# A CYCLICALLY CATALYTIC SUPER-BROWNIAN MOTION ${ }^{1}$ 

By Klaus Fleischmann and Jie Xiong

Weierstrass Institute and University of Tennessee
In generalization of the mutually catalytic super-Brownian motion in $R$ of Dawson and Perkins and Mytnik, a function-valued cyclically catalytic model $X$ is constructed as a strong Markov solution to a martingale problem. Starting with a finite population $X_{0}$, each pair of neighboring types will globally segregate in the long-term limit (noncoexistence of neighboring types). Also finer extinction-survival properties depending on $X_{0}$ are studied in the spirit of Mueller and Perkins. In fact, $X_{0}$ can be chosen in such a way that all types survive for all finite times. On the other hand, sufficient conditions on $X_{0}$ are stated for the following situation: given a type $k$ and a positive time $t$, the $k$ th subpopulation $X^{k}$ dies by time $t$ with a large probability, provided that its initial value $X_{0}^{k}$ was sufficiently small.

## Contents

## 1. Introduction

1.1. Background and motivation
1.2. Rough description of the model
2. Results
2.1. Preliminaries: notations
2.2. Existence of $X$ and basic properties of all solutions
2.3. Global segregation of neighboring types
2.4. Finite time survival-extinction
3. Construction (proof of Theorem 3)
3.1. Construction of a solution to the martingale problem
3.2. Some properties of all martingale problem solutions ( $X, P_{r, x}$ )
3.3. The mapping $(r, x) \mapsto \mathscr{P}_{r, x}$
3.4. Selection of a strong Markov solution
4. Global segregation of neighboring types (proof of Theorem 4)
5. Finite time behavior (proof of Theorem 5).
5.1. Finite time survival of all types [proof of (a)]
5.2. Finite time extinction of a type [proof of (b)]

Acknowledgment
References

[^0]
## 1. Introduction.

1.1. Background and motivation. Recently Dawson and Perkins [11] and Mytnik [27] introduced and studied a mutually catalytic super-Brownian motion in R. This is a function-valued diffusion of two types of materials (species) where the small portions of mass ("particles") move chaotically in R but additionally branch (split or die) with a locally and temporally given rate proportional to the density of mass of the other type. Thus, each type serves as a catalyst for the other type's branching. This true interaction of types destroys the usual independence assumption in branching theory; in particular, this model is not a superprocess in its standard definition (for superprocesses, see, for instance, [14]). For a recent survey on catalytic and mutually catalytic branching models, we refer to [8, 9].

It is a natural desire to extend the mutually catalytic model to $K \geq 2$ types $A^{0}, \ldots, A^{K-1}$ of materials (as a rule, we write the index referring to the type as an upper index: please do not misunderstand the index as a power). We restrict to a cyclic situation, as often met in epidemics (see, for instance, [25]), networks of neurons (see, e.g., [17]), or biological competition models (see, e.g., [13]),

$$
\begin{equation*}
A^{k}+A^{k+1} \longrightarrow A^{k}, \quad k \in \mathrm{~K}, \tag{1}
\end{equation*}
$$

where $\mathrm{K}=\{0, \ldots, K-1\}$ denotes the cyclic group of size $K \geq 2$ (the additive group modulo $K$ ).

For treatments of cyclic reactions in terms of interacting particle systems, see [2], and [12]. Related to noise-induced transport phenomena, see [16], and in terms of deterministic equations, see [1, 22, 24, 30] (for instance).
1.2. Rough description of the model. A bit more precisely, we consider the following stochastic equation:

$$
\begin{equation*}
d X_{t}^{k}(a)=\frac{\sigma^{2}}{2} \Delta X_{t}^{k}(a) d t+\sqrt{\gamma^{k} X_{t}^{k}(a) X_{t}^{k+1}(a)} d W_{t}^{k}(a), \tag{2}
\end{equation*}
$$

$t>0,(k, a) \in \mathrm{K} \times \mathrm{R}$. Here the one-dimensional Laplacian $\Delta$ acts on the realvalued variable $a$, and ( $\left.\sigma^{2} / 2\right) \Delta$ reflects the chaotic motion of particles with diffusion constant $\sigma>0$. Moreover, the constants $\gamma^{k}>0$ are the interaction rates, and $d W$ denotes a standard white noise on $\mathbf{R}_{+} \times \mathrm{K} \times \mathrm{R}$.

The quantity $X_{t}^{k}(a)$ can be interpreted as the density of mass of type $k$ at time $t$ at site $a$. Intuitively, the subpopulation $X^{k}$ of $X$ of type $k$ evolves as a super-Brownian motion in R but with branching rate $\gamma^{k} X_{t}^{k+1}(a)$ changing with time $t$ and site $a$. Hence, the subpopulation $X^{k+1}$ serves as a catalyst for the branching of $X^{k}$, for each $k \in \mathrm{~K}$. Recall again that by this cyclic interaction over all the types, the basic independence assumption in branching theory is violated, so that $X$ is not a superprocess according to the usual definition.

Of course, in the special case $K=2$ we get the mutually catalytic branching model in R of [11] and [27]. (For further results on mutually catalytic models, see also [3-7,10, 26].)

Intuitively, a solution to (2) should be a (time-homogeneous) Markov process $X$. The first purpose of the paper is to establish that a weak solution $X$ to (2) exists which is a strong Markov process (see Theorem 3 below). Unfortunately, uniqueness of solutions remains open at this stage. The main obstacle for this is that as opposed to the mutually catalytic model, for $K \geq 3$ a self-duality [27] does not hold, and we also have not been able to find any other dual (or approximate dual) process for $X$. Nevertheless, each strong Markov solution $X$ to (2) we call a cyclically catalytic super-Brownian motion (SBM) in $\mathrm{K} \times \mathrm{R}$ (see also Definition 2 below). Besides the construction, we start the investigation of the survival-extinction behavior of cyclically catalytic SBMs in the case of finite populations (Theorems 4 and 5).

One expects that also a strong Markov $Z^{d}$-version of the cyclically catalytic model exists just as in the mutually catalytic case of [11]. The long-time results presented for the present cyclically catalytic SBMs in R should hold also for cyclically catalytic simple super-random walks in $Z^{d}$ (for Theorem 4: restrict to $d \leq 2$ ). The existence of a $\mathrm{R}^{2}$-version, however, remains open at this stage (note that for $K \geq 3$, moment dual processes or moment equations for $\varepsilon \mathrm{Z}^{2}$-approximations are much more complicated compared with the $K=2$ case, so that it is not clear how methods from [6, 7] could be extended).

## 2. Results.

2.1. Preliminaries: notations. With $c$ we always denote a positive constant which might vary from place to place. A $c$ with some additional mark (as $\bar{c}$ or $c_{1}$ ) will, however, denote a specific constant. A constant of the form $c_{(\#)}$ means, this constant first occurred related to formula line (\#).

For $\lambda \in \mathrm{R}$, introduce the reference function $\phi_{\lambda}$,

$$
\begin{equation*}
\phi_{\lambda}(a):=e^{-\lambda|a|}, \quad a \in \mathrm{R} \tag{3}
\end{equation*}
$$

(as usual, the colon attached to an equality sign " $=$ " refers to the side of the introduced notation). For $f: \mathrm{K} \times \mathrm{R} \rightarrow \mathrm{R}$, put

$$
\begin{equation*}
|f|_{\lambda}:=\sup _{k \in \mathrm{~K}, a \in \mathrm{R}}\left|f^{k}(a)\right| / \phi_{\lambda}(a), \quad \lambda \in \mathrm{R} . \tag{4}
\end{equation*}
$$

(Note that compared with [11] we reversed the sign in the definition of $\phi_{\lambda}$ but we kept it in the definition of $|\cdot|_{\lambda}$, and, concerning this, we use the same conventions as in [31].)

At some places we will need also a smoothed version $\tilde{\phi}_{\lambda}$ of $\phi_{\lambda}$. For this purpose, introduce the mollifier

$$
\begin{equation*}
\rho(a):=c_{(5)} \mathbf{1}_{\{|a|<1\}} \exp \left[-1 /\left(1-a^{2}\right)\right], \quad a \in \mathrm{R}, \tag{5}
\end{equation*}
$$

with $c_{(5)}$ the normalizing constant such that $\int d a \rho(\alpha)=1$. For $\lambda \in \mathrm{R}$, set

$$
\begin{equation*}
\tilde{\phi}_{\lambda}(a):=\int d b \phi_{\lambda}(b) \rho(b-a), \quad a \in \mathrm{R} . \tag{6}
\end{equation*}
$$

Note that to each $\lambda \in \mathrm{R}$ and $m \geq 0$ there are positive constants $\underline{c}_{\lambda, m}$ and $\bar{c}_{\lambda, m}$ such that

$$
\begin{equation*}
\underline{c}_{\lambda, m} \phi_{\lambda}(a) \leq\left|\frac{d^{m}}{d a^{m}} \tilde{\phi}_{\lambda}(a)\right| \leq \bar{c}_{\lambda, m} \phi_{\lambda}(a), \quad a \in \mathrm{R} \tag{7}
\end{equation*}
$$

(cf. [23], (2.1)).
For $\lambda \in \mathrm{R}$, let $\mathscr{C}_{\lambda}=\mathscr{C}_{\lambda}(\mathrm{K} \times \mathrm{R})$ denote the set of all continuous (real-valued) functions $f$ on $\mathrm{K} \times \mathrm{R}$ such that $|f|_{\lambda}$ is finite, and such that $f^{k}(\alpha) / \phi_{\lambda}(a)$ has a finite limit as $|a| \uparrow \infty$, for each $k \in \mathrm{~K}$. Introduce the spaces

$$
\begin{equation*}
\mathscr{C}_{\text {tem }}=\mathscr{C}_{\text {tem }}(\mathrm{K} \times \mathrm{R}):=\bigcap_{\lambda>0} \mathscr{C}_{-\lambda}, \quad \mathscr{C}_{\text {rap }}=\mathscr{C}_{\text {rap }}(\mathrm{K} \times \mathrm{R}):=\bigcap_{\lambda>0} \mathscr{C}_{\lambda} \tag{8}
\end{equation*}
$$

of tempered and rapidly decreasing functions, respectively. (Roughly speaking, the functions in $\mathscr{C}_{\text {tem }}$ are allowed to have a subexponential growth, whereas the ones in $\mathscr{C}_{\text {rap }}$ decay faster than exponentially.) Write $\mathscr{C}_{\text {rap }}^{(m)}=\mathscr{C}_{\text {rap }}^{(m)}(\mathrm{K} \times \mathrm{R})$ if we additionally require that all partial derivatives $\partial^{m} / \partial a^{m}$ up to the order $m \geq 1$ belong to $\mathscr{C}_{\text {rap }}$.

For each $\lambda \in$ R, the linear space $\mathscr{C}_{\lambda}$ equipped with the norm $|\cdot|_{\lambda}$ is a separable Banach space. The spaces $\mathscr{C}_{\text {tem }}$ and $\mathscr{C}_{\text {rap }}$ are topologized by the metrics

$$
\begin{align*}
& d_{\mathrm{tem}}(f, g):=\sum_{n=1}^{\infty} 2^{-n}\left(|f-g|_{-1 / n} \wedge 1\right), \quad f, g \in \mathscr{C}_{\mathrm{tem}}  \tag{9}\\
& d_{\mathrm{rap}}(f, g):=\sum_{n=1}^{\infty} 2^{-n}\left(|f-g|_{1 / n} \wedge 1\right), \quad f, g \in \mathscr{C}_{\mathrm{rap}} \tag{10}
\end{align*}
$$

making them Polish spaces. Similarly, we also define in $\mathscr{C}_{\text {rap }}^{(m)}, m \geq 1$, metrics in the obvious way to make them Polish.

Write $\Omega:=\mathscr{C}\left(\mathrm{R}_{+}, \mathscr{C}_{\text {tem }}^{+}\right)$for the set of all continuous paths $t \mapsto \omega_{t} \in \mathscr{C}_{\text {tem }}$. Equipped with the metric

$$
\begin{equation*}
d_{\Omega}\left(\omega, \omega^{\prime}\right):=\sum_{n=1}^{\infty} 2^{-n}\left(\sup _{0 \leq t \leq n} d_{\mathrm{tem}}\left(\omega_{t}, \omega_{t}^{\prime}\right) \wedge 1\right), \quad \omega, \omega^{\prime} \in \Omega \tag{11}
\end{equation*}
$$

$\Omega$ is a Polish space. The $\sigma$-field of all Borel subsets of $\Omega$ is denoted by $\mathscr{T}$.
If $E$ is a topological space, a measure on $E$ is meant to be a measure defined on the $\sigma$-field of all Borel subsets of $E$.

Let $\mathscr{P}$ denote the set of all probability measures on $\Omega$. Endowed with the Prohorov metric $d_{\mathscr{P}}$, we get a Polish space ([15], Theorem 3.1.7). Write com( $\mathscr{P}$ ) for the collection of all compact subsets of $\mathscr{P}$, equipped with the metric

$$
d_{\mathrm{com}}\left(K_{1}, K_{2}\right):=\inf \left\{\varepsilon>0: K_{1} \subseteq K_{2}^{\varepsilon} \text { and } K_{2} \subseteq K_{1}^{\varepsilon}\right\}, \quad K_{1}, K_{2} \in \operatorname{com}(\mathscr{P})
$$

where $K^{\varepsilon}$ is the $\varepsilon$-neighborhood of $K$ (based on $d_{\mathscr{P}}$ ). Then the metric space $\left(\operatorname{com}(\mathscr{P}), d_{\text {com }}\right)$ is separable ([32], Lemma 12.1.1).

As a rule, the processes $X=\left\{X_{t}: t \geq 0\right\}$ considered in this paper are $\mathscr{C}_{\text {tem }}^{+}-$ valued, continuous and presented in their canonical form. That is, we identify each process $X$ with a probability law $P$ on $\Omega=\mathscr{C}\left(\mathrm{R}_{+}, \mathscr{C}_{\text {tem }}^{+}\right)$, in other words,
with the probability space $(\Omega, \mathscr{F}, P)$. More precisely, we always consider $(\Omega, \mathscr{F}, P)$ as a filtered probability space, using the usual filtration $\left\{\mathscr{F}_{t}: t \geq 0\right\}$. Write $\mathscr{F}^{r}$ for the sub- $\sigma$-field of $\mathscr{F}$ generated by the coordinate maps $\omega \mapsto \omega_{t}$, for $t \geq r$.

Let $d k$ denote the counting measure (Haar measure) on the cyclic group K [that is, $\int_{\mathrm{K}} d k f(k)=\sum_{k \in \mathrm{~K}} f(k)$ for all functions $f: \mathrm{K} \rightarrow \mathrm{R}_{+}$). For functions $f, g$ on $\mathrm{K} \times \mathrm{R}$ or R , we write $\langle f, g\rangle$ for the integral of $f \cdot g$ with respect to $d k d a$ or $d a$, respectively (if the integral makes sense). As opposed to the notation $|\cdot|_{\lambda}$ introduced in (4), for functions $f \geq 0$ on $\mathrm{K} \times \mathrm{R}$ or R we define

$$
\begin{equation*}
\|f\|_{\lambda}:=\left\langle f \cdot \tilde{\phi}_{-\lambda}, 1\right\rangle, \quad \lambda \in \mathrm{R} \tag{12}
\end{equation*}
$$

[with the smoothed reference function $\tilde{\phi}_{-\lambda}$ from (6)]. Set $\|f\|:=\|f\|_{0}$ for the "total mass" of the (density) function $f$.

Let $p$ denote the heat kernel in R related to $\left(\sigma^{2} / 2\right) \Delta$,

$$
\begin{equation*}
p_{t}(a):=\left(2 \pi \sigma^{2} t\right)^{-1 / 2} \exp \left[-\frac{|a|^{2}}{2 \sigma^{2} t}\right], \quad t>0, \quad a \in \mathrm{R} \tag{13}
\end{equation*}
$$

and $\left\{S_{t}: t \geq 0\right\}$ the corresponding heat flow semigroup. Write $\xi=\left(\xi, \Pi_{a}\right)$ for the related Brownian motion in R , with $\Pi_{a}$ denoting the law of $\xi$ if $\xi_{0}=a \in \mathrm{R}$.

The (usually upper) index + on a set of real-valued functions will refer to the collection of all nonnegative members of this set, similarly to our notation $\mathrm{R}_{+}=[0, \infty)$. The Kronecker symbol is denoted by $\delta_{k, l}$.
2.2. Existence of $X$ and basic properties of all solutions. A more precise formulation of the stochastic equation (2) can be given in terms of the following martingale problem $\mathbf{M P}_{x}$. Recall that $K \geq 2, \sigma>0$, and $\gamma^{h}>0, k \in \mathrm{~K}$.

Definition 1 (Martingale problem MP $\mathbf{M P}_{x}$ ). Fix $x \in \boldsymbol{b}_{\text {tem }}^{+}=b_{\text {tem }}^{+}(\mathrm{K} \times \mathrm{R})$. We say a stochastic process $X=\left\{X_{t}: t \geq 0\right\}$ with law $P_{x}$ on $\Omega=6\left(\mathrm{R}_{+}, \varnothing_{\text {tem }}^{+}\right)$ is a solution to the martingale problem $\mathbf{M P}_{x}$ if $P_{x}\left(X_{0}=x\right)=1$ and, for test functions $\varphi \in \mathscr{C}_{\text {rap }}^{(2)}=\mathscr{\sigma}_{\text {rap }}^{(2)}(\mathrm{K} \times \mathrm{R})$, setting

$$
\begin{equation*}
M_{t}^{k}\left(\varphi^{k}\right):=\left\langle X_{t}^{k}, \varphi^{k}\right\rangle-\left\langle x^{k}, \varphi^{k}\right\rangle-\int_{0}^{t} d s\left\langle X_{s}^{k}, \frac{\sigma^{2}}{2} \Delta \varphi^{k}\right\rangle \tag{14}
\end{equation*}
$$

$t \geq 0, k \in \mathrm{~K}$, one has orthogonal continuous square-integrable martingales $M^{k}\left(\varphi^{k}\right), k \in \mathrm{~K}$, starting from $M_{0}^{k}\left(\varphi^{k}\right) \equiv 0$, and with square functions

$$
\begin{equation*}
\left\langle\left\langle M^{k}\left(\varphi^{k}\right)\right\rangle_{t}=\gamma^{k} \int_{0}^{t} d s \int_{\mathrm{R}} d a X_{s}^{k}(a) X_{s}^{k+1}(a)\left[\varphi^{k}(a)\right]^{2}\right. \tag{15}
\end{equation*}
$$

$t \geq 0, k \in \mathrm{~K}$.
Now the definition of our basic object of interest follows.
Definition 2 (Cyclically catalytic SBM $X$ ). If ( $X, P_{x}, x \in \sigma_{\text {tem }}^{+}$) is a (time-homogeneous) strong Markov process such that ( $X, P_{x}$ ) is a solution to the martingale problem $\mathbf{M P}_{x}$ of Definition 1, for each $x \in \mathscr{b}_{\text {tem }}^{+}$, then it is called a cyclically catalytic super-Brownian motion (SBM) in $\mathrm{K} \times \mathrm{R}$ with diffusion constant $\sigma$ and interaction rate $\gamma=\left(\gamma^{k}\right)_{k \in \mathrm{~K}}$.

Here is our first result.
Theorem 3 (Cyclically catalytic SBM $X$ ).
(a) (Existence of $X$ ). To each $K \geq 2, \sigma>0$ and vector $\gamma>0$, there exists a cyclically catalytic super-Brownian motion ( $X, P_{x}, x \in \mathscr{b}_{\text {tem }}^{+}$) in $\mathrm{K} \times \mathrm{R}$ with diffusion constant $\sigma$ and interaction rate $\gamma$ according to Definition 2.
(b) (Finite moments). Each cyclically catalytic SBM X has finite moments of all orders: for fixed $c_{0}, T, q>0$ and $\lambda^{\prime}, \lambda \in \mathrm{R}$ with $q \lambda^{\prime}<\lambda$,

$$
\sup _{x \in E_{\text {tem }}^{+},\left.|x|\right|_{-\lambda^{\prime}} \leq c_{0}} P_{x} \sup _{0 \leq t \leq T} \sum_{k \in \mathrm{~K}}\left\langle\left(X_{t}^{k}\right)^{q}, \phi_{\lambda}\right\rangle<\infty .
$$

The expectation of $X$ is given by

$$
P_{x} X_{t}^{k}(a)=S_{t} x^{k}(a), \quad x \in \mathscr{b}_{\text {tem }}^{+}, \quad(t, k, a) \in \mathrm{R}_{+} \times \mathrm{K} \times \mathrm{R},
$$

and the covariance by

$$
\begin{aligned}
& \operatorname{Cov}_{x}\left(X_{t_{1}}^{k_{1}}\left(a_{1}\right), X_{t_{2}}^{k_{2}}\left(a_{2}\right)\right)= \gamma^{k_{1}} \delta_{k_{1}, k_{2}} \int_{0}^{t_{1} \wedge t_{2}} d s \int_{\mathrm{R}} d b S_{s} x^{k_{1}}(b) S_{s} x^{k_{1}+1}(b) \\
& \times \mathrm{p}_{t_{1}-s}\left(a_{1}-b\right) \mathrm{p}_{t_{2}-s}\left(a_{2}-b\right), \\
& x \in \mathscr{C}_{\text {tem }}^{+}, t_{1}, t_{2} \geq 0, k_{1}, k_{2} \in \mathrm{~K}, a_{1}, a_{2} \in \mathrm{R} .
\end{aligned}
$$

Note that the covariance vanishes only if $x=0, k_{1} \neq k_{2}$, or $t_{1} \wedge t_{2}=0$. In particular, the process $X$ is nondegenerate.

