

Regularization of the superposition principle: Potential theory meets Fokker-Planck equations

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Abstract

For a (probability measure valued) solution to a (possibly nonlinear) Fokker-Planck equation (FPE) the powerful superposition principle renders a probability measure on path space with one dimensional time marginals equal to this solution, and additionally solving the martingale problem for the (possibly nonlinear) Kolmogorov operator given by the FPE. The superposition principle thus reveals that such parabolic PDEs have a probabilistic counter part. The aim of this work is to go a substantial further step and, by exploiting the superposition principle, construct a full fledged Markov process, i.e. a family of path space measures for a large set of space time starting points connected by the Markov property, associated to the (linearized) FPE in the above way. In fact, under very general (merely measurability) conditions on the coefficients of the FPE this is achieved in this paper in such a way that the resulting process is a right process, which is a particularly useful class of Markov processes, enjoying among other regularity properties the strong Markov property, which is fundamental for the analysis of the underlying FPE as a (nonlinear) parabolic PDE by probabilistic tools. For example, as two main applications we construct fundamental flow solutions for the FPE and we prove a well-posedness result for the parabolic Dirichlet problem through probabilistic means for more general coefficients than could be treated in the existing literature. Furthermore, we introduce a Choquet capacity for such FPEs using the corresponding right process. The validity of the strong Markov property in the context of the superposition principle was an open problem even in the linear case. In this paper we solve this also in the nonlinear case, i.e. for path laws of solutions to McKean-Vlasov SDEs with Nemytskii type coefficients. A main application here is the FPE given by the generalized porous media equation and its corresponding McKean-Vlasov SDE. For the reader's convenience, an introduction to the theory of right process and its potential theory is given in the last section of the paper.

Keywords: Nonlinear Fokker–Planck–Kolmogorov equation; McKean–Vlasov stochastic differential equation; Superposition principle; Right process; Nonlinear Markov process; Porous media equation

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1 Introduction

Let $(\mu_t)_{t \in [0, T]} \subset \mathcal{P}(\mathbb{R}^d)$ be a weakly continuous solution to the time-dependent (non)linear Fokker-Planck equation

$$\frac{d}{dt} \mu_t = L_{t, \mu_t}^* \mu_t, \quad t \in (0, T), \tag{1.1}$$

with L_{t, μ_t} given by

$$L_{t, \mu_t} f(t, x) := \sum_{i=1}^d b_i(t, x, \mu_t) \frac{\partial}{\partial x_i} f(t, x) + \frac{1}{2} \sum_{i, j=1}^d a_{i, j}(t, x, \mu_t) \frac{\partial^2}{\partial x_i \partial x_j} f(t, x), \quad t \in (0, T), x \in \mathbb{R}^d,$$

where $f \in C_c^\infty((0, \infty) \times \mathbb{R}^d)$. Then, by [4], under mild assumptions on the coefficients a and b , the well-known superposition principle ensures the existence of a probability measure η on the paths-space $C([0, T]; \mathbb{R}^d)$ which is a solution to the (non)linear martingale problem associated to L_{t, μ_t} with one-dimensional time marginals given by $\mu_t, t \geq 0$, or equivalently, which renders a weak solution to the McKean-Vlasov SDE

$$dX(t) = b(t, x, \mathcal{L}_{X(t)}) dt + \sigma(t, x, \mathcal{L}_{X(t)}) dW(t), \quad \mathcal{L}_{X(t)} = \mu_t, t \in [0, T], \tag{1.2}$$

where W is a d -dimensional standard Brownian motion whilst $\mathcal{L}_{X(t)}$ denotes the law of $X(t), t \geq 0$. The converse implication, namely that the one-dimensional time marginals of a weak solution to the McKean-Vlasov SDE (1.2) solves the (non)linear Fokker-Planck equation (1.1), is just a simple consequence of Ito's formula. Thus, under mild assumptions on the coefficients, there is a well-understood correspondence between

(1.1) and (1.2). The hard implication is, of course, the one that comes from the superposition principle, which was first obtained in the linear case (i.e. the coefficients do not depend on μ_t) by [1] and [45], followed up by [64], [35], as well as by other contributions. The extension of the superposition principle to the nonlinear setting and the corresponding McKean-Vlasov SDEs was realized in [3], [4]. We refer to [36] for a comprehensive study of the theory of linear Fokker-Planck-Kolmogorov equations, to [39] for a detailed presentation of McKean-Vlasov SDEs and their applications, and to [6] for the analysis of nonlinear Fokker-Planck equations and their probabilistic counterpart which are nonlinear Markov processes in the sense of McKean; see also the seminal paper [54], as well as the very recent work [59] in which a general recipe is developed how to prove that the family of laws of solutions to a McKean-Vlasov stochastic differential equation with merely measurable coefficients (also only measurable in the measure variable) indexed by a large enough set of initial distributions form a nonlinear Markov process in the sense of McKean. This applies to a large number of examples (see [6] and the references therein), in particular, to also those in [8], [7], where the underlying nonlinear Fokker-Planck equation is the parabolic p -Laplace equation and the Leibenson equation respectively, with corresponding non-linear Markov process being the p -Brownian motion and the Leibenson process. See also [2] for a gradient flow theory in metric spaces and the superposition principle for the continuity equation.

Another fundamental feature of the superposition principle is that, under a weak uniqueness assumption, the martingale solution denoted above by η enjoys the "simple" Markov property; see Section 2.1, more precisely (2.10). Thus, under very mild assumptions on the coefficients, the correspondence of (non)linear Fokker-Planck equations and McKean Vlasov SDEs could be analyzed through the lens of Markov processes. However, it is well known that the "simple" Markov property (2.10) is typically not enough in order to benefit from the rich theory of Markov processes, since additional regularity properties of the process, like the strong Markov property, are necessary. This leads us to the formulation of the first main goal which we achieve in this paper:

- *We show that, under reasonable conditions (see Theorem 2.12, Theorem 2.15), one can construct a homogenized right (hence strong Markov) process X on a subset $E \subset [0, \infty) \times \mathbb{R}^d$ with $\mu_t(\{x : (t, x) \in E\}) = 1, t \geq 0$, which is associated to a given weakly continuous solution $(\mu_t)_{t \geq 0}$ to the linear Fokker-Planck equation (1.1) in such a way that the family of one-dimensional time marginals of the spatial component of X coincides with the given solution $(\mu_t)_{t \geq 0}$.*

As indicated above, a right process automatically has the strong Markov property, in particular we answer the question on the validity of this property for the superposed measure η mentioned above, and which was left open in [64, Remark 2.10]. Moreover, we deduce the strong Markov property for the nonlinear Markov processes associated with nonlinear Fokker-Planck equation from [6].

Furthermore, a right process has associated a transition semigroup with regularity properties that substantially generalize the usual ones valid for Feller transition semigroups. However, the class of right processes is a much larger class than the well studied Feller processes. For a concise introduction to the theory of right processes, we refer to Section 6.

We here recall that right processes can be immediately constructed starting from a Feller transition semigroup on a locally compact space with countable base, that is a Markov transition semigroup which is also a C_0 -semigroup on the space of continuous functions vanishing at infinity; see e.g. [33, Theorem 4.4]. However, the right processes we shall construct will not have a Feller transition semigroup.

As already mentioned, the interest in constructing right processes is motivated by the fact that they enjoy plenty of useful properties which are otherwise not ensured if a process satisfies merely the simple Markov property (2.10). Some of these key properties that a right process enjoys are briefly listed below (for details we refer to Section 6 and to the classical work [33] and [60]):

- The underlying filtration satisfies the usual hypotheses.
- The hitting times of Borel sets are stopping times, and they can be approximated from above by hitting times of compact sets.
- The strong Markov property holds.
- The Blumenthal zero-one law is satisfied.
- The Dirichlet problem can be analyzed by means of hitting distributions.

- Analytic concepts such as polar sets and Choquet capacities have a full probabilistic interpretations.
- Many useful transformations such as killing, time change, subordination, Girsanov transforms, concatenation, perturbation by kernels, are possible for right processes and leave this class invariant.

Concerning the methods and techniques of proof we would like to indicate the following:

- *In order to construct the desired right process we first treat linear time-dependent Fokker-Planck equations, and the approach we adopt mixes fundamental tools from the theory of generalized Dirichlet forms, the theory of Fokker-Planck-Kolmogorov equations, and potential theory. It consists of three main parts which are systematically treated in Section 2.2, Section 2.3, and Section 2.4, and which can be summarized as follows:*

- I (see Section 2.2) Starting from a solution $(\mu_t)_{t \geq 0}$ to the linear Fokker-Planck equation (2.3), for each $N > 0$ we use the theory of generalized Dirichlet forms for time-dependent operators on L^1 , as developed in [61] (see also [62] and [50], as well as the classical work [47] and [53]), in order to construct a time-homogenized right process on a subset of the product space $(0, N) \times \mathbb{R}^d$, which is of full $\mu_t(dx)dt$ -measure.
- II (see Section 2.4.2) We use a convenient restricted uniqueness assumption, the standard superposition principle from Theorem 2.5, and potential theoretic tools, in order to construct a right process on a subset of $(0, \infty) \times \mathbb{R}^d$ of full $\mu_t(dx)dt$ -measure, by taking the projective limit of the processes constructed in I for each $N > 0$ using [61].
- III (see Section 2.4.3) Finally, we use again the superposition principle together with potential theoretic tools developed in [21] in order to extend the process constructed in II to a conservative right process with continuous paths defined on a larger set $E \subset [0, \infty) \times \mathbb{R}^d$ such that $\mu_t(E) = 1, t \in [0, \infty)$, and such that, if we start from $\delta_0 \otimes \mu_0$ then the family of one-dimensional time marginals of the spatial component of the right process coincides with the given $(\mu_t)_{t \geq 0}$.

Among the main consequences of our results Theorem 2.12 and Theorem 2.15, we would like to mention the following:

- *We deduce several new results in the theory of time-dependent PDEs with only measurable coefficients, namely:*
 - We show the existence of a *fundamental flow solution* to the time-dependent linear Fokker-Planck equation (2.13); see Corollary 3.2.
 - We derive maximal tail estimates for the path-space law obtained by the superposition of the fundamental flow solution to the time-dependent linear Fokker-Planck equation (2.13) by means of Lyapunov functions; see Proposition 3.9.
 - We solve the backward Kolmogorov PDE associated to the time-dependent linear operators (2.1) by solving the parabolic Dirichlet problem through the hitting distribution of the homogenized right process obtained in Theorem 2.12-Theorem 2.15; see Proposition 3.4.
 - We solve the time-dependent linear Fokker-Planck equation perturbed by a nonlocal operator, and simultaneously the corresponding martingale problem, through a method of adding bounded jumps to a right process; see Section 3.4.

As already mentioned, the above results are valid for solutions to linear Fokker-Planck equations. By the linearization technique (in the spirit of [4]), we then lift these results to non-linear Fokker-Planck equations. More precisely, we achieve the second fundamental goal of this paper:

- *In Section 5, see Theorem 5.4, we construct a right processes X on a subset $E \subset [0, \infty) \times \mathbb{R}^d$ with $\mu_t(\{x : (t, x) \in E\}) = 1, t \geq 0$, in such a way that the family of one-dimensional time marginals of the path-law η of the spatial component of X started at $\delta_0 \otimes \mu_0$, coincides with the given solution $(\mu_t)_{t \geq 0}$ of the nonlinear Fokker-Planck equation (thus realizing McKean vision from [54]), or equivalently, η coincides with the path-law of the weak solution to the McKean-Vlasov SDE (5.4). In particular, we derive the strong*

Markov property for the weak solution to the McKean-Vlasov SDE, which, as mentioned before, was an open problem even for the linear case.

In the same Section 5 we achieve two more goals:

- We construct a Choquet capacity for the solution to the nonlinear Fokker-Planck equation through the corresponding hitting distribution of the solution to the McKean-Vlasov SDE (5.4); see Corollary 5.6.

- As our main application in the nonlinear case, we investigate sufficient explicit conditions on the coefficients of a class of nonlinear Fokker-Planck equations, more precisely, generalized porous media equations, for which all the hypotheses in Theorem 5.4 are fulfilled and hence our theory fully applies; see Section 5.2.

Many techniques in constructing and manipulating the right processes mentioned above are of potential theoretic nature. Thus, to make this work as self-contained as possible:

- *In Section 6 we give a concise introduction to the theory of probabilistic potential theory, with special emphasis on those tools that are fundamentally used in order to derive the main results from Section 2 and Section 3. A reader who is interested in the technical details of this work is encouraged to have a preliminary look at Section 6 before diving into the details of the main results.*

- We would like to point out that, besides well known results, in Section 6 we also include new results on right processes that are of general nature and which are, as well, needed in the main part concerning Fokker-Planck equations, namely Section 2 and Section 3; see e.g. Lemma 6.23, Proposition 6.24, Proposition 6.39, Proposition 6.44, Proposition 6.56, as well as those from Section 6.5.

