=1}^{d}\left[f\left(n-e_{j}\right)-f(n)\right] n_{j} d_{j}(n) .
$$

Under Assumption (3.2), setting $W(n)=|n|$, we have

$$
\frac{L W(n)}{W(n)^{1+\eta}} \rightarrow-\infty .
$$

This classically entails that

$$
\begin{equation*}
\sup _{n \in \mathbb{N}^{r}} \mathbb{E}_{n}\left(\left|X_{1}\right|\right)<+\infty \tag{3.3}
\end{equation*}
$$

(The argument is very similar to the one used for Theorem 2.4.) In particular, the process is non-explosive and $\tau_{\partial}$ is finite almost surely. Therefore, if we define for all $m \geq 1$

$$
O_{m}=U_{m}=\left\{x \in \mathbb{N}^{d}:|x| \leq m\right\},
$$

it is clear that the process $X$ satisfies the conditions of Section 2.1 and is regularly absorbed. We can now state the main result of the section.

Theorem 3.1. Under Assumption 5, the multi-dimensional competitive birth and death process $\left(X_{t}, t \geq 0\right)$ absorbed when one of its coordinates hits 0 admits a unique quasi-stationnary distribution $\nu_{Q S D}$ and there exist constants $C, \gamma>0$ such that, for all probability measure $\mu$ on $\mathbb{N}^{d}$,

$$
\left\|\mathbb{P}_{\mu}\left(X_{t} \in \cdot \mid t<\tau_{\partial}\right)-\nu_{Q S D}\right\|_{T V} \leq C e^{-\gamma t}, \quad \forall t \geq 0
$$

It is easy to check that Assumption 5 is satisfied in the general LotkaVolterra birth and death process of the introduction. Indeed, we clearly have $\bar{d}(k) \leq C k^{2}$ and

$$
\begin{aligned}
\underline{d}(k) & \geq \inf _{n \in \mathbb{N}^{d},|n|=k} \sum_{i=1}^{d}\left(d_{i}(n)-b_{i}(n)\right)-\sup _{n \in \mathbb{N}^{d},|n|=k} \sum_{i=1}^{n} d_{i}(n) \mathbb{1}_{n_{i}=1} \\
& \geq-C k+\inf _{n \in \mathbb{N}^{d},|n|=k} \sum_{i, j=1}^{d} n_{i}\left(c_{i j}-\gamma_{i j}\right) n_{j}
\end{aligned}
$$

for $C=\sup _{i}\left(b_{i}+d_{i}\right)+\sup _{i j} c_{i j}$. Under the assumptions of Theorem 1.1, there exists $C^{\prime}>0$ such that, for all $n \in \mathbb{N}^{d}$,

$$
\sum_{i, j=1}^{d} n_{i}\left(c_{i j}-\gamma_{i j}\right) n_{j} \geq C^{\prime}|n|^{2}
$$

This entails Assumption 5 .
Proof. We are going to use the criterion of Corollary 2.8 , Using (3.3) and copying the arguments of [9, Sec.4.1.1 and Thm.4.1], one deduces that Assumption 2 is satisfied with $\nu=\delta_{(1, \ldots, 1)}$ and Assumptions 3 and 4 are also satisfied.

Hence we only have to find a couple of admissible functions $(V, \varphi)$ satisfying Conditions (a) and (b) of Proposition 2.7 and satisfying (2.9) (note that the fact that $\varphi$ is positive on $\mathbb{N}^{d}$ entails immediately that $\left.\inf _{x \in O_{m}} \varphi(x)>0\right)$. This couple of functions is given for all $n \in \mathbb{Z}_{+}^{d}$ by

$$
V(n)= \begin{cases}\sum_{k=1}^{|n|} \frac{1}{k^{\alpha}} & \text { if } n \in \mathbb{N}^{d} \\ 0 & \text { if } n \in \partial\end{cases}
$$

and

$$
\varphi(n)= \begin{cases}\sum_{k=|n|+1}^{+\infty} \frac{1}{k^{\beta}} & \text { if } n \in \mathbb{N}^{d} \\ 0 & \text { if } n \in \partial\end{cases}
$$

for appropriate choices of $\alpha, \beta>1$. Note that the two functions are bounded, nonnegative and positive of $\mathbb{N}^{d}$. Note also that, $\operatorname{since}^{\inf }{ }_{n \in \mathbb{N}^{d}} V(n)>0$, Conditions (2.3) and (2.9) are obviously satisfied. Note also that (2.5) is clear for non-explosive Markov processes on discrete state spaces, and that (2.4) is true from the definition of $V$ and $\varphi$.

Hence, we only have to check Conditions (a) and (b) since they imply the last condition needed for the couple of functions $(V, \varphi)$ to be admissible: the fact that $L V$ is bounded from above and $L \varphi$ bounded from below. So we compute

$$
\begin{aligned}
L \varphi(n) & =-\sum_{i=1}^{d} \frac{n_{i} b_{i}(n)}{(1+|n|)^{\beta}}+\sum_{i=1}^{d} \frac{n_{i} \mathbb{1}_{n_{i} \neq 1} d_{i}(n)}{|n|^{\beta}}-\sum_{i=1}^{d} \mathbb{1}_{n_{i}=1} d_{i}(n) \sum_{k=|n|+1}^{+\infty} \frac{1}{k^{\beta}} \\
& \geq \frac{1}{|n|^{\beta}}\left[\underline{d}(|n|)-\frac{\bar{d}(|n|)}{\beta-1}\right],
\end{aligned}
$$

where we used the fact that

$$
\sum_{k=x+1}^{+\infty} \frac{1}{k^{\beta}} \leq \int_{x}^{+\infty} \frac{d y}{y^{\beta}}=\frac{1}{(\beta-1) x^{\beta-1}}
$$

Hence it follows from Assumption (3.1) that there exists $\beta>1$ large enough such that $L \varphi(n) \geq 0$ for all $|n|$ large enough. This entails Condition (a).