Recall that the novelty of this theorem concerns the case $K \geq 3$, since $K=2$ is due to [11], and that uniqueness remains unsolved if $K \geq 3$.

The proof of Theorem 3 will be provided in Section 3 below. There we will start from an approximating system of processes where on small time intervals we consider $K$ conditionally independent catalytic super-Brownian motions with frozen, smoothed and truncated branching rate functions (catalyst). Then tightness will be shown by an adoption of a method used in [11] which was based on [31]. This then yields the existence of solutions to the martingale problem $\mathbf{M P}_{x}$ of Definition 1. Note that the existence of a weak solution to the stochastic equation (2) (on an enlarged probability space) then follows from the standard martingale representation theorem [35]. We mention also that the convolution form of (2) is given in (57) below.

Since uniqueness in the martingale problem is not established, some more efforts are needed to construct a Markov solution to the martingale problem. Moreover, since the topic of continuous dependence on the initial data of constructed solutions is also a delicate unsolved problem in the present model, we could not follow the usual route to deduce the strong Markov property from a Feller property. Nevertheless, by an adoption of methods developed in [32] for finite-dimensional diffusions, we succeeded in selecting a time-homogeneous strong Markov process from the set of all solutions of the family of martingale problems.

Part (a) of Theorem 3 is implied by Theorem 23 below, whereas (b) follows from Corollary 16.
2.3. Global segregation of neighboring types. Now we restrict our attention to any cyclically catalytic SBM $X$ as introduced in Definition 2 (which exists by Theorem 3), and fix its initial state $X_{0}$.

In the mutually catalytic model (in R), the self-duality is a powerful tool not only for establishing the uniqueness in the martingale problem, but also to get results on the long-term behavior [11]. In fact, the total mass process $t \mapsto\left(\left\|X_{t}^{0}\right\|,\left\|X_{t}^{1}\right\|\right)$ in the case of finite initial masses $\left\|X_{0}\right\|=X_{0}^{0}(\mathrm{R})+X_{0}^{1}(\mathrm{R})$ is a nonnegative martingale, and its a.s. limit $\left(\left\|X_{\infty}^{0}\right\|,\left\|X_{\infty}^{1}\right\|\right)$, say, can be identified in relatively simple terms. Indeed, it coincides in law with the state $B_{\tau}$ of a Brownian motion $B$ in $\mathrm{R}_{+}^{2}$ in its first hitting time $\tau$ of the boundary $\partial \mathrm{R}_{+}^{2}$ of $\mathrm{R}_{+}^{2}$, if $B$ was started from $B_{0}=\left(\left\|X_{0}^{0}\right\|,\left\|X_{0}^{1}\right\|\right)$ (see the proof of Theorem $1.2(\mathrm{a})$ in [11]). In particular, the limit population is nondegenerate and has full expectation (persistence).

Of course, the present cyclically catalytic model also has that convergence property,

$$
\begin{equation*}
\lim _{t \uparrow \infty}\left(\left\|X_{t}^{0}\right\|, \ldots,\left\|X_{t}^{K-1}\right\|\right)=:\left(\left\|X_{\infty}^{0}\right\|, \ldots,\left\|X_{\infty}^{K-1}\right\|\right) \quad \text { exists a.s. } \tag{16}
\end{equation*}
$$

provided that $\left\|X_{0}\right\|<\infty$. But we have not been able to identify the limit (16). An obstacle is that the random time-change argument of [11] is not as powerful, since it leads to $K$ Brownian motions which run with different clocks, as opposed to the K = 2 case. In other words, in the terminology of Swart [33], the $K \geq 3$ case is an anisotropic situation, which is much more delicate than the isotropic $K=2$ case. Nevertheless, we are able to verify the following "global segregation" (noncoexistence) of neighboring types in the limit, which in the $K=2$ case is a simple consequence of a property of the hitting state $B_{\tau}$, namely that $B_{\tau}^{0} B_{\tau}^{1}=0$. Recall that $K \geq 2$.

Theorem 4 (Global segregation of neighboring types). Start any cyclically catalytic super-Brownian motion $X$ with a finite initial mass $\left\|X_{0}\right\|$. Then, for each $k \in \mathrm{~K}$,

$$
\begin{equation*}
\lim _{t \uparrow \infty}\left\|X_{t}^{k}\right\| \cdot\left\|X_{t}^{k+1}\right\|=0 \quad \text { a.s. } \tag{17}
\end{equation*}
$$

Consequently, for each pair of neighboring types, only one of them has the chance to survive in the limit.

The proof of Theorem 4 will be given in Section 4. It will be based on a modification of arguments of [11], adapted to the aforementioned case of different clocks if $K \geq 3$. The strategy of proof is as follows. Set

$$
\begin{equation*}
Z_{t}:=\sum_{k \in \mathrm{~K}} \gamma^{k}\left\|X_{t}^{k}\right\| \cdot\left\|X_{t}^{k+1}\right\|, \quad t \geq 0 . \tag{18}
\end{equation*}
$$

Assuming now that

$$
\begin{equation*}
\inf _{t \geq 0} Z_{t}>0 \text { with positive probability, } \tag{19}
\end{equation*}
$$

our task is to construct stopping times $T_{1}, T_{2}, \ldots$ such that $Z_{T_{n}} \rightarrow 0$ as $n \uparrow \infty$ on the event in (19). But this is an obvious contradiction. Then the almost sure convergence of the martingales $t \mapsto\left\|X_{t}^{k}\right\|, k \in \mathrm{~K}$, will yield the claim (17).

As opposed to the mutually catalytic case, Theorem 4 in particular leaves open whether for $K \geq 3$ the limit ( $\left\|X_{\infty}^{0}\right\|, \ldots,\left\|X_{\infty}^{K-1}\right\|$ ) in (16) is nondegenerate, and whether it has full expectation (persistence).
2.4. Finite time survival-extinction. For the mutually catalytic model in $\mathbf{Z}^{d}$ (also established and investigated in [11]), a recent paper [26] addresses the following questions: Is it possible that, depending on the finite initial state, both types survive all finite times a.s., or that one of the types dies in a given finite time with high probability? The following results on our cyclic model are in that spirit. Recall the reference function $\phi_{\lambda}$ introduced in (3).

Theorem 5 (Finite time behavior). Fix again any cyclically catalytic SBM $X$ with $X_{0} \in \mathscr{b}_{\text {tem }}^{+}$satisfying $\left\|X_{0}\right\|<\infty$.
(a) (Finite time survival of all types). Assume that $\prod_{k \in K}\left\|X_{0}^{k}\right\|>0$, and that there is a $T>0$ such that

$$
\begin{equation*}
\max _{k \in \mathrm{~K}} \liminf _{|a| \uparrow \infty} \frac{S_{t}\left[X_{0}^{k}\right]^{2}(a) S_{t}\left[X_{0}^{k+1}\right]^{2}(a)}{\left[S_{t} X_{0}^{k}(a)\right]^{4}}=0, \quad t \geq T . \tag{20}
\end{equation*}
$$

Then

$$
\begin{equation*}
\prod_{k \in \mathrm{~K}}\left\|X_{t}^{k}\right\|>0 \quad \text { for all } t>0 \text { a.s. } \tag{21}
\end{equation*}
$$

(b) (Finite time extinction of a type with high probability). Fix a type $k_{0} \in$ K. For $i=0,1,2$, consider positive constants $c_{i}, \lambda_{i}$, and $c_{1}^{\prime}, \lambda_{1}^{\prime}$ with

$$
\begin{equation*}
\lambda_{0}>\lambda_{1}^{\prime}>\lambda_{1}, \quad 2 \lambda_{1}^{\prime}<\lambda_{1}+\lambda_{2} \quad \text { and } \quad c_{1}^{\prime} \leq c_{1} . \tag{22}
\end{equation*}
$$

Then the following statement holds. For $\varepsilon \in(0,1]$ and $T>0$ fixed, $c_{0}$ can be chosen so small that if the initial state $X_{0}=x \in \mathscr{b}_{\text {tem }}^{+}$is such that

$$
\begin{equation*}
x^{k_{0}} \leq c_{0} \phi_{\lambda_{0}} \tag{23}
\end{equation*}
$$

as well as

$$
\begin{equation*}
|x|_{\lambda_{1}} \leq c_{1}, \quad x^{k_{0}+2} \leq c_{2} \phi_{\lambda_{2}} \quad \text { and } \quad x^{k_{0}+1} \geq c_{1}^{\prime} \phi_{\lambda_{1}^{\prime}}, \tag{24}
\end{equation*}
$$

then

$$
\begin{equation*}
P_{x}\left(X_{t}^{k_{0}}=0 \text { for } t \geq T\right) \geq 1-\varepsilon . \tag{25}
\end{equation*}
$$

The proof of Theorem 5 in Section 5 below uses ideas of [26]. Of course, the condition (20) in Theorem 5(a) looks a bit complicated, so we have to discuss it. Roughly speaking, it is, for instance, satisfied if the initial states of each pair of neighboring types are separated in different half axes and have sufficiently large mass tails. This will now be made more precise in the following example.

Example 6 (Starting from separated neighbors with large tails). Assume $K \geq 2$ is even and that

$$
\begin{equation*}
X_{0}^{2 k}:=\phi_{1} \mathbf{1}_{\mathrm{R} \backslash \mathrm{R}_{+}} \quad \text { and } \quad X_{0}^{2 k+1}:=\phi_{1} \mathbf{1}_{\mathrm{R}_{+}}, \quad k \in \mathrm{~K} \tag{26}
\end{equation*}
$$

[with the reference function $\phi_{1}$ from (3)], ignoring the discontinuity at $0 \in R$, which can simply be overcome by a smoothing procedure, for instance using the mollifier $\rho$ from (5). Then the simultaneous finite time survival as claimed in (21) holds. In fact,

$$
\begin{align*}
S_{t} X_{0}^{2 k}(a) & =\mathscr{N}\left(\frac{-a-t}{\sqrt{t}}\right) e^{a+t / 2},  \tag{27a}\\
S_{t}\left[X_{0}^{2 k}\right]^{2}(a) & =\mathscr{N}\left(\frac{-a-2 t}{\sqrt{t}}\right) e^{2 a+2 t},  \tag{27b}\\
S_{t}\left[X_{0}^{2 k+1}\right]^{2}(a) & =\mathscr{N}\left(\frac{a-2 t}{\sqrt{t}}\right) e^{-2 a+2 t}, \tag{27c}
\end{align*}
$$

with $\mathscr{N}$ denoting the distribution function of the standard normal law on R . As $a \downarrow-\infty$, the $\mathscr{N}$-expressions in (27a) and (27b) tend to 1 , for fixed $t>0$. Therefore, the ratio in assumption (20) with $k$ replaced by $2 k$ is of order $\mathscr{N}(a / \sqrt{t}) / e^{4 a}$ as $a \downarrow-\infty$, hence converges to zero by l'Hôpital's rule. On the other hand, if we shift the type by one, then we get the same order of decay if $a \uparrow \infty$ instead. Altogether, assumption (20) is fulfilled, hence (21) holds.

The philosophy behind the proof of Theorem 5(b) is as follows. Since 0 is an absorbing state for the subprocess $t \mapsto X_{t}^{k_{0}}$, it suffices to consider $X$ on a possibly smaller time interval $[0, T]$. Moreover, because initially the catalyst $X^{k_{0}+2}$ for $X^{k_{0}+1}$ is not too large by assumption, $X^{k_{0}+1}$ should not be very small on $[0, T]$, and since $X^{k_{0}+1}$ serves as the catalyst for $X^{k_{0}}$, the latter should have some chance to die by time $T$. Actually, we want to bound $\left\|X_{T}^{k_{0}}\right\|_{\lambda}$ from above (for an appropriate $\lambda$ and on a suitable new time scale) by a supercritical Feller's branching diffusion, which of course dies by a given time with a positive probability. Making then its initial state $\left\|x^{k_{0}}\right\|_{\lambda}$ sufficiently small, this extinction probability can be forced to be sufficiently close to one.

Remark 7 (Property of all solutions). As can be seen from the proof in Section 5 , Theorem 5 actually applies for any family $\left\{\left(X, P_{x}\right): x \in \mathscr{C}_{\text {tem }}^{+}\right\}$of processes such that $P_{x}$ solves $\mathbf{M P}_{x}, x \in \mathscr{C}_{\text {tem }}^{+}$, that is, without requiring the Markov property.

Unfortunately, we do not have results on the long-time behavior of $X$ for infinite initial populations (for $K \geq 3$ ). Note that the study of the long-term behavior of the mutually catalytic model in the case of infinite initial populations (see $[3,4,5,7,11]$ ) also relies heavily on the self-duality of the model. In fact, via self-duality, it is based on the long-term behavior of the finite mass system.
3. Construction (proof of Theorem 3). We will start with proving the existence of a solution $X$ to the martingale problem $\mathbf{M P}_{x}$ of Definition 1. Then a time-homogeneous strong Markov solution will be selected from the set of all solutions to the family of martingale problems $\mathbf{M P} \mathbf{P}_{x}, x \in \mathscr{C}_{\text {tem }}^{+}$, as needed for Theorem 3.
3.1. Construction of a solution to the martingale problem. In this subsection, we want to verify that the martingale problem $\mathbf{M P}_{x}$ of Definition 1 has a solution. To this aim, we will start from some approximations (Definition 10), and will verify some properties of them (Lemmas 11-13), which turn out to be owned also by all solutions of the martingale problem $\mathbf{M P}_{x}$ (see Lemma 14 in Section 3.2).

To prepare for the selection of a strong Markov solution, a timeinhomogeneous point of view will be convenient to use: we start the process at times $r \geq 0$ (the model will still be time-homogeneous). On the other hand, for the sake of working with a single path space, we formally extend the paths backward by assuming that they are constant in the interval $[0, r]$.

DEFINITION 8 (Martingale problem MP $\mathbf{M P}_{r, x}$ ). Fix $(r, x) \in \mathrm{R}_{+} \times \mathscr{C}_{\text {tem }}^{+}$. We say a stochastic process $X=\left\{X_{t}: t \geq 0\right\}$ with law $P_{r, x}$ on $\Omega=\mathscr{C}\left(\mathrm{R}_{+}, \mathscr{b}_{\text {tem }}^{+}\right)$is $a$ solution to the martingale problem $\mathbf{M P}_{r, x}$ if the following three conditions hold:
(i) $P_{r, x}\left(X_{t}=x\right.$ for all $\left.t \leq r\right)=1$.
(ii) For test functions $\varphi \in \mathscr{C}_{\text {rap }}^{(2)}$, setting

$$
M_{r, t}^{k}\left(\varphi^{k}\right):=\left\langle X_{t}^{k}, \varphi^{k}\right\rangle-\left\langle x^{k}, \varphi^{k}\right\rangle-\int_{r}^{t} d s\left\langle X_{s}^{k}, \frac{\sigma^{2}}{2} \Delta \varphi^{k}\right\rangle
$$

$t \geq r, k \in \mathrm{~K}$, one has orthogonal continuous square-integrable martingales $t \mapsto M_{r, t}^{k}\left(\varphi^{k}\right), k \in \mathrm{~K}$, (after time $r$ ) starting from $M_{r, r}^{k}\left(\varphi^{k}\right) \equiv 0$.
(iii) The square functions satisfy

$$
\left\langle\left\langle M_{r, .}^{k}\left(\varphi^{k}\right)\right\rangle\right\rangle_{t}=\gamma^{k} \int_{r}^{t} d s \int_{\mathrm{R}} d a X_{s}^{k}(a) X_{s}^{k+1}(a)\left[\varphi^{k}(a)\right]^{2}
$$

$t \geq r, k \in \mathrm{~K}$.
PROPOSITION 9 (Existence of a solution to the martingale problem $\mathbf{M P}_{r, x}$ ). For each $(r, x) \in \mathrm{R}_{+} \times \mathscr{C}_{\text {tem }}^{+}$, there exists a solution $\left(X, P_{r, x}\right)$ to the martingale problem $\mathbf{M P}_{r, x}$ of Definition 8.

For the verification of this proposition, we partly borrow ideas from the proof of Theorem 6.1 in [11] (starting from page 1133), which are in part based on [31], Appendix. So first we will introduce in Definition 10 below an approximating sequence $\left\{{ }^{n} X: n \geq 1\right\}$ of continuous $\mathscr{C}_{\text {tem }}^{+}$-valued processes. ${ }^{n} X$ has the property that on small time periods $[i / n,(i+1) / n), i \geq 0$, given ${ }^{n} X_{i / n}$ the single subpopulations ${ }^{n} X^{k}$ behave as independent continuous catalytic
super-Brownian motions in R with frozen, smoothed and bounded branching rate function (catalyst) given by $\gamma^{k}\left(S_{1 / n}{ }^{n} X_{i / n}^{k+1} \wedge n\right)$ (see [11]). (The additional smoothing with $S_{1 / n}$-recall that $S$ denotes the heat flow semigroup-will help us to make working a Gronwall's inequality argument in the proof of Lemma 12 below.). Then we pass to a pointwise stochastic equation (Lemma 11) and use it to derive some moment estimates (Lemmas 12 and 13). After these preparations, Proposition 9 then easily follows by an application of [31], Lemma 6.3(ii).

We start with introducing the system $\left\{{ }^{n} X: n \geq 1\right\}$ of approximating $\mathscr{C}_{\text {tem }}^{+}-$ valued processes.

DEFINITION 10 (Martingale problem $\mathbf{M P}_{r, x}^{n}$ ). For $n \geq 1,(r, x) \in \mathrm{R}_{+} \times \mathscr{b}_{\text {tem }}^{+}$, let ( ${ }^{n} X, P_{r, x}^{n}$ ) denote the unique (in law) process with the following two properties. First, ${ }^{n} X_{t} \equiv x$ for $t \leq r$. On the other hand, for $\varphi \in \mathscr{C}_{\text {rap }}^{(2)}$, setting

$$
\begin{equation*}
{ }^{n} M_{r, t}^{k}\left(\varphi^{k}\right):=\left\langle{ }^{n} X_{t}^{k}, \varphi^{k}\right\rangle-\left\langle x^{k}, \varphi^{k}\right\rangle-\int_{r}^{t} d s\left\langle{ }^{n} X_{s}^{k}, \frac{\sigma^{2}}{2} \Delta \varphi^{k}\right\rangle \tag{28}
\end{equation*}
$$

$t \geq r, k \in \mathrm{~K}$, one has orthogonal continuous square-integrable martingales $t \mapsto{ }^{n} M_{r, t}^{k}\left(\varphi^{k}\right), k \in \mathrm{~K}$, starting from ${ }^{n} M_{r, r}^{k}\left(\varphi^{k}\right)=0$ and with square functions

$$
\begin{equation*}
\left\langle\left\langle{ }^{n} M_{r, \cdot}^{k}\left(\varphi^{k}\right)\right\rangle\right\rangle_{t}=\gamma^{k} \int_{r}^{t} d s \int_{\mathrm{R}} d b^{n} X_{s}^{k}(b)\left(S_{1 / n}^{n} X_{[n s] / n}^{k+1}(b) \wedge n\right)\left[\varphi^{k}(b)\right]^{2} \tag{29}
\end{equation*}
$$

$t \geq r, k \in \mathrm{~K}$.
Note that the uniqueness in law can be proved via log-Laplace representations of this system of "piecewise independent superprocesses."

The family (28) of martingales extends [35] to orthogonal square-integrable martingale measures ${ }^{n} M_{r}^{k}={ }^{n} M_{r}^{k}(d(s, b))$ and to the usual class of predictable integrands. Moreover, for $t \geq r$ and $\varphi \in \mathscr{C}_{\text {rap }}$, as well as $k \in \mathrm{~K}$ fixed, the function $(s, a) \mapsto S_{t-s} \varphi^{k}(a)$ on $[r, t] \times \mathrm{R}$ can be included as integrand of the stochastic integrals. Then ${ }^{n} X$ can be shown to satisfy the following stochastic equation:

$$
\begin{equation*}
\left\langle{ }^{n} X_{t}^{k}, \varphi^{k}\right\rangle=\left\langle x^{k}, S_{t-r} \varphi^{k}\right\rangle+\int_{[r, t] \times \mathrm{R}}{ }^{n} M_{r}^{k}(d(s, b)) S_{t-s} \varphi^{k}(b) \tag{30}
\end{equation*}
$$

$P_{r, x}^{n}$-a.s., for $t \geq r, k \in \mathrm{~K}, \varphi \in \mathscr{C}_{\text {rap }}$. This, in particular, immediately implies the following moment formulas:

$$
\begin{align*}
P_{r, x}^{n}{ }^{n} X_{t}^{k}(a) & =S_{t-r} x^{k}(a)  \tag{31}\\
P_{r, x}^{n}{ }^{n} X_{t}^{k}(a)^{n} X_{t}^{l}(a) & =S_{t-r} x^{k}(a) S_{t-r} x^{l}(a) \tag{32}
\end{align*}
$$

for all $t \geq r, a \in \mathrm{R}$, and $k \neq l$. Moreover, replacing $\varphi^{k}$ by $\mathrm{p}_{\varepsilon}(\cdot-a)$ in (30), where $0<\varepsilon \leq 1$, and $a \in \mathrm{R}$ are fixed, gives

$$
\begin{equation*}
S_{\varepsilon}{ }^{n} X_{t}^{k}(a)=S_{\varepsilon+t-r} x^{k}(a)+\int_{[r, t] \times \mathrm{R}}{ }^{n} M_{r}^{k}(d(s, b)) \mathrm{p}_{\varepsilon+t-s}(b-a) \tag{33}
\end{equation*}
$$

$P_{r, x}^{n}$-a.s. We want to let $\varepsilon \downarrow 0$.