The paper is structured as follows:

In Section 2 we deal with time-dependent linear Fokker-Planck equations. In Section 2.1 we briefly recall the well-known superposition principle, a restricted uniqueness assumption and its relation with the "simple" Markov property. In Section 2.2 we place our problem of constructing more regular Markov processes in the context given by generalized Dirichlet forms, recall some fundamental notions and results, and discuss two problems which have remained unsolved up to now (OP1 and OP2), and which are the main obstacles that we overcome in this work. In Section 2.3 we present the main results of Section 2. Section 2.4 is entirely devoted to the proof of these main results, namely Theorem 2.12 and Theorem 2.12, for which we allocate separate subsections.

In Section 3 we use the right processes constructed in Section 2 in order to deduce new results about time-dependent linear PDEs: Construction of fundamental flow solutions, solving parabolic Dirichlet problems by probabilistic means, deriving maximal tail estimates for the path-space law obtained by the superposition of the fundamental flow solution to the time-dependent linear Fokker-Planck equation (2.13) by means of Lyapunov functions, and solving Fokker-Planck equations perturbed by nonlocal operators.

In Section 4 we present several results of independent interest establishing that the square root of the density of solutions to Fokker-Planck equations with non-degenerate diffusion coefficients belongs locally to H^1 . This regularity result plays a central role in the present work, as it shows that one of our main assumptions, namely (2.14), is in fact satisfied under very mild conditions; it significantly generalizes [36, Theorem 7.4.1].

In Section 5 we extend to nonlinear Fokker-Planck equations the results presented above for the linear case. We construct Choquet capacities for such nonlinear equations by probabilistic means and we investigate explicit sufficient conditions on the coefficients for a class of generalized porous media equations to which our theory applies.

In Section 6 we offer a concise introduction to the theory of right processes and their potential theory, and we also include here some new results with full proofs, which we need in the previous sections.

2 Regularization of the superposition principle for linear time-dependent Fokker-Planck equations

Let us consider the time-dependent diffusion operator

$$\mathbb{L}_t f(t, x) := \sum_{i=1}^d b_i(t, x) \frac{\partial}{\partial x_i} f(t, x) + \frac{1}{2} \sum_{i,j=1}^d a_{i,j}(t, x) \frac{\partial^2}{\partial x_i \partial x_j} f(t, x), \quad (2.1)$$

for $(t, x) \in (0, \infty) \times \mathbb{R}^d$, which acts on Borel measurable functions $f : (0, \infty) \times \mathbb{R}^d \rightarrow \mathbb{R}$ which are twice continuously differentiable in the second argument, whilst

$$b = (b_i)_{1 \leq i \leq d} : [0, \infty) \times \mathbb{R}^d \rightarrow \mathbb{R}^d, \quad a = (a_{ij})_{1 \leq i, j \leq d} : [0, \infty) \times \mathbb{R}^d \rightarrow \mathbb{R}^{d \times d} \quad (2.2)$$

are Borel measurable coefficients, such that $a(t, x)$ is symmetric and non-negative definite for all $(t, x) \in [0, \infty) \times \mathbb{R}^d$.

Throughout, for a measurable space (E, \mathcal{B}) we denote by $\mathcal{M}(E)$ the set of all finite (non-negative) measures on E , whilst by $\mathcal{P}(E)$ we denote the subset of $\mathcal{M}(E)$ containing all probability measures on E . Furthermore, by $p\mathcal{B}$ (resp. $b\mathcal{B}$) we denote the space of all non-negative (resp. bounded), real valued, and \mathcal{B} -measurable functions defined on E ; $bp\mathcal{B} := p\mathcal{B} \cap b\mathcal{B}$. If E is a topological spaces then by $\mathcal{B}(E)$ we denote the Borel σ -algebra, whilst by $C_b(E)$ we denote the space of all real-valued, bounded, and continuous functions defined on E . Also, by $C_c(E)$ we denote the subset of $C_b(E)$ consisting of functions which are compactly supported.

For $\mu \in \mathcal{M}(\mathbb{R}^d)$ and $f \in p\mathcal{B}$ we often use the notation

$$\mu(f) := \int_E f d\mu.$$

Let $(\mu_n)_n \subset \mathcal{M}(\mathbb{R}^d)$ and $\mu \in \mathcal{M}(\mathbb{R}^d)$. We say that $(\mu_n)_{n \geq 1}$ converges weakly (resp. vaguely) to μ if

$$\lim_n \mu_n(f) = \mu(f) \quad \text{for all } f \in C_b(\mathbb{R}^d) \quad (\text{resp. } C_c(\mathbb{R}^d)).$$

Throughout, we use $B(a, r)$ to denote the Euclidean ball in \mathbb{R}^d centered at $a \in \mathbb{R}^d$ and with radius $r > 0$.

Definition 2.1. Let $0 \leq s < T < \infty$.

(i) A family $(\mu_t)_{t \in (s, T)} \subset \mathcal{M}(\mathbb{R}^d)$ is called a (weak, or distributional) solution on (s, T) to the linear Fokker-Planck equation

$$\frac{d}{dt} \mu_t = \mathsf{L}_t^* \mu_t, \quad t \in (s, T), \quad (2.3)$$

if $(\mu_t)_{t \in (s, T)}$ is a Borel curve in $\mathcal{M}(\mathbb{R}^d)$ endowed with the Borel σ -algebra $\mathcal{B}(\mathcal{M}(\mathbb{R}^d))$ associated to the topology of weak convergence,

$$\int_s^T \int_{B(0, R)} \|b(t, x)\| + \|a(t, x)\| \mu_t(dx) dt < \infty, \quad R > 0, \quad (2.4)$$

and

$$\int_s^T \int_{\mathbb{R}^d} \left[\frac{d}{dt} f(t, x) + \mathsf{L}_t f(t, x) \right] \mu_t(dx) dt = 0, \quad f \in C_c^\infty((s, T) \times \mathbb{R}^d). \quad (2.5)$$

(ii) We say that $(\mu_t)_{t \in [s, T)} \subset \mathcal{M}(\mathbb{R}^d)$ is a solution to the linear Fokker-Planck equation (2.3) (with initial condition μ_s) if $(\mu_t)_{t \in (s, T)}$ is a solution on (s, T) as in (i), and $\lim_{t \searrow s} \mu_t = \mu_s$ weakly.

(iii) We say that $(\mu_t)_{t \in [s, T)} \subset \mathcal{M}(\mathbb{R}^d)$ is a weakly (resp. vaguely) continuous solution to the linear Fokker-Planck equation (2.3) if it is a solution on (s, T) in the sense of (i), and the curve of measures $[s, T) \ni t \mapsto \mu_t \in \mathcal{M}(\mathbb{R}^d)$ is (resp. vaguely) continuous; this definition is obviously extended when $(\mu_t)_{t \in [s, T]}$ is merely right-continuous.

Remark 2.2. Let $(\mu_t)_{t \in [s, T)} \subset \mathcal{M}(\mathbb{R}^d)$ be a vaguely (right-)continuous solution to the linear Fokker-Planck equation (2.3). For $0 \leq s < T < \infty$, if we take f of the form $f(r, x) = g(r)h(x)$ with $g \in C_c^\infty((s, T))$ and $h \in C_c^\infty(\mathbb{R}^d)$, then using first Fubini theorem and then integration by parts, (2.5) gives

$$- \int_s^T g'(r) \int_{\mathbb{R}^d} h d\mu_r dr = \int_s^T g(r) \int_{\mathbb{R}^d} \mathsf{L}_t h d\mu_r dr. \quad (2.6)$$

Let $s < s' < t' < T$, $g_n \in C_c^\infty(s, T)$, $n \geq 1$ such that $\lim_n g_n = 1_{[s', t']}$ a.e., and $\lim_n g'_n = \delta_{t'} - \delta_{s'}$ in the sense of weak convergence of measures; note that this can be done e.g. by convoluting $1_{[s', t]}$ with smooth and compactly supported mollifiers ρ_n , namely taking $g_n := \rho_n * 1_{[s', t]}$, $n \geq 1$. Thus, replacing g with g_n in (2.6), using that $r \mapsto \int_{\mathbb{R}^d} h \, d\mu_r$ is bounded and continuous, and letting n tend to infinity, we obtain

$$\int_{\mathbb{R}^d} h \, d\mu_{t'} - \int_{\mathbb{R}^d} h \, d\mu_{s'} = \int_{s'}^{t'} \int_{\mathbb{R}^d} \mathbb{L}_r h \, d\mu_r \, dr, \quad s < s' < t' < T, h \in C_c^\infty(\mathbb{R}^d). \quad (2.7)$$

It is easy to see that one can also pass from (2.7) to (2.5), hence the two relations are in fact equivalent.

2.1 The superposition principle and the "simple" Markov property

For $0 \leq s < T < \infty$ we define

$$\Pi_s(t, \cdot) : C([s, T]; \mathbb{R}^d) \rightarrow \mathbb{R}^d, \quad \Pi_s(t, \omega) := \omega(t), \quad \omega \in C([s, T]; \mathbb{R}^d), t \in [s, T].$$

As usual, we endow $C([s, T]; \mathbb{R}^d)$ with the Borel σ -algebra rendered by the topology of uniform convergence, which coincides with the σ -algebra generated by $(\Pi_s(t, \cdot))_{t \in [s, T]}$.

If $\eta \in \mathcal{P}(C([s, T]; \mathbb{R}^d))$, then by η_t , $t \in [s, T]$ we denote the one-dimensional time marginals, namely

$$\eta_t := \eta \circ (\Pi_s(t, \cdot))^{-1} \in \mathcal{P}(\mathbb{R}^d), \quad t \in [s, T].$$

Definition 2.3. Let $0 \leq s < T < \infty$. A probability $\eta \in \mathcal{P}(C([s, T]; \mathbb{R}^d))$ is called a solution to the canonical martingale problem (canonical MP) associated to $(\mathbb{L}_t)_{t \in [s, T]}$, with initial distribution $\nu_s \in \mathcal{P}(\mathbb{R}^d)$ if

$$\int_s^T \int_{B(0, R)} \|b(t, x)\| + \|a(t, x)\| \eta_t(dx) dt < \infty, \quad R > 0,$$

and for all $f \in C_c^{1,2}((s, T) \times \mathbb{R}^d)$ the process $(M(t))_{t \in [s, T]}$ given by

$$M(t) := f(t, \Pi_s(t)) - f(s, \Pi_s(s)) - \int_s^t \bar{\mathbb{L}}f(r, \Pi_s(r)) \, dr, \quad t \in [s, T]$$

is an $(\mathcal{F}_s(t))$ -martingale with respect to η , where $\mathcal{F}_s(t) := \sigma(\Pi_s(r), s \leq r \leq t)$, $t \in [s, T]$.

Remark 2.4. Just by taking expectations, it follows immediately that if η is a solution to the martingale problem as in Definition 2.3, then $(\eta_t)_{t \in [s, T]}$ is a weakly continuous solution to the linear Fokker-Planck equation (2.13). The following highly nontrivial result, known as the superposition principle, gives the converse implication.

Theorem 2.5 (The superposition principle; [35], [64]). Let $(\mu_t)_{t \in [s, T]}$ be a weakly continuous solution to (2.3), and suppose further that

$$\int_s^T \int_{B(0, R)} \frac{|\langle b(t, x), x \rangle| + \|a(t, x)\|}{(1 + \|x\|)^2} \mu_t(dx) dt < \infty. \quad (2.8)$$

Then, according to [35, Theorem 1.1], there exists $\eta \in \mathcal{P}(C([s, T]; \mathbb{R}^d))$ a solution to the canonical martingale problem associated to $(\mathbb{L}_t)_{t \in [s, T]}$ on $[s, T]$ such that the one-dimensional time marginals $(\eta_t)_{t \in [s, T]}$ satisfy $\eta_t = \mu_t$, $t \in [s, T]$. Moreover, by [64, Proposition 2.8], for each $x \in \mathbb{R}^d$ there exists a probability $\eta^x \in \mathcal{P}(C([s, T]; \mathbb{R}^d))$ such that

(i) For every $A \in \mathcal{B}(C([s, T]; \mathbb{R}^d))$

$$\text{the mapping } \mathbb{R}^d \ni x \mapsto \eta^x(A) \text{ is Borel measurable and } \eta(A) = \int_{\mathbb{R}^d} \eta^x(A) \mu_s(dx).$$

(ii) For μ_s -a.e. $x \in \mathbb{R}^d$, η^x is a solution to the martingale problem associated to $(\mathbf{L}_t)_{t \in [s, T]}$ on $[s, T]$, with initial distribution δ_x .