We fix such a value of $\beta$. Using that

$$
\sup _{n \in \mathbb{N}^{d}} V(n)=\sum_{k=1}^{+\infty} \frac{1}{k^{\alpha}} \leq 1+\int_{1}^{\infty} \frac{d x}{x^{\alpha}}=\frac{\alpha}{\alpha-1}
$$

and

$$
\varphi(n) \geq \int_{|n|+1}^{\infty} \frac{d x}{x^{\beta}}=\frac{(1+|n|)^{1-\beta}}{\beta-1} \geq \frac{|n|^{1-\beta}}{2(\beta-1)}
$$

for $|n|$ large enough, we compute for such $n$

$$
\begin{aligned}
L V(n)+\frac{V^{1+\varepsilon}(n)}{\varphi^{\varepsilon}(n)} & \leq \sum_{i=1}^{d} \frac{n_{i} b_{i}(n)}{(|n|+1)^{\alpha}}-\sum_{i=1}^{d} \frac{n_{i} d_{i}(n) \mathbb{1}_{n_{i} \neq 1}}{|n|^{\alpha}}+C|n|^{\varepsilon(\beta-1)} \\
& \leq-\frac{d(|n|)}{|n|^{\alpha}}+C|n|^{\varepsilon(\beta-1)},
\end{aligned}
$$

where $C=[\alpha /(\alpha-1)]^{1+\varepsilon}[2(\beta-1)]^{\varepsilon}$. Choosing $\alpha=1+\eta / 2$ and $\varepsilon=\eta /[2(\beta-$ 1)], Assumption (3.2) implies that $L V(n)+\frac{V^{1+\varepsilon}(n)}{\varphi^{\varepsilon}(n)} \leq 0$ for $n \notin O_{m}$ with $m$ large enough. Since $\inf _{n \in O_{m}} \varphi(n)>0$, we have proved Condition (b).

## 4 Application to multidimensional Feller diffusions absorbed when one of the coordinates hits 0

We consider a general multitype Feller diffusion $\left(X_{t}, t \geq 0\right)$ in $\mathbb{R}_{+}^{d}$, solution to the stochastic differential equation

$$
\begin{equation*}
d X_{t}^{i}=\sqrt{\gamma_{i} X_{t}^{i}} d B_{t}^{i}+X_{t}^{i} r_{i}\left(X_{t}\right) d t, \quad 1 \leq i \leq d \tag{4.1}
\end{equation*}
$$

where ( $B_{t}^{i}, t \geq 0$ ) are independent standard Brownian motions, $\gamma_{i}$ are positive constants and $r_{i}$ are measurable maps from $\mathbb{R}_{+}^{d}$ to $\mathbb{R}$. From the biological point of view, $r_{i}(x)$ represents the growth rate per individual of species $i$ in a population of size vector $x \in \mathbb{R}_{+}^{d}$. We shall make the following assumption.

Assumption 6. Assume that, for all $i \in\{1, \ldots, d\}, r_{i}$ is locally Hölder on $\mathbb{R}_{+}^{d}$ and that there exist $a>0$ and $0<\eta<1$ such that

$$
\begin{equation*}
r_{i}(x) \leq a^{\eta}-x_{i}^{\eta} \tag{4.2}
\end{equation*}
$$

and there exist constants $B_{a}>a, C_{a}>0$ and $D_{a}>0$ such that

$$
\begin{equation*}
\sum_{i=1}^{d} \mathbb{1}_{x_{i} \geq B_{a}} r_{i}(x) \leq C_{a}\left(\sum_{i=1}^{d} \mathbb{1}_{x_{i} \leq a} r_{i}(x)+D_{a}\right), \quad \forall x \in \mathbb{R}_{+}^{d} \tag{4.3}
\end{equation*}
$$

Assumptions (4.2) and (4.3) correspond to some form of competitivity for the system of interacting Feller diffusions. Moreover, Assumption (4.2) entail that each coordinate of the process can be upper bounded by the (strong) solution of the one-dimensional Feller diffusion

$$
\begin{equation*}
d \hat{X}_{t}^{i}=\sqrt{\gamma_{i} \hat{X}_{t}^{i}} d B_{t}^{i}+a^{\eta} \hat{X}_{t}^{i} d t, \quad 1 \leq i \leq d \tag{4.4}
\end{equation*}
$$

with initial value $\hat{X}_{0}^{i}=X_{0}^{i}$. To prove this, one can apply the transformation $Y_{t}^{i}=2 \sqrt{X_{t}^{i} / \gamma_{i}}$ to obtain a SDE of the form

$$
d Y_{t}^{i}=d B_{t}^{i}+\hat{r}_{i}\left(Y_{t}\right) d t, \quad 1 \leq i \leq d
$$

with $\hat{r}_{i}$ measurable, locally Hölder, hence locally bounded in $(0,+\infty)^{d}$, and (possibly) singular at $\partial:=\mathbb{R}_{+}^{d} \backslash(0,+\infty)^{d}$, and apply classical comparison results to solutions of SDEs with constant diffusion coefficient. Since the SDE (4.4) is non-explosive $(+\infty$ is a natural boundary and 0 is an exit boundary for this diffusion by classical criteria, cf. e.g. [15]), the SDE (4.1) is non explosive. Similarly, Assumption (4.2) implies that $X_{t}^{i} \leq \bar{X}_{t}^{i}$ where