LEMMA 11 (Pointwise equation for ${ }^{n} X$ ). For $n \geq 1,(r, x) \in \mathrm{R}_{+} \times \mathscr{C}_{\text {tem }}^{+}, t \geq r$ and $(k, a) \in \mathrm{K} \times \mathrm{R}$ fixed,

$$
\begin{equation*}
{ }^{n} X_{t}^{k}(a)=S_{t-r} x^{k}(a)+\int_{[r, t] \times \mathrm{R}}{ }^{n} M_{r}^{k}(d(s, b)) \mathrm{p}_{t-s}(b-a), \quad P_{r, x}^{n}-a . s \tag{34}
\end{equation*}
$$

(reading the integral term as 0 if $t=r$ ).
Proof. Fix $n, r, x, t, k, a$ as in the lemma. To check that the stochastic integrals in (33) converge in $L^{2}$ as $\varepsilon \downarrow 0$ to the one in (34), consider

$$
\begin{array}{r}
P_{r, x}^{n}\left(\int_{[r, t] \times \mathrm{R}}{ }^{n} M_{r}^{k}(d(s, b))\left[p_{\varepsilon+t-s}(b-a)-p_{t-s}(b-a)\right]\right)^{2} \\
=\gamma^{k} P_{r, x}^{n} \int_{r}^{t} d s \int_{\mathrm{R}} d b^{n} X_{s}^{k}(b)\left(S_{1 / n}{ }^{n} X_{[n s] / n}^{k+1}(b) \wedge n\right)  \tag{35}\\
\times\left[p_{\varepsilon+t-s}(b-a)-p_{t-s}(b-a)\right]^{2}
\end{array}
$$

(which holds by the well-known isometry properties of stochastic integration). By the mixed moment formula (32) we may continue with

$$
\begin{align*}
\leq c \int_{r}^{t} d s & \int_{\mathrm{R}} d b S_{s-r} x^{k}(b) S_{[1+n s] / n-r} x^{k+1}(b)  \tag{36}\\
\times & {\left[p_{\varepsilon+t-s}(b-a)-p_{t-s}(b-a)\right]^{2} }
\end{align*}
$$

By definition, for fixed $x \in \mathscr{C}_{\text {tem }}^{+}$and $\lambda>0$,

$$
\begin{equation*}
x^{k} \leq c \phi_{-\lambda} \tag{37}
\end{equation*}
$$

On the other hand, for fixed $T>0$ and $\lambda \in \mathrm{R}$ there are constants $\underline{c}$ and $\bar{c}$ such that

$$
\begin{equation*}
\underline{c} \phi_{\lambda} \leq S_{s} \phi_{\lambda} \leq \bar{c} \phi_{\lambda}, \quad 0 \leq s \leq T \tag{38}
\end{equation*}
$$

(cf. [31], Lemma 6.2(ii)). Hence, by Cauchy-Schwarz, the estimate (36) can be continued with

$$
\begin{align*}
\leq c & \left(\int_{0}^{t} d s \int_{\mathrm{R}} d b\left[p_{\varepsilon+t-s}(b-a)-p_{t-s}(b-a)\right]^{2}\right)^{1 / 2}  \tag{39}\\
& \times\left(\int_{0}^{t} d s \int_{\mathrm{R}} d b\left[p_{\varepsilon+t-s}^{2}(b-a)+p_{t-s}^{2}(b-a)\right] \phi_{-4 \lambda}(b)\right)^{1 / 2}
\end{align*}
$$

But the first term can be bounded by $c \varepsilon^{1 / 4}$ (cf. [31], Lemma 6.2(i)), whereas for the second term we use again (38) to get the bound

$$
\begin{equation*}
c \phi_{-2 \lambda}(a)\left(\int_{0}^{t} d s\left[(\varepsilon+t-s)^{-1 / 2}+(t-s)^{-1 / 2}\right]\right)^{1 / 2} \leq c \phi_{-2 \lambda}(a) \tag{40}
\end{equation*}
$$

Altogether, (35) tends to zero as $\varepsilon \downarrow 0$, for fixed $k, r, t, a$ and $x$, uniformly in $n$. Thus (33) implies (34), completing the proof of the lemma.

Since for the later selection of a strong Markov solution we will need some measurable dependence on the initial data $(r, x)$, in the construction we will already allow that the ${ }^{n} X$ additionally depend on some varying initial data ( $r_{m}, x_{m}$ ), and we will write ${ }^{m, n} X$ instead of ${ }^{n} X$.

To be more precise, consider $\left(r_{m}, x_{m}\right) \in \mathrm{R}_{+} \times \mathscr{b}_{\text {tem }}^{+}, m \geq 1$. Fix now $m, n \geq 1$, hence $\left(r_{m}, x_{m}\right) \in \mathrm{R}_{+} \times \mathscr{b}_{\text {tem }}^{+}$for the moment. By definition, ( ${ }^{m, n} X, P_{r_{m}, x_{m}}^{n}$ ) is the unique solution to the martingale problem $\mathbf{M P}_{r_{m}, x_{m}}^{n}$ of Definition 10, and for the related martingale measures we write now ${ }^{m, n} M_{r_{m}}^{k}$ instead of ${ }^{n} M_{r}^{k}$. Recall the notation $|\cdot|_{\lambda}$ from (4).

LEMMA 12 (Uniformly bounded moments for ${ }^{m, n} X$ ). For fixed $c_{0}, T, q>0$ and $\lambda^{\prime}, \lambda \in \mathrm{R}$ with $2 q \lambda^{\prime}<\lambda$,

$$
\begin{equation*}
\sup _{\substack{m, n \geq 1, r_{m}, t \in[0, T] \\ x_{m} \in \mathscr{C}_{\text {tem }}^{+},\left|x_{m}\right|_{-\lambda^{\prime}} \leq c_{0}}} \sum_{k \in \mathrm{~K}} P_{r_{m}, x_{m}}^{n}\left\langle\left(m, n X_{t}^{k}\right)^{2 q}, \phi_{\lambda}\right\rangle<\infty \tag{41}
\end{equation*}
$$

Proof. Fix $c_{0}, T, q, \lambda^{\prime}, \lambda$ as in the lemma, where without loss of generality we may assume that $q>5$. We may also restrict our attention to $r_{m} \leq t$. In order to handle later the imposed time-partitioning in a Gronwall's inequality argument, we include now the approximating equation (33) in our consideration. Let $0 \leq \varepsilon \leq 1$. Using equations (33) and (34) as well as Burkholder-Davis-Gundy's inequality applied to the martingale,

$$
\begin{equation*}
t \mapsto \int_{\left[r_{m}, t\right] \times \mathrm{R}}{ }^{m, n} M_{r_{m}}^{k}(d(s, b)) p_{\varepsilon+t^{\prime}-s}(b-a), \quad r_{m} \leq t \leq t^{\prime} \tag{42}
\end{equation*}
$$

gives the inequality

$$
\begin{align*}
P_{r_{m}, x_{m}}^{n}\left(S_{\varepsilon}^{m, n} X_{t}^{k}(a)\right)^{2 q} \leq & c\left(S_{\varepsilon+t-r_{m}} x_{m}^{k}(a)\right)^{2 q}+c P_{r_{m}, x_{m}}^{n} \\
& \times\left(\int_{r_{m}}^{t} d s \int_{\mathrm{R}} d b p_{\varepsilon+t-s}^{2}(b-a)^{m, n} X_{s}^{k}(b)\right.  \tag{43}\\
& \left.\times\left(S_{1 / n}^{m, n} X_{[n s] / n}^{k+1}(b) \wedge n\right)\right)^{q}
\end{align*}
$$

Using the presupposed bound $c_{0} \phi_{-\lambda^{\prime}}$ for $x_{m}$, and the heat flow estimate (38), the first term at the right-hand side of (43) has the bound $c \phi_{-2 q \lambda^{\prime}}(a)$, which paired with $\phi_{\lambda}$ leads to a finite expression within (41), independent of $m, n$, $k, r_{m}, t, x_{m}$.

Hence it remains to deal with the remaining second term at the right-hand side of (43). First of all, in the integrand of the double integral we may additionally introduce $\phi_{-\lambda / q}(b) \phi_{\lambda / q}(b) \equiv \mathbf{1}$. Moreover, we decompose the square term by using $2=\left(2-\frac{2}{q}\right)+\frac{2}{q}$. Then by Hölder's inequality with $p$ such that $\frac{1}{p}+\frac{1}{q}=1$, the $q$ th power of the double integral in (43) can be estimated from
above by

$$
\begin{align*}
& \left(\int_{r_{m}}^{t} d s \int_{\mathrm{R}} d b p_{\varepsilon+t-s}^{2}(b-a) \phi_{-\lambda p / q}(b)\right)^{q / p}  \tag{44a}\\
& \quad \times \int_{r_{m}}^{t} d s \int_{\mathrm{R}} d b p_{\varepsilon+t-s}^{2}(b-a) \phi_{\lambda}(b)  \tag{44b}\\
& \quad \times\left({ }^{m, n} X_{s}^{k}(b)\left({S_{1 / n}}^{m, n} X_{[n s] / n}^{k+1}(b) \wedge n\right)\right)^{q}
\end{align*}
$$

where we used that $\left(2-\frac{2}{q}\right) p=2$. In the double integral in (44a), we estimate one of the $p$-factors by $c(t-s)^{-1 / 2}$, and apply (38) in order to get for the $\frac{q}{p}$ th power of that integral the bound

$$
\begin{equation*}
c \phi_{-\lambda}(a)\left(\int_{0}^{t} d s(t-s)^{-1 / 2}\right)^{q / p} \leq c \phi_{-\lambda}(a) \tag{45}
\end{equation*}
$$

with $c$ again independent of $m, n, k, r_{m}, t, x_{m}$, which cancels the $\phi_{\lambda}(\alpha)$ in (41). On the other hand, in the double integral in (44b) we split $p^{2}$ as before, but use this time that the remaining $p$ can be paired with 1 in (41). Altogether we found

$$
\begin{align*}
& \int_{\mathrm{R}} d a \phi_{\lambda}(a) P_{r_{m}, x_{m}}^{n}\left(S_{\varepsilon}^{m, n} X_{t}^{k}(a)\right)^{2 q} \\
& \quad \leq c+c \int_{r_{m}}^{t} d s(\varepsilon+t-s)^{-1 / 2} \int_{\mathrm{R}} d b \phi_{\lambda}(b)  \tag{46}\\
& \quad \times P_{r_{m}, x_{m}}^{n}\left({ }^{m, n} X_{s}^{k}(b){S_{1 / n}}^{m, n} X_{[n s] / n}^{k+1}(b)\right)^{q}
\end{align*}
$$

with the constants $c$ independent of $m, n, k, r_{m}, t, x_{m}$. The latter term can further be estimated by using the elementary inequality

$$
\begin{equation*}
(u v)^{q} \leq u^{2 q}+v^{2 q}, \quad u, v \geq 0, \quad q>0 \tag{47}
\end{equation*}
$$

Also, we may sum these inequalities over $k \in \mathrm{~K}$. Setting (for fixed $m, n, r_{m}, x_{m}$ )

$$
\begin{equation*}
f_{\varepsilon}(t):=\sum_{k \in \mathrm{~K}} \int_{\mathrm{R}} d a \phi_{\lambda}(a) P_{r_{m}, x_{m}}^{n}\left(S_{\varepsilon}^{m, n} X_{t}^{k}(a)\right)^{2 q} \tag{48}
\end{equation*}
$$

$0 \leq \varepsilon \leq 1, r_{m} \leq t \leq T$, we thus obtained the following two estimates:

$$
\begin{aligned}
f_{0}(t) & \leq \underline{c}+\underline{c} \int_{r_{m}}^{t} d s(t-s)^{-1 / 2}\left[f_{0}(s)+f_{1 / n}([n s] / n)\right] \\
f_{1 / n}([n t] / n) & \leq \underline{c}+\underline{c} \int_{r_{m}}^{[n t] / n} d s(1 / n+[n t] / n-s)^{-1 / 2}\left[f_{0}(s)+f_{1 / n}([n s] / n)\right]
\end{aligned}
$$

with the constant $\underline{c}$ independent of $m, n, r_{m}, t, x_{m}$. Using in the latter integral first $1 / n+[n t] / n \geq t$ and then $[n t] / n \leq t$, we see that $g(t):=f_{0}(t)+$ $f_{1 / n}([n t] / n)$ satisfies

$$
\begin{equation*}
g(t) \leq \underline{c}+\underline{c} \int_{r_{m}}^{t} d s(t-s)^{-1 / 2} g(s), \quad r_{m} \leq t \leq T \tag{49}
\end{equation*}
$$

with $\underline{c}$ independent of $m, n, r_{m}, t, x_{m}$. Then Gronwall's inequality implies $g(t) \leq c, 0 \leq r_{m} \leq t \leq T$, with $c$ independent of $m, n, r_{m}, t, x_{m}$ (see [20], page 138). Hence $f_{0}(t) \leq g(t)$ gives the claim (41), completing the proof.

Next we want to deal with moments of time increments of the integral part
${ }^{m, n} Y_{t}^{k}(a):=\int_{\left[r_{m}, t\right] \times \mathrm{R}}{ }^{m, n} M_{r_{m}}^{k}(d(s, b)) p_{t-s}(b-a), \quad t \geq r_{m}, k \in \mathrm{~K}, a \in \mathrm{R}$ in (34).

LEMMA 13 (Moments of increments). For constants $c_{0}, T, p, q>0$ with $\frac{1}{p}+$ $\frac{1}{q}=1$ and $q>5$, and $\lambda^{\prime}, \lambda \in \mathrm{R}$ with $2 q \lambda^{\prime}<\lambda$, we have

$$
\begin{gather*}
\sup _{\substack{m, n \geq 1, r_{m} \in[0, T] \\
x_{m} \in \epsilon_{\mathrm{tem}}^{+},\left|x_{m}\right|_{-\lambda^{\prime}} \leq c_{0}}} \sum_{k \in \mathrm{~K}} P_{r_{m}, x_{m}}^{n}\left|m, n Y_{t^{\prime}}^{k}\left(a^{\prime}\right)-m, n Y_{t}^{k}(a)\right|^{2 q}  \tag{50}\\
\leq c\left(\left|t^{\prime}-t\right|^{1 / 2}+\left|a^{\prime}-a\right|\right)^{q / p} \phi_{-\lambda}(a)
\end{gather*}
$$

whenever $t, t^{\prime} \in[0, T], a, a^{\prime} \in \mathrm{R}$, and $\left|a-a^{\prime}\right| \leq 1$.
Proof. We may assume that $r_{m} \leq t \leq t^{\prime}$. Let

$$
N_{r}:=\int_{\left[r_{m}, r\right] \times \mathrm{R}}{ }^{m, n} M_{r_{m}}^{k}(d(s, b))\left(p_{t^{\prime}-s}\left(b-a^{\prime}\right)-p_{t-s}(b-a)\right), \quad r_{m} \leq r \leq t^{\prime}
$$

Then $r \mapsto N_{r}$ is a martingale with square function

$$
\langle N\rangle_{r}=\int_{r_{m}}^{r} d s \int_{\mathrm{R}} d b\left[p_{t^{\prime}-s}\left(b-a^{\prime}\right)-p_{t-s}(b-a)\right]^{m, n} X_{s}^{k}(b){S_{1 / n}}^{m, n} X_{[n s] / n}^{k+1}(b)
$$

Note that

$$
\begin{equation*}
{ }^{m, n} Y_{t^{\prime}}^{k}\left(a^{\prime}\right)-{ }^{m, n} Y_{t}^{k}(a)=N_{t^{\prime}} \tag{51}
\end{equation*}
$$

where we used the convention

$$
\begin{equation*}
p_{s}:=0 \quad \text { if } s<0 \tag{52}
\end{equation*}
$$

Again by Burkholder-Davis-Gundy's inequality, we will deal with

$$
\begin{aligned}
& P_{r_{m}, x_{m}}^{n}\left(\int_{r_{m}}^{t^{\prime}} d s \int_{\mathrm{R}} d b\right. {\left[p_{t^{\prime}-s}\left(b-a^{\prime}\right)-p_{t-s}(b-a)\right]^{2} } \\
&\left.\times{ }^{m, n} X_{s}^{k}(b) S_{1 / n}^{m, n} X_{[n s] / n}^{k+1}(b)\right)^{q} .
\end{aligned}
$$

As we derived (44a) and (44b), we get the bound

$$
\begin{gathered}
c\left(\int_{0}^{t^{\prime}} d s \int_{\mathrm{R}} d b\left[p_{t^{\prime}-s}\left(b-a^{\prime}\right)-p_{t-s}(b-a)\right]^{2}\right)^{q / p} \\
\quad \times \int_{r_{m}}^{t^{\prime}} d s \int_{\mathrm{R}} d b\left[p_{t^{\prime}-s}^{2}\left(b-a^{\prime}\right)+p_{t-s}^{2}(b-a)\right] \\
\quad \times P_{r_{m}, x_{m}}^{n}\left({ }^{m, n} X_{s}^{k}(b) S_{1 / n}^{m, n} X_{[n s] / n}^{k+1}(b)\right)^{q}
\end{gathered}
$$

By [31], Lemma 6.2 (i),

$$
\begin{equation*}
\int_{0}^{t^{\prime}} d s \int_{\mathrm{R}} d b\left[p_{t^{\prime}-s}\left(b-a^{\prime}\right)-p_{t-s}(b-a)\right]^{2} \leq c\left(\left|t^{\prime}-t\right|^{1 / 2}+\left|a^{\prime}-a\right|\right)^{q / p} \tag{53}
\end{equation*}
$$

$t, t^{\prime} \geq 0, a, a^{\prime} \in \mathrm{R}$. Therefore, the first double integral leads to the desired right-hand side in (50), except for the $\phi_{-\lambda}(a)$. So it remains to show that the second double integral can uniformly be bounded by $c \phi_{-\lambda}(a)$. For this purpose we may assume that $t^{\prime}=t$ [recall the convention (52)]. In the integrand we additionally introduce $\phi_{-\lambda}(b) \phi_{\lambda}(b) \equiv 1$, and apply the Cauchy-Schwarz inequality to the $d b$-integral. This gives the bound

$$
\begin{align*}
& 2\left(\int_{\mathrm{R}} d b p_{t-s}^{4}(b-a) \phi_{-2 \lambda}(b)\right)^{1 / 2}  \tag{54a}\\
& \quad \times\left(\int_{\mathrm{R}} d b \phi_{2 \lambda}(b) P_{r_{m}, x_{m}}^{n}\left({ }^{m, n} X_{s}^{k}(b) S_{1 / n}^{m, n} X_{[n s] / n}^{k+1}(b)\right)^{2 q}\right)^{1 / 2} \tag{54b}
\end{align*}
$$

for the $d b$-integral. The factor in (54b) is uniformly bounded. In fact, use once more (47) and (38), to get expressions of the type as in Lemma 12 with $q, \lambda$ replaced by $2 q, 2 \lambda$. On the other hand, in the factor in (54a) we split $p^{4}=p p^{3}$, and use $p_{t-s}^{3}(b-a) \leq c(t-s)^{-3 / 2}$ which, after taking the $1 / 2$ power, has a bounded $d s$-integral. In fact, the $d b$-integral of the remaining quantities gives $\phi_{-\lambda}(a)$ by (38), uniformly in $(t-s)$ and $r_{m}, x_{m}$. Thus the proof of Lemma 13 is complete.

Completion of the proof of Proposition 9. Fix $\left(r_{m}, x_{m}\right)$ converging in $\mathrm{R}_{+} \times$ $\sigma_{\text {tem }}^{+}$to $(r, x)$ as $m \uparrow \infty$, as well as $p, q>0$ as in Lemma 13 , implying $q / p>4$. Since $T, \lambda^{\prime}, \lambda$ in Lemma 13 are arbitrary, from (50) and Lemma 6.3(ii) in [31], we see that the sequence of the laws of ${ }^{m, n} Y$ with respect to $P_{r_{m}, x_{m}}^{n}$ is tight in $\mathscr{P}$. [Note that at the right-hand side of the condition (6.5) in [31] the factor $e^{\lambda|x|}$ has to be added, and that the laws of the initial states ${ }^{n} X_{0} \in \mathscr{C}_{\text {tem }}, n \geq 1$, in Lemma 6.3(ii) should be tight by assumption.] So is that of the laws of ${ }^{m, n} X$ [recall the definition of ${ }^{m, n} Y$ before Lemma 13, and (34)]. Let $X$ denote any limit point (in law) of the ${ }^{m, n} X$ as $m \uparrow \infty$ and $n \uparrow \infty$. Since ${ }^{m, n} X$ satisfies the martingale problem $\mathbf{M P}_{r_{m}, x_{m}}^{n}$ of Definition 8, for each ( $m, n$ ), it follows from a standard limiting argument that $X$ satisfies the martingale problem $\mathbf{M P}_{r, x}$. This finishes the proof of Proposition 9.
3.2. Some properties of all martingale problem solutions $\left(X, P_{r, x}\right)$. After we have constructed a solution to our basic martingale problem, we now want to collect some properties of all the solutions (that is, not only of the constructed ones).

For this purpose, we redefine ( ${ }^{m, n} X, P_{r_{m}, x_{m}}^{n}$ ) introduced before Lemma 12 as any solutions to the martingale problem $\mathbf{M P}_{r_{m}, x_{m}}$ of Definition 8 (instead of $\left.\mathbf{M P}_{r_{m}, x_{m}}^{n}\right)$ for each $n \geq 1$. In particular, in the case $\left(r_{m}, x_{m}\right) \equiv\left(r_{1}, x_{1}\right)=:(r, x)$, we have a whole sequence $\left\{{ }^{1, n} X=:{ }^{n} X: n \geq 1\right\}$ of solutions to $\mathbf{M P}_{r, x}$.