As a consequence, if $\nu_s \in \mathcal{P}(\mathbb{R}^d)$ satisfies $\nu_s \leq c\mu_s$ for some constant $c \in (0, \infty)$, then

$$\eta^{\nu_s}(\cdot) := \int_{\mathbb{R}^d} \eta^x(\cdot) \nu_s(dx)$$

is also a solution to the martingale problem associated to $(\mathbf{L}_t)_{t \in [s, T]}$ on $[s, T]$, with initial distribution ν_s .

Remark 2.6. We point out that in [64], the superposition principle is derived with (2.8) replaced by the stronger condition

$$\int_{(s, T) \times \mathbb{R}^d} \|b(t, x)\| + \|a(t, x)\| \mu_t(dx) dt < \infty. \quad (2.9)$$

However, there are relevant examples for which (2.9) is not fulfilled, but for which the relaxed condition (2.8) is valid. Other key results in [64] remain valid if we assume (2.8) instead of (2.9), this being the case of [64, Lemma 2.12] which we shall use below in the proof of Proposition 2.7.

Let $N > 0$ and $(\mu_t)_{t \in [0, N]} \subset \mathcal{P}(\mathbb{R}^d)$ be a given weakly continuous solution to the linear Fokker-Planck equation (2.3). With respect to such a fixed solution, we consider the following uniqueness hypothesis:

H₁ μ For every $0 \leq s < N < \infty$ and $c > 0$ we have that whenever $(\nu_t)_{t \in [s, N]}, (\tilde{\nu}_t)_{t \in [s, N]} \subset \mathcal{P}(\mathbb{R}^d)$ are two weakly continuous solutions to the linear Fokker-Planck equation (2.3), we have:

$$\text{If } \nu_t, \tilde{\nu}_t \leq c\mu_t, t \in [s, N], \text{ and } \nu_s = \tilde{\nu}_s, \text{ then } (\nu_t)_{t \in [s, N]} = (\tilde{\nu}_t)_{t \in [s, N]}.$$

Proposition 2.7 (Well-posedness of the canonical MP). *Let $N > 0$ and $(\mu_t)_{t \in [0, N]} \subset \mathcal{P}(\mathbb{R}^d)$ be a given weakly continuous solution to the linear Fokker-Planck equation (2.3). If **H₁ μ** is satisfied, $0 \leq s < N$, $c \in (0, \infty)$, and $(\nu_t)_{t \in [s, N]} \subset \mathcal{P}(\mathbb{R}^d)$ is a weakly continuous solution to the linear Fokker-Planck equation (2.3) on $[s, N]$ satisfying $\nu_t \leq c\mu_t, t \in [s, N]$, then there exists a unique $\eta \in \mathcal{P}(C([s, N]; \mathbb{R}^d))$ which is a solution to the canonical MP associated to $(\mathbf{L}_t)_{t \in [s, N]}$ and such that $\eta_t = \nu_t$ for every $t \in [s, N]$.*

Proof. The statement follows by [64, Lemma 2.12], choosing $T = N$ and

$$\mathcal{R}_{[s, T]} := \left\{ \nu = (\nu_t)_{t \in [s, T]} \subset \mathcal{P}(\mathbb{R}^d) : \nu \text{ is a weakly continuous solution to (2.3) and} \right. \\ \left. \text{there exists } c < \infty \text{ such that } \nu_t \leq c\mu_t, t \in [s, T] \right\}, \quad 0 \leq s < T.$$

□

Another result from [64] that remains valid if we assume (2.8) instead of (2.9) is the following:

Proposition 2.8 ([64, Remark 2.3]). *Let $(\mu_t)_{t \in (s, T)}$ be a solution to the linear Fokker-Planck equation as in Definition 2.1, which in addition satisfies*

$$(\mu_t)_{t \in (s, T)} \subset \mathcal{P}(\mathbb{R}^d) \text{ as well as (2.8).}$$

Then $(\mu_t)_{t \in (s, T)}$ admits a unique dt-version $(\tilde{\mu}_t)_{t \in [s, T]}$ which is weakly continuous.

However, let us point that the proof of Proposition 2.8 requires a light adaptation. More precisely, recall first that [64, Remark 2.3] is fundamentally [2, Lemma 8.1.2]. The latter result, remains nevertheless valid if we replace the assumption (2.9) from [64] with the relaxed one (2.8) from [35], with the mention that, under (2.8), the tightness of $(\mu_t)_{t \in [s, T]}$ required in [2] is ensured by [35, Proposition 2.2 and Proof of Theorem 1.1]; also, note that [35, Proposition 2.2] is valid without the assumption that $(\mu_t)_t$ is weakly continuous, since the latter was used only for applying the Gronwall's lemma for continuous functions, but the Gronwall's lemma is valid for merely measurable functions as well, see e.g. [44, Theorem 5.1], and which can be employed in [35, Proposition 2.2].

The "simple" Markov property Let $(\mu_t)_{t \in [0, T]}$ be a weakly continuous solution to (2.3) such that (2.8) and \mathbf{H}_μ hold. Then, one can easily deduce:

The solution $\eta \in \mathcal{P}(C([0, T]; \mathbb{R}^d))$ provided by Theorem 2.5 and Proposition 2.7 has the simple Markov property, namely for every $f \in b\mathcal{B}(\mathbb{R}^d)$ we have

$$\mathbb{E}_\eta \{f(\Pi_0(t+s)) \mid \mathcal{F}_0(t)\} = \mathbb{E}_\eta \{f(\Pi_0(t+s)) \mid \Pi_0(t)\}, \quad \eta\text{-a.s.} \quad t \in [0, T], s \in [0, T-t], \quad (2.10)$$

where \mathbb{E}_η denotes the expectation with respect to the probability η on $C([0, T]; \mathbb{R}^d)$.

Recall that the filtration $(\mathcal{F}_0(t))_{t \in [0, T]}$ is not right continuous, and there is a priori no guarantee that the simple Markov property (2.10) is preserved relative to $(\mathcal{F}_0(t+))_{t \in [0, T]}$.

For the reader convenience let us rapidly prove (2.10), following a standard argument from e.g. [44, Theorem 4.2]: Let $t \in [0, T)$ and $F \in \mathcal{F}_0(t)$ with $\eta(F) > 0$ be fixed and consider the following probabilities

$$\eta^{(1)}(A) = \frac{\mathbb{E}_\eta \{1_F \mathbb{E}[A \mid \mathcal{F}_0(t)]\}}{\eta(F)}, \quad \eta^{(2)}(A) = \frac{\mathbb{E}_\eta \{1_F \mathbb{E}[A \mid \Pi_0(t)]\}}{\eta(F)}, \quad A \in \mathcal{F}_0(T).$$

Then, it follows precisely as in [44, Theorem 4.2] that the probability measures

$$\eta^{t, (i)} := \eta^{(i)} \circ (\Pi_0(t + \cdot))^{-1} \in \mathcal{P}(C([t, T]; \mathbb{R}^d)), \quad i \in \{1, 2\},$$

are both solutions to the martingale problem associated to $(L_r)_{r \in [t, T]}$, having the distribution at time t

$$\eta_t^{t, (1)}(B) = \frac{\eta(\Pi_0(t)^{-1}(B) \cap F)}{\eta(F)} = \eta_t^{t, (2)}(B), \quad B \in \mathcal{B}(\mathbb{R}^d).$$

By Remark 2.4 we get that $(\eta_s^{t, (i)})_{s \in [t, T]}$ is weakly continuous solution to the linear Fokker-Planck equation (2.3), $i \in \{1, 2\}$. Moreover, we have

$$\eta_s^{t, (i)} := \eta^{(i)} \circ (\Pi_0(s))^{-1} \leq \frac{\eta \circ (\Pi_0(s))^{-1}}{\eta(F)} = \frac{1}{\eta(F)} \mu_s, \quad r \in [t, T], \quad i \in \{1, 2\}.$$

Therefore, by \mathbf{H}_μ we deduce that $\eta_s^{t, (1)} = \eta_s^{t, (2)}$, hence for $B \in \mathcal{B}(\mathbb{R}^d)$

$$\begin{aligned} \mathbb{E}_\eta \{1_F \mathbb{E}[\Pi_{t+r} \in B \mid \mathcal{F}_0(t)]\} &= \eta(F) \eta_{t+r}^{t, (1)}(B) = \eta(F) \eta_{t+r}^{t, (2)}(B) \\ &= \mathbb{E}_\eta \{1_F \mathbb{E}[\Pi_{t+r} \in B \mid \Pi_t]\}, \quad r \in [0, T-t], \end{aligned}$$

from which the simple Markov property follows.

Remark 2.9. Note that in [44, Theorem 4.2], the simple Markov property is deduced under the much stronger assumption requiring that the martingale problem has at most one solution for every given initial distribution which even implies that the strong Markov property holds. We stress that above we proved the simple Markov property (2.10) under assumption \mathbf{H}_μ which is much weaker than the mentioned uniqueness assumption in [44, Theorem 4.2]. However, the strong Markov property can no longer be deduced under this weaker uniqueness assumption by the arguments in [44, Theorem 4.2]. In this paper we are fundamentally interested in proving the strong Markov property as well as other useful regularity properties for the solution the martingale problem associated to the linear Fokker-Planck equation (2.3), not assuming we have uniqueness for the martingale problem under any initial condition. Moreover, we extend these results to the non-linear case as well.

2.2 Towards more regularity through generalized Dirichlet forms

In this section we approach the superposition principle from the perspective of generalized Dirichlet forms for time-dependent operators on L^1 , as developed in [61]. As we shall point out, this approach has the key advantage of providing a "full" Hunt (in particular strong) Markov process which is associated to the solution of the linear Fokker-Planck equation. However, two important open problems show up, which we explain below in detail, and systematically solve afterwards.

Throughout we are given a solution $\mu := (\mu_t)_{t \geq 0}$ to the linear Fokker-Planck equation on $[0, \infty)$, i.e. on each $[0, N)$, $N > 0$, which consists of absolutely continuous probability measures on \mathbb{R}^d , namely: There exists $u : [0, \infty) \times \mathbb{R}^d \rightarrow \mathbb{R}_+$ jointly measurable, such that

$$\mu_t := u(t, x) dx, \quad \mu_t(\mathbb{R}^d) = 1, \quad t \geq 0, \quad (2.11)$$

$$\int_0^N \int_{B(0, R)} (\|b(t, x)\| + \|a(t, x)\|) u(t, x) dx dt < \infty, \quad N > 0, R > 0, \quad (2.12)$$

and

$$\int_0^\infty \int_{\mathbb{R}^d} \left[\frac{d}{dt} f(t, x) + \mathbf{L}_t f(t, x) \right] u(t, x) dx dt = 0, \quad f \in C_c^\infty((0, \infty) \times \mathbb{R}^d). \quad (2.13)$$

Further, we consider the measures

$$\bar{\mu} := u(t, x) dt \otimes dx \text{ on } [0, \infty) \times \mathbb{R}^d, \quad \bar{\mu}_N := u(t, x) dt \otimes dx \text{ on } [0, N) \times \mathbb{R}^d, N > 0,$$

so that, by setting

$$\bar{\mathbf{L}}f(t, x) := \frac{d}{dt} f(t, x) + \mathbf{L}_t f(t, x), \quad f \in C_c^\infty((0, \infty) \times \mathbb{R}^d),$$

(2.13) rewrites as

$$\int_0^\infty \int_{\mathbb{R}^d} \bar{\mathbf{L}}f d\bar{\mu} = 0, \quad f \in C_c^\infty((0, \infty) \times \mathbb{R}^d).$$

That is, $\bar{\mu}$ is an *infinitesimally invariant* σ -finite measure for $\bar{\mathbf{L}}$.

Below, we will fix such a measure $\bar{\mu}$ having certain properties, and call it a *reference measure* and its density a *reference density*.

Thus, we now perfectly fit in the framework of [61], which we further follow and from which we take over the next additional assumptions on the coefficients and on $\bar{\mu}$:

H_a For every $N, r > 0$ we have

$$\partial_{x_j} a_{i,j} \in L^2([0, N) \times B(0, r); \bar{\mu}_N), \quad 1 \leq i, j \leq d,$$

and there exists a constant $C = C(r, N) > 0$ such that

$$C^{-1} \|\xi\|^2 \leq \langle a(t, x)\xi, \xi \rangle \leq C \|\xi\|^2, \quad \xi \in \mathbb{R}^d, (t, x) \in [0, N) \times B(0, r).$$

H_b For every $N, r > 0$ we have

$$b \in L^2([0, N) \times B(0, r); \bar{\mu}_N).$$

H \sqrt{u} For every $h \in C_c^\infty(\mathbb{R}^d)$ and $N > 0$ we have

$$uh \in L^\infty([0, N) \times \mathbb{R}^d), \quad \text{and} \quad \sqrt{u}h \in L^2([0, N); H^1(\mathbb{R}^d)). \quad (2.14)$$

We say that condition **H_{a,b} \sqrt{u}** holds if all the above three conditions hold.