$$
\begin{equation*}
d \bar{X}_{t}^{i}=\sqrt{\gamma_{i} \bar{X}_{t}^{i}} d B_{t}^{i}+\bar{X}_{t}^{i}\left[a^{\eta}-\left(\bar{X}_{t}^{i}\right)^{\eta}\right] d t, \quad 1 \leq i \leq d \tag{4.5}
\end{equation*}
$$

for some constant $B$ and with initial value $\bar{X}_{0}^{i}=X_{0}^{i}$. Since $\bar{X}^{i}$ is a diffusion on $\mathbb{R}_{+}$for which $+\infty$ is an entrance boundary and 0 an exit boundary, we deduce that $\tau_{\partial}<+\infty$ almost surely.

We also deduce from the fact that $r_{i}$ is locally bounded and from (4.2) that there is strong existence and pathwise uniqueness for (4.1). The ideas to prove this seem quite standard, but we couldn't find a proof of this fact in our situation, so we explain the main steps below. Assume that $X_{0}^{i}>0$ for all $1 \leq i \leq d$ (otherwise, we will see below what to do). Since the drifts $\hat{r}_{i}$ of $Y^{i}$ are locally bounded in $(0,+\infty)^{d}$, we can use the well-known result of Veretennikov [25] to prove local strong existence and pathwise uniqueness
until possible explosion or the first exit time $\tau_{\text {exit }}$ from all compact subsets of $(0,+\infty)^{d}$. As explained above, the process is non-explosive. In addition, on the event $\left\{\tau_{\text {exit }}<\infty\right\}$, the process admits a left limit at the time $\tau_{\text {exit }}$ since its drift and diffusion coefficients are bounded on each compact subset of $\mathbb{R}_{+}^{d}$. Hence, we can define $X_{\tau_{\text {exit }}}$ as the right limit of $X_{t}$ at this time, which has at least one zero coordinate, i.e. $\tau_{\partial}=\tau_{\text {exit. }}$. We then set $X_{t}^{j}=0$ for all $t \geq \tau_{\partial}$ and all $j$ such that $X_{\tau_{\partial}}^{j}=0$. We can then proceed as above for the set of coordinates such that $X_{\tau_{\partial}}^{i}>0$. This gives the existence of a global strong solution.

For pathwise uniqueness, we know that it holds locally between the times of absorption of a coordinate, but only for the process formed by the nonzero coordinates, so we need to prove that no solution of (4.1) starting with one zero coordinate, say $X_{0}^{i}=0$, can have $X_{t}^{i}>0$ for some positive time $t$. This can be proved using again the comparison between $X^{i}$ and the Feller diffusions (4.4), since 0 is absorbing (since it is an exit boundary) for this diffusion.

Strong existence and pathwise uniqueness imply well-posedness of the martingale problem, hence the strong Markov property hold on the canonical space with respect to the natural filtration (see e.g. [23]). Since the paths of $X$ are continuous, the hitting times of closed subsets of $\mathbb{R}_{+}^{d}$ are stopping times for this filtration. Hence, defining $O_{n}=\left\{x \in(0,+\infty)^{d}: 1 /(2 n)<\right.$ $\left.x_{i}<2 n, \forall 1 \leq i \leq d\right\}$ and $U_{n}=\left\{x \in(0,+\infty)^{d}: 1 /(2 n-1) \leq x_{i} \leq\right.$ $2 n-1, \forall 1 \leq i \leq d\}$ for all $n \geq 1$, it is then clear that the process $X$ satisfies the conditions of Section 2.1 and, in view of the previous construction of the process, it is regularly absorbed. In addition, it follows from Itô's formula and the local boundedness of the coefficients of the SDE that any measurable function $f: \mathbb{R}_{+}^{d} \rightarrow \mathbb{R}$ twice continuously differentiable on $(0,+\infty)$ belongs to the domain of the weakened generator of $X$ and

$$
L f(x)=\sum_{i=1}^{d} \frac{\gamma_{i} x_{i} b_{i}(x)}{2} \frac{\partial f}{\partial x_{i}}(x)+\sum_{i=1}^{d} x_{i} \frac{\partial^{2} f}{\partial x_{i}^{2}}(x), \quad \forall x \in(0,+\infty)^{d} .
$$

We can now state the main result of the section.
Theorem 4.1. Under Assumption 6, the multi-dimensional Feller diffusion process ( $X_{t}, t \geq 0$ ) absorbed when one of its coordinates hits 0 admits a unique quasi-stationnary distribution $\nu_{Q S D}$ and there exist constants $C, \gamma>$ 0 such that, for all probability measure $\mu$ on $\mathbb{N}^{d}$,

$$
\left\|\mathbb{P}_{\mu}\left(X_{t} \in \cdot \mid t<\tau_{\partial}\right)-\nu_{Q S D}\right\|_{T V} \leq C e^{-\gamma t}, \quad \forall t \geq 0
$$

It is straightforward to check that Assumption 6 is satisfied in the competitive Lotka-Volterra case, that is when $r_{i}(x)=r_{i}-\sum_{i=1}^{d} c_{i j} x_{j}$ with $c_{i j} \geq 0$ and $c_{i i}>0$ for all $1 \leq i, j \leq d$. Hence Theorem 1.2 is an immediate corollary of Theorem 4.1. Assumption 6 allows for other biologically relevant models. For instance, one can consider ecosystems where the competition among individuals only begins when the population size reaches a level $K>0$, which leads for instance to the SDE

$$
d X_{t}^{i}=\sqrt{\gamma_{i} X_{t}^{i}} d B_{t}^{i}+X_{t}^{i}\left(r-\sum_{j=1}^{d} c_{i j}\left(X_{t}^{j}-K\right)_{+}\right) d t, X_{0}^{i}>0
$$

where $c_{i j}$ are non-negative constants and $c_{i i}>0$ for all $i, j \in\{1, \ldots, d\}$. Note also that a similar approach (i.e. using Theorem 2.4 with Lyapunov functions of the form $\prod_{i=1}^{d} h\left(x_{i}\right)$ ) can also be used to handle diffusion processes evolving in bounded boxes.