With that system $\left\{\left({ }^{m, n} X, P_{r_{m}, x_{m}}^{n}\right): m, n \geq 1\right\}$ we now repeat all the constructions in Section 3.1. Then, in particular, analogs of Lemmas 11-13 are true, and once more by tightness, any limit point ( $X, P_{r, x}$ ) of that new system again satisfies the martingale problem $\mathbf{M} \mathbf{P}_{r, x}$ of Definition 8.

Lemma 14 (Properties of all martingale problem solutions). For $m, n \geq 1$, take $\left(r_{m}, x_{m}\right) \in \mathrm{R}_{+} \times \mathcal{C}_{\text {tem }}^{+}$and let $\left({ }^{m, n} X, P_{r_{m}, x_{m}}^{n}\right)$ denote any solution to the martingale problem $\mathbf{M P}_{r_{m}, x_{m}}$ of Definition 8. Then the following statements hold:
(a) (Pointwise equation). For $t \geq r_{m}$, and $(k, a) \in \mathrm{K} \times \mathrm{R}$ fixed, $P_{r_{m}, x_{m}}^{n}$-almost surely,

$$
{ }^{m, n} X_{t}^{k}(a)=S_{t-r_{m}} x_{m}^{k}(a)+\int_{\left[r_{m}, t\right] \times \mathrm{R}}{ }^{m, n} M_{r_{m}}^{k}(d(s, b)) p_{t-s}(b-a)
$$

(with ${ }^{m, n} M_{r_{m}}^{k}$ denoting the related martingale measure).
(b) (Uniformly bounded moments). For fixed $c_{0}, T, q>0$ and $\lambda^{\prime}, \lambda \in \mathrm{R}$ with $2 q \lambda^{\prime}<\lambda$,

$$
\sup _{\substack{m, n \geq 1, r_{m}, t \in\left[0, T \backslash \\ x_{m} \in \epsilon_{\text {tem }}\left|x_{m}\right| x_{n} \mid \lambda \leq c_{0}\right.}} \sum_{k \in \mathrm{~K}} P_{r_{m}, x_{m}}^{n}\left\langle\left(^{m, n} X_{t}^{k}\right)^{2 q}, \phi_{\lambda}\right\rangle<\infty .
$$

(c) (Moments of increments). For constants $c_{0}, T, p, q>0$ with $\frac{1}{p}+\frac{1}{q}=1$ and $q>5$, and $\lambda^{\prime}, \lambda \in \mathrm{R}$ with $2 q \lambda^{\prime}<\lambda$, we have (with the notation ${ }^{m, n} Y$ introduced before Lemma 22)

$$
\begin{gathered}
\left.\sup _{\substack{m, n \geq 1, r_{|l|} \\
x_{m} \in \epsilon_{\text {tem }}^{+},\left|x_{m}\right| T_{-} \leq T_{1} \leq 0}} \sum_{k \in \mathrm{~K}} P_{r_{m}, x_{m}}^{n}\right|^{m, n} Y_{t^{\prime}}^{k}\left(a^{\prime}\right)-\left.{ }^{m, n} Y_{t}^{k}(a)\right|^{2 q} \\
\leq c\left(\left|t^{\prime}-t\right|^{1 / 2}+\left|a^{\prime}-a\right|\right)^{q / p} \phi_{-\lambda}(a),
\end{gathered}
$$

whenever $t, t^{\prime} \in[0, T], a, a^{\prime} \in \mathrm{R}$, and $\left|a-a^{\prime}\right| \leq 1$.
(d) (Limit points). Assume that $\left(r_{m}, x_{m}\right)$ converges in $\mathrm{R}_{+} \times \mathscr{C}_{\mathrm{tem}}^{+}$to ( $r, x$ ) as $m \uparrow \infty$. Then any limit point $\left(X, P_{r, x}\right)$ of $\left\{\left({ }^{m, n} X, P_{r_{m}, x_{m}}^{n}\right): m, n \geq 1\right\}$ as $m \uparrow \infty$ and $n \uparrow \infty$ satisfies the martingale problem $\mathbf{M P}_{r, x}$ of Definition 8.

We also need the following property.
Corollary 15 (Uniformly bounded moments for each solution $X$ ). Fix $c_{0}$, $T, q>0$ and $\lambda^{\prime}, \lambda$ in R with $2 q \lambda^{\prime}<\lambda$. Let $\left(X, P_{r, x}\right)$ be any solution to the martingale problem $\mathbf{M P}_{r, x}$ of Definition 8. Then,

$$
\begin{equation*}
\sup _{\substack{r \in|, T|, x \in \operatorname{tem} \\|x|-\left.\right|_{\text {tom }} ^{+}}} P_{r, x} \sup _{\substack{0 \leq t \leq T}} \sum_{k \in K}\left\langle\left(X_{t}^{k}\right)^{2 q}, \phi_{\lambda}\right\rangle<\infty . \tag{55}
\end{equation*}
$$

Proof. We specialize in Lemma 14(c) to ${ }^{m, n} X \equiv{ }^{1,1} X=: X$. Using the Banach space $L^{2 q}\left(\mathrm{R}, \phi_{\lambda}(a) d a\right)$ with $q$ sufficiently large, from the proof of Theorem 1.2.1 in [28], and Lemma 14(c) we get (55) with $X$ replaced by
$Y \equiv{ }^{m, n} Y$. But by Lemma 14(a) and the heat flow estimate (38), claim (55) also holds for $X$. Finally, (55) is then true for all $q>0$, completing the proof.

The special case ${ }^{m, n} X \equiv X$ also gives the following result.
Corollary 16 (Pointwise equation for each solution). For each solution $\left(X, P_{r, x}\right)$ to the martingale problem $\mathbf{M P}_{r, x}$ of Definition 8, the family of martingales extends to orthogonal square-integrable martingale measures $M_{r}^{k}=$ $M_{r}^{k}(d(s, b))$ such that, for the usual predictable functions $f:[r, \infty] \times \mathrm{K} \times \mathrm{R} \times \Omega \rightarrow$ R in their domain,

$$
\begin{align*}
& \left\langle\left\langle\int_{[r, \cdot] \times \mathrm{R}} M_{r}^{k}(d(s, b)) f_{s}^{k}(b)\right\rangle_{t}\right.  \tag{56}\\
& \quad=\gamma^{k} \int_{r}^{t} d s \int_{\mathrm{R}} d a X_{s}^{k}(a) X_{s}^{k+1}(a)\left[f_{s}^{k}(b)\right]^{2}, \quad t \geq r, \quad k \in \mathrm{~K}
\end{align*}
$$

Moreover, for $t \geq r$ and $(k, a) \in \mathrm{K} \times \mathrm{R}$ fixed,

$$
\begin{equation*}
X_{t}^{k}(a)=S_{t-r} x^{k}(a)+\int_{[r, t] \times \mathrm{R}} M_{r}^{k}(d(s, b)) p_{t-s}(b-a), \quad P_{r, x}-a . s . \tag{57}
\end{equation*}
$$

In particular, the expectation formula

$$
\begin{equation*}
P_{r, x} X_{t}^{k}(a)=S_{t-r} x^{k}(a), \quad t \geq r, \quad k \in \mathrm{~K}, \quad a \in \mathrm{R} \tag{58}
\end{equation*}
$$

and the covariance formula

$$
\begin{align*}
& \operatorname{Cov}_{r, x}\left(X_{t_{1}}^{k_{1}}\left(a_{1}\right), X_{t_{2}}^{k_{2}}\left(a_{2}\right)\right) \\
& =\gamma_{k_{1}} \delta_{k_{1}, k_{2}} \int_{r}^{t_{1} \wedge t_{2}} d s \int_{\mathrm{R}} d b S_{s-r} x^{k_{1}}(b) S_{s-r}  \tag{59}\\
& \quad \times x^{k_{1}+1}(b) p_{t_{1}-s}\left(a_{1}-b\right) p_{t_{2}-s}\left(a_{2}-b\right),
\end{align*}
$$

$t_{1}, t_{2} \geq r, k_{1}, k_{2} \in \mathrm{~K}, a_{1}, a_{2} \in \mathrm{R}$, are valid.
Note that (55), (58) and (59) yield already the moment formulas in Theorem 3(b).
3.3. The mapping $(r, x) \mapsto \mathscr{P}_{r, x} . \quad$ For $(r, x) \in \mathrm{R}_{+} \times \mathscr{b}_{\text {tem }}^{+}$, let $\mathscr{P}_{r, x} \subseteq$ $\mathscr{P}$ denote the set of all solutions $P_{r, x}$ to the martingale problem $\mathbf{M P}_{r, x}$ of Definition 8. Note that $\mathscr{P}_{r, x} \neq \varnothing$ by Proposition 9. Recall the metric space $\left(\operatorname{com}(\mathscr{P}), d_{\text {com }}\right)$ introduced in Section 2.1.

LEMMA 17 (Set of all solutions). $\quad(r, x) \mapsto \mathscr{P}_{r, x} \neq \varnothing$ is a measurable mapping of $\mathrm{R}_{+} \times \mathscr{C}_{\text {tem }}^{+}$into $\operatorname{com}(\mathscr{P})$.

Proof. From the special case $\left(r_{m}, x_{m}\right) \equiv(r, x)$, that is ${ }^{m, n} X \equiv{ }^{1, n} X$, in Lemma 14(d) we see that the set $\mathscr{P}_{r, x}$ of all solutions to the martingale problem $\mathbf{M P}_{r, x}$ of Definition 8 is compact. On the other hand, the special case ${ }^{m, n} X \equiv{ }^{m, 1} X$ shows that the map $(r, x) \mapsto \mathscr{P}_{r, x}$ is measurable (see [32], Lemma 12.1.8). This completes the proof.

We continue with a time-homogeneity property of the family

$$
\begin{equation*}
\left\{\mathscr{P}_{r, x}:(r, x) \in \mathrm{R}_{+} \times \mathscr{C}_{\text {tem }}^{+}\right\} \tag{60}
\end{equation*}
$$

of all solutions to the martingale problems of Definition 8. For this purpose, for $r \geq 0$, we introduce the shift operator $\Phi_{r}$ on $\Omega=\mathscr{C}\left(\mathrm{R}_{+}, \mathscr{C}_{\text {tem }}^{+}\right)$by

$$
\begin{equation*}
\left(\Phi_{r} \omega\right)_{t}:=\omega_{(t-r) \vee 0}, \quad \omega \in \Omega, \quad t \geq 0 \tag{61}
\end{equation*}
$$

(producing a constant initial piece).
LEMMA 18 (Time homogeneity). The $\operatorname{map}(r, x) \mapsto \mathscr{P}_{r, x}$ istime-homogeneous, that is,

$$
\begin{equation*}
\mathscr{P}_{r, x}=\mathscr{P}_{0, x} \circ \Phi_{r}^{-1}, \quad(r, x) \in \mathrm{R}_{+} \times \mathscr{C}_{\text {tem }}^{+} \tag{62}
\end{equation*}
$$

(with the obvious notation).
Proof. The proof is quite elementary and shows that there is a one-to-one correspondence between the solution to the martingale problems $\mathbf{M P}_{r, x}$ and $\mathbf{M P}_{0, x}$ (compare with [32], Lemma 6.5.1). Fix $(r, x) \in \mathrm{R}_{+} \times \mathscr{C}_{\text {tem }}^{+}$and $P \in \mathscr{P}$.

Step $1^{\circ}$ (Constancy). By the notation (61),

$$
\begin{equation*}
\left\{\left(\Phi_{r} X\right)_{t}=x, \quad t \leq r\right\}=\left\{X_{0}=x\right\} . \tag{63}
\end{equation*}
$$

Hence, the process with law $P \circ \Phi_{r}^{-1}$ equals constantly $x$ up to time $r$ if and only if $P\left(X_{0}=x\right)=1$.

Step $2^{\circ}$ (Martingale property). For $\varphi \in \digamma_{\text {rap }}^{(2)}, t \geq r$, and $k \in \mathrm{~K}$, again by definition of the shift operator,

$$
\begin{align*}
& \left\langle\left(\Phi_{r} X\right)_{t}^{k}, \varphi^{k}\right\rangle-\left\langle x^{k}, \varphi^{k}\right\rangle-\int_{r}^{t} d s\left\langle\left(\Phi_{r} X\right)_{s}^{k}, \frac{\sigma^{2}}{2} \Delta \varphi^{k}\right\rangle \\
& \quad=\left\langle X_{t-r}^{k}, \varphi^{k}\right\rangle-\left\langle x^{k}, \varphi^{k}\right\rangle-\int_{0}^{t-r} d s\left\langle X_{s}^{k}, \frac{\sigma^{2}}{2} \Delta \varphi^{k}\right\rangle  \tag{64}\\
& \quad=M_{0, t-r}^{k}\left(\varphi^{k}\right)
\end{align*}
$$

Thus, by the martingale problem $\mathbf{M P}_{r, x}$ of Definition 8 , the map $t \mapsto M_{r, t}^{k}\left(\varphi^{k}\right)$ with respect to the law $P \circ \Phi_{r}^{-1}$ is a martingale after time $r$ starting from 0 if and only if $t \mapsto M_{0, t-r}^{k}\left(\varphi^{k}\right)$ with respect to $P$ is a martingale after 0 starting from 0 .

Step $3^{\circ}$ (Square function). Similarly, by Definition 8(iii),
(65) $t \mapsto\left[M_{r, t}^{k}\left(\varphi^{k}\right)\right]^{2}-\gamma^{k} \int_{r}^{t} d s \int_{\mathrm{R}} d a X_{s}^{k}(a) X_{s}^{k+1}(a)\left[\varphi^{k}(a)\right]^{2}=: N_{r, t}^{k}\left(\varphi^{k}\right)$
with respect to $P \circ \Phi_{r}^{-1}$ is a martingale after time $r$ if and only if the same is true for $r=0$.

Step $4^{\circ}$ (Conclusion). Putting Steps $1^{\circ}-3^{\circ}$ together, the claim in the lemma follows.

We finish this section with an optional stopping argument which we need later for an integrability statement in the proof of selection of a strong Markov solution. By (7), for fixed $\lambda \in \mathrm{R}$, there is a constant $c_{(66)}$ such that for the smoothed reference function $\tilde{\phi}_{\lambda}$,

$$
\begin{equation*}
\left|\frac{\sigma^{2}}{2} \Delta \tilde{\phi}_{\lambda}\right| \leq c_{(66)} \tilde{\phi}_{\lambda} . \tag{66}
\end{equation*}
$$

Recall also the notation $\|\cdot\|_{\lambda}$ introduced in (12).
Lemma 19 (A conditional moment estimate). Let the law $P_{r, x}$ belong to $\mathscr{P}_{r, x}$, for $(r, x) \in \mathrm{R}_{+} \times \mathscr{C}_{\text {tem }}^{+}$. Consider $T \geq r$ and $[r, T]$-valued stopping times $\eta \leq \vartheta$. Then

$$
\begin{equation*}
P_{r, x}\left\{\left\|X_{\vartheta}^{k}\right\|_{-\lambda}^{n} \mid \mathscr{F}_{\eta}\right\} \leq e^{n c_{(66)}(T-\eta)}\left\|X_{\eta}^{k}\right\|_{-\lambda}^{n}<\infty, \quad P_{r, x}-a . s . \tag{67}
\end{equation*}
$$

for all $k \in \mathrm{~K}, n \geq 1, \lambda>0$.
Proof. Fix $P_{r, x}, k, n$, and $\lambda$. From the martingale problem $\mathbf{M P}_{r, x}$ and Itô's formula for $t \geq r$,

$$
\begin{align*}
d\left(e^{-n c_{(66)} t}\left\|X_{t}^{k}\right\|_{-\lambda}^{n}\right)= & n e^{-n c_{(66)} t}\left\|X_{t}^{k}\right\|_{-\lambda}^{n-1},\left\langle X_{t}^{k},\left(\frac{\sigma^{2}}{2} \Delta-c_{(66)}\right) \tilde{\phi}_{\lambda}\right) d t  \tag{68}\\
& +n e^{-n c_{(66)} t}\left\|X_{t}^{k}\right\|_{-\lambda}^{n-1} d M_{r, t}^{k}\left(\tilde{\phi}_{\lambda}\right) .
\end{align*}
$$

Hence, by the definition (66) of $c_{(66)}$, the process $t \mapsto e^{-n c_{(66)} t}\left\|X_{t}^{k}\right\|_{-\lambda}^{n}$ is a $P_{r, x}$-supermartingale after time $r$. Then the claim (67) immediately follows from [19], Theorem 1.1.39.
3.4. Selection of a strong Markov solution. We now want to select a timehomogeneous strong Markov version $(r, x) \mapsto P_{r, x}$ from $(r, x) \mapsto \mathscr{P}_{r, x}$. The idea behind this is an optimization procedure as in [32], which goes back to [21] and which uses an extremal property which is well behaved under conditioning and weak convergence.

To make this precise, we first recall some notation taken from [32]. We stress the fact that all those results we quote from [32] are valid also in our present case of the set $\Omega=\mathscr{6}\left(\mathrm{R}_{+}, \mathscr{b}_{\text {tem }}^{+}\right)$of paths in an infinite-dimensional space. Recall that $\mathscr{F}^{r}$ denotes the $\sigma$-field generated by the coordinate process at times $t \geq r$.

Notation 20 (Composition I). For fixed $\omega \in \Omega, r \geq 0$, and a law $P$ on ( $\Omega, \mathscr{F}^{r}$ ) with the property that $P\left(X_{r}=\omega_{r}\right)=1$, let $\delta_{\omega} \otimes_{r} P$ denote the unique law on $\Omega$ satisfying

$$
\begin{equation*}
\delta_{\omega} \otimes_{r} P\left(X_{t}=\omega_{t} \text { for } t \leq r\right)=1 \quad \text { and } \quad \delta_{\omega} \otimes_{r} P=P \text { on } \mathscr{\mathscr { T }}^{r} \tag{69}
\end{equation*}
$$

(see [32], Lemma 6.1.1). Roughly speaking, the irrelevant history of the process $(X, P)$ up to time $r$ is replaced by the one of $\omega$ yielding the process $\left(X, \delta_{\omega} \otimes_{r}\right.$ $P)$. For fixed $x \in \mathscr{C}_{\text {tem }}^{+}$, let the notation $\delta_{x} \otimes_{r} P$, however, refer to the special case $\omega_{t}=x$ for $t \leq r$ (constant initial piece).

We also need the following notation.
Notation 21 (Composition II). Consider a probability measure $P$ on $\Omega$, a stopping time $\tau$ on $\Omega$, and a mapping $\omega \mapsto Q_{\omega}$ of $\Omega$ into $\mathscr{P}$ satisfying the following two conditions:
(a) $\omega \mapsto Q_{\omega}$ is $\mathscr{F}_{\tau}$-measurable.
(b) $Q_{\omega}\left(X_{\tau(\omega)}=\omega_{\tau(\omega)}\right)=1$ for all $\omega \in \Omega$.

Then let $P \otimes_{\tau} Q$ denote the unique probability measure on $\Omega$ which has the following two properties:
(c) $P \otimes_{\tau} Q$ equals $P$ on $\mathscr{F}_{\tau}$.
(d) $\omega \mapsto \delta_{\omega} \otimes_{\tau(\omega)} Q_{\omega}$ (recall Notation 20) is a regular conditional probability distribution of $P \otimes_{\tau} Q$ given $\mathscr{F}_{\tau}$ (see [32], Theorem 6.1.2).

Roughly speaking, the process $\left(X, P \otimes_{\tau} Q\right)$ has the law $P$ until the random time $\tau$, and its conditional law after time $\tau$ is given by the family $Q$.

DEFINITION 22 (Strong Markov solution). $\quad\left\{P_{r, x}:(r, x) \in \mathrm{R}_{+} \times \mathscr{C}_{\text {tem }}^{+}\right\} \subseteq \mathscr{P}$ is said to be a time-homogeneous strong Markov solution to the family

$$
\begin{equation*}
\mathbf{M P}:=\left\{\mathbf{M P}_{r, x}:(r, x) \in \mathbf{R}_{+} \times \mathscr{C}_{\text {tem }}^{+}\right\} \tag{70}
\end{equation*}
$$

of martingale problems of Definition 8 , if $(r, x) \mapsto P_{r, x}$ is a measurable map of $\mathrm{R}_{+} \times \mathscr{C}_{\text {tem }}^{+}$into $\mathscr{P}$, and if for each $(r, x) \in \mathrm{R}_{+} \times \mathscr{C}_{\text {tem }}^{+}$:
(a) $P_{r, x} \in \mathscr{P}_{r, x}$.
(b) $P_{r, x}=P_{0, x} \circ \Phi_{r}^{-1}$ (time homogeneity).
(c) For each stopping time $\tau \geq r$ on $\Omega$ and each regular conditional probability distribution $\omega \mapsto P_{\omega}$ of $P_{r, x}$ given $\mathscr{F}_{\tau}$, there is a $P_{r, x}$-null set $\mathrm{N} \in \mathscr{F}_{\tau}$ such that

$$
P_{\omega}=\delta_{\omega} \otimes_{\tau(\omega)} P_{\tau(\omega), \omega_{\tau(\omega)}}, \quad \omega \notin \mathrm{N}
$$

(recall Notation 20).
In other words, solutions $P_{r, x}$ of $\mathbf{M P}_{r, x}$ are selected in such a way that

$$
\begin{equation*}
\left(X, P_{r, x},(r, x) \in \mathrm{R}_{+} \times \mathscr{C}_{\mathrm{tem}}^{+}\right) \tag{71}
\end{equation*}
$$

is a time-homogeneous strong Markov process.
The existence statement on a cyclically catalytic SBM as claimed in Theorem 3(a) can now be restated as follows.

THEOREM 23 (Existence of a strong Markov solution). There exists a timehomogeneous strong Markov solution to the martingale problem MP according to Definition 22.

The verification of this theorem (in the end of this subsection) needs some further preparation. Recall the second part of Notation 20 and Notation 21.