According to [61, Theorem 1.7, Remark 1.18, Theorem 1.9, Proposition 1.10, Remark 1.11, and Remark 1.12] we have the following collection of fundamental results:

Theorem 2.10 (cf. [61]). *Let u be a reference density (i.e. satisfying (2.11)-(2.13)) such that condition **H_{a,b} \sqrt{u}** is fulfilled. Then the following assertions hold for each $N > 0$:*

(i) *There exists a maximal extension $(\bar{\mathbf{L}}_1^N, \mathbf{D}(\bar{\mathbf{L}}_1^N))$ of $(\bar{\mathbf{L}}, C_c^\infty((0, N) \times \mathbb{R}^d))$ on $L^1([0, N) \times \mathbb{R}^d; \bar{\mu}_N)$ satisfying the following properties:*

(i.1) *It generates a C_0 -semigroup $(\bar{\mathbf{T}}_t^N)_{t \geq 0}$ on $L^1([0, N) \times \mathbb{R}^d; \bar{\mu}_N)$, which can be extended (keeping the same notation) as a C_0 -semigroup of contractions on every $L^p([0, N) \times \mathbb{R}^d; \bar{\mu}_N)$, $p \in [1, \infty)$, whose corresponding generator is further denoted by $(\bar{\mathbf{L}}_p^N, \mathbf{D}(\bar{\mathbf{L}}_p^N))$.*

(i.2) $(\bar{T}_t^N)_{t \geq 0}$ is of evolution type, namely for every $t \geq 0$

$$\bar{T}_t^N f(s, x) = 1_{[0, N-t)}(s) \bar{T}_t^N (f(s+t, \cdot))(s, x), \quad (2.15)$$

a.e. $(s, x) \in [0, N) \times \mathbb{R}^d$ and $f \in L^p([0, N) \times \mathbb{R}^d; \bar{\mu}_N)$, $1 \leq p \leq \infty$.

(ii) There exists a right (hence strong Markov) process (see Definition 6.10)

$$X^N = (\Omega^N, \mathcal{F}^N, (\mathcal{F}^N(t))_{t \geq 0}, (X^N(t))_{t \geq 0}, \mathbb{P}_{s,x}^N, (s, x) \in [0, N) \times \mathbb{R}^d),$$

with lifetime ζ^N and shift operators $(\theta^N(t))_t$, such that the following are satisfied:

(ii.1) Its transition function $(\bar{P}_t^N)_{t \geq 0}$ and corresponding resolvent $\bar{U}^N := (\bar{U}_\alpha^N)_{\alpha > 0}$ are such that for each $f \in b\mathcal{B}([0, N) \times \mathbb{R}^d)$, $\bar{U}_\alpha^N f$ is a $\bar{\mu}_N$ -version of $(\alpha - \bar{L}_1^N)^{-1} f$, $\alpha > 0$.

(ii.2) For $(s, x) \in [0, N) \times \mathbb{R}^d$ $\bar{\mu}_N$ -a.e. we have

$$\bar{\mathbb{P}}_{s,x}^N([0, \zeta^N) \ni t \mapsto X^N(t) \in [0, N) \times \mathbb{R}^d \text{ is continuous}) = 1. \quad (2.16)$$

(ii.3) If we set $X^N(t) = (X_1^N(t), X_2^N(t))$, $t \geq 0$, then for $(s, x) \in [0, N) \times \mathbb{R}^d$ $\bar{\mu}_N$ -a.e. we have

$$\bar{\mathbb{P}}_{s,x}^N(X_1^N(t) \neq t + s; t < \zeta^N) = 0. \quad (2.17)$$

(ii.4) For every $N > 0$, and every $\bar{\nu} \in \mathcal{P}([0, N) \times \mathbb{R}^d)$ which satisfies $\bar{\nu} \leq c \bar{\mu}_N$ for some constant $c > 0$, the process X^N solves the $\mathbb{P}_{\bar{\nu}}^N$ -martingale problem associated with $(\bar{\mathbb{L}}, C_c^\infty((0, \infty) \times \mathbb{R}^d))$ in the sense that

$$f(X^N(t \wedge \zeta^N)) - f(X^N(0)) - \int_0^{t \wedge \zeta^N} \bar{\mathbb{L}}f(X^N(s)) ds, \quad t \geq 0 \quad (2.18)$$

is an $(\mathcal{F}^N(t))$ -martingale under $\mathbb{P}_{\bar{\nu}}^N := \int_{[0, N) \times \mathbb{R}^d} \mathbb{P}_{s,x}^N \bar{\nu}(ds, dx)$, for all $f \in C_c^\infty((0, \infty) \times \mathbb{R}^d)$.

Remark 2.11. Recall that $(\bar{P}_t)_{t \geq 0}$ is represented by X^N through

$$\bar{P}_t f(s, x) = \mathbb{E}_{s,x}^N \{f(X^N(t)); t < \zeta^N\}, \quad t \geq 0, f \in b\mathcal{B}([0, N) \times \mathbb{R}^d).$$

In particular, $(\bar{P}_t)_{t \geq 0}$ is a sub-Markovian transition function which is not Markovian, as $\bar{\mathbb{P}}_{\bar{\mu}_N}^N(\zeta^N \leq N) = 1$ by (2.17). Throughout the paper we shall often consider the distribution $\mathbb{P}_{\bar{\nu}}^N \circ (X(t))^{-1} := \bar{\nu} \circ \bar{P}_t$ on $[0, N) \times \mathbb{R}^d$, $t \geq 0$, for some $\bar{\nu} \in \mathcal{P}([0, N) \times \mathbb{R}^d)$. However, since $\zeta^N \leq N$, $\mathbb{P}_{\bar{\nu}}^N \circ (X(t))^{-1}$ is merely as a sub-probability on $[0, N) \times \mathbb{R}^d$. We shall later on prove that under some conditions, we have $\zeta^N = N - s$ \mathbb{P}_{s, μ_s} -a.s., which will be crucial in order to show that the second components X_2^N , $N > 0$ admit a projective limit with infinite lifetime.

There are two unsolved problems that concern the processes X^N , $N > 0$ constructed in Theorem 2.10, (ii), which we explain below and solve afterwards (see Theorem 2.12 and Theorem 2.15 for the solutions to OP1 and OP2, respectively):

Open problem 1 (OP1): Is there a global Markov process X on $[0, \infty) \times \mathbb{R}^d$ which is obtained as the projective limit of the processes X^N , $N > 0$, constructed in Theorem 2.10, (ii)? And if yes, is it conservative, right, or Hunt? This issue is highly nontrivial since it is related to the uniqueness property of the martingale problems (2.18), or equivalently, to the corresponding linear Fokker-Planck equation (2.13). This is a sensitive issue since, in general, uniqueness is expected to hold only in a restrictive sense, as it will be discussed later on.

Open problem 2 (OP2): The second question is technically much more delicate than the first one, and it concerns both the locally constructed process X^N in Theorem 2.10, (ii), i.e. on each time interval $[0, N)$,

and the global one X which is the projective limit of X^N , $N > 0$, and which exists if **Q1** has an affirmative answer. It goes as follows: Fix some $N > 0$, take $f(t, x) = h(x)$, $(t, x) \in [0, N] \times \mathbb{R}^d$ with $h \in C_c^\infty(\mathbb{R}^d)$, and $\bar{\nu} \in \mathcal{P}([0, N] \times \mathbb{R}^d)$ with $\bar{\nu} \leq c\bar{\mu}_N$ for some $c > 0$. Then, using Theorem 2.10, (ii.3) and (2.18), it follows that

$$h(X_2^N(t \wedge \zeta^N)) - h(X_2^N(0)) - \int_0^{t \wedge \zeta^N} \bar{\mathbf{L}}_{s+X_1^N(0)} h(X_2^N(s)) ds, \quad t \geq 0$$

is a $(\mathcal{F}^N(t))$ -martingale under $\mathbb{P}_{\bar{\nu}}^N$ which is continuous on $[0, \zeta^N)$. Thus, if we are allowed to consider $\mathbb{P}_{\bar{\nu}}^N \circ (X^N(0))^{-1} = \delta_0 \otimes \nu_0$ for some $\nu_0 \in \mathcal{P}(\mathbb{R}^d)$, then the *spatial* component X_2^N solves the martingale problem associated with the time-dependent linear operator \mathbf{L}_t on $0 \leq t < \zeta^N$, starting at $t = 0$ from ν_0 . But this is not possible since the initial distribution of $\mathbb{P}_{\bar{\nu}}^N \circ (X_1^N(0))^{-1}$ is assumed to be absolutely continuous with respect to dt . This issue is a deep one and in fact ubiquitous in the theory of (generalized) Dirichlet forms. It essentially arises from the fact that (the resolvent of the) Markov process X^N is in fact constructed on the state space $[0, N] \times \mathbb{R}^d$, uniquely except on some abstract set \mathcal{N} which is $\bar{\mu}_N$ -inessential (or $\bar{\mu}_N$ -negligible and $\bar{\mu}_N$ -polar), and on \mathcal{N} it is trivially constructed so that if it starts from $x \in \mathcal{N}$ then it remains in x forever. As a consequence, such a Markov process set to be stuck at any point in \mathcal{N} will not solve the martingale problem starting from an initial distribution that charges \mathcal{N} . In the case of the process X^N furnished by Theorem 2.10, due to the fact that the first component is the uniform motion to the right, one can easily see that the set $\{0\} \times \mathbb{R}^d$ is never hit by the process, namely it is polar; hence, the process X^N given by Theorem 2.10, uniquely in law on $[0, T] \times \mathbb{R}^d$ except some abstract set \mathcal{N} as above, is in particular not guaranteed to solve the martingale problem (2.18) with initial distribution $\delta_0 \times \mu_0$. As a conclusion, *the second question is if, and under which conditions, one can construct the process X so that it can solve the martingale problem (2.18) starting at $t = 0$ from generic distributions of the type $\delta_0 \times \nu_0$* . By the same reason, this question concerns also the global process X from **Q1**, once we show it can indeed be constructed.

2.3 The main results: The corresponding right (in particular strong) Markov process

This subsection is devoted to answering the two fundamental questions raised above. Throughout we assume that the basic settings (2.11)-(2.13) are in force. Further, we introduce the following notation:

$$\begin{aligned} \mathcal{A}_{\leq u}^{s,N} &:= \{v : [s, N] \times \mathbb{R}^d \rightarrow [0, \infty) \text{ measurable} : v \leq u \text{ dt} \otimes dx\text{-a.e.}\} \\ \mathcal{A}_{\lesssim u}^{s,N} &:= \{v : [s, N] \times \mathbb{R}^d \rightarrow [0, \infty) \text{ measurable} : \exists c > 0 \text{ such that } v \leq cu \text{ dt} \otimes dx\text{-a.e.}\} \end{aligned} \quad (2.19)$$

We introduce the following key hypothesis that concerns the given solution $(\mu_t := u(t, x)dx)_{t \geq 0}$ of the linear Fokker-Planck equation (2.3):

$\mathbf{H}_{\leq u}$ $(\mu_t = u(t, x)dx)_{t \in [0, \infty)}$ is weakly continuous and for every $0 \leq s < N < \infty$ we have that whenever $v \in \mathcal{A}_{\leq u}^{s,N}$ induces a weakly right-continuous and vaguely left-continuous solution $(v(t, x)dx)_{t \in [s, N)}$ (of sub-probabilities) to the linear Fokker-Planck equation (2.3) on $[s, N)$ with initial condition $v(s, x)dx = \mu_s$, then we have the identification $(v(t, x)dx)_{t \in [s, N)} = (\mu_t)_{t \in [s, N)}$.

We introduce now a uniqueness hypothesis which is stronger than $\mathbf{H}_{\leq u}$:

$\mathbf{H}_{\lesssim u}$ We have that $\mathbf{H}_{\leq u}$ holds and for every $0 \leq s < N < \infty$ and $c > 0$ we have that whenever $(v(t, x)dx)_{t \in [s, N)}$ and $(\tilde{v}(t, x)dx)_{t \in [s, N)}$ are two weakly continuous curves with values in $\mathcal{P}(\mathbb{R}^d)$ which are solutions to the linear Fokker-Planck equation (2.3) on $[s, N)$ such that $\tilde{v}, v \in \mathcal{A}_{\lesssim u}^{s,N}$ and $v(s, x)dx = \tilde{v}(s, x)dx$, then the two solutions must coincide, namely $(v(t, x)dx)_{t \in [s, N)} = (\tilde{v}(t, x)dx)_{t \in [s, N)}$.