Proof. Up to a linear scaling of the coordinates, we can assume without loss of generality that $\gamma_{i}=2$ for all $1 \leq i \leq d$, so we will only consider this case from now on. Note that Assumption 6 is not modified by the rescaling.

We shall make use of Corollary 2.8. We divide the proof into five steps, respectively devoted to the construction of a function $\varphi$ satisfying Conditions (a), of a function $V$ satisfying Condition (b), to the proof of (2.9), to the proof of a local Harnack inequality, and to check Assumption 2, Assumption 3 and Assumption 4.

## Step 1: construction of a function $\varphi$ satisfying Condition (a).

Recall the definition of the constants $a>0$ and $B_{a}>a$ from Assumption 6. We use the following lemma, proved at the end of this section.
Lemma 4.2. There exists $M>0$ such that, for all $\beta \geq M$, there exists a function $h_{\beta}: \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}$twice continuously differentiable on $(0,+\infty)$ such that

$$
h_{\beta}(x)= \begin{cases}4 x^{2} / a^{2} & \text { if } x \in[0, a / 2], \\ B_{a}^{\beta}(2 x)^{-\beta} & \text { if } x \geq B_{a},\end{cases}
$$

$h_{\beta}(x) \geq 1$ for all $x \in[a / 2, a], h_{\beta}$ is nonincreasing and convex on $[a,+\infty)$,
$M^{\prime}:=\sup _{\beta \geq M} \sup _{x \in[a / 2, a]}\left|h_{\beta}^{\prime}(x)\right|<+\infty \quad$ and $\quad M^{\prime \prime}:=\sup _{\beta \geq M} \sup _{x \in[a / 2, a]}\left|h_{\beta}^{\prime \prime}(x)\right|<+\infty$.
We set $\beta=M+\left(2 \vee a M^{\prime}\right) / C_{a}+1$,

$$
\varphi(x)=\prod_{i=1}^{d} h_{\beta}\left(x_{i}\right), \quad \forall x \in \mathbb{R}_{+}^{d}
$$

and we shall prove that $\varphi$ satisfies Condition (a). We have

$$
\frac{L \varphi(x)}{\varphi(x)}=\sum_{i=1}^{d} \frac{x_{i} h_{\beta}^{\prime}\left(x_{i}\right) r_{i}(x)+x_{i} h_{\beta}^{\prime \prime}(x)}{h_{\beta}\left(x_{i}\right)} .
$$

Now, it follows from the properties of $h_{\beta}$ and Assumptions 6 that, for all $x \in \mathbb{R}_{+}$and all $1 \leq i \leq d$,

$$
\frac{x_{i} h_{\beta}^{\prime}\left(x_{i}\right) r_{i}(x)+x_{i} h_{\beta}^{\prime \prime}(x)}{h_{\beta}\left(x_{i}\right)} \geq \begin{cases}0 & \text { if } x_{i} \geq a \\ -\beta r_{i}(x) & \text { if } x_{i} \geq B_{a} \\ 2 r_{i}(x)+\frac{2}{x_{i}} & \text { if } x_{i} \leq a / 2 \\ a M^{\prime}\left(r_{i}(x)-a^{\eta}\right)-a^{1+\eta} M^{\prime}-a M^{\prime \prime} & \text { if } a / 2 \leq x_{i} \leq a\end{cases}
$$

where we used in the last inequality the fact that $r_{i}(x)-a^{\eta} \leq 0$ for all $x$. Using once again this property, we deduce that, for some constant $B$ independent of $\beta \geq M$,

$$
\frac{x_{i} h_{\beta}^{\prime}\left(x_{i}\right) r_{i}(x)+x_{i} h_{\beta}^{\prime \prime}(x)}{h_{\beta}\left(x_{i}\right)} \geq \begin{cases}0 & \text { if } x_{i} \geq a,  \tag{4.6}\\ -\beta r_{i}(x) & \text { if } x_{i} \geq B_{a}, \\ \left(2 \vee a M^{\prime}\right) r_{i}(x)+\frac{2}{x_{i}}-B & \text { if } x_{i} \leq a .\end{cases}
$$

Hence, for all $x \in \mathbb{R}_{+}^{d}$,

$$
\frac{L \varphi(x)}{\varphi(x)} \geq-\beta \sum_{i=1}^{d} \mathbb{1}_{x_{i} \geq B_{a}} r_{i}(x)+\sum_{i=1}^{d} \mathbb{1}_{x_{i} \leq a}\left(\left(2 \vee a M^{\prime}\right) r_{i}(x)+\frac{2}{x_{i}}\right)-d B
$$

This and Assumption (4.3) imply that

$$
\begin{aligned}
\frac{L \varphi(x)}{\varphi(x)} & \geq-\sum_{i=1}^{d} \mathbb{1}_{x_{i} \geq B_{a}} r_{i}(x)+2 \sum_{i=1}^{d} \mathbb{1}_{x_{i} \leq a} \frac{1}{x_{i}}-d B-\left(2 \vee a M^{\prime}\right) D_{a} \\
& \geq \sum_{i=1}^{d}\left(x_{i}^{\eta}+\frac{2}{x_{i}}\right)-B^{\prime},
\end{aligned}
$$

for some constant $B^{\prime}$, where we used Assumption (4.2) in the last inequality.
Hence, there exists $n \geq 1$ and a constant $C>0$ such that

$$
L \varphi(x) \geq-C \mathbb{1}_{x \in O_{n}} .
$$

This ends the proof that $\varphi$ satisfies Condition (a).