Lemma 24 (Compositions). Fix $x \in \mathscr{b}_{\text {tem }}^{+}, P \in \mathscr{P}_{0, x}$, and a finite stopping time $\tau$ on $\Omega=\mathscr{C}\left(\mathrm{R}_{+}, \mathscr{b}_{\text {tem }}^{+}\right)$.
(a) (Composition I). If $\omega \mapsto P_{\omega}$ is a regular conditional probability distribution of $P$ given $\mathscr{F}_{\tau}$, then there is a $P$-null set $\mathrm{N} \in \mathscr{F}_{\tau}$ such that

$$
\delta_{\omega_{\tau(\omega)}} \otimes_{\tau(\omega)} P_{\omega} \in \mathscr{P}_{\tau(\omega), \omega_{\tau(\omega)}}, \quad \omega \notin \mathrm{N}
$$

(b) (Composition II). If $\omega \mapsto Q_{\omega}$ is an $\mathscr{F}_{\tau}$-measurable map of $\Omega$ into $\mathscr{P}$ such that

$$
\delta_{\omega_{\tau(\omega)}} \otimes_{\tau(\omega)} Q_{\omega} \in \mathscr{P}_{\tau(\omega), \omega_{\tau(\omega)}}, \quad \omega \in \Omega
$$

then $P \otimes_{\tau} Q$ belongs to $\mathscr{P}_{0, x}$.
Proof. Fix $x, P, \tau$ as in the lemma.
(a) Let $\omega \mapsto P_{\omega}$ be a regular conditional probability distribution of $P$ given $\mathscr{F}_{\tau}$. Denote by $\mathscr{A}$ a countable dense subset of $\mathscr{C}_{\text {rap }}^{(2)}$. Fix $\varphi \in \mathscr{A}$ and $k \in \mathrm{~K}$ for a while. First note that by Definition 8(ii),

$$
\begin{equation*}
M_{\tau(\omega), t}^{k}\left(\varphi^{k}\right)=M_{0, t}^{k}\left(\varphi^{k}\right)-M_{0, \tau(\omega)}^{k}\left(\varphi^{k}\right), \quad t \geq \tau(\omega) \tag{72}
\end{equation*}
$$

Hence, by the last part of Theorem 1.2.10 in [32] (applied to $\theta(t)=M_{0, t}^{k}\left(\varphi^{k}\right)$ and $s=0$ ), there exists a $P$-null set $\mathrm{N}_{\varphi} \in \mathscr{F}_{\tau}$ such that for all $\omega \notin \mathrm{N}_{\varphi}$,

$$
\begin{equation*}
t \mapsto M_{\tau(\omega), t}^{k}\left(\varphi^{k}\right) \quad \text { is a } P_{\omega} \text {-martingale after time } \tau(\omega) \tag{73}
\end{equation*}
$$

Recalling the notation $N_{r, t}^{k}\left(\varphi^{k}\right)$ introduced in (65), by Definition 8(iii) the following identity holds:

$$
\begin{equation*}
N_{\tau(\omega), t}^{k}\left(\varphi^{k}\right)=N_{0, t}^{k}\left(\varphi^{k}\right)-N_{0, \tau(\omega)}^{k}\left(\varphi^{k}\right)-2 M_{\tau(\omega), t}^{k}\left(\varphi^{k}\right) M_{0, \tau(\omega)}^{k}\left(\varphi^{k}\right) \tag{74}
\end{equation*}
$$

Appealing again to the same theorem in [32] and combining with (73), we may redefine the $P$-null set $\mathrm{N}_{\varphi} \in \mathscr{F}_{\tau}$ such that for $\omega \notin \mathrm{N}_{\varphi}$ additionally,

$$
\begin{equation*}
t \mapsto N_{\tau(\omega), t}^{k}\left(\varphi^{k}\right) \quad \text { is a } P_{\omega} \text {-martingale after time } \tau(\omega) \tag{75}
\end{equation*}
$$

Introduce the $P$-null set $\mathrm{N}:=\bigcap_{\varphi \in \mathscr{A}} \mathrm{N}_{\varphi} \in \mathscr{F}_{\tau}$. We may additionally assume that N is independent of $k \in \mathrm{~K}$. To $\varphi \in \mathscr{C}_{\text {rap }}^{(2)}$ we now choose $\varphi_{n} \in \mathscr{A}$ converging in $\mathscr{C}_{\text {rap }}^{(2)}$ to $\varphi$. Then, from (73) and (75) we conclude that for $\omega \notin \mathrm{N}$,
(76) $\quad M_{\tau(\omega),}^{k} \cdot\left(\varphi^{k}\right)$ and $N_{\tau(\omega),}^{k} \cdot\left(\varphi^{k}\right)$ are $P_{\omega}$-martingales after time $\tau(\omega)$.

Since $k$ and $\varphi$ are arbitrary, and N does not depend on them, claim (a) is true.
(b) Let $\omega \mapsto Q_{\omega}$ be an $\mathscr{F}_{\tau}$-measurable map as presupposed in (b). First of all, $P \otimes_{\tau} Q$ makes sense, according to Notation 21. Trivially, $P \otimes_{\tau} Q$ has the right initial state,

$$
\begin{equation*}
P \otimes_{\tau} Q\left(X_{0}=x\right)=P\left(X_{0}=x\right)=x \tag{77}
\end{equation*}
$$

Fix $\varphi \in \mathscr{C}_{\text {rap }}^{(2)}$ and $k \in \mathrm{~K}$. Next we want to show that

$$
\begin{equation*}
\left\{M_{0, t}^{k}\left(\varphi^{k}\right): t \geq 0\right\} \quad \text { is a } P \otimes_{\tau} Q \text {-martingale. } \tag{78}
\end{equation*}
$$

By the last part of Theorem 6.1.2 in [32] [again with $\theta(t)=M_{0, t}^{k}\left(\varphi^{k}\right)$ and $s=0]$ it suffices to show that

$$
\begin{equation*}
M_{0, t}^{k}\left(\varphi^{k}\right) \text { is } P \otimes_{\tau} Q \text {-integrable, } \quad t \geq 0 \tag{79}
\end{equation*}
$$

that

$$
\begin{equation*}
\left\{M_{0, t \wedge \tau}^{k}\left(\varphi^{k}\right): t \geq 0\right\} \quad \text { is a } P \text {-martingale } \tag{80}
\end{equation*}
$$

and that

$$
\begin{equation*}
\left\{M_{0, t}^{k}\left(\varphi^{k}\right)-M_{0, t \wedge \tau(\omega)}^{k}\left(\varphi^{k}\right): t \geq 0\right\} \quad \text { is a } Q_{\omega} \text {-martingale, } \quad \omega \in \Omega . \tag{81}
\end{equation*}
$$

In order to check the integrability statement (79), we fix $T \geq t \vee 1$ and $\lambda>0$, as well as a constant $c_{(82)}>0$ such that

$$
\begin{equation*}
\left|\varphi^{k}\right|+\frac{\sigma^{2}}{2}\left|\Delta \varphi^{k}\right| \leq c_{(82)} \tilde{\phi}_{\lambda} \tag{82}
\end{equation*}
$$

[recall (7)]. Then by the martingale definition (14),

$$
\begin{equation*}
P \otimes_{\tau} Q\left|M_{0, t}^{k}\left(\varphi^{k}\right)\right| \leq 2 c_{(82)} T \sup _{s \leq T} P \otimes_{\tau} Q\left\|X_{s}^{k}\right\|_{-\lambda} \tag{83}
\end{equation*}
$$

[recall notation (12)]. Conditioning on $\mathscr{F}_{\tau}$ in the latter expectation expression, by Notation 21 we get

$$
\begin{equation*}
P \otimes_{\tau} Q\left\|X_{s}^{k}\right\|_{-\lambda}=\int_{\Omega} P(d \omega) \delta_{\omega} \otimes_{\tau(\omega)} Q_{\omega}\left\|X_{s}^{k}\right\|_{-\lambda} . \tag{84}
\end{equation*}
$$

First we restrict in the latter integral additionally to $\tau(\omega)>s$. Then concerning the internal expectation, $\left\|X_{s}^{k}\right\|_{-\lambda}$ equals the deterministic value $\left\|\omega_{s}^{k}\right\|_{-\lambda}$, just by notation (69). Hence, for the considered first part of (84) we found the bound

$$
\begin{equation*}
\int_{\Omega} P(d \omega)\left\|\omega_{s}^{k}\right\|_{-\lambda}=P\left\|X_{s}^{k}\right\|_{-\lambda} \leq e^{c_{(66)} T}\left\|x^{k}\right\|_{-\lambda} \tag{85}
\end{equation*}
$$

where we used Lemma 19.
Under the restriction $\tau(\omega) \leq s$, however, again by notation (69),

$$
\begin{equation*}
\delta_{\omega} \otimes_{\tau(\omega)} Q_{\omega}=Q_{\omega}=\delta_{\omega_{\tau(\omega)}} \otimes_{\tau(\omega)} Q_{\omega} \text { on } \mathscr{F}^{\tau(\omega)} . \tag{86}
\end{equation*}
$$

Since by assumption, $\delta_{\omega_{\tau(\omega)}} \otimes_{\tau(\omega)} Q_{\omega}$ satisfies the martingale problem $\mathbf{M P}_{\tau(\omega), \omega_{\tau(\omega)}}$ (except for $\omega$ in the null set $\mathrm{N}_{\tau}$ ), we may apply Lemma 19 to get

$$
\begin{equation*}
\delta_{\omega_{\tau(\omega)}} \otimes_{\tau(\omega)} Q_{\omega}\left\|X_{s}^{k}\right\|_{-\lambda} \leq e^{c_{(66} T}\left\|\omega_{\tau(\omega)}^{k}\right\|_{-\lambda} . \tag{87}
\end{equation*}
$$

Thus, for the part of (84) under consideration, we got the bound

$$
\begin{equation*}
e^{c_{(66)} T} \int_{\Omega} P(d \omega)\left\|\omega_{\tau(\omega)}^{k}\right\|_{-\lambda}=e^{c_{(66)} T} P\left\|X_{\tau}^{k}\right\|_{-\lambda} \leq e^{2 c_{(66)} T}\left\|x^{k}\right\|_{-\lambda}, \tag{88}
\end{equation*}
$$

where in the last step we exploited once more our conditional moment estimate in Lemma 19.

Altogether we obtained

$$
\begin{equation*}
\sup _{s \leq T} P \otimes_{\tau} Q\left\|X_{s}^{k}\right\|_{-\lambda} \leq 2 e^{2 c_{(66)} T}\left\|x^{k}\right\|_{-\lambda}<\infty . \tag{89}
\end{equation*}
$$

Thus, by (83), the integrability claim (79) is verified.
Statement (80) is immediately clear, and we turn to (81). Fix $\omega \in \Omega$. Since $\delta_{\omega_{\tau(\omega)}} \otimes_{\tau(\omega)} Q_{\omega} \in \mathscr{P}_{\tau(\omega), \omega_{\tau(\omega)}}$ by assumption,
(90) $\quad t \mapsto M_{\tau(\omega), t}^{k}\left(\varphi^{k}\right) \quad$ is a $\delta_{\omega_{\tau(\omega)}} \otimes_{\tau(\omega)} Q_{\omega}$-martingale after time $\tau(\omega)$,
by definition. But $M_{\tau(\omega), t}^{k}\left(\varphi^{k}\right)$ is $\mathscr{F}^{\tau(\omega)}$-measurable; thus in (90) we may replace $\delta_{\omega_{\tau(\omega)}} \otimes_{\tau(\omega)} Q_{\omega}$ by $Q_{\omega}$. Hence, by the martingale identity (72),
(91) $t \mapsto M_{0, t}^{k}\left(\varphi^{k}\right)-M_{0, \tau(\omega)}^{k}\left(\varphi^{k}\right) \quad$ is a $Q_{\omega}$-martingale after time $\tau(\omega)$,
and (81) follows.
Exploiting again [32], Theorem 6.1.2, analogously to (78) it can be shown that

$$
\begin{equation*}
\left\{N_{0, t}^{k}\left(\varphi^{k}\right): t \geq 0\right\} \quad \text { is a } P \otimes_{\tau} Q \text {-martingale. } \tag{92}
\end{equation*}
$$

Together with (78), the claim follows, completing the proof of Lemma 24.
Now it will be convenient for us to consider the map $(r, x) \mapsto \mathscr{P}_{r, x}$ introduced in the beginning of Section 3.3 also from a more general point of view.

Definition 25 (Nice family). A family $\left\{\mathscr{P}_{r, x} \neq \varnothing:(r, x) \in \mathrm{R}_{+} \times \sigma_{\text {tem }}^{+}\right\}$of subsets of the set $\mathscr{P}$ of all probability laws on $\Omega=\mathscr{C}\left(\mathrm{R}_{+}, \mathscr{b}_{\text {tem }}^{+}\right)$is said to be nice, if it is measurable of $\mathrm{R}_{+} \times \mathscr{b}_{\text {tem }}^{+}$into $\operatorname{com}(\mathscr{P})$ as in Lemma 17, timehomogeneous as in Lemma 18, and if it has the composition properties as in Lemma 24.

Now we are ready to verify the existence Theorem 23 .
Proof of Theorem 23. By the Lemmas 17, 18 and 24, we already know that our family $(r, x) \mapsto \mathscr{P}_{r, x}$ of all solutions to the martingale problem MP in (70) is nice according to the previous definition. By a successive optimization procedure, we would like to shrink down the sets $\mathscr{P}_{r, x}$ to single point sets $\left\{P_{r, x}\right\}$.

Let $\mathscr{A}_{1}, \mathscr{A}_{2}$, and $\mathscr{A}_{3}$ denote countable dense subsets of $(0, \infty), \mathscr{C}_{0}$ and $\mathscr{C}_{\text {rap }}$, respectively, where $\mathscr{C}_{0}:=\mathscr{C}_{0}(\mathrm{R})$ is the separable Banach space of all functions $f: R \rightarrow R$ vanishing at infinity, equipped with the supremum norm of uniform convergence. Let $\left\{\left(\theta_{n}, f_{n}, \varphi_{n}\right): n \geq 1\right\}$ denote an enumeration of $\mathscr{A}_{1} \times \mathscr{A}_{2} \times \mathscr{A}_{3}$.

Fix $(r, x) \in \mathrm{R}_{+} \times \mathscr{C}_{\text {tem }}^{+}$for the moment. For $P \in \mathscr{P}_{r, x}$, set

$$
\begin{equation*}
L_{r, x}^{n}(P):=\int_{0}^{\infty} d t e^{-\theta_{n} t} P f_{n}\left(\left\langle X_{r+t}, \varphi_{n}\right\rangle\right), \quad n \geq 1 . \tag{93}
\end{equation*}
$$

Define inductively

$$
\begin{equation*}
\mathscr{P}_{r, x}^{n+1}:=\left\{P \in \mathscr{P}_{r, x}^{n}: L_{r, x}^{n}(P)=\sup _{P^{\prime} \in \mathscr{P}_{r, x}^{n}} L_{r, x}^{n}\left(P^{\prime}\right)\right\}, \quad n \geq 1 \tag{94}
\end{equation*}
$$

where $\mathscr{P}_{r, x}^{1}:=\mathscr{P}_{r, x}$. Then, by [32], Lemma 12.2.2,

$$
\begin{equation*}
\mathscr{P}^{n}:=\left\{\mathscr{P}_{r, x}^{n}:(r, x) \in \mathrm{R}_{+} \times \mathscr{C}_{\mathrm{tem}}^{+}\right\} \tag{95}
\end{equation*}
$$

is again a nice family, for each $n \geq 1$. Moreover, as in the proof of Theorem 12.2.3 in [32], also the monotone limits

$$
\begin{equation*}
\mathscr{P}_{r, x}^{\infty}:=\bigcap_{n \geq 1} \mathscr{P}_{r, x}^{n}, \quad(r, x) \in \mathrm{R}_{+} \times \mathscr{C}_{\text {tem }}^{+} \tag{96}
\end{equation*}
$$

form a nice family.
Fix again $(r, x) \in \mathrm{R}_{+} \times \mathscr{C}_{\text {tem }}^{+}$, and consider $P, P^{\prime} \in \mathscr{P}_{r, x}^{\infty}$. In order to finish the proof of the theorem, it remains to show that $P=P^{\prime}$. By construction,

$$
\begin{equation*}
\int_{0}^{\infty} d t e^{-\theta t} P f\left(\left\langle X_{r+t}, \varphi\right\rangle\right)=\int_{0}^{\infty} d t e^{-\theta t} P^{\prime} f\left(\left\langle X_{r+t}, \varphi\right\rangle\right) \tag{97}
\end{equation*}
$$

for all $(\theta, f, \varphi) \in \mathscr{A}_{1} \times \mathscr{A}_{2} \times \mathscr{A}_{3}$. Since $\mathscr{A}_{1}$ is dense in $(0, \infty)$ by assumption, by the uniqueness theorem of Laplace transforms and by the integrands' continuity in $t$, we get

$$
\begin{equation*}
P f\left(\left\langle X_{r+t}, \varphi\right\rangle\right)=P^{\prime} f\left(\left\langle X_{r+t}, \varphi\right\rangle\right), \quad t \geq 0,(f, \varphi) \in \mathscr{A}_{2} \times \mathscr{A}_{3} \tag{98}
\end{equation*}
$$

But also $\mathscr{A}_{2}$ and $\mathscr{A}_{3}$ are dense in $\mathscr{C}_{0}$ and $\mathscr{C}_{\text {rap }}$, respectively, and we conclude that the martingale problem solutions $(X, P)$ and $\left(X, P^{\prime}\right)$ have the same one-dimensional distributions. Thus the laws $P$ and $P^{\prime}$ coincide (cf. [15], Theorem 4.4.2), completing the proof of Theorem 23.
4. Global segregation of neighboring types (Proof of Theorem 4). Fix $X_{0}=x \in \mathscr{C}_{\text {tem }}^{+}$and suppose it has a finite total mass $\|x\|$. In accordance with (18), set

$$
\begin{equation*}
Z_{T}=Z_{T}(x):=\sum_{k \in \mathrm{~K}} \gamma^{k}\left\|X_{T}^{k}\right\| \cdot\left\|X_{T}^{k+1}\right\| \tag{99}
\end{equation*}
$$

The strategy of the following proof of Theorem 4 is to construct a contradiction by assuming that

$$
\begin{equation*}
P_{x}\left(\inf _{T \geq 0} Z_{T}>0\right)>0 \tag{100}
\end{equation*}
$$

Then $Z_{T} \rightarrow 0 P_{x}$-a.s. will follow [and therefore the claim (17)], since 0 is an absorbing state for the process $Z$, and since for the continuous nonnegative martingales $T \mapsto\left\|X_{T}^{k}\right\|$ we have that

$$
\begin{equation*}
\lim _{T \uparrow \infty}\left\|X_{T}^{k}\right\|=:\left\|X_{\infty}^{k}\right\| \text { exists in } \mathrm{R}_{+}, \quad k \in \mathrm{~K} \tag{101}
\end{equation*}
$$

The contradiction will arise when under the event in (100) we construct finite stopping times $T_{n} \uparrow \infty$ such that $Z_{T_{n}} \rightarrow 0$. This construction requires some preparation.

Step $1^{\circ}$. First, for $T \geq 0$, we introduce the "global clock"

$$
\begin{equation*}
A_{T}:=\int_{0}^{T} d t \sum_{k \in \mathrm{~K}} \gamma^{k} \int_{\mathrm{R}} d a X_{t}^{k}(a) X_{t}^{k+1}(a) \underset{T \uparrow \infty}{\nearrow} \text { some } A_{\infty}<\infty \tag{102}
\end{equation*}
$$

Indeed, $A_{\infty}$ is finite $P_{x}$-a.s., since by (15) and orthogonality, $A$ is the square function of the nonnegative (hence convergent) martingale $T \mapsto\left\|X_{T}\right\|$. We want to decompose $A_{T}$ by using the pointwise equation (57) (with $r=0$ ). For this purpose, put

$$
\begin{align*}
{ }^{0} A_{T} & :=\int_{0}^{T} d t \sum_{k \in \mathrm{~K}} \gamma^{k} \int_{\mathrm{R}} d a S_{t} x^{k}(a) S_{t} x^{k+1}(a),  \tag{103a}\\
{ }^{1} N_{T}(T) & :=\int_{0}^{T} d t \sum_{k \in \mathrm{~K}} \gamma^{k} \int_{\mathrm{R}} d a S_{t} x^{k}(a) N_{t}^{k+1}(t, a),  \tag{103b}\\
{ }^{2} N_{T}(T) & :=\int_{0}^{T} d t \sum_{k \in \mathrm{~K}} \gamma^{k} \int_{\mathrm{R}} d a N_{t}^{k}(t, a) S_{t} x^{k+1}(a),  \tag{103c}\\
{ }^{3} N_{T}(T) & :=\int_{0}^{T} d t \sum_{k \in \mathrm{~K}} \gamma^{k} \int_{\mathrm{R}} d a N_{t}^{k}(t, a) N_{t}^{k+1}(t, a), \tag{103d}
\end{align*}
$$

where, for $a \in \mathrm{R}$,

$$
\begin{equation*}
N_{r}^{k}(t, a):=\int_{[0, r] \times \mathrm{R}} M^{k}(d(s, b)) p_{t-s}(b-a), \quad 0<r \leq t \tag{104}
\end{equation*}
$$

with the martingale measures $M^{k}:=M_{0}^{k}$ from Corollary 16. Note that all quantities make sense by the uniform moment estimates Corollary 15. Then, by our equation (57) (recall that there we have continuity in $t, a$ ), we get the decomposition

$$
\begin{equation*}
A_{T}={ }^{0} A_{T}+{ }^{1} N_{T}(T)+{ }^{2} N_{T}(T)+{ }^{3} N_{T}(T)=:{ }^{0} A_{T}+N_{T}(T) \tag{105}
\end{equation*}
$$

Step $2^{\circ}$. In analyzing the fluctuating part $N_{T}(T)$ of $A_{T}$, a little care has to be taken since $T \mapsto N_{T}(T)$ [or the single terms $T \mapsto{ }^{i} N_{T}(T)$ ] are not martingales.