We have now all the ingredients so give a positive answer to the open question OP1:

Theorem 2.12 (Answer to OP1). *Let u be a reference density such that that $\mathbf{H}_{\mathbf{a}, \mathbf{b}}^{\sqrt{u}}$ and $\mathbf{H}_{\lesssim u}$ hold. Then the following assertions hold:*

(i) There exists a maximal extension $(\bar{\mathbb{L}}_1, \mathbb{D}(\bar{\mathbb{L}}_1))$ of $(\bar{\mathbb{L}}, C_c^\infty((0, \infty) \times \mathbb{R}^d))$ on $L^1([0, \infty) \times \mathbb{R}^d; \bar{\mu})$. It generates a C_0 -resolvent $\bar{U} = (\bar{U}_\alpha)_{\alpha > 0}$ on $L^1([0, \infty) \times \mathbb{R}^d; \bar{\mu})$, which, keeping the same notation, induces a C_0 -resolvent of Markov contractions \bar{U} on every $L^p([0, \infty) \times \mathbb{R}^d; \bar{\mu})$, and whose corresponding generator is further denoted by $(\bar{\mathbb{L}}_p, \mathbb{D}(\bar{\mathbb{L}}_p))$, $p \in [1, \infty)$. Moreover, for every $N > 0$ we have

$$\bar{U}_\alpha (f 1_{[0, N) \times \mathbb{R}^d})|_{[0, N) \times \mathbb{R}^d} = \bar{U}_\alpha^N (f|_{[0, N) \times \mathbb{R}^d}), \quad f \in L^p([0, \infty) \times \mathbb{R}^d; \bar{\mu}),$$

where \bar{U}_α^N is given in Theorem 2.10.

(ii) There exists a set $E \in \mathcal{B}([0, \infty) \times \mathbb{R}^d)$ and a conservative right (hence strong Markov) process on E

$$X = (\Omega, \mathcal{F}, \mathcal{F}_t, X(t), \mathbb{P}_{s,x}, (s, x) \in E, t \geq 0) \text{ such that:}$$

(ii.1) $\bar{\mu}([0, \infty) \times \mathbb{R}^d \setminus E) = 0$ and $\mu_s(\{x \in \mathbb{R}^d, (s, x) \in E\}) = 1$, $s > 0$.

(ii.2) The transition function of X denoted by $(P_t)_{t \geq 0}$, and the corresponding resolvent $\mathcal{U} := (U_\alpha)_{\alpha > 0}$, are such that

- $P_t(f|_E)$ is finely continuous for every $f \in C_b([0, \infty) \times \mathbb{R}^d)$
- $U_\alpha(f|_E)$ is a finely continuous $\bar{\mu}$ -version of $(\alpha - \bar{\mathbb{L}}_p)^{-1} f$, $\alpha > 0$, $f \in p\mathcal{B}([0, \infty) \times \mathbb{R}^d) \cap L^p([0, \infty) \times \mathbb{R}^d; \bar{\mu})$.

(ii.3) For $(s, x) \in E$ we have

$$\mathbb{P}_{s,x}([0, \infty) \ni t \mapsto X(t) \in E \subset [0, \infty) \times \mathbb{R}^d \text{ is continuous}) = 1.$$

(ii.4) If we set $X(t) = (X_1(t), X_2(t))$, $t \geq 0$, then for $(s, x) \in E$ we have

$$\mathbb{P}_{s,x}(X_1(t) = t + s; 0 \leq t < \infty) = 1.$$

Moreover, we have

$$\mathbb{P}_{s,\mu_s} \circ (X_2(t))^{-1} = \mu_{t+s} \quad \text{for every } s > 0 \text{ and } t \geq 0.$$

(ii.5) For every $(s, x) \in E$ and every $f \in C_c^2([0, \infty) \times \mathbb{R}^d)$ we have that

$$\mathbb{E}_{(s,x)} \left\{ \int_0^t |\bar{\mathbb{L}}f(X(r))| dr \right\} < \infty, \quad t \geq 0,$$

and

$$f(X(t)) - f(X(0)) - \int_0^t \bar{\mathbb{L}}f(X(r)) dr, \quad t \geq 0$$

is a continuous $(\mathcal{F}(t))$ -martingale under $\mathbb{P}_{(s,x)}$.

(ii.6) For every $s \in (0, \infty)$, and every $\nu_s \in \mathcal{P}(\mathbb{R}^d)$ which satisfies $\nu_s \leq c\mu_s$ for some constant $c > 0$, we have that

$$\nu_{t+s} := \mathbb{P}_{\delta_s \otimes \nu_s} \circ (X_2(t))^{-1}, \quad t \geq 0,$$

is the unique weakly continuous solution to the LFPE (2.3) on $[s, N)$ which consists of probability measures, belongs to $\mathcal{A}_{<u}^{s,N}$, and whose initial condition (at time s) is ν_s , for all $N > s$. Moreover, the (strong Markov) process X is the unique solution to the martingale problem on $[s, N)$ associated with $(\bar{\mathbb{L}}, C_c^\infty((0, \infty) \times \mathbb{R}^d))$, with initial distribution $\delta_s \otimes \nu_s$ and with $(X_2(t))_{t \geq 0}$ having the one-dimensional marginals in $\mathcal{A}_{<u}^{s,N}$, for every $N > s$.

Proof. The proof is given in Section 2.4.2. □

Remark 2.13. Note that the martingale problem in Theorem 2.12, (ii.5), may be solved for every $f \in b\mathcal{B}(E)$ such that $f \in D(\bar{\mathbb{L}}_1)$, in the following sense: Let \tilde{f} be a finely continuous $\bar{\mu}$ -version of f , and $E' \in \mathcal{B}(E)$ such that $E \setminus E'$ is $\bar{\mu}$ -inessential and \tilde{f} is bounded on E' ; \tilde{f} always exists. By considering the restriction of X to E' , taking g to be any $\mathcal{B}(E)$ -measurable $\bar{\mu}$ -version of $\bar{\mathbb{L}}_1 f$, and using Lemma 6.54 and Proposition 6.56, it follows that there exists a second set $E'' \in \mathcal{B}(E)$ such that $E \setminus E''$ is $\bar{\mu}$ -inessential and

$$\tilde{f}(X(t)) - \tilde{f}(X(0)) - \int_0^t g(X(r)) dr, \quad t \geq 0$$

is a càdlàg (\mathcal{F}_t) -martingale under $\mathbb{P}_{(s,x)}$ for every $(s,x) \in E''$.

In order to answer positively to the open question OP2, we need additional regularity properties for the resolvent \mathcal{U} constructed in Theorem 2.12, at points of the type $(0,x)$, namely at $s=0$. To this end, the key needed assumption turns out to be the following:

\mathbf{H}_0^{TV} The following hold:

(i) We have

$$\lim_{t \rightarrow 0} \mu_t = \mu_0 \quad \text{in total variation, or equivalently, } \lim_{t \rightarrow 0} u(t, \cdot) = u(0, \cdot) \text{ in } L^1(dx). \quad (2.20)$$

(ii) There exists a family \mathcal{A}_0 of probability densities which separates the measures on $\text{supp}(\mu_0)$, with the additional properties that if $\nu_0 \in \mathcal{A}_0$ then there exists a constant c such that $\nu_0 := \nu_0 dx \leq c u(0, x) dx$, and the corresponding solution $(\eta_t^{1, \nu_0})_{0 \leq t \leq 1}$, given by Theorem 2.5, (iii), is continuous at $t=0$ in total variation, or equivalently

$$\lim_{t \rightarrow 0} \frac{d\eta_t^{1, \nu_0}}{dx} = \nu_0 \quad \text{in } L^1(dx);$$

see also Remark 2.14, below.

Remark 2.14. Assume that $\mathbf{H}_{\lesssim u}$ holds. Then any weakly continuous $(\nu_t)_{0 \leq t \leq 1} \subset \mathcal{P}(\mathbb{R}^d)$ which is a solution to the linear Fokker-Planck equation (2.3) such that $\nu_t \leq c \mu_t, 0 \leq t \leq 1$ for some constant c , satisfies

$$\nu_t = \eta_t^{1, \nu_0}, \quad 0 \leq t \leq 1.$$

Thus, under \mathbf{H}_0^{TV} and if $\nu_0 \in \mathcal{A}_0$, we have

$$\lim_{t \rightarrow 0} \frac{d\nu_t}{dx} = \frac{d\nu_0}{dx} \quad \text{in } L^1(dx).$$

Theorem 2.15 (Answer to OP2). Let u be a reference density such that $\mathbf{H}_{\mathbf{a}, \mathbf{b}}^{\sqrt{u}}, \mathbf{H}_{\lesssim u}$ hold. If, in addition, \mathbf{H}_0^{TV} holds, then in assertions (ii.1), (ii.4), and (ii.6) from Theorem 2.12 we can include the case $s=0$ as well; more precisely, the three assertions can be respectively replaced with the following stronger ones:

(ii.1') $\bar{\mu}([0, \infty) \times \mathbb{R}^d) \setminus E = 0$ and $\mu_s(\{x \in \mathbb{R}^d, (s, x) \in E\}) = 1, s \geq 0$.

(ii.4') If we set $X(t) = (X_1(t), X_2(t)), t \geq 0$, then for $(s, x) \in E$ we have

$$\mathbb{P}_{s,x}(X_1(t) = t + s; 0 \leq t < \infty) = 1.$$

Moreover, we have

$$\mathbb{P}_{s, \mu_s} \circ (X_2(t))^{-1} = \mu_{t+s} \quad \text{for every } s \geq 0 \text{ and } t \geq 0.$$

(ii.6') For every $s \in [0, \infty)$, and every $\nu_s \in \mathcal{P}(\mathbb{R}^d)$ which satisfies $\nu_s \leq c \mu_s$ for some constant $c > 0$, we have that

$$\nu_{t+s} := \mathbb{P}_{\delta_s \otimes \nu_s} \circ (X_2(t))^{-1}, \quad t \geq 0,$$

is the unique weakly continuous solution to the LFPE (2.3) on $[s, N)$ which consists of probability measures, belongs to $\mathcal{A}_{\lesssim u}^{s, N}$, and whose initial condition (at time s) is ν_s , for all $N > s$. Moreover, the (strong Markov) process X is the unique solution to the martingale problem on $[s, N)$ associated with $(\bar{\mathbb{L}}, C_c^\infty((0, \infty) \times \mathbb{R}^d))$, with initial distribution $\delta_s \otimes \nu_s$ and with $(X_2(t))_{t \geq 0}$ having the one-dimensional marginals in $\mathcal{A}_{\lesssim u}^{s, N}$, for every $N > s$.

Proof. The proof is given in Section 2.4.3. □

Remark 2.16. Let \tilde{X} be given by the trivial extension (see Definition 6.41) of X given in Theorem 2.12 or in Theorem 2.15, from E to $[0, \infty) \times \mathbb{R}^d$. Then \tilde{X} is a right process by Proposition 6.40, and note that X is the restriction (see Definition 6.36) of \tilde{X} from $[0, \infty) \times \mathbb{R}^d$ to E , whilst all points in $[0, \infty) \times \mathbb{R}^d \setminus E$ are trap points for \tilde{X} . Consequently, \tilde{X} is a conservative norm-continuous right process on $[0, \infty) \times \mathbb{R}^d$.

2.4 Proofs of the main results

This subsection has three parts. In the first part we prepare some technical results that shall be used in the proofs of Theorem 2.12 and Theorem 2.15, which are given in great detail in the last two parts, respectively.

2.4.1 Some technical lemmas needed in the proofs

Lemma 2.17. Assume that $\mathbf{H}_{\mathbf{a}, \mathbf{b}}^{\sqrt{u}}$ holds, let $N > 0$ and $v^N \in \mathcal{A}_{\lesssim u}^{0, N}$ with $c > 0$ as in (2.19), such that $v^N(t, x)dx \in \mathcal{P}(\mathbb{R}^d)$ for every $t \in [0, N]$. Then there exists a measurable set $S \subset [0, N]$ of full Lebesgue measure such that for every $s \in S$ there exists $v_s^N : [0, \infty) \times \mathbb{R}^d \rightarrow [0, \infty)$ measurable satisfying the following:

(i) For every $t \in [0, \infty)$ we have

$$\mathbb{P}_{\delta_s \otimes v^N(s, x)dx}^N \circ (X_2^N(t))^{-1} = v_s^N(t, x)dx,$$

in particular, $v_s^N(t, x)dx$ is a sub-probability on \mathbb{R}^d .

(ii) We have dt-a.e. that

$$v_s^N(t, x)dx \leq c1_{[0, N]}(s+t)u(s+t, x)dx, \quad \text{hence} \quad v_s^N(\cdot - s, \cdot) \in \mathcal{A}_{\lesssim u}^{s, N}.$$

(iii) The family of sub-probability measures $(v_s^N(t-s, x)dx)_{t \in [s, N]}$ is a weakly right-continuous solution to the linear Fokker-Planck equation (2.3) on $[s, N]$ with $v_s^N(0, x)dx = v^N(s, x)dx$. Moreover, $(s, N) \ni t \mapsto v_s^N(t-s, x)dx$ is also vaguely left-continuous.