Step 2: construction of a function $V$ satisfying Condition (b) and verification that $(V, \varphi)$ is a couple of admissible functions.

For $V$, we define

$$
V(x)=\prod_{i=1}^{d} g\left(x_{i}\right), \quad \forall x \in \mathbb{R}_{+}^{d},
$$

where the function $g: \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}$is twice continuously differentiable on $(0,+\infty)$, increasing concave and such that

$$
g(x)= \begin{cases}x^{\gamma} & \text { if } x \leq 1 \\ \delta-x^{-\eta / 2} & \text { if } x \geq 2\end{cases}
$$

for some constants $\gamma<1$ and $\delta>0$ and where $\eta$ is defined in Assumption (4.2). Since $g^{\prime}(1)=\gamma$ and $g^{\prime}(2)=\eta 2^{-2-\eta / 2}$, it is possible to find $\delta>0$ such that such a function $g$ exists as soon as $\eta 2^{-2-\eta / 2}<\gamma$. Hence, we shall assume that $\gamma$ belongs to the non-empty interval $\left(\eta 2^{-2-\eta / 2}, 1\right)$. We have

$$
\frac{L V(x)}{V(x)}=\sum_{i=1}^{d} \frac{x_{i} g^{\prime}\left(x_{i}\right) r_{i}(x)+x_{i} g^{\prime \prime}\left(x_{i}\right)}{g\left(x_{i}\right)}
$$

and

$$
\frac{x_{i} g^{\prime}\left(x_{i}\right) r_{i}(x)+x_{i} g^{\prime \prime}\left(x_{i}\right)}{g\left(x_{i}\right)} \leq \begin{cases}\gamma r_{i}(x)-\frac{\gamma(1-\gamma)}{x_{i}} & \text { if } x_{i} \leq 1 \\ 2 a^{\eta} \sup _{1 \leq x \leq 2} g^{\prime}(x) & \text { if } 1 \leq x_{i} \leq 2 \\ \frac{\eta r_{i}(x) x_{i}^{-\eta / 2}}{2\left(\delta-x_{i}^{-\eta / 2}\right)} & \text { if } x_{i} \geq 2\end{cases}
$$

We deduce from Assumptions (4.2) that there exist constants $B^{\prime}, B^{\prime \prime}>0$ such that

$$
\frac{x_{i} g^{\prime}\left(x_{i}\right) r_{i}(x)+x_{i} g^{\prime \prime}(x)}{g\left(x_{i}\right)} \leq B^{\prime}- \begin{cases}\frac{\gamma(1-\gamma)}{x_{i}} & \text { if } x_{i} \leq 1  \tag{4.7}\\ 0 & \text { if } 1 \leq x_{i} \leq 2 \\ B^{\prime \prime} x_{i}^{\eta / 2} & \text { if } x_{i} \geq 2\end{cases}
$$

Thus, since $h_{\beta}\left(x_{i}\right) \geq D_{\beta}\left(x_{i}^{2} \wedge x_{i}^{-\beta}\right)$ for some constant $D_{\beta}>0$ and since

$$
\begin{aligned}
& g\left(x_{i}\right) \leq \delta \\
& \frac{L V(x)}{V(x)}+\frac{V(x)^{\varepsilon}}{\varphi(x)^{\varepsilon}} \leq B^{\prime} d-\gamma(1-\gamma) \sum_{i=1}^{d} \frac{\mathbb{1}_{x_{i} \leq 1}}{x_{i}}-B^{\prime \prime} \sum_{i=1}^{d} \mathbb{1}_{x_{i} \geq 2} x_{i}^{\eta / 2} \\
&+\left(\frac{\delta}{D_{\beta}}\right)^{d \varepsilon} \prod_{i=1}^{d}\left(x_{i}^{\varepsilon \beta} \vee x_{i}^{-2 \varepsilon}\right) \\
& \leq B^{\prime} d+\gamma(1-\gamma)+B^{\prime \prime} 2^{\eta / 2}-\gamma(1-\gamma)\left(\inf _{i} x_{i}\right)^{-1}-B^{\prime \prime}\left(\sup _{i} x_{i}\right)^{\eta / 2} \\
&+\left(\frac{\delta}{D_{\beta}}\right)^{d \varepsilon}\left[\left(\sup _{i} x_{i}\right)^{\beta d \varepsilon}+\left(\inf _{i} x_{i}\right)^{-2 d \varepsilon}\right]
\end{aligned}
$$

Therefore, choosing $\varepsilon>0$ such that $d(1+\alpha) \varepsilon<1$ and $\beta d \varepsilon<\eta / 2, L V(x)+$ $V(x)^{1+\varepsilon} / \varphi(x)^{\varepsilon} \leq 0$ for all $x \in \mathbb{R}_{+}^{d}$ such that $\inf _{i} x_{i}$ is small enough or $\sup _{i} x_{i}$ is big enough. Since $L V$ is bounded from above by 4.7) and since $V$ and $\varphi$ are positive continuous on any compact subset of $(0,+\infty)^{d}$, we have proved Condition (b).