Interchanging the order of integration in (103b) gives

$$
\begin{equation*}
{ }^{1} N_{T}(T)=\sum_{k} \gamma^{k} \int_{[0, T] \times \mathrm{R}} M^{k+1}(d(s, b)) \int_{s}^{T} d t S_{2 t-s} x^{k}(b) \tag{106}
\end{equation*}
$$

We generalize now the notation $T \mapsto{ }^{1} N_{T}(T)$ by putting

$$
\begin{equation*}
{ }^{1} N_{r}(T):=\sum_{k} \gamma^{k} \int_{[0, r] \times \mathrm{R}} M^{k+1}(d(s, b)) \int_{s}^{T} d t S_{2 t-s} x^{k}(b), \quad 0 \leq r \leq T \tag{107}
\end{equation*}
$$

As opposed to $T \mapsto{ }^{1} N_{T}(T)$, for fixed $T>0$, the process $r \mapsto{ }^{1} N_{r}(T), r \in$ [ $0, T$ ], is a martingale. Analogously, we can define the martingale $r \mapsto{ }^{2} N_{r}(T)$.

Integrating by parts in the third fluctuation term (103d) gives

$$
\begin{aligned}
{ }^{3} N_{T}(T)=\int_{0}^{T} d t \sum_{k} \gamma^{k} \int_{\mathrm{R}} d a[ & \int_{[0, t] \times \mathrm{R}} M^{k}(d(s, b)) p_{t-s}(b-a) N_{s}^{k+1}(t, a) \\
& \left.+\int_{[0, t] \times \mathrm{R}} M^{k+1}(d(s, b)) p_{t-s}(b-a) N_{s}^{k}(t, a)\right]
\end{aligned}
$$

which we write as ${ }^{31} N_{T}(T)+{ }^{32} N_{T}(T)$ in the obvious correspondence. Interchanging the order of integration yields

$$
\begin{align*}
{ }^{31} N_{T}(T)= & \sum_{k} \gamma^{k} \int_{[0, T] \times \mathrm{R}} M^{k}(d(s, b))  \tag{108}\\
& \times \int_{s}^{T} d t \int_{\mathrm{R}} d a p_{t-s}(b-a) N_{s}^{k+1}(t, a) .
\end{align*}
$$

Put

$$
\begin{align*}
{ }^{31} N_{r}(T):= & \sum_{k} \gamma^{k} \int_{[0, r] \times \mathrm{R}} M^{k}(d(s, b))  \tag{109}\\
& \times \int_{s}^{T} d t \int_{\mathrm{R}} d a p_{t-s}(b-a) N_{s}^{k+1}(t, a),
\end{align*}
$$

$0 \leq r \leq T$, getting again a martingale $t \mapsto{ }^{31} N_{r}(T)$. Similarly, we define the martingale $r \mapsto{ }^{32} N_{r}(T)$.

Altogether, in generalization of the notations (103b)-(103d) and (105), for $T>0$ fixed, we defined the martingales ${ }^{i} N .(T), i \in\{1,2,31,32\}$, and

$$
\begin{equation*}
r \mapsto N_{r}(T):={ }^{1} N_{r}(T)+{ }^{2} N_{r}(T)+{ }^{31} N_{r}(T)+{ }^{32} N_{r}(T), \quad 0 \leq r \leq T \tag{110}
\end{equation*}
$$

Step $3^{\circ}$. Let us next mention the idea behind the following construction of a contradiction. It is relatively easy to see that for the deterministic part ${ }^{0} \mathrm{~A}$ of $A$ in the decomposition (105) we have

$$
\begin{equation*}
{ }^{0} A_{T} \geq c Z_{0} \quad \text { for large } T \tag{111}
\end{equation*}
$$

[see (127) below]. On the other hand, the martingale

$$
\begin{equation*}
r \mapsto N_{r}(T) \text { from (110) has a square function bounded by } A_{T} \tag{112}
\end{equation*}
$$

[see (123) below]. Since $B_{t} \approx \sqrt{t}$ for Brownian motion in R, the martingale representation theorem "yields"

$$
\begin{equation*}
\left|N_{T}(T)\right| \leq \sup _{0 \leq r \leq T}\left|N_{r}(T)\right| \leq c \sqrt{A_{T}} \tag{113}
\end{equation*}
$$

hence

$$
\begin{equation*}
N_{T}(T) \geq-c \sqrt{A_{T}} \quad \text { for large } T \tag{114}
\end{equation*}
$$

Combining with the decomposition (105) and the estimate (111) gives

$$
\begin{equation*}
A_{T}+c \sqrt{A_{T}} \geq c Z_{0} \quad \text { for a large } T_{1} \tag{115}
\end{equation*}
$$

Hence, there is a continuous function $h$ with $h(0)=0$ such that $h\left(A_{T_{1}}\right) \geq Z_{0}$. By our assumption (100), $Z_{T_{1}}$ is different from 0 with positive probability. Starting at time $T_{1}$ anew, we will find $T_{2}>T_{1}$ such that $h\left(A_{T_{2}}-A_{T_{1}}\right) \geq Z_{T_{1}}$, as so on. But $A_{T_{n+1}}-A_{T_{n}} \rightarrow 0$ as $n \uparrow \infty$ by (102) [provided that $T_{n} \uparrow \infty$ ]; therefore $Z_{T_{n}} \rightarrow 0$, which contradicts (100), as desired.

Step $4^{\circ}$. In order to make precise the previous ideas, we will control the random expressions ${ }^{i} N_{T}(T)$ in terms of $A_{T}$, as needed for (112). By orthogonality, for the square function of the martingale ${ }^{1} N(T)$ as defined in (107) we get

$$
\left\langle\left.\right|^{1} N(T)\right\rangle_{r}=\sum_{k} \gamma^{k} \int_{0}^{r} d s \int_{\mathrm{R}} d b \gamma^{k+1} X_{s}^{k+1}(b) X_{s}^{k+2}(b)\left[\int_{s}^{T} d t S_{2 t-s} x^{k}(b)\right]^{2} .
$$

Setting

$$
\begin{equation*}
q_{t}(a):=\int_{0}^{t} d s p_{s}(a), \quad t \geq 0, a \in \mathrm{R} \tag{116}
\end{equation*}
$$

for each constant $0 \leq \theta \leq 1$ we obtain

$$
\begin{equation*}
\int_{s}^{T} d t S_{2 t-s-\theta s} X_{\theta s}^{k}(b) \leq \frac{1}{2} q_{2 T}(0) \max _{k \in K, 0 \leq s \leq T}\left\|X_{s}^{k}\right\| \tag{117}
\end{equation*}
$$

Applying this for $\theta=0$, and denoting by $\bar{\gamma}$ the maximum of the $\gamma^{k}$, we find

$$
\begin{equation*}
\left\langle\left\langle{ }^{1} N(T)\right\rangle\right\rangle_{r} \leq \frac{1}{4} \bar{\gamma} q_{2 T}^{2}(0) \max _{k \in \mathrm{~K}, 0 \leq s \leq T}\left\|X_{s}^{k}\right\|^{2} A_{T}, \quad 0 \leq r \leq T . \tag{118}
\end{equation*}
$$

The same estimate is true for $\left.\left\langle 2^{2} N(T)\right\rangle\right\rangle_{r}$.
The square function value $\left\langle\left\langle{ }^{31} N(T)\right\rangle\right\rangle_{r}$ of the martingale ${ }^{31} N(T)$ of (109) equals

$$
\begin{align*}
& \sum_{k} \gamma^{k} \int_{0}^{r} d s \int_{\mathrm{R}} d b \gamma^{k} X_{s}^{k}(b) X_{s}^{k+1}(b) \\
& \times {\left[\int_{s}^{T} d t \int_{\mathrm{R}} d a p_{t-s}(b-a) N_{s}^{k+1}(t, a)\right]^{2} . } \tag{119}
\end{align*}
$$

But

$$
\begin{equation*}
N_{s}^{k+1}(t, a)=S_{t-s} N_{s}^{k+1}(s, \cdot)(a), \quad 0 \leq s \leq t, \tag{120}
\end{equation*}
$$

hence, by (57),

$$
\begin{equation*}
\left|N_{s}^{k+1}(t, a)\right| \leq S_{t-s} X_{s}^{k+1}(a)+S_{t} X_{0}^{k+1}(a) . \tag{121}
\end{equation*}
$$

Thus, the expression under the square brackets in (119) can in absolute value be bounded from above by

$$
\begin{equation*}
\int_{s}^{T} d t\left[S_{2(t-s)} X_{s}^{k+1}(b)+S_{2 t-s} X_{0}^{k+1}(b)\right] \leq q_{2 T}(0) \max _{k \in \mathrm{~K}, 0 \leq s \leq T}\left\|X_{s}^{k}\right\|, \tag{122}
\end{equation*}
$$

where in the last step we used (117) (for $\theta=1$ and $\theta=0$ ). Hence, for $\left\langle\left\langle{ }^{31} N(T)\right\rangle\right\rangle_{r}$ we get the same bound as in (118), except for the factor $\frac{1}{4}$. Moreover, for $\left\langle\left\langle{ }^{32} N(T)\right\rangle\right\rangle_{r}$ we get the same bound as for $\left\langle\left\langle{ }^{31} N(T)\right\rangle\right\rangle_{r}$.

Altogether, by the Kunita-Watanabe inequality, for the martingale $N(T)$ as defined in (110) we get, as announced in (112), the square function estimate

$$
\begin{equation*}
\langle\langle N(T)\rangle\rangle_{r} \leq c_{(123)} \bar{\gamma} q_{2 T}^{2}(0) A_{T} \max _{k \in \mathrm{~K}, 0 \leq s \leq T}\left\|X_{s}^{k}\right\|^{2}, \quad 0 \leq r \leq T \tag{123}
\end{equation*}
$$

where $c_{(123)}$ is a (universal) constant.
Step $5^{\circ}$. Now we want to derive a lower estimate for the deterministic term ${ }^{0} A_{T}$ from (103a), as announced in (111). For this purpose, for the fixed initial state $x$, choose a constant $L=L(x) \geq 1$ such that

$$
\begin{equation*}
\left\langle x^{k}, \mathbf{1}_{[-L / 2, L / 2]}\right\rangle \geq \frac{1}{2}\left\|x^{k}\right\|, \quad k \in \mathrm{~K} \tag{124}
\end{equation*}
$$

Also, there is a (universal) constant $c_{(125)}$ such that

$$
\begin{equation*}
q_{2 T}(L) \geq c_{(125)} q_{2 T}(0), \quad T \geq L^{2}=: T_{1}(x) \geq 1 \tag{125}
\end{equation*}
$$

Then in the identity

$$
\begin{equation*}
{ }^{0} A_{T}=\frac{1}{2} \sum_{k} \gamma^{k} \int_{\mathrm{R}} d a \int_{\mathrm{R}} d b x^{k}(a) x^{k+1}(b) q_{2 T}(b-a) \tag{126}
\end{equation*}
$$

we first restrict the integration domains in order to use (125), getting

$$
{ }^{0} A_{T} \geq \frac{c_{(125)}}{2} \sum_{k} \gamma^{k} \int_{-L / 2}^{L / 2} d a \int_{-L / 2}^{L / 2} d b x^{k}(a) x^{k+1}(b) q_{2 T}(0), \quad T \geq T_{1}(x)
$$

Then (124) yields the estimate

$$
\begin{equation*}
{ }^{0} A_{T} \geq \frac{c_{(125)}}{8} q_{2 T}(0) \sum_{k} \gamma^{k}\left\|x^{k}\right\| \cdot\left\|x^{k+1}\right\|=\frac{c_{(125)}}{8} q_{2 T}(0) Z_{0} \tag{127}
\end{equation*}
$$

for $T \geq T_{1}(x)$, with $Z_{0}=Z_{0}(x)$ from (99).
Now we modify our definition of $T_{1}(x)$ from (125): if $Z_{0}=0$, we set $T_{1}(x):=$ $\infty$. Otherwise we may enlarge $T_{1}(x)$ from the former definition (125) to a finite value by requiring that additionally,

$$
\begin{equation*}
\frac{c_{(125)}}{8} q_{2 T}(0) Z_{0} \geq 2, \quad T \geq T_{1}(x) \tag{128}
\end{equation*}
$$

Assume $T \in\left[T_{1}(x), \infty\right)$ for a while. From (127) and (128) we already know that

$$
\begin{equation*}
{ }^{0} A_{T} \geq \frac{c_{(125)}}{8} q_{2 T}(0) Z_{0} \geq 2 \tag{129}
\end{equation*}
$$

Step $6^{\circ}$. Next we want to bound below the probability $P_{x}\left(A_{T} \geq 1\right)$. Recall that $T_{1}(x) \leq T<\infty$. If we assume for the moment that $A_{T}<1$, then (129) and (105) imply that

$$
\begin{equation*}
N_{T}(T) \leq-\frac{c_{(125)}}{16} q_{2 T}(0) Z_{0} . \tag{130}
\end{equation*}
$$

Consequently,

$$
\begin{equation*}
P_{x}\left(A_{T}<1\right)=P_{x}\left(A_{T}<1, N_{T}(T) \leq-\frac{c_{(125)}}{16} q_{2 T}(0) Z_{0}\right) . \tag{131}
\end{equation*}
$$

Let

$$
\begin{equation*}
R=R(x) \geq \max _{k \in \mathrm{~K}}\left\|x^{k}\right\| . \tag{132}
\end{equation*}
$$

Distinguishing between

$$
\begin{equation*}
\max \left\{\left\|X_{t}^{k}\right\|: k \in \mathrm{~K}, 0 \leq t \leq T\right\}>2 R \tag{133}
\end{equation*}
$$

and the opposite, identity (131) can be continued with

$$
\begin{align*}
\leq & \sum_{k} P_{x}\left(A_{T} \leq 1, \max _{0 \leq t \leq T}\left\|X_{t}^{k}\right\| \geq 2 R\right) \\
& +P_{x}\left(A_{T} \leq 1, N_{T}(T) \leq-\frac{c_{(125)}}{16} q_{2 T}(0) Z_{0}, \quad \max _{\substack{k \in \kappa T \\
0 \leq \leq \leq T}}\left\|X_{t}^{k}\right\| \leq 2 R\right) . \tag{134}
\end{align*}
$$

For the first term in (134) we use that by (57) the process $t \mapsto\left\|X_{t}^{k}\right\|-\left\|x^{k}\right\|$ equals in law to a one-dimensional standard Brownian motion $\left(B, \Pi_{0}\right)$ (starting from 0 ) running with a clock bounded by $t \mapsto A_{t}$. Hence, for the first term in (134) we get the bound

$$
\begin{equation*}
\sum_{k} \Pi_{0}\left(\left\|x^{k}\right\|+\max _{0 \leq t \leq 1} B_{t} \geq 2 R\right) \leq K \Pi_{0}\left(\max _{0 \leq t \leq 1} B_{t} \geq R\right) \tag{135}
\end{equation*}
$$

where in the last step we used the definition (132) of $R$. By the reflection principle of Brownian motion, and an elementary estimate for the normal law,

$$
\begin{equation*}
\Pi_{0}\left(\max _{0 \leq t \leq 1} B_{t} \geq R\right) \leq 2 \Pi_{0}\left(B_{1} \geq R\right) \leq \frac{2}{R} e^{-R^{2} / 2} \tag{136}
\end{equation*}
$$

Consequently, for the first term in (134) we got the bound ( $2 K / R$ ) $e^{-R^{2} / 2}$.
For the second term in (134) we use the square function estimate (123) to obtain the bound

$$
\begin{equation*}
P_{x}\left(N_{T}(T) \leq-\frac{c_{(125)}}{16} q_{2 T}(0) Z_{0},\langle\langle N(T)\rangle\rangle_{T} \leq 4 c_{(123)} \bar{\gamma} q_{2 T}^{2}(0) R^{2}\right) . \tag{137}
\end{equation*}
$$

But the law of $r \mapsto N_{r}(T)$ coincides with the distribution of $B$ running with the clock $r \mapsto\left\langle\langle N(T)\rangle_{r}\right.$ (for a finite time). Hence, (137) is bounded from above
by

$$
\begin{align*}
& \Pi_{0}\left(\min \left\{B_{t}: 0 \leq t \leq 4 c_{(123)} \bar{\gamma} q_{2 T}^{2}(0) R^{2}\right\} \leq-\frac{c_{(125)}}{16} q_{2 T}(0) Z_{0}\right) \\
& \quad=\Pi_{0}\left(\min \left\{B_{t}: 0 \leq t \leq 1\right\} \leq-\frac{c_{(125)} Z_{0}}{32 \sqrt{c_{(123)}} \bar{\gamma} R}\right) \tag{138}
\end{align*}
$$

where in the last step we used Brownian scaling. Changing from $B$ to $-B$, again by the first part of (136) we may continue with

$$
\begin{equation*}
\leq 1-\Pi_{0}\left(\left|B_{1}\right| \leq \frac{c_{(125)} Z_{0}}{32 \sqrt{c_{(123)} \bar{\gamma}} R}\right) \leq 1-\frac{c_{(139)} Z_{0}}{R} \tag{139}
\end{equation*}
$$

where the constant $c_{(139)}$ does not depend on $x$ and $T$.
Altogether

$$
\begin{equation*}
P_{x}\left(A_{T} \geq 1\right) \geq-\frac{2 K}{R(x)} e^{-R^{2}(x) / 2}+\frac{c_{(139)} Z_{0}(x)}{R(x)}=: f(x) \tag{140}
\end{equation*}
$$

provided that $T \in\left[T_{1}(x), \infty\right)$.
Step $7^{\circ}$. Now we will make more precise our choice of $R(x)$ in (132). In fact, for the $x$ considered in this proof, set

$$
R(x):= \begin{cases}\sqrt{2 \log 2 K} \vee \max _{k \in \mathrm{~K}}\left\|x^{k}\right\|, & \text { if } c_{(139)} Z_{0}(x) \geq 2  \tag{141}\\ \sqrt{2 \log \frac{4 K}{c_{(139)} Z_{0}(x)}} \vee \max _{k \in \mathrm{~K}}\left\|x^{k}\right\|, & \text { otherwise }\end{cases}
$$

Note that then

$$
\begin{equation*}
f(x) \geq \frac{c_{(139)} Z_{0}(x)}{2 R(x)} \geq 0 \tag{142}
\end{equation*}
$$

Moreover, setting

$$
\begin{equation*}
V_{\delta, C}:=\left\{x: Z_{0}(x) \geq \delta, \max _{k \in \mathrm{~K}}\left\|x^{k}\right\| \leq C\right\}, \quad 0<\delta<C<\infty \tag{143}
\end{equation*}
$$

our choice (141) of $R$ yields
(144) $R\left(V_{\delta, C}\right)$ is a relatively compact subset of $(0, \infty), \quad 0<\delta<C<\infty$.

Step $8^{\circ}$. Setting $T_{0}:=0$, and recalling the definition of $T_{1}$ around (128), define inductively the stopping times

$$
T_{n+1}:= \begin{cases}T_{n}+T_{1}\left(X_{T_{n}}\right), & \text { if } T_{n}<\infty  \tag{145}\\ \infty, & \text { otherwise }\end{cases}
$$

$n \geq 1$. Note that $T_{n} \geq n$ for all $n$. Recalling that almost surely,

$$
\begin{equation*}
0=A_{0} \leq A_{t} \uparrow A_{\infty}<\infty \quad \text { as } t \uparrow \infty \tag{146}
\end{equation*}
$$

by the strong Markov property we have

$$
\begin{align*}
P_{x}\left\{A_{T_{n+1}}-A_{T_{n}} \geq 1 \mid \mathscr{F}_{T_{n}}\right\} & =\mathbf{1}_{\left\{T_{n}<\infty\right\}} P_{X_{T_{n}}}\left(A_{T_{1}} \geq 1\right)  \tag{147}\\
& \geq \mathbf{1}_{\left\{T_{n}<\infty\right\}} f\left(X_{T_{n}}\right) .
\end{align*}
$$

Hence, by the conditional version of Borel-Cantelli (see [36], 12.15),

$$
\begin{equation*}
\left\{A_{T_{n+1}}-A_{T_{n}} \geq 1 \text { infinitely often }\right\} \supseteq\left\{\sum_{n=1}^{\infty} \mathbf{1}_{\left\{T_{n}<\infty\right\}} f\left(X_{T_{n}}\right)=\infty\right\} \tag{148}
\end{equation*}
$$

$P_{x}$-a.s. But by (146), the left-hand side of (148) must be a null set. Hence,

$$
\begin{equation*}
\sum_{n=1}^{\infty} \mathbf{1}_{\left\{T_{n}<\infty\right\}} f\left(X_{T_{n}}\right)<\infty, \quad P_{x} \text {-a.s. } \tag{149}
\end{equation*}
$$

Since $f \geq 0$, on the set $\left\{T_{n}<\infty: n \geq 1\right\}$ we have ( $P_{x}$-a.s.)

$$
\begin{equation*}
\lim _{n \uparrow \infty} f\left(X_{T_{n}}\right)=0 \quad \text { hence } \lim _{n \uparrow \infty} \frac{Z_{T_{n}}}{R\left(X_{T_{n}}\right)}=0 \tag{150}
\end{equation*}
$$

the latter by (142).
Step $9^{\circ}$. Suppose now that (100) is valid, and we want to derive a contradiction. By (100), there exist constants $\delta>0$ and $\varepsilon \in\left(0, \frac{1}{2}\right)$ such that for our fixed $x$,

$$
\begin{equation*}
P_{x}\left(\inf _{T \geq 0} Z_{T} \geq \delta\right) \geq 2 \varepsilon \tag{151}
\end{equation*}
$$

On the other hand, from the martingale convergence (101) we conclude for the existence of a constant $C>\delta$ such that

$$
\begin{equation*}
P_{x}\left(\sup _{k \in K, T \geq 0}\left\|X_{T}^{k}\right\| \leq C\right) \geq 1-\varepsilon \tag{152}
\end{equation*}
$$

Introduce the event

$$
\begin{equation*}
\Omega_{\delta, C}:=\left\{\omega: \inf _{T \geq 0} Z_{T} \geq \delta, \sup _{k \in \mathrm{~K}, T \geq 0}\left\|X_{T}^{k}\right\| \leq C\right\} \tag{153}
\end{equation*}
$$

Then from (151) and (152),

$$
\begin{equation*}
P_{x}\left(\Omega_{\delta, C}\right) \geq \varepsilon \tag{154}
\end{equation*}
$$

Note that for $\omega \in \Omega_{\delta, C}$ we have $T_{n}<\infty$ for all $n$, hence, by (150),

$$
\begin{equation*}
P_{x}\left(\omega \in \Omega_{\delta, C}, \lim _{n \uparrow \infty} \frac{Z_{T_{n}}}{R\left(X_{T_{n}}\right)}=0\right) \geq \varepsilon>0 \tag{155}
\end{equation*}
$$

But $X_{T_{n}} \in V_{\delta, C}$ for each $n$ on the event $\Omega_{\delta, C}$, implying $Z_{T_{n}} \rightarrow 0$ as $n \uparrow \infty$ by the relative compactness in (144), which contradicts $\inf _{T \geq 0} Z_{T} \geq \delta>0$ in the definition of $\Omega_{\delta, C}$. Therefore the statement (100) cannot be true, and the claim in Theorem 4 follows as already explained in the beginning of this subsection. This completes the proof of Theorem 4.
5. Finite time behavior (Proof of Theorem 5). Finally, the finite time behavior Theorem 5 will be proved in the following two subsections.
5.1. Finite time survival of all types [proof of (a)]. As a preparation for the proof, for convenience we give the following variance estimate.