Proof. First of all, recall that $\bar{\mu}_N$ is (\bar{P}_t^N) -sub-invariant, so that the dual (w.r.t. $\bar{\mu}_N$) semigroup of kernels $\bar{P}_t^{N, *}, t \geq 0$ is sub-Markovian, namely $\bar{P}_t^{N, *}1 \leq 1, t \geq 0$. Moreover, property (2.15) transfers by duality as

$$\bar{P}_t^{N, *} f(s, x) = \bar{P}_t^{N, *} (f(s-t, \cdot))(s, x)1_{[t, N]}(s) \quad (s, x) \in [0, N] \times \mathbb{R}^d \text{ a.s.}$$

Let $g, h \in C([0, N]; \mathbb{R})$ and $f \in C_b(\mathbb{R}^d)$. Then

$$\begin{aligned} & \int_0^N g(t) \int_0^N h(s) \mathbb{E}_{\delta_s \otimes v^N(s, x)dx}^N f(X_2^N(t)) \, ds dt \\ &= \int_0^N g(t) \int_{[0, N] \times \mathbb{R}^d} h(s) \bar{P}_t^N f(s, x) v^N(s, x) \, dx ds dt \\ &= \int_0^N g(t) \int_{[0, N] \times \mathbb{R}^d} h(s) \frac{v^N(s, x)}{u(s, x)} \bar{P}_t^N f(s, x) \bar{\mu}_N(ds, dx) dt \\ &= \int_0^N g(t) \int_{[0, N] \times \mathbb{R}^d} f(x) \bar{P}_t^{N, *} \left(h(\cdot) \frac{v^N(\cdot, \cdot)}{u(\cdot, \cdot)} \right) (s, x) \bar{\mu}_N(ds, dx) dt \\ &= \int_0^N g(t) \int_{[0, N] \times \mathbb{R}^d} f(x) h(s-t) 1_{[t, N]}(s) \bar{P}_t^{N, *} \left(\frac{v^N}{u} \right) (s, x) \bar{\mu}_N(ds, dx) dt \\ &= \int_0^N g(t) \int_0^N h(s) \int_{\mathbb{R}^d} f(x) 1_{[0, N-t]}(s) \bar{P}_t^{N, *} \left(\frac{v^N}{u} \right) (s+t, x) u(s+t, x) \, dx ds dt. \end{aligned}$$

Consequently, by setting

$$\tilde{v}_s^N(t, x) := 1_{[0, N)}(s+t) \bar{P}_t^{N,*} \left(\frac{v^N}{u} \right) (s+t, x) u(s+t, x), \quad (s, t, x) \in [0, N) \times [0, \infty) \times \mathbb{R}^d,$$

and since f, g, h are arbitrarily chosen, it follows that

$$\mathbb{P}_{\delta_s \otimes v^N(s, x) dx}^N \circ (X_2(t))^{-1} = \tilde{v}_s^N(t, x) dx \quad \text{for } (s, t) \in [0, N) \times [0, \infty) \text{ } ds dt\text{-a.e.} \quad (2.21)$$

Also, notice that because $v^N \in \mathcal{A}_{\lesssim u}^{0, N}$ and because $\bar{P}_t^{N,*} 1 \leq 1$, we get that

$$\tilde{v}_s^N(t, x) dx \leq c 1_{[0, N)}(s+t) u(s+t, x) dx \quad ds dt\text{-a.e.} \quad (2.22)$$

Now we shall construct a modification of $\tilde{v}_s^N(t, x) dx$ which is weakly continuous with respect to $t \geq 0$, ds -a.e. To this end, let $D \subset [0, \infty)$ be dense (from the right) and countable, and $\tilde{S} \subset [0, N)$ be measurable and of full Lebesgue measure such that for every $(s, t) \in \tilde{S} \times D$ we have

$$\mathbb{P}_{\delta_s \otimes v(s, x) dx}^N \circ (X_2(t))^{-1} = \tilde{v}_s^N(t, x) dx \quad \text{and} \quad \tilde{v}_s^N(t, x) dx \leq c 1_{[0, N)}(s+t) u(s+t, x) dx, \quad (2.23)$$

which is possible due to (2.21) and (2.22). Let $s \in \tilde{S}$ be fixed. From the last inequality we deduce that $(\tilde{v}_s^N(t, x) dx)_{t \in D}$ is a family of sub-probability measures which is tight, hence by Prohorov's theorem, it is sequentially relatively compact with respect to the weak convergence. Further, let $t \in [0, \infty)$, and $(t_n)_{n \geq 1} \subset D$ such that $\lim_n t_n = t$ decreasingly. By passing to a subsequence, we may also assume that there exists a sub-probability measure $m_s^N(t)$ on \mathbb{R}^d such that

$$\lim_n \tilde{v}_s^N(t_n, x) dx = m_s^N(t), \quad \text{weakly.}$$

But with $s \in \tilde{S}$ fixed as above, $t \mapsto \mathbb{P}_{\delta_s \otimes v(s, x) dx}^N \circ (X_2(t))^{-1}$ is (right) weakly continuous, hence (2.23) and the above equality entail

$$\mathbb{P}_{\delta_s \otimes v^N(s, x) dx}^N \circ (X_2(t))^{-1} = m_s^N(t) = \lim_{D \ni r \searrow t} \tilde{v}_s^N(r, x) dx, \quad \text{weakly.}$$

Moreover, by the inequality in (2.23) we also have that for every $s \in \tilde{S}$ there exists $v_s^N(t, x)$, $(t, x) \in [0, \infty) \times \mathbb{R}^d$ such that

$$m_s^N(t) = v_s^N(t, x) dx \leq c 1_{[0, N)}(s+t) u(s+t, x) dx, \quad t \in [0, \infty). \quad (2.24)$$

It is then straightforward to see that for each $s \in \tilde{S}$ we have that $v_s^N(\cdot, \cdot)$ is jointly measurable on $[0, \infty) \times \mathbb{R}^d$. Thus, we have proved (i) and (ii).

To prove (iii), let us note that by Remark 2.11 and the above construction we have for all $s \in \tilde{S}$ that

$$\int_{\mathbb{R}^d} f(x) v_s^N(t, x) dx = \int_{\mathbb{R}^d} v^N(s, x) \mathbb{E}_{s, x}^N \{ f(X_2^N(t)); t < \zeta^N \} dx, \quad f \in C_b(\mathbb{R}^d),$$

hence using the path-continuity property (2.16) we get that

$$[0, \infty) \ni t \mapsto v_s^N(t, x) dx \quad \text{is weakly right continuous.}$$

Further, let us show that one can choose $S \subset \tilde{S}$ of full measure on $(0, N)$ such that for every $s \in \tilde{S}$ we have that $(v_s^N(t, x) dx)_{t \in (0, N-s)}$ is a solution to the linear Fokker-Planck equation (2.3). We proceed as follows: First of all, notice that since \bar{P}_t^N and \bar{T}_t^N coincide as operators on $L^p([0, N) \times \mathbb{R}^d; \bar{\mu}_N)$ for all $t \geq 0$, we get that $(\bar{P}_t^N)_{t \geq 0}$ solves the Cauchy equation

$$\bar{P}_t^N f(s, x) - f(s, x) = \int_0^t \bar{P}_r^N \bar{\Gamma} f(s, x) dr \bar{\mu}_N\text{-a.e.}, \quad 0 < t < \infty, f \in C_c^\infty(\mathbb{R}^d). \quad (2.25)$$

Let $g, h \in C_c^\infty((0, N); \mathbb{R})$ and $f \in C_c^\infty(\mathbb{R}^d)$, so that using the evolution character of (\overline{T}_t^N) given by (2.15), we have

$$\begin{aligned}
& \int_0^N g(s) \int_s^N h'(t) \int_{\mathbb{R}^d} f(x) v_s^N(t-s, x) dx dt ds = \int_0^N g(s) \int_s^N h'(t) \int_{\mathbb{R}^d} \overline{P}_{t-s} f(s, x) v(s, x) dx dt ds \\
&= \int_0^N g(s) \int_{\mathbb{R}^d} \int_s^N h'(t) \overline{P}_{t-s} f(s, x) v(s, x) dx dt ds \\
&= \int_0^N g(s) \int_{\mathbb{R}^d} h(s) f(x) v(s, x) dx ds - \int_0^N g(s) \int_s^N h(t) \int_{\mathbb{R}^d} \overline{P}_{t-s} \mathbb{L} f(s, x) v^N(s, x) dx dt ds \\
&= \int_0^N g(s) h(s) \int_{\mathbb{R}^d} f(x) v(s, x) dx ds - \int_0^N g(s) \int_s^N h(t) \int_{\mathbb{R}^d} \mathbb{L}_t f(x) v_s^N(t-s, x) dx dt ds.
\end{aligned}$$

Thus

$$\int_0^N g(s) \int_s^N \int_{\mathbb{R}^d} (h'(t) f(x) + h(t) \mathbb{L}_t f(x) v_s^N(t-s, x)) dx dt ds = \int_0^N g(s) h(s) \int_{\mathbb{R}^d} f(x) v(s, x) dx ds.$$

Since g is arbitrarily chosen and using a density argument, it is easy to see that one can find a measurable subset $S' \subset \tilde{S}$ of full Lebesgue measure on $(0, N)$ such that

$$\int_s^N \int_{\mathbb{R}^d} \left(\frac{d}{dt} f(t, x) + \mathbb{L}_t f(t, x) v_s^N(t-s, x) \right) dx dt = \int_{\mathbb{R}^d} f(s, x) v(s, x) dx$$

for all $s \in S'$ and $f \in C_c^\infty((0, N) \times \mathbb{R}^d)$. Thus, for every $s \in S'$ and $f \in C_c^\infty((s, N) \times \mathbb{R}^d)$ we obtain that

$$\int_s^N \int_{\mathbb{R}^d} \left(\frac{d}{dt} f(t, x) + \mathbb{L}_t f(t, x) v_s^N(t-s, x) \right) dx dt = 0,$$

which means that $(v_s^N(t-s, x) dx)_{t \in [s, N]}$ is a weakly right-continuous solution to the linear Fokker-Planck equation (2.3).

To finish the proof of (iii) it remains to show that $(v_s^N(t-s, x) dx)_{t \in (s, N)}$ is vaguely left-continuous. To this end, note that if $\mathcal{C} \subset C_c^\infty((0, N) \times \mathbb{R}^d)$ be countable and dense in $C_c((0, N) \times \mathbb{R}^d)$ with respect to the uniform convergence, so that by (2.25) we get for every $t > 0$, ds -a.e. on $(0, N)$ that

$$\int_{\mathbb{R}^d} \overline{P}_t f(s, x) v^N(s, x) dx = \int_{\mathbb{R}^d} 1_{[0, N)}(t+s) f(s+t, x) v_s^N(t, x) dx, \quad f \in \mathcal{C}.$$

and

$$\int_{\mathbb{R}^d} \overline{P}_t f(s, x) v^N(s, x) dx = \int_{\mathbb{R}^d} f(s, x) v^N(s, x) dx + \int_{\mathbb{R}^d} v^N(s, x) \int_0^t \overline{P}_r \mathbb{L} f(s, x) dr dx, \quad f \in \mathcal{C}.$$

Thus, on the one hand, there exists $S \subset S'$ of full Lebesgue measure on $(0, N)$ such that for every $s \in S'$

$$\int_{\mathbb{R}^d} \overline{P}_t f(s, x) v^N(s, x) dx = \int_{\mathbb{R}^d} f(s, x) v^N(s, x) dx + \int_{\mathbb{R}^d} v^N(s, x) \int_0^t \overline{P}_r \mathbb{L} f(s, x) dr dx, \quad f \in \mathcal{C}, t \in \mathbb{Q}_+,$$

in particular the above integrals are well defined. On the other hand,

$$[0, \infty) \ni t \mapsto \int_{\mathbb{R}^d} \overline{P}_t f(s, x) v^N(s, x) dx = \int_{\mathbb{R}^d} f(s, x) v^N(s, x) dx + \int_{\mathbb{R}^d} v^N(s, x) \int_0^t \overline{P}_r \mathbb{L} f(s, x) dr dx$$

is right continuous for every $f \in C_c^\infty((0, N) \times \mathbb{R}^d)$ and $s \in S$, hence we must have

$$\int_{\mathbb{R}^d} \overline{P}_t f(s, x) v^N(s, x) dx = \int_{\mathbb{R}^d} f(s, x) v^N(s, x) dx + \int_{\mathbb{R}^d} v^N(s, x) \int_0^t \overline{P}_r \mathbb{L} f(s, x) dr dx, \quad f \in \mathcal{C}, t \in [0, \infty).$$

This means that for every $s \in S$ and $f \in \mathcal{C}$ we have

$$[0, \infty) \ni t \mapsto \int_{\mathbb{R}^d} \bar{P}_t f(s, x) v^N(s, x) dx = \int_{\mathbb{R}^d} 1_{[0, N)}(t + s) f(s + t, x) v_s^N(t, x) dx \in \mathbb{R} \quad \text{is continuous.}$$

By a density argument we can further assume that \mathcal{C} consists of functions with separate variables, so taking $f \in \mathcal{C}$ of the form $f(r, x) = g(r)h(x)$, $(r, x) \in (0, N) \times \mathbb{R}^d$, we get for every $s \in S$ that

$$[0, N - s) \ni t \mapsto \int_{\mathbb{R}^d} h(x) v_s^N(t, x) dx \in \mathbb{R} \quad \text{is continuous.}$$

Since we assumed that the class of functions h for which the above holds is dense in $C_c(\mathbb{R}^d)$ with respect to the uniform convergence, we get that

$$[0, N - s) \ni t \mapsto v_s^N(t, x) dx \in \mathcal{M}(\mathbb{R}^d) \quad \text{is vaguely continuous.}$$

□

Remark 2.18. We emphasize that although X^N has continuous paths a.s. on $[0, \zeta^N)$ as stated in (2.16), it does not mean that the curve of sub-probability measures $\mathbb{P}_{s, x} \circ (X_2^N(t))^{-1}$ on \mathbb{R}^d is weakly left-continuous on $(s, N - s)$ because the process $t \mapsto 1_{[0, \zeta^N)}(t)$ is not guaranteed to be stochastically left continuous in t . In other words, based only on the results obtained in [61], it is not clear that ζ^N is a continuous random variable, at least on the event $\zeta^N < N - s$ for s fixed; such a property would have solved the left-continuity issue in $t \in (s, N - s)$, since in that case, the jump event $\zeta^N = t$ would have zero probability.