We can now check that $(V, \varphi)$ is a couple of admissible functions. First, $V$ and $\varphi$ are both bounded, positive on $(0,+\infty)$ and vanishing on $\partial$. They both belong to the domain of the weakened infinitesimal generator of $X$. Since the function $g / h_{\beta}$ is positive continuous on $(0,+\infty)$ and

$$
\frac{g(x)}{h_{\beta}(x)}= \begin{cases}x^{\gamma-2} & \text { if } x \leq 1 \wedge a / 2 \\ \left(\delta-x^{-\eta / 2}\right)(2 x)^{\beta} / B_{a}^{\beta} & \text { if } x \geq B_{a} \vee 2\end{cases}
$$

we deduce that $\inf _{x \in(0,+\infty)} g(x) / h_{\beta}(x)>0$, hence 2.3 holds true, and 2.4 is also clear. Since $V$ is continuous on $\mathbb{R}_{+}^{d}$ and $X_{T_{n}} \rightarrow X_{\tau_{\partial}}$ almost surely, 2.5) also holds true. Finally, since $L V$ is bounded from above and $L \varphi$ is bounded from below, we have proved that $(V, \varphi)$ is a couple of admissible functions.
Step 3: proof of 2.9 .
Using the upper bound $X_{t}^{i} \leq \bar{X}_{t}^{i}$ for all $t \geq 0$ and $1 \leq i \leq d$, where $\bar{X}^{i}$ is solution to the $\operatorname{SDE}$ (4.5), and noting that the processes $\left(\bar{X}^{i}\right)_{1 \leq i \leq d}$ are independent, we have for all $x \in \mathbb{R}_{+}^{d}$ and all $t_{2}>0$,

$$
\mathbb{P}_{x}\left(t_{2}<\tau_{\partial}\right) \leq \prod_{i=1}^{d} \mathbb{P}_{x_{i}}\left(\bar{X}_{t_{2}}^{i}>0\right)
$$

Now, there exist constants $D$ and $D^{\prime}$ such that

$$
\begin{equation*}
\mathbb{P}_{x_{i}}\left(\bar{X}_{t_{2}}^{i}>0\right) \leq\left(D x_{i}\right) \wedge 1 \leq D^{\prime} g\left(x_{i}\right) \quad \text { for all } x_{i}>0 \tag{4.8}
\end{equation*}
$$

To prove this, we can consider a scale function $s$ of the diffusion $\bar{X}^{i}$ such that $s(0)=0$. Using the expression of the scale function and the speed measure (see e.g. [23, V.52]), one easily checks that $s\left(x_{i}\right) \sim x_{i}$ when $x_{i} \rightarrow 0$ and that Proposition 4.9 of [8] is satisfied, so that $\mathbb{P}_{x_{i}}\left(\bar{X}_{t_{2}}^{i}>0\right) \leq M s\left(x_{i}\right)$ for some $M>0$. Since $s\left(x_{i}\right) \sim x_{i}$ when $x_{i} \rightarrow 0,4.8$ is proved and hence (2.9) holds true.

Step 4: Harnack inequality for $u$. Consider a bounded measurable function $f$ and define the application $u:(t, x) \in \mathbb{R}_{+} \times E \mapsto \mathbb{E}_{x}\left(f\left(X_{t}\right)\right)$. Our aim is to prove that, for all $m \geq 1$, there exist two constants $N_{m}>0$ and $\delta_{m}>0$, which do not depend on $f$, such that
$u\left(\delta_{m}+\delta_{m}^{2}, x\right) \leq N_{m} u\left(\delta_{m}+2 \delta_{m}^{2}, y\right)$, for all $x, y \in O_{m}$ such that $|x-y| \leq \delta_{m}$.

First, the function $u$ is continuous on $\mathbb{R}_{+}^{*} \times E$. Indeed, since the SDE (4.1) has Hölder coefficients in $O_{m}$, using the same approach as in the proof of [24, Thm 7.2.4], one deduces that, for all $x \in E$ and for all $\varepsilon>0$, there exists $\delta>0$ such that, for all $y \in B(x, \delta)$, all $s, t \in[\delta, 2 \delta]$ and all bounded measurable function $\varphi: E \rightarrow \mathbb{R}$,

$$
\left|\mathbb{E}_{x}\left(\varphi\left(X_{s}\right)\right)-\mathbb{E}_{y}\left(\varphi\left(X_{t}\right)\right)\right| \leq \varepsilon\|\varphi\|_{\infty}
$$

This result applied to the function $\varphi: x \mapsto \mathbb{E}_{x}\left(f\left(X_{t-\delta}\right)\right)$ for $t \geq \delta$ and the Markov property entails the continuity of $u$ on $\mathbb{R}_{+}^{*} \times E$.