Lemma 26 (Variance estimate). For $x \in \mathscr{b}_{\text {tem }}^{+}$and $(t, k, a) \in \mathrm{R}_{+} \times \mathrm{K} \times \mathrm{R}$,

$$
\begin{equation*}
\operatorname{Var}_{x} X_{t}^{k}(a) \leq \gamma^{k} \sqrt{\frac{2 t}{\pi}} \sqrt{S_{t}\left[x^{k}\right]^{2}(a) S_{t}\left[x^{k+1}\right]^{2}(a)} \tag{156}
\end{equation*}
$$

Proof. By the covariance formula in Theorem 3(b), $X_{t}^{k}(a)$ has the following variance:

$$
\begin{equation*}
\operatorname{Var}_{x} X_{t}^{k}(a)=\gamma^{k} \int_{0}^{t} d s \int_{\mathrm{R}} d b S_{s} x^{k}(b) S_{s} x^{k+1}(b) p_{t-s}^{2}(a-b) \tag{157}
\end{equation*}
$$

Estimate one of the $p$-factors by $p_{t-s}(0)=1 / \sqrt{2 \pi(t-s)}$ and use CauchySchwarz to get the upper bound

$$
\begin{equation*}
\gamma^{k} \int_{0}^{t} d s \frac{1}{\sqrt{2 \pi(t-s)}} \prod_{i=0}^{1}\left(\int_{\mathrm{R}} d b p_{t-s}(a-b)\left[S_{s} x^{k+i}(b)\right]^{2}\right)^{1 / 2} \tag{158}
\end{equation*}
$$

By Jensen's inequality, $\left[S_{s} x^{k}(b)\right]^{2} \leq S_{s}\left[x^{k}\right]^{2}(b)$, and altogether we get the desired variance estimate (156).

COMPLETION OF THE PROOF OF Theorem 5(a). Fix $X_{0}=x$, and $t \geq T>0$ as in the theorem, and let $P_{x}$ be any solution to the martingale problem $\mathbf{M P}_{x}$ (that is, the Markov property is not needed for this proof). Let $k \in \mathrm{~K}$ and $\varepsilon>0$. Then by our assumption (20), there is an $a \in \mathrm{R}$ such that

$$
\begin{equation*}
\frac{4 \gamma^{k} \sqrt{2 t}}{\sqrt{\pi}} \frac{\sqrt{S_{t}\left[x^{k}\right]^{2}(a) S_{t}\left[x^{k+1}\right]^{2}(a)}}{\left[S_{t} x^{k}(a)\right]^{2}}<\varepsilon \tag{159}
\end{equation*}
$$

Fix this $a$. By the continuity of states, $\left\|X_{t}^{k}\right\|=0$ implies that $X_{t}^{k}(a)=0$. Hence,

$$
\begin{equation*}
P_{x}\left(\left\|X_{t}^{k}\right\|=0\right) \leq P_{x}\left(\left|S_{t} x^{k}(a)-X_{t}^{k}(a)\right| \geq \frac{1}{2} S_{t} x^{k}(a)\right) \tag{160}
\end{equation*}
$$

By the expectation formula in Theorem 3(b), Chebyshev's inequality, and the variance estimate in Lemma 26,

$$
\begin{equation*}
P_{x}\left(\left\|X_{t}^{k}\right\|=0\right) \leq \frac{4 \gamma^{k} \sqrt{2 t}}{\sqrt{\pi}} \frac{\sqrt{S_{t}\left[x^{k}\right]^{2}(a) S_{t}\left[x^{k+1}\right]^{2}(a)}}{\left[S_{t} x^{k}(a)\right]^{2}}<\varepsilon \tag{161}
\end{equation*}
$$

the latter by our choice (159) of $a$. Since $\varepsilon$ is arbitrary, we arrive at

$$
\begin{equation*}
P_{x}\left(\left\|X_{t}^{k}\right\|=0\right)=0, \quad t \geq T \tag{162}
\end{equation*}
$$

Denote by $\tau \leq \infty$ the hitting time of 0 of the nonnegative continuous martingale $t \mapsto\left\|X_{t}^{k}\right\|$. Assume for the moment that $P_{x}(\tau<\infty)>0$ is true. Then also $P_{x}(\tau<t)>0$ holds for some $t \geq T$. But the state 0 is a trap of that martingale, and we get $P_{x}\left(\left\|X_{t}^{k}\right\|=0\right)>0$ for that $t$, which contradicts (162). Hence, $P_{x}(\tau<\infty)=0$, and since $k$ is arbitrary, the claim (21) follows, finishing the proof of Theorem 5(a).
5.2. Finite time extinction of a type [proof of (b)]. Recall that $k_{0} \in \mathrm{~K}$ is fixed, where we may assume without loss of generality that $k_{0}=0$. Recall also that we have positive constants $c_{i}$ and $\lambda_{i}, 0 \leq i \leq 2$, as well as $c_{1}^{\prime}, \lambda_{1}^{\prime}$ which are related by assumption (22). Additionally, fix positive constants $\alpha, \beta^{\prime}, \beta, \gamma$ such that

$$
\begin{equation*}
\lambda_{0}>\alpha>\lambda_{1}^{\prime}, \quad \lambda_{1}>\beta, \quad \lambda_{2}>\gamma \quad \text { and } \quad 2 \lambda_{1}^{\prime}<2 \beta^{\prime}<\beta+\gamma \tag{163}
\end{equation*}
$$

Once and for all, fix $\varepsilon \in(0,1]$, and a constant $p>9$. Without loss of generality, we may consider a terminal time $T \in(0,1]$.

Let $I_{c_{0}}$ denote the set of all initial states $x \in \mathscr{l}_{\text {tem }}^{+}$with $\|x\|<\infty$ and which satisfy (23) and (24). We may assume that $c_{0} \leq 1$, implying $I_{c_{0}} \subseteq I_{1}$.

Fix for the moment $c_{0}$ and $x \in I_{c_{0}}$. Let $P_{x}$ be any solution to the martingale problem MP ${ }_{x}$ (that is, the Markov property will again not be used).

Recall the smooth reference functions $\tilde{\phi}_{\lambda}$ introduced in (6) and the related notation $\|\cdot\|_{\lambda}$ from (12). Introduce the stopping times

$$
\begin{equation*}
\tau_{L}^{k, p, \lambda}:=\inf \left\{t \geq 0:\left\|\left(X_{t}^{k}\right)^{p}\right\|_{\lambda} \geq L\right\}, \quad L>0, \quad \lambda>0 \tag{164}
\end{equation*}
$$

and

$$
\begin{equation*}
\tau_{L}(t):=t \wedge \tau_{L}^{1,4 p, 4 p \beta} \wedge \tau_{L}^{2,4 p, 4 p \gamma}, \quad t \leq T \tag{165}
\end{equation*}
$$

Recall the martingales $r \mapsto N_{r}^{k}(t, a), r \leq t$, from (104), for each $(k, a) \in \mathrm{K} \times \mathrm{R}$. Set

$$
\begin{equation*}
N_{t}^{k}(a):=N_{t}^{k}(t, a) \tag{166}
\end{equation*}
$$

Since the claim of Theorem 5(b) highly depends on the interplay of "sizes" of the different types, some efforts are needed to control them. Here is our first result in this direction.

LEMMA 27 (A moment estimate for some fluctuation increments). For $0 \leq$ $t^{\prime} \leq t \leq T, a, a^{\prime} \in \mathrm{R}$, and $L>0$, there is a constant $c_{(167)}$ such that

$$
\begin{align*}
& \sup _{x \in I_{1}} P_{x}\left|N_{t}^{1}(a)-N_{t^{\prime}}^{1}\left(a^{\prime}\right)\right|^{2 p} \mathbf{1}\left\{t \leq \tau_{L}^{1,4 p, 4 p \beta} \wedge \tau_{L}^{2,4 p, 4 p \gamma}\right\}  \tag{167}\\
& \quad \leq c_{(167)} \sqrt{L}\left(\left|a-a^{\prime}\right|+\left|t-t^{\prime}\right|^{1 / 2}\right)^{p-1} T^{1 / 4}\left(\phi_{p(\beta+\gamma)}\left(a^{\prime}\right)+\phi_{p(\beta+\gamma)}(a)\right)
\end{align*}
$$

Proof. By definition (104) of $N_{r}^{k}(t, a)$, the moment expression in (167) can be written as

$$
P_{x}\left|\int_{[0, t] \times \mathrm{R}} M^{1}(d(s, b))\left[p_{t-s}(b-a)-p_{t^{\prime}-s}\left(b-a^{\prime}\right)\right] \mathbf{1}\left\{t \leq \tau_{L}^{1,4 p, 4 p \beta} \wedge \tau_{L}^{2,4 p, 4 p \gamma}\right\}\right|^{2 p}
$$

Under the restriction as in the indicator, for the upper integration bound we may use that $t=\tau_{L}(t)$, by definition (165). Hence, the latter moment expression can be estimated from above by

$$
\begin{equation*}
P_{x}\left|\int_{\left[0, \tau_{L}(t)\right] \times \mathrm{R}} M^{1}(d(s, b))\left[p_{t-s}(b-a)-p_{t^{\prime}-s}\left(b-a^{\prime}\right)\right]\right|^{2 p} \tag{168}
\end{equation*}
$$

By virtue of the Burkholder-Davis-Gundy inequality, this can be further bounded by

$$
\begin{equation*}
c P_{x}\left|\int_{0}^{\tau_{L}(t)} d s \int_{\mathrm{R}} d b\left[p_{t-s}(b-a)-p_{t^{\prime}-s}\left(b-a^{\prime}\right)\right]^{2} X_{s}^{1}(b) X_{s}^{2}(b)\right|^{p} \tag{169}
\end{equation*}
$$

Writing $2=\left(2-\frac{2}{p}\right)+\frac{2}{p}$, by Hölder's inequality the latter double integral can be estimated from above by
(170a) $\left|\int_{0}^{\tau_{L}(t)} d s \int_{\mathrm{R}} d b\left[p_{t-s}(b-a)-p_{t^{\prime}-s}\left(b-a^{\prime}\right)\right]^{2}\right|^{p-1}$

$$
\begin{equation*}
\times \int_{0}^{\tau_{L}(t)} d s \int_{\mathrm{R}} d b\left[p_{t-s}(b-a)-p_{t^{\prime}-s}\left(b-a^{\prime}\right)\right]^{2}\left[X_{s}^{1}(b)\right]^{p}\left[X_{s}^{2}(b)\right]^{p} \tag{170b}
\end{equation*}
$$

For the first factor (170a) we use $\tau_{L}(t) \leq t$ and then the heat kernel estimate (53) to get the bound

$$
\begin{equation*}
c\left(\left|a-a^{\prime}\right|+\left|t-t^{\prime}\right|^{1 / 2}\right)^{p-1} \tag{171}
\end{equation*}
$$

On the other hand, in the second factor (170b) we introduce $\tilde{\phi}_{p(\beta+\gamma)}(b) \times$ $\tilde{\phi}_{-p(\beta+\gamma)}(b)$ [which is bounded away from 0, recall (7)], and use CauchySchwarz to get for the internal integral in (170b) the bound

$$
\begin{align*}
& \left|\int_{\mathrm{R}} d b\left[p_{t-s}(b-a)+p_{t^{\prime}-s}\left(b-a^{\prime}\right)\right]^{4} \tilde{\phi}_{2 p(\beta+\gamma)}(b)\right|^{1 / 2}  \tag{172a}\\
& \quad \times\left|\int_{\mathrm{R}} d b\left[X_{s}^{1}(b)\right]^{2 p}\left[X_{s}^{2}(b)\right]^{2 p} \tilde{\phi}_{-2 p(\beta+\gamma)}(b)\right|^{1 / 2} \tag{172b}
\end{align*}
$$

In the new first factor (172a), estimate three of the $p$-factor pairs by $2 p_{t^{\prime}-s}(0)$, and use the heat kernel estimate (38) to get the bound

$$
\begin{equation*}
c p_{t^{\prime}-s}^{3 / 2}(0)\left[\phi_{p(\beta+\gamma)}(a)+\phi_{p(\beta+\gamma)}\left(a^{\prime}\right)\right] \tag{173}
\end{equation*}
$$

for (172a). In the second new factor (172b) apply once more Cauchy-Schwarz to get

$$
\begin{equation*}
\leq\left\langle\left[X_{s}^{1}\right]^{4 p}, \phi_{-4 p \beta}^{1}\right\rangle^{1 / 4}\left\langle\left[X_{s}^{2}\right]^{4 p}, \tilde{\phi}_{-4 p \gamma}\right)^{1 / 4} \leq \sqrt{L}, \tag{174}
\end{equation*}
$$

where in the last step we used $s \leq \tau_{L}(t) \leq \tau_{L}^{1,4 p, 4 p \beta} \wedge \tau_{L}^{2,4 p, 4 p \gamma}$ [recall (165)]. Consequently, for (172a) and (172b) we have the bound

$$
\begin{equation*}
c p_{t^{\prime}-s}^{3 / 2}(0)\left[\phi_{p(\beta+\gamma)}(a)+\phi_{p(\beta+\gamma)}\left(a^{\prime}\right)\right] \sqrt{L} . \tag{175}
\end{equation*}
$$

Inserting (171) and (175) into (170a) and 170b) gives the bound
(176) $c \sqrt{L}\left(\left|a-a^{\prime}\right|+\left|t-t^{\prime}\right|^{1 / 2}\right)^{p-1}\left[\phi_{p(\beta+\gamma)}(a)+\phi_{p(\beta+\gamma)}\left(a^{\prime}\right)\right] \int_{0}^{t} d s p_{t^{\prime}-s}^{3 / 2}(0)$,
since $\tau_{L}(t) \leq t$. This is clearly bounded by the right-hand side of (167), completing the proof of Lemma 27.

We continue with the proof of Theorem 5(b). For the purpose of establishing a further control of the states of our process, for each $n \geq 1$, we consider the equidistant grid,

$$
\begin{equation*}
G_{n}:=\left\{\left(t_{n, i}, a_{n, j}\right): t_{n, i}:=i 2^{-n} T, a_{n, j}:=j 2^{-n}, 0 \leq i \leq 2^{n}, j \in Z\right\} \tag{177}
\end{equation*}
$$

partitioning $[0, T] \times \mathrm{R}$.
The idea is now to show that $X_{t}^{1}(a)$ is "not too small." As

$$
\begin{equation*}
X_{t}^{1}(a)=S_{t} x^{1}(a)+N_{t}^{1}(a) \tag{178}
\end{equation*}
$$

[recall (57), (104) and (166)], we first will show that for the fluctuation part $N^{1}$ with a large probability,

$$
\begin{equation*}
\left|N_{t}^{1}(a)\right| \leq \frac{1}{2} S_{t} x^{1}(a), \quad 0 \leq t \leq T, \quad a \in \mathrm{R}, \tag{179}
\end{equation*}
$$

[see (192) below]. In fact, using Lemma 27, we can estimate the increments of $N_{t}^{1}(a)$ for $(t, a)$ in the union $G=\cup_{n} G_{n}$ of grids with a large probability. Thus, for any $(t, a) \in[0, T] \times \mathrm{R}$, we can approximate $N_{t}^{1}(a)$ by the sum of the aforementioned increments towards the boundary $\{0\} \times \mathrm{R}$ in order to obtain (179). Then by (178),

$$
\begin{equation*}
X_{t}^{1}(a) \geq \frac{1}{2} S_{t} x^{1}(a), \quad 0 \leq t \leq T, \quad a \in \mathrm{R} \tag{180}
\end{equation*}
$$

with a large probability [see (193)].
Here are the details. Let $n \geq 1$. Two points $g=(t, b)$ and $g^{\prime}=\left(t^{\prime}, b^{\prime}\right)$ in the grid $G_{n}$ are called neighboring points, if one of their coordinates coincide and the other one are neighbors in the obvious meaning. For $0<\varepsilon_{1} \leq 1$, and neighboring points $g, g^{\prime} \in G_{n}$ with $g \geq g^{\prime}$, introduce the event
(181) $A_{\varepsilon_{1},{ }_{n}^{g}}^{g, g^{\prime}}:=\left\{\left|N_{t}^{1}(a)-N_{t^{\prime}}^{1}\left(a^{\prime}\right)\right| \geq 2^{-n / p} \varepsilon_{1} \phi_{\beta^{\prime}}(a), t \leq \tau_{L}^{1,4 p, 4 p \beta} \wedge \tau_{L}^{2,4 p, 4 p \gamma}\right\}$,
and denote by $A_{\varepsilon_{1}}$ the union of $A_{\varepsilon_{1}, n}^{g, g^{\prime}}$, with all these $g, g^{\prime} \in G_{n}$ and $n \geq 1$. By Markov's inequality and Lemma 27,

$$
\begin{aligned}
& \sup _{x \in I_{1}} P_{x}\left(A_{\varepsilon_{1}, n}^{g, g^{\prime}}\right) \\
& \quad \leq 2^{2 n} \varepsilon_{1}^{-2 p} \phi_{-2 p \beta^{\prime}}(a) \sup _{x \in I_{1}} P_{x}\left|N_{t}^{1}(a)-N_{t^{\prime}}^{1}\left(a^{\prime}\right)\right|^{2 p} \mathbf{1}\left\{t \leq \tau_{L}^{1,4 p, 4 p \beta} \wedge \tau_{L}^{2,4 p, 4 p \gamma}\right\} \\
& \quad \leq 2^{2 n} \varepsilon_{1}^{-2 p} \phi_{-2 p \beta^{\prime}}(a) c_{(167)} \sqrt{L}\left(2^{-n / 2}\right)^{p-1} T^{1 / 4}\left(\phi_{p(\beta+\gamma)}\left(a^{\prime}\right)+\phi_{p(\beta+\gamma)}(a)\right) \\
& \quad=c_{(167)} \sqrt{L} T^{1 / 4} 2^{-n(p-5) / 2} \varepsilon_{1}^{-2 p} \phi_{-2 p \beta^{\prime}}(a)\left(\phi_{p(\beta+\gamma)}\left(a^{\prime}\right)+\phi_{p(\beta+\gamma)}(a)\right)
\end{aligned}
$$

Using our assumption (163) on $\beta^{\prime}, \beta, \gamma$, having in mind $\alpha^{\prime}=a-2^{-n}$,

$$
\begin{equation*}
\sum_{a \in G_{n}(t)} \phi_{-2 p \beta^{\prime}}(a)\left(\phi_{p(\beta+\gamma)}\left(a-2^{-n}\right)+\phi_{p(\beta+\gamma)}(a)\right) \leq c 2^{n} \tag{182}
\end{equation*}
$$

where $G_{n}(t)$ denotes the section of $G_{n}$ with a fixed $t$ from the grid. Hence, since there are $2^{n}+1$ different $t$ in $G_{n}$,

$$
\begin{align*}
\sup _{x \in I_{1}} P_{x}\left(A_{\varepsilon_{1}}\right) & \leq \sum_{n \geq 1} c \sqrt{L} T^{1 / 4} 2^{-n(p-5) / 2} \varepsilon_{1}^{-2 p}\left(2^{n}+1\right) c 2^{n}  \tag{183}\\
& \leq c_{(183)} \sqrt{L} T^{1 / 4} \varepsilon_{1}^{-2 p},
\end{align*}
$$

since we assumed $p>9$. Similarly to [34], page 295, for all $0 \leq t \leq T$ and $a \in \mathrm{R}$ we then obtain

$$
\begin{equation*}
\left|N_{t}^{1}(a)\right| \leq c_{(184)} \varepsilon_{1}(1+T) \phi_{\beta^{\prime}}(a) \quad \text { on }\left\{T \leq \tau_{L}^{1,4 p, 4 p \beta} \wedge \tau_{L}^{2,4 p, 4 p \gamma}\right\} \cup A_{\varepsilon_{1}}^{\mathrm{c}} \tag{184}
\end{equation*}
$$

for some constant $c_{(184)}$ only depending on $p$ and $\beta^{\prime}$. Recalling (24) and (38),

$$
\begin{equation*}
\frac{1}{2} S_{t} x^{1} \geq \frac{1}{2} c_{1}^{\prime} S_{t} \phi_{\lambda_{1}^{\prime}} \geq c \phi_{\lambda_{1}^{\prime}} \geq c_{(185)} \phi_{\beta^{\prime}} \tag{185}
\end{equation*}
$$