Lemma 2.19. Assume that conditions $\mathbf{H}_{\mathbf{a}, \mathbf{b}}^{\sqrt{\mathbf{u}}}$ and $\mathbf{H}_{\leq \mathbf{u}}$ are fulfilled. Then, the following assertions hold for every $0 < s < N < \infty$.

(i) We have the identification

$$\mathbb{P}_{\delta_s \otimes \mu_s}^N \circ (X_2^N(t))^{-1} = \mu_{t+s}, \quad 0 \leq t < N - s. \quad (2.26)$$

(ii) We have the following refinement of Theorem 2.10, (ii.3):

$$\mathbb{P}_{\delta_s \otimes \mu_s}^N (X_1^N(t) \neq t + s; t < \zeta^N) = 0. \quad (2.27)$$

(iii) The lifetime ζ^N satisfies

$$\mathbb{P}_{s, x}^N (\zeta^N = N - s) = 1 \quad \mu_s\text{-a.e. } x \in \mathbb{R}^d, \quad (2.28)$$

Proof. (i)-(ii). Since $\mathbf{H}_{\mathbf{a}, \mathbf{b}}^{\sqrt{\mathbf{u}}}$ is in force we apply Lemma 2.17 for $v^N(s, x) = u(s, x)$, $(s, x) \in [0, \infty) \times \mathbb{R}^d$, noticing that we trivially have $u \in \mathcal{A}_{\leq \mathbf{u}}^{0, N}$ for all $N > 0$. Thus, for every $N > 0$ there exists $S \subset [0, N)$ of full Lebesgue measure such that for every $s \in S$ there exists a family $(v_s^N(t, x) dx)_{t \geq 0}$ having all the properties derived in Lemma 2.17. In particular, $(v_s^N(t - s, x) dx)_{t \in [s, N)}$ is a weakly right-continuous and vaguely left-continuous solution to the linear Fokker-Planck equation (2.3) on $[s, N)$ with initial condition $u(s, x) dx = \mu_s$. Consequently, by $\mathbf{H}_{\leq \mathbf{u}}$ we have

$$\mathbb{P}_{\delta_s \otimes u(s, x) dx}^N \circ (X_2^N(t - s))^{-1} = v_s^N(t - s, x) dx = u(t, x) dx, \quad s \in S, t \in [s, N),$$

or equivalently, that

$$\mathbb{P}_{\delta_s \otimes \mu_s}^N \circ (X_2^N(t))^{-1} = v_s^N(t, x) dx = u(t + s, x) dx = \mu_{t+s}, \quad s \in S, t \in [0, N - s). \quad (2.29)$$

Furthermore, by (2.10), (ii.2), and integrating (2.16) with respect to $\bar{\mu}_N$, we get

$$\int_0^N \bar{\mathbb{P}}_{\delta_s \otimes \mu_s}^N (X_1^N(t) \neq t + s; t < \zeta^N) ds = 0,$$

which means that by passing to a subset of S , we may assume without loss of generality that

$$\bar{\mathbb{P}}_{\delta_s \otimes \mu_s}^N (X_1^N(t) \neq t + s; t < \zeta^N) = 0, \quad s \in S. \quad (2.30)$$

To prove that (2.29) and (2.30) from above hold for all $0 < s < N$, we use a bootstrap argument: Let $0 < s < N$ and $s' \in S$ such that $s' < s$. Then

$$\begin{aligned} \delta_s \otimes u(s, x) dx &= \delta_{s'+s-s'} \otimes u(s' + s - s', x) dx = \delta_{s'+s-s'} \otimes v_{s'}(s - s', x) dx \\ &= \mathbb{P}_{\delta_{s'} \otimes \mu_{s'}}^N \circ (X_1^N(s - s'), X_2^N(s - s'))^{-1}, \end{aligned}$$

where the last equality follows from (2.29) and (2.30). Hence, by the Markov property of X^N we get

$$\begin{aligned} \mathbb{P}_{\delta_s \otimes \mu_s}^N \circ (X_1^N(t), X_2^N(t))^{-1} &= \mathbb{P}_{\delta_{s'} \otimes \mu_{s'}}^N \circ (X_1^N(s-s'), X_2^N(s-s'))^{-1} \circ (X_1^N(t), X_2^N(t))^{-1} \\ &= \mathbb{P}_{\delta_{s'} \otimes \mu_{s'}}^N \circ (X_1^N(s - s' + t), X_2^N(s - s' + t))^{-1}, \end{aligned}$$

and thus for all $t < N - s$ we have that $t + s - s' < N - s'$, so by (2.29) we obtain

$$\begin{aligned} \mathbb{P}_{\delta_s \otimes u(s, x) dx}^N \circ (X_2^N(t))^{-1} &= \mathbb{P}_{\delta_{s'} \otimes u(s', x) dx}^N \circ (X_2^N(s - s' + t))^{-1} \\ &= u(s + t, x) dx, \end{aligned}$$

so (i)-(ii) are completely proved.

To prove (iii) we make use of (i)-(ii) to deduce that for every $0 < s < N < \infty$ we have

$$\mathbb{P}_{\delta_s \otimes \mu_s}^N \circ (X_1^N(t), X_2^N(t))^{-1} = \delta_{s+t} \otimes \mu_{t+s}, \quad 0 \leq t < N - s,$$

hence for $0 \leq t < N - s$

$$\mathbb{E}_{\delta_s \otimes \mu_s} \{ \mathbf{1}_{[0, N-s] \times \mathbb{R}^d} (X^N(t)); t < \zeta^N \} = \int_{[0, N-s] \times \mathbb{R}^d} \mathbf{1} d\mathbb{P}_{\delta_s \otimes \mu_s}^N \circ (X_1^N(t), X_2^N(t))^{-1} = 1,$$

in other words $\mathbb{P}_{\delta_s \otimes \mu_s} \{ t < \zeta^N \} = 1$, $0 \leq t < N - s$, or $\mathbb{P}_{\delta_s \otimes \mu_s} \{ N - s \leq \zeta^N \} = 1$. At the same time, by (2.27) we clearly have that $\mathbb{P}_{\delta_s \otimes \mu_s}^N \{ N - s < \zeta^N \} = 0$, since $\mathbb{P}_{\delta_s \otimes \mu_s}^N$ -a.s. on the event $N - s < \zeta^N$ we have $X_1^N(N - s) = N$ which entails the contradiction $X^N(N - s) \notin [0, N] \times \mathbb{R}^d$. In conclusion, we must have (2.28). \square

Lemma 2.20. *Assume that conditions $\mathbf{H}_{\mathbf{a}, \mathbf{b}}^{\sqrt{\mathbf{u}}}$ and $\mathbf{H}_{\leq \mathbf{u}}$ are fulfilled. Also, let $(\nu_t := v(t, x) dx)_{t \in [0, N]} \subset \mathcal{P}(\mathbb{R}^d)$ be a weakly continuous solution to the linear Fokker-Planck equation (2.3) on $[0, N]$ such that $v \in \mathcal{A}_{\leq \mathbf{u}}^{0, N}$. Then for every $0 < s < N < \infty$ we have the identification*

$$\mathbb{P}_{\delta_s \otimes \nu_s}^N \circ (X_2^N(t))^{-1} = \nu_{t+s}, \quad 0 \leq t < N - s. \quad (2.31)$$

Moreover, there exists a unique solution $\eta^N \in \mathcal{P}(C([0, N]; \mathbb{R}^d))$ to the canonical MP associated to $(\mathbf{L}_t)_{t \in [0, T]}$ such that for every $\varepsilon \in (0, N)$

$$\eta^N|_{C([\varepsilon, N]; \mathbb{R}^d)} = \mathbb{P}_{\delta_\varepsilon \otimes \nu_\varepsilon}^N \circ \left[(X_2^N(t - \varepsilon))_{t \in [\varepsilon, N]} \right]^{-1} \quad \text{as distributions on } C([\varepsilon, N]; \mathbb{R}^d).$$

Proof. First of all, since $\mathbf{H}_{\leq \mathbf{u}}$ is stronger than $\mathbf{H}_{\mathbf{a}, \mathbf{b}}$ we have that the conclusions of Lemma 2.19 are valid. In particular, using Lemma 2.19, (iii), applied to $v^N = v$, we have that the solutions $(v_s^N(t, x) dx)_{0 \leq t < N-s}$ for $s \in S$ provided by Lemma 2.17 are weakly right-continuous curves in $\mathcal{P}(\mathbb{R}^d)$. Thus, by Proposition 2.8, they are weakly continuous, and the uniqueness condition $\mathbf{H}_{\leq \mathbf{u}}$ can be used to prove (2.31) by following the same lines as in the proof of (i)-(ii) in Lemma 2.19.

Now, let us show the second part of the statement. To this end, let $\varepsilon \in (0, T)$, $f \in C_c^\infty([\varepsilon, N] \times \mathbb{R}^d)$, and define the process

$$M(t) := f(X^N(t)) - f(X^N(0)) - \int_0^t \bar{\mathbf{L}}f(X^N(r)) dr, \quad t \in [0, N - \varepsilon].$$

Recall that $h(\Delta) = 0$ for any function $h : [0, N] \times \mathbb{R}^d \rightarrow \mathbb{R}$, where Δ is an abstract cemetery points; see the paragraph on right processes from Section 6. Also, note that by Lemma 2.19 we have

$$\mathbb{P}_{\delta_\varepsilon \otimes \nu_\varepsilon}^N (\zeta^N = N - \varepsilon) = 1.$$

Furthermore, by Lemma 2.19, (ii), (2.31), and (2.12), we have

$$\mathbb{E}_{\delta_\varepsilon \otimes \nu_\varepsilon}^N \int_0^t |\bar{\Gamma}f| (X^N(r)) dr = \int_0^t \nu_{r+\varepsilon} \left(|\bar{\Gamma}f| (r + \varepsilon, \cdot) \right) dr \leq c \int_0^N \mu_r \left(|\bar{\Gamma}f| (r, \cdot) \right) dr < \infty, \quad t < N - \varepsilon.$$

Consequently, under $\mathbb{P}_{\delta_\varepsilon \otimes \nu_\varepsilon}^N$, the process $M(t)$, $\varepsilon \leq t \leq N$, is well-defined and integrable. Moreover, since X^N is a right process, it is a standard fact that M is a continuous additive functional, namely

$$M(t+s) = M(s) + M(t) \circ \theta^N(s), \quad 0 \leq s \leq N - \varepsilon - t, \quad \mathbb{P}_{\delta_\varepsilon \otimes \nu_\varepsilon}^N \text{-a.e.}$$

Consequently, $(M(t))_{t \in [0, N-\varepsilon]}$ is a $\mathbb{P}_{\delta_\varepsilon \otimes \nu_\varepsilon}^N$ -martingale as soon as we show that

$$\mathbb{E}_{\delta_\varepsilon \otimes \nu_\varepsilon}^N \{M(t)\} = 0, \quad t \in [0, N - \varepsilon].$$

But this can be easily verified since again by (2.31) and Lemma 2.19 we have

$$\begin{aligned} \mathbb{E}_{\delta_\varepsilon \otimes \nu_\varepsilon}^N \{M(t)\} &= \nu_{t+\varepsilon}(f(t + \varepsilon, \cdot)) - \nu_\varepsilon(f(\varepsilon, \cdot)) - \int_0^t \nu_{r+\varepsilon}(\bar{\Gamma}f(r + \varepsilon, \cdot)) dr \\ &= \nu_{t+\varepsilon}(f(t + \varepsilon, \cdot)) - \nu_\varepsilon(f(\varepsilon, \cdot)) - \int_\varepsilon^{t+\varepsilon} \nu_r(\bar{\Gamma}f(r, \cdot)) dr \\ &= 0, \quad t < N - \varepsilon, \end{aligned}$$

where the last equality follows from the fact that $(\nu_t)_{t \in [0, N]}$ is assumed to be a solution to the linear Fokker-Planck equation (2.3).