Second, fix $m \geq 1$ and let $K$ be a compact set with $C^{\infty}$ boundary such that $O_{m} \subset K \subset E$ and such that $d\left(O_{m}, \partial K\right)>0$. We set $\delta=$ $d\left(O_{m}, \partial K\right) / 3$. Let $\left(h_{n}\right)_{n \geq 1}$ be a sequence of bounded $C^{\infty}$ functions from $\left(\left[\delta, \infty[\times \partial K) \cup(\{\delta\} \times K)\right.\right.$ to $\mathbb{R}$ such that $\left(h_{n}\right)_{n \geq 1}$ converges uniformly locally toward $u$ restricted to $([\delta, \infty[\times \partial K) \cup(\{\delta\} \times K)$ (such a sequence exists because of the continuity of $u)$. For all $n \geq 1$, we let $u_{n}:[\delta,+\infty[\times K \rightarrow \mathbb{R}$ be the solution to the linear parabolic equation

$$
\left\{\begin{array}{l}
\partial_{t} u_{n}=\sum_{i=1}^{d} x_{i} \frac{\partial^{2} u_{n}}{\partial x_{i}^{2}}+x_{i} r_{i}(x) \frac{\partial u_{n}}{\partial x_{i}}, \\
u_{n}(t, x)=h_{n}(t, x), \forall(t, x) \in([\delta, \infty[\times \partial K) \cup(\{\delta\} \times K) .
\end{array}\right.
$$

By [18, Thm 5.1.15], for all $n \geq 1, u_{n}$ is of regularity $C^{1,2}$. Moreover, using the Harnack inequality provided by [16, Theorem 1.1] (with $\theta=2$ and $R=\delta$ ), we deduce that there exists a constant $N>0$ which does not
depend on $f$ nor on $n$ such that

$$
\begin{equation*}
u_{n}\left(\delta+\delta^{2}, x\right) \leq N u_{n}\left(\delta+2 \delta^{2}, y\right), \text { for all } x, y \in O_{m} \text { such that }|x-y| \leq \delta \tag{4.10}
\end{equation*}
$$

In particular, applying Itô's formula to $s \mapsto u_{n}\left(\delta+t-s, X_{s}\right)$ at time $\tau_{K} \wedge t$ and taking the expectation, one deduces that, for all $t \geq 0$,

$$
\begin{aligned}
u_{n}(\delta+t, x) & =\mathbb{E}_{x}\left[u_{n}\left(\delta+t-\tau_{K} \wedge t, X_{\tau_{K} \wedge t}\right)\right] \\
& =\mathbb{E}_{x}\left(\mathbb{1}_{\tau_{K} \leq t} u_{n}\left(t+\delta-\tau_{K}, X_{\tau_{K}}\right)\right)+\mathbb{E}_{x}\left(\mathbb{1}_{\tau_{K}>t} u_{n}\left(\delta, X_{t}\right)\right) \\
& =\mathbb{E}_{x}\left(\mathbb{1}_{\tau_{K} \leq t} h_{n}\left(t+\delta-\tau_{K}, X_{\tau_{K}}\right)\right)+\mathbb{E}_{x}\left(\mathbb{1}_{\tau_{K}>t} h_{n}\left(\delta, X_{t}\right)\right) .
\end{aligned}
$$

By Lebesgue's theorem, the last quantity converges when $n \rightarrow+\infty$ to

$$
\begin{aligned}
\mathbb{E}_{x}\left(\mathbb{1}_{\tau_{K} \leq t} u(t+\right. & \left.\left.\delta-\tau_{K}, X_{\tau_{K}}\right)\right)+\mathbb{E}_{x}\left(\mathbb{1}_{\tau_{K}>t} u\left(\delta, X_{t}\right)\right) \\
& =\mathbb{E}_{x}\left(\mathbb{1}_{\tau_{K} \leq t} \mathbb{E}_{X_{\tau_{K}}}\left(f\left(X_{t+\delta-\tau_{K}}\right)\right)+\mathbb{E}_{x}\left(\mathbb{1}_{\tau_{K}>t} \mathbb{E}_{X_{t}}\left(f\left(X_{\delta}\right)\right)\right)\right. \\
& =\mathbb{E}_{x}\left(f\left(X_{\delta+t}\right)\right)=u(\delta+t, x),
\end{aligned}
$$

where we used the strong Markov property at time $\tau_{K}$.
We deduce from this convergence and the Harnack inequalities 4.10 that there exist two constants $N_{m}>0$ and $\delta_{m}>0$, which do not depend on $f$, such that (4.9) holds true.

Step 5 : proof that Assumptions 2 , 3 and 4 are satisfied. Fix $x_{1} \in O_{1}$ and let $\nu$ denote the conditional law $\mathbb{P}_{x_{1}}\left(X_{\delta_{1}+\delta_{1}^{2}} \in \cdot \mid \delta_{1}+\delta_{1}^{2}<\tau_{\partial}\right)$. Then the Harnack inequality (4.9) entails that, for all $x \in O_{1}$ such that $\left|x-x_{1}\right| \leq \delta_{1}$,

$$
\mathbb{P}_{x}\left(X_{\delta_{1}+2 \delta_{1}^{2}} \in \cdot\right) \geq \frac{\mathbb{P}_{x_{1}}\left(\delta_{1}+\delta_{1}^{2}<\tau_{\partial}\right)}{N_{1}} \nu
$$

Replacing $O_{1}$ by its intersection with the open ball with center $x_{1}$ and radius $\delta_{1}$, this implies that Assumption 2 is satisfied for $k_{0}=1$ (the second part of this assumption is a classical property for locally elliptic diffusion processes).

Assumption 3 is a direct consequence of the domination by solutions to (4.5), since these solutions come down from infinity and hit 0 in finite time almost surely (cf. e.g. [2]).