Make $\varepsilon_{1}>0$ now so small that

$$
\begin{equation*}
c_{(185)} \geq c_{(184)} \varepsilon_{1}(1+T) \tag{186}
\end{equation*}
$$

implying

$$
\begin{equation*}
\frac{1}{2} S_{t} x^{1} \geq c_{(184)} \varepsilon_{1}(1+T) \phi_{\beta^{\prime}} \tag{187}
\end{equation*}
$$

Then, by (184),

$$
\begin{align*}
& P_{x}\left(\left|N_{t}^{1}(a)\right| \leq \frac{1}{2} S_{t} x^{1}(a), 0 \leq t \leq T, a \in \mathrm{R}\right) \\
& \quad \geq P_{x}\left(\left|N_{t}^{1}(a)\right| \leq c_{(184)} \varepsilon_{1}(1+T) \phi_{\beta^{\prime}}(a), 0 \leq t \leq T, a \in \mathrm{R}\right)  \tag{188}\\
& \quad \geq 1-P_{x}\left(T>\tau_{L}^{1,4 p, 4 p \beta}\right)-P_{x}\left(T>\tau_{L}^{2,4 p, 4 p \gamma}\right)-P_{x}\left(A_{\varepsilon_{1}}\right)
\end{align*}
$$

By the definition (164) of $\tau_{L}^{1,4 p, 4 p \beta}$,

$$
\begin{equation*}
P_{x}\left(T>\tau_{L}^{1,4 p, 4 p \beta}\right) \leq L^{-1} \sup _{x \in \zeta_{\text {tem }}^{+},|x|_{\lambda_{1}} \leq c_{1}} P_{x} \sup _{0 \leq t \leq 1}\left\|\left(X_{t}^{1}\right)^{4 p}\right\|_{-4 p \beta} \tag{189}
\end{equation*}
$$

since $T \leq 1$. Hence, because we assumed $\lambda_{1}>\beta$, by the uniform moment bounds in Corollary 15,

$$
\begin{equation*}
\sup _{x \in I_{1}} P_{x}\left(T>\tau_{L}^{1,4 p, 4 p \beta}\right) \leq c L^{-1} \leq \frac{\varepsilon}{6} \tag{190}
\end{equation*}
$$

where for the latter estimate we made finally $L$ sufficiently large (recall that we fixed $\varepsilon$ in the beginning of the proof). Similarly, we may assume that also

$$
\begin{equation*}
\sup _{x \in I_{1}} P_{x}\left(T>\tau_{L}^{2,4 p, 4 p \gamma}\right) \leq \frac{\varepsilon}{6} \tag{191}
\end{equation*}
$$

Now we further redefine our $T \in(0,1]$ by making it additionally so small that the right-hand side in (183) gets smaller than $\varepsilon / 6$. Then from the chain (188) of inequalities, from (190), (191) and (183) we obtain

$$
\begin{equation*}
\inf _{x \in I_{1}} P_{x}\left(\left|N_{t}^{1}(a)\right| \leq \frac{1}{2} S_{t} x^{1}(a), 0 \leq t \leq T, a \in \mathrm{R}\right) \geq 1-\varepsilon / 2 \tag{192}
\end{equation*}
$$

[as announced in (179)]. Then by (178),

$$
\begin{equation*}
\inf _{x \in I_{1}} P_{x}\left(X_{t}^{1}(a) \geq \frac{1}{2} S_{t} x^{1}(a), 0 \leq t \leq T, a \in \mathrm{R}\right) \geq 1-\varepsilon / 2 \tag{193}
\end{equation*}
$$

[as announced in (180)].
As $X^{1}$ is now seen to be not too small with a high probability, and since it is the catalyst for $X^{0}$, it will kill $X^{0}$ by time $T$ with a large probability. This idea we want to make precise by comparing $t \mapsto\left\|X_{t}^{0}\right\|_{\alpha}$, after an appropriate random time change, with a supercritical Feller's branching diffusion [see (216) below].

Let

$$
\begin{equation*}
\kappa:=\inf \left\{t \geq 0: X_{t}^{1}(a)<\frac{1}{2} S_{t} x^{1}(a) \text { for some } a \in \mathrm{R}\right\} . \tag{194}
\end{equation*}
$$

Then from (193) we already know that

$$
\begin{equation*}
\inf _{x \in I_{1}} P_{x}(\kappa \geq T) \geq 1-\varepsilon / 2 \tag{195}
\end{equation*}
$$

By Corollary 16 [recall our notation (12)], for $t \geq 0$,

$$
\begin{equation*}
\left\|X_{t}^{0}\right\|_{\alpha}=\left\|x^{0}\right\|_{\alpha}+\int_{0}^{t} d s\left\langle X_{s}^{0}, \frac{\sigma^{2}}{2} \Delta \tilde{\phi}_{-\alpha}\right\rangle+M_{t}^{0}\left(\tilde{\phi}_{-\alpha}\right) \tag{196}
\end{equation*}
$$

with the stochastic integral

$$
\begin{equation*}
M_{t}^{0}\left(\tilde{\phi}_{-\alpha}\right):=\int_{[0, t] \times \mathrm{R}} M_{0}^{0}(d(s, b)) \tilde{\phi}_{-\alpha}(b) \tag{197}
\end{equation*}
$$

satisfying

$$
\begin{equation*}
d\left\langle\left\langle M^{0}\left(\tilde{\phi}_{-\alpha}\right)\right\rangle\right\rangle_{t}=\gamma^{0}\left\langle X_{t}^{0} X_{t}^{1},\left(\tilde{\phi}_{-\alpha}\right)^{2}\right\rangle d t \tag{198}
\end{equation*}
$$

But for $t \leq \kappa$ we have by the definition (194) of $\kappa$,

$$
\begin{equation*}
X_{t}^{1} \geq \frac{1}{2} S_{t} x^{1} \geq c \phi_{\lambda_{1}^{\prime}} \tag{199}
\end{equation*}
$$

[as in (185)]. Moreover,

$$
\begin{equation*}
\phi_{\lambda_{1}^{\prime}} \tilde{\phi}_{-\alpha} \geq c \tag{200}
\end{equation*}
$$

since $\alpha>\lambda_{1}^{\prime}$ by assumption. Hence,

$$
\begin{equation*}
d\left\langle\left\langle M_{\cdot}^{0}\left(\tilde{\phi}_{-\alpha}\right)\right\rangle\right\rangle_{t} \geq c_{(201)}\left\|X_{t}^{0}\right\|_{\alpha} d t \quad \text { on }[0, \kappa] \tag{201}
\end{equation*}
$$

uniformly for $x \in I_{1}$.
Note that $\left\|X_{t}^{0}\right\|_{\alpha}=0$ if and only if $\left\|X_{t}^{0}\right\|=0$, and recall that the state 0 is absorbing for the continuous martingale $t \mapsto\left\|X_{t}^{0}\right\|$. Thus, for our further proof we may assume that $\left\|x^{0}\right\|_{\alpha}>0$.

Set

$$
\begin{equation*}
A_{t}:=\int_{0}^{t} d\left\langle\left\langle M^{0}\left(\tilde{\phi}_{-\alpha}\right)\right\rangle\right\rangle_{s} \frac{1}{\left\|X_{t}^{0}\right\|_{\alpha}}(\leq \infty), \quad t \in[0, \infty] \tag{202}
\end{equation*}
$$

We introduce the new time scale,

$$
\begin{equation*}
t \mapsto \vartheta_{t}:=\inf \left\{r \geq 0: A_{r}>t\right\} \tag{203}
\end{equation*}
$$

(on which $A$ grows linearly) and the process

$$
\begin{equation*}
U_{t}:=\left\|X_{\vartheta_{t}}^{0}\right\|_{\alpha}, \quad t<A_{\infty} \tag{204}
\end{equation*}
$$

This $U$ we want to bound by a supercritical Feller's branching diffusion.
By (196), we have

$$
\begin{equation*}
U_{t}=\left\|x^{0}\right\|_{\alpha}+\int_{0}^{\vartheta_{t}} d s\left\langle X_{s}^{0}, \frac{\sigma^{2}}{2} \Delta \tilde{\phi}_{-\alpha}\right\rangle+M_{t}, \quad t<A_{\kappa} \tag{205}
\end{equation*}
$$

where

$$
\begin{equation*}
t \mapsto M_{t}:=M_{\vartheta_{t}}^{0}\left(\tilde{\phi}_{-\alpha}\right), \quad t<A_{\kappa} \tag{206}
\end{equation*}
$$

is a continuous local martingale such that

$$
\begin{equation*}
d\langle\langle M\rangle\rangle_{t}=U_{t} d t \tag{207}
\end{equation*}
$$

In fact, from (202),

$$
\begin{equation*}
d\left\langle\left\langle M_{.}^{0}\left(\tilde{\phi}_{-\alpha}\right)\right\rangle_{t}=\left\|X_{t}^{0}\right\|_{\alpha} d A_{t}\right. \tag{208}
\end{equation*}
$$

hence, by (206), and a change of variables [see, e.g., [28], Proposition (0.4.9)],

$$
\begin{equation*}
\left\langle\langle M\rangle_{t}=\int_{0}^{\vartheta_{t}} d A_{s}\left\|X_{s}^{0}\right\|_{\alpha}=\int_{0}^{t} d s\left\|X_{\vartheta_{s}}^{0}\right\|_{\alpha}=\int_{0}^{t} d s U_{s}\right. \tag{209}
\end{equation*}
$$

By (205)-(207), and the martingale representation theorem (see, e.g., [18], Theorem 2.7.1), passing to an enlarged probability space $\left(\Omega^{\prime}, \mathscr{F}^{\prime}, \mathscr{P}\right)$, there is a (standard) Brownian motion $B$ in R such that

$$
\begin{equation*}
U_{t}=\left\|x^{0}\right\|_{\alpha}+\int_{0}^{\vartheta_{t}} d s\left\langle X_{s}^{0}, \frac{\sigma^{2}}{2} \Delta \tilde{\phi}_{-\alpha}\right\rangle+\int_{0}^{t} d B_{s} \sqrt{U_{s}}, \quad t<A_{\kappa} \tag{210}
\end{equation*}
$$

Again by a change of variables,

$$
\begin{equation*}
\int_{0}^{\vartheta_{t}} d s\left\langle X_{s}^{0}, \frac{\sigma^{2}}{2} \Delta \tilde{\phi}_{-\alpha}\right\rangle=\int_{0}^{t} d \vartheta_{s}\left\langle X_{\vartheta_{s}}^{0}, \frac{\sigma^{2}}{2} \Delta \tilde{\phi}_{-\alpha}\right\rangle . \tag{211}
\end{equation*}
$$

Recall from (66) that

$$
\begin{equation*}
\frac{\sigma^{2}}{2} \Delta \tilde{\phi}_{-\alpha} \leq c \tilde{\phi}_{-\alpha} \tag{212}
\end{equation*}
$$

and from (201) and (202) that

$$
\begin{equation*}
d A_{t} \geq c_{(201)} d t \quad \text { on }[0, \kappa] \tag{213}
\end{equation*}
$$

implying

$$
\begin{equation*}
d \vartheta_{s} \leq c d s \quad \text { on }\left[0, A_{\kappa}\right] \tag{214}
\end{equation*}
$$

Inserting (211), (212) and (214) into (210), we get

$$
\begin{equation*}
U_{t} \leq\left\|x^{0}\right\|_{\alpha}+c_{(215)} \int_{0}^{t} d s U_{s}+\int_{0}^{t} d B_{s} \sqrt{U_{s}}, \quad 0 \leq t \leq A_{\kappa} \tag{215}
\end{equation*}
$$

with $c_{(215)}$ uniform in $x \in I_{1}$. Thus, by comparison (see [29], V.43.1),

$$
\begin{equation*}
U \leq \widehat{U} \quad \text { on }\left[0, A_{\kappa}\right], \tag{216}
\end{equation*}
$$

where $\widehat{U}$ is the pathwise unique solution to

$$
\begin{equation*}
\widehat{U}_{t}=\left\|x^{0}\right\|_{\alpha}+c_{(215)} \int_{0}^{t} d s \widehat{U}_{s}+\int_{0}^{t} d B_{s} \sqrt{\widehat{U}_{s}}, \quad t \geq 0 \tag{217}
\end{equation*}
$$

In other words, $\widehat{U}$ is a certain supercritical Feller's branching diffusion.
Now

$$
\begin{equation*}
P_{x}\left(X_{t}^{0}=0, t \geq T\right) \geq P_{x}\left(\left\|X_{T}^{0}\right\|_{\alpha}=0, \kappa \geq T\right) \tag{218}
\end{equation*}
$$

But $T=\vartheta_{A_{T}}$ by the definitions (202) and (203). Hence, by definition (204) of $U$, we may continue inequality (218) with
(219) $\geq P_{x}\left(U_{A_{T}}=0, \kappa \geq T\right) \geq \mathscr{P}\left(\widehat{U}_{A_{T}}=0 \mid \widehat{U}_{0}=\left\|x^{0}\right\|_{\alpha}\right)-P_{x}(\kappa<T)$,
where we also used (216). But $A_{T} \geq c_{(213)} T$ by (213), and by assumption (23),

$$
\begin{equation*}
\left\|x^{0}\right\|_{\alpha} \leq c_{0}\left\|\phi_{\lambda_{0}}\right\|_{\alpha}=c_{0} c_{(220)} \tag{220}
\end{equation*}
$$

[recall (163)]. Thus, by the branching property of Feller's branching diffusion and the estimate (195), we may continue (219) with

$$
\begin{equation*}
\geq\left(\mathscr{P}\left(\widehat{U}_{c_{(213)} T}=0 \mid \widehat{U}_{0}=1\right)\right)^{c_{0} c_{(220)}}-\varepsilon / 2 \tag{221}
\end{equation*}
$$

But the latter probability expression is positive, thus the right-hand side in (221) can be made greater than or equal to $1-\varepsilon$ by choosing $c_{0}>0$ sufficiently small. This completes the proof of Theorem 5(b).

Acknowledgment. We are grateful to Don Dawson who suggested to us the cyclically catalytic model. We also thank Achim Klenke and an anonymous referee for carefully reading the manuscript and suggestions for an improvement of the exposition. Most of this project was done when the second author visited the Weierstrass Institute, respectively the first author the University of Tennessee. Both thank these institutions for hospitality and financial support.

## REFERENCES

[1] Boerlijst, M. C. and Hogeweg P. (1995). Spatial gradients enhance persistence of hypercycles. Phys. D 88 29-39.
[2] Bramson, M. and Griffeath, D. (1989). Flux and fixation in cyclic particle systems. Ann. Probab. 17 26-45.
[3] Cox, J. T., Dawson, D. A., Klenke, A. and Perkins, E. A. (2000). On the growth of one-type blocks in the mutually catalytic branching model on $Z^{d}$. Unpublished manuscript.
[4] Cox, J. T. and Klenke, A. (2000). Recurrence and ergodicity of interacting particle systems. Probab. Theory Related Fields 116 239-255.
[5] Cox, J. T., Klenke, A. and Perkins, E. A. (2000). Convergence to equilibrium and linear systems duality. In Stochastic Models (L. G. Gorostiza and B. G. Ivanoff, eds.) 41-66. Amer. Math. Soc., Providence, RI.
[6] Dawson, D. A., Etheridge, A. M., Fleischmann, K., Mytnik, L., Perkins, E. A. and Xiong, J. (2000). Mutually catalytic branching in the plane: finite measure states. Preprint 615, WIAS Berlin.
[7] Dawson, D. A., Etheridge, A. M., Fleischmann, K., Mytnik, L., Perkins, E. A. and XIONG, J. (2001). Mutually catalytic branching in the plane: infinite measure states. WIAS Berlin, Preprint 633.
[8] Dawson, D. A. and Fleischmann, K. (2000). In Infinite Dimensional Stochastic Analysis 145-170. Royal Netherlands Academy of Arts and Sciences, Amsterdam.
[9] Dawson, D. A. and Fleischmann, K. (2000). Catalytic and mutually catalytic superBrownian motions. WIAS Berlin, Preprint 546.
[10] Dawson, D. A., Fleischmann, K., Mytnik, L., Perkins, E. A. and Xiong, J. (2001). Mutually catalytic branching in the plane: uniqueness. Preprint 641, WIAS Berlin.
[11] Dawson, D. A. and Perkins, E. A. (1998). Long-time behavior and coexistence in a mutually catalytic branching model. Ann. Probab. 26 1088-1138.
[12] Durrett, R. (1995). Ten lectures on particle sytems. Lecture Notes in Math. 1608 97-201. Springer, Berlin.
[13] Durrett, R. and Levin S. A. (1998). Spatial aspects of interspecific competition. Theoret. Population Biol. 53 30-43.
[14] Dynkin, E. B. (1994). An Introduction to Branching Measure-valued Processes. Amer. Math. Soc., Providence, RI.
[15] Eithier, S. N. and Kurtz, T. G. (1986). Markov Process: Characterization and Convergence. Wiley, New York.
[16] Freund, J. and Schimansky-Geier, L. (1999). Diffusion in discrete ratchets. Phys. Rev. E 60 1304-1309.
[17] Gravner, J. and Griffeath, D. (1998). Cellular automaton growth on $\mathbf{z}^{2}$ : theorems, examples, and problems. Adv. in Appl. Math. 21 241-304.
[18] Ikeda, N. and Watanabe, S. (1981). Stochastic Differential Equations and Diffusion Processes. North-Holland, Amsterdam.
[19] Jacod, J. and Shiryaev, A. N. (1987). Limit Theorems for Stochastic Processes. Springer, Berlin.
[20] Kallianpur, G. and Xiong, J. (1995). Stochastic Differential Equations in InfiniteDimensional Spaces. IMS, Hayward, CA.
[21] Krylov, N. V. (1993). The selection of a Markov process from a Markov system of processes Izv. Akad. Nauk SSSR Ser. Mat. 37 691-708 (in Russian).
[22] Merino, S. (1996). Cyclic competition of three species in the time periodic and diffusive case. J. Math. Biol. 34 789-809.
[23] Mitoma, I. (1985). An $\infty$-dimensional inhomogeneous Langevin equation. J. Funct. Anal. 61 342-359.
[24] Molina, A., Gonzalez, J. and Lopez-Tenes, M. (1998). Application of the superposition principle to the study of CEC, CE, EC and catalytic mechanisms in cyclic chronopotentiometry III. J. Math. Chem. 23 277-296.
[25] Mollison, D., ed. (1995). Epidemics Models: Their Structure and Relation to Data. Cambridge Univ. Press.
[26] Mueller, C. and Perkins, E. A. (2000). Extinction for two parabolic stochastic PDEs on the lattice. Ann. Inst. H. Poincaré Probab. Statist., 36 301-338.
[27] MyTNIK, L. (1998). Uniqueness for a mutually catalytic branching model. Probab. Theory Related Fields 112 245-253.
[28] Revuz, D. and Yor, M. (1991). Continuous Martingales and Brownian Motion. Springer, Berlin.
[29] Rogers, L. C. G. and Williams, D. (1987). Diffusions, Markov Processes, and Martingales 2. Itô Calculus. Wiley, Chichester.
[30] Rujigrok, Th. and Rujigrok, M. (1997). A reaction-diffusion equation for a cyclic system with three components. J. Statist. Phys. 87 1145-1164.
[31] Shiga, T. (1994). Two contrasting properties of solutions for one-dimensional stochastic partial differential equations. Canad. J. Math. 46 415-437.
[32] Stroock, D. W. and Varadhan, S. R. S. (1979). Multidimensional Diffusion Processes. Springer, Berlin.
[33] Swart, J. M. (1999). Clustering of linearly interacting diffusions and universality of their long-time distribution. Preprint, Kathol. Univ. Nijmegen. Available at http://www-math.sci.kun.nl/math/onderzoek/reports/reports_nl.html.
[34] Tribe, R. (1995). Large time behavior of interface solutions to the heat equation with FisherWright white noise. Probab. Theory Related Fields 102 289-311.
[35] WALSH, J. B. (1986). An introduction to stochastic partial differential equations. École d’Été de Probabilités de Saint-Flour XIV. Lecture Notes in Math. 1180 266-439. Springer, Berlin.
[36] Williams, D. (1991). Probability with Martingales. Cambridge Univ. Press.

| Weierstrass Institute <br> for Applied Analysis and Stochastics | University of Tennessee, Department of Mathematics |
| :---: | :---: |
| Mohrenstr. 39 | Knoxville, Tennessee 37996-1300 |
| D-10117 Berlin | E-MAIL: jxiong@math.utk.edu |
| Germany |  |
| E-mALL: fleischmann@wias-berlin.de |  |


[^0]:    Received December 1999; revised September 2000.
    ${ }^{1}$ Supported in part by the DFG.
    AMS 2000 subject classifications. Primary 60K35; secondary 60H15, 60 J 80.
    Key words and phrases. Catalyst, reactant, superprocess, duality, martingale problem, cyclic reaction, global segregation of neighboring types, finite time survival, extinction, strong Markov selection, stochastic equation.