Finally, Assumption 4 is a consequence of (4.9). Indeed, for any fixed $m$ and for all $t \geq \delta_{m}+2 \delta_{m}^{2}$, this equation applied to $f(x)=\mathbb{P}_{x}\left(t-\delta_{m}-2 \delta_{m}^{2}<\right.$ $\left.\tau_{\partial}\right)$ and the Markov property entails that
$\mathbb{P}_{x}\left(t-\delta_{m}^{2}<\tau_{\partial}\right) \leq N_{m} \mathbb{P}_{y}\left(t<\tau_{\partial}\right)$, for all $x, y \in O_{m}$ such that $|x-y| \leq \delta_{m}$.

Since $s \mapsto \mathbb{P}_{x}\left(s<\tau_{\partial}\right)$ is non-increasing, we deduce that

$$
\mathbb{P}_{x}\left(t<\tau_{\partial}\right) \leq N_{m} \mathbb{P}_{y}\left(t<\tau_{\partial}\right), \text { for all } x, y \in O_{m} \text { such that }|x-y| \leq \delta_{m} .
$$

Since $O_{m}$ has a finite diameter and is connected, we deduce that there exists $N_{m}^{\prime}$ such that, for all $t \geq \delta_{m}+2 \delta_{m}^{2}$,

$$
\mathbb{P}_{x}\left(t<\tau_{\partial}\right) \leq N_{m}^{\prime} \mathbb{P}_{y}\left(t<\tau_{\partial}\right), \text { for all } x, y \in O_{m}
$$

Now, for $t \leq \delta_{m}+2 \delta_{m}^{2}$, we simply use the fact that $x \mapsto \mathbb{P}_{x}\left(\delta_{m}+2 \delta_{m}^{2}<\tau_{\partial}\right)$ is uniformly bounded from below on $O_{m}$ by a constant $1 / N_{m}^{\prime \prime}>0$. In particular,
$\mathbb{P}_{x}\left(t<\tau_{\partial}\right) \leq 1 \leq N_{m}^{\prime \prime} \mathbb{P}_{y}\left(\delta_{m}+2 \delta_{m}^{2}<\tau_{\partial}\right) \leq N_{m}^{\prime \prime} \mathbb{P}_{y}\left(t<\tau_{\partial}\right)$, for all $x, y \in O_{m}$.
As a consequence, Assumption 4 is satisfied.
Finally, we deduce from Steps 1, 2, 3 and 5 that all the assumptions of Theorem 2.8 are satisfied. This concludes the proof of Theorem 1.2 .

Proof of Lemma 4.2. We consider a $C^{\infty}$ function $\varphi: \mathbb{R} \rightarrow[0,1]$ such that $\varphi(x)=0$ for all $x \leq a / 2$ and $\varphi(x)=1$ for all $x \geq a$. For any $\beta>0$, we define the function $h_{\beta}$ as

$$
h_{\beta}(x)= \begin{cases}4 x^{2} / a^{2} & \text { if } x \in[0, a / 2], \\ P_{1}(x) & \text { if } x \in[a / 2, a], \\ P_{2}(x) & \text { if } x \in\left[a, B_{a}\right], \\ B_{a}^{\beta}(2 x)^{-\beta} & \text { if } x \geq B_{a},\end{cases}
$$

where

$$
P_{1}(x)=4 x^{2} / a^{2}(1-\varphi(x))+\varphi(x) P_{2}(x)
$$

and

$$
P_{2}(x)=2^{-\beta}-\frac{\beta 2^{-\beta}}{B_{a}}\left(x-B_{a}\right)+\frac{\beta(\beta+1) 2^{-\beta-1}}{B_{a}^{2}}\left(x-B_{a}\right)^{2}+C_{\beta}\left(x-B_{a}\right)^{4},
$$

with

$$
C_{\beta}=\frac{-2^{-\beta}+\beta 2^{-\beta}\left(a / B_{a}-1\right)-\beta(\beta+1) 2^{-\beta-1}\left(a / B_{a}-1\right)^{2}+1}{\left(a-B_{a}\right)^{4}} .
$$

Note that the coefficients of the polynomial $P-2$ have been chosen so that $h_{\beta}$ is $C^{2}$ at $B_{a}$ and $P_{2}(a)=1$.

There exists $M>0$ such that $C_{\beta}>0$ for all $\beta \geq M$, so that the first and second derivative of $P_{2}$ satisfy, for all $\beta \geq M$,

$$
\begin{aligned}
& P_{2}^{\prime}(x)=-\frac{\beta 2^{-\beta}}{B_{a}}+\frac{\beta(\beta+1) 2^{-\beta}}{B_{a}^{2}}\left(x-B_{a}\right)+4 C_{\beta}\left(x-B_{a}\right)^{3} \leq 0, \forall x \leq B_{a}, \\
& P_{2}^{\prime \prime}(x)=\frac{\beta(\beta+1) 2^{-\beta}}{B_{a}^{2}}+12 C_{\beta}\left(x-B_{a}\right)^{2} \geq 0, \forall x \in \mathbb{R} .
\end{aligned}
$$

In particular, $P_{2}$ is decreasing and convex on $\left(-\infty, B_{a}\right]$. Moreover, since $P_{2}(a)=1$, we deduce that $P_{2}(x) \geq 1$ for all $x \leq a$.

Finally, one easily deduces from the above definitions and properties that, for all $\beta \geq M, h_{\beta}$ is of regularity $C^{2}$, that $h_{\beta}(x) \geq 1$ for all $x \in[a / 2, a]$, that $h_{\beta}$ is decreasing and convex on $[a,+\infty)$. The finiteness of $M^{\prime}$ and $M^{\prime \prime}$ are simple consequences of the fact that the coefficients of $P_{2}$ are uniformly bounded in $\beta \geq M$. This concludes the proof of Lemma 4.2.

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