Now, let us notice that $(M(t))_{t \in [0, N-\varepsilon]}$ being a $\mathbb{P}_{\delta_\varepsilon \otimes \nu_\varepsilon}^N$ -martingale actually means that

$$f(t + \varepsilon, X_2^N(t)) - f(\varepsilon, X_2^N(0)) - \int_0^t \bar{\Gamma}f(r + \varepsilon, X_2^N(r)) dr, \quad t \in [0, N - \varepsilon],$$

is a $\mathbb{P}_{\delta_\varepsilon \otimes \nu_\varepsilon}^N$ -martingale, hence, by a change of variable, that

$$f(t, X_2^N(t - \varepsilon)) - f(\varepsilon, X_2^N(0)) - \int_\varepsilon^t \bar{\Gamma}f(r, X_2^N(r - \varepsilon)) dr, \quad t \in [\varepsilon, N],$$

is a $\mathbb{P}_{\delta_\varepsilon \otimes \nu_\varepsilon}^N$ -martingale. Consequently,

$$\mathbb{P}_{\delta_\varepsilon \otimes \nu_\varepsilon}^N \circ \left[(X_2^N(t - \varepsilon))_{t \in [\varepsilon, N]} \right]^{-1} \in \mathcal{P}(C([\varepsilon, N]; \mathbb{R}^d))$$

is a solution to the canonical MP associated to $(\mathbb{L}_t)_{t \in [\varepsilon, N]}$, with initial distribution ν_ε . Finally, the second part of Lemma 2.20 follows by Proposition 2.7. \square

2.4.2 Proof of Theorem 2.12

Let

$$X^N = \left(\Omega^N, \mathcal{F}^N, (\mathcal{F}_t^N)_{t \geq 0}, \mathbb{P}_{s,x}^N, (s, x) \in [0, N] \times \mathbb{R}^d \right) \text{ with lifetime } \zeta^N, \quad N > 0.$$

be given by Theorem 2.10. In order to construct the projective limit of the above family of right processes, the idea is to first construct the projective limit of the analytic objects $(\bar{P}_t^N, t \geq 0)_{N > 0}$ and $(\bar{U}^N)_{N > 0}$. One this first part is achieved, we proceed to the construction of the corresponding stochastic process which shall stand as the projective limit of $(X^N)_{N > 0}$. The proof is a bit involved, so let us split it in several steps.

Step OP1.1: We claim that for all $N > 0$ we have

$$\bar{P}_t^N((s, x), \cdot) = \bar{P}_t^{N+1}((s, x), \cdot) \Big|_{[0, N] \times \mathbb{R}^d} \quad \bar{\mu}_N\text{-a.e. on } [0, N] \times \mathbb{R}^d, t \geq 0,$$

and thus, by taking the Laplace transform with respect to the t variable,

$$\bar{U}_\alpha^N((s, x), \cdot) = \bar{U}_\alpha^{N+1}((s, x), \cdot) \Big|_{[0, N] \times \mathbb{R}^d} \quad \bar{\mu}_N\text{-a.e. on } [0, N] \times \mathbb{R}^d, \alpha > 0.$$

To prove the claim, let $h \in pb\mathcal{B}([0, \infty) \times \mathbb{R}^d)$, so that $hu \in \mathcal{A}_{\lesssim u}^{0, N}$, $N > 0$ for $c := \|h\|_\infty$. Further, set

$$v(t, x) := \frac{h(t, x)u(t, x)}{\int_{\mathbb{R}^d} h(t, x)u(t, x)dx}, \quad t \geq 0, x \in \mathbb{R}^d,$$

so that Lemma 2.17 as well as Lemma 2.19 can be applied to deduce: For each $N > 0$ there exists a set $S \subset (0, N + 1)$ of full Lebesgue measure such that for every $s \in S$ one can construct the family of probability measures $(v_s^{N+1}(t - s, x)dx)_{t \in [s, N+1]}$ which is a weakly continuous solution to the linear Fokker-Planck equation (2.3) on $[s, N + 1]$ with $v_s^{N+1}(0, x)dx = v(s, x)dx$. Thus, for $s \in S, t < N - s$, by Lemma 2.19 and then Lemma 2.20 we have

$$\begin{aligned} (\delta_s \otimes v(s, x)dx) \circ \bar{P}_t^N &= \mathbb{P}_{\delta_s \otimes v(s, x)dx}^N \circ (X_1^N(t), X_2^N(t))^{-1} \\ &= \delta_{t+s} \otimes \mathbb{P}_{\delta_s \otimes v(s, x)dx}^N \circ (X_2^N(t))^{-1} \\ &= \delta_{t+s} \otimes v_s^{N+1}(t - s, x)dx, \end{aligned}$$

as well as

$$\begin{aligned} (\delta_s \otimes v(s, x)dx) \circ \bar{P}_t^{N+1} &= \mathbb{P}_{\delta_s \otimes v(s, x)dx}^{N+1} \circ (X_1^{N+1}(t), X_2^{N+1}(t))^{-1} \\ &= \delta_{t+s} \otimes \mathbb{P}_{\delta_s \otimes v(s, x)dx}^{N+1} \circ (X_2^{N+1}(t))^{-1} \\ &= \delta_{t+s} \otimes v_s^{N+1}(t - s, x)dx. \end{aligned}$$

Consequently, for every $s \in S$ and $t < N - s$ we have

$$(\delta_s \otimes v(s, x)dx) \circ \bar{P}_t^N = (\delta_s \otimes v(s, x)dx) \circ \bar{P}_t^{N+1},$$

or equivalently,

$$\int_{\mathbb{R}^d} h(s, x)u(s, x)\bar{P}_t^N((s, x), \cdot) dx = \int_{\mathbb{R}^d} h(s, x)u(s, x)\bar{P}_t^{N+1}((s, x), \cdot) dx,$$

where the dot indicates that the equality is for the corresponding measures on $[0, N] \times \mathbb{R}^d$. As a matter of fact, the above equality holds for any $t \geq 0$ since by Lemma 2.19 it is easy to see that as measures on $[0, N] \times \mathbb{R}^d$,

$$(\delta_s \otimes v(s, x)dx) \circ \bar{P}_t^N = 0 = (\delta_s \otimes v(s, x)dx) \circ \bar{P}_t^{N+1}, \quad t \geq N - s.$$

Finally, integrating also with respect to s we obtain

$$\int_{[0, N] \times \mathbb{R}^d} h(s, x)\bar{P}_t^N((s, x), \cdot) \bar{\mu}_N(ds, dx) = \int_{[0, N] \times \mathbb{R}^d} h(s, x)\bar{P}_t^{N+1}((s, x), \cdot) \bar{\mu}_N(ds, dx).$$

The claim now follows since h was arbitrarily chosen.

Step OP1.2: For $\alpha > 0$, $0 \leq f \in L^p([0, \infty) \times \mathbb{R}^d; \bar{\mu})$, $1 \leq p \leq \infty$, set

$$\bar{U}_\alpha f := \sup_N \left(1_{[0, N] \times \mathbb{R}^d} \bar{U}_\alpha^N(f|_{[0, N] \times \mathbb{R}^d}) \right).$$

We claim that $\bar{U} := (\bar{U}_\alpha)_{\alpha>0}$ is a C_0 -resolvent of Markov operators on $L^p([0, \infty) \times \mathbb{R}^d; \bar{\mu})$, $1 \leq p < \infty$, and $\bar{\mu}$ is sub-invariant with respect to \bar{U} .

To prove this claim, let $0 \leq f \in L^p([0, \infty) \times \mathbb{R}^d; \bar{\mu})$, $1 \leq p < \infty$ and note first that by **Step Q1.1** from above we have

$$1_{[0, N) \times \mathbb{R}^d} \bar{U}_\alpha^N (f|_{[0, N) \times \mathbb{R}^d}) \text{ is increasing with respect to } N > 0.$$

Therefore, by monotone convergence

$$\bar{\mu} \left((\alpha \bar{U}_\alpha f)^p \right) = \sup_N \bar{\mu}_N \left(\left(\alpha \bar{U}_\alpha^N (f|_{[0, N) \times \mathbb{R}^d}) \right)^p \right) \leq \sup_N \bar{\mu}_N \left((f|_{[0, N) \times \mathbb{R}^d})^p \right) = \bar{\mu}(f^p), \quad \alpha > 0.$$

Furthermore, since by Lemma 2.19, (iii) we get $\bar{P}_t^N 1(s, x) = 1_{[0, N)}(s+t) \bar{\mu}_N$ -a.e., $t \geq 0$, we deduce that $\alpha \bar{U}_\alpha 1 = 1$. Also, $\alpha \bar{U}_\alpha$ is clearly positivity preserving, consequently it is a Markov operator on $L^p(\bar{\mu})$ for every $1 \leq p \leq \infty$, and it is straightforward to check that $(\bar{U}_\alpha)_\alpha$ satisfies the resolvent equation. To conclude this step it remains to show that

$$\lim_{\alpha \rightarrow \infty} \alpha \bar{U}_\alpha = \text{Id} \quad \text{in } L^p(\bar{\mu}),$$

which by a density argument boils down to testing the above convergence on smooth functions f with compact support in $[0, N) \times \mathbb{R}^d$, for some arbitrarily fixed $N > 0$. But for such f , by **Step Q1.1** and Lemma 2.19, (ii), (iii), we have

$$\alpha \bar{U}_\alpha f = \alpha \bar{U}_\alpha^N (f|_{[0, N) \times \mathbb{R}^d}) 1_{[0, N) \times \mathbb{R}^d}, \quad (2.32)$$

which converges to f in $L^p(\bar{\mu})$ as $\alpha \rightarrow \infty$ by Theorem 2.10.

Step OP1.3: Proof of Theorem 2.12, (i) Let us denote by $(\bar{L}_p, D(\bar{L}_p))$ the infinitesimal generators of the C_0 -resolvent \bar{U} on $L^p([0, \infty) \times \mathbb{R}^d; \bar{\mu})$, that is

$$D(\bar{L}_p) := \bar{U}_1 (L^p([0, \infty) \times \mathbb{R}^d; \bar{\mu})), \quad \bar{L}_p(\bar{U}_1 f) := \bar{U}_1 f - f, \quad f \in L^p([0, \infty) \times \mathbb{R}^d; \bar{\mu}), 1 \leq p < \infty.$$

Since $\alpha \bar{U}_\alpha$ is a contraction on $L^p([0, \infty) \times \mathbb{R}^d; \bar{\mu})$ for every $\alpha > 0$, we have that $(\bar{L}_p, D(\bar{L}_p))$ is a maximal dissipative operator on $L^p([0, \infty) \times \mathbb{R}^d; \bar{\mu})$, $1 \leq p < \infty$. Let us show that $(\bar{L}_1, D(\bar{L}_1))$ is an extension of $(\bar{L}, C_c^\infty((0, \infty) \times \mathbb{R}^d))$ on $L^1([0, \infty) \times \mathbb{R}^d; \bar{\mu})$: Let $f \in C_c^\infty((0, \infty) \times \mathbb{R}^d)$ so that $f \in C_c^\infty((0, N) \times \mathbb{R}^d)$ for some $N > 0$ that we fix. Thus, by Theorem 2.10, $f \in D(\bar{L}_1^N)$ and there exists $f_0 \in L^1([0, N) \times \mathbb{R}^d; \bar{\mu}_N)$ such that

$$f = \bar{U}_1^N f_0 \quad \text{and} \quad \bar{L} f = \bar{U}_1^N f_0 - f_0 \quad \bar{\mu}_N\text{-a.e.}$$

Extending f_0 by zero on $[N, \infty) \times \mathbb{R}^d$ and using (2.32), we get

$$\bar{L} f = \bar{U}_1^N f_0 - f_0 = \bar{U}_1 f_0 - f_0 = \bar{L}_1 \bar{U}_1 f_0 = \bar{L}_1 f \quad \bar{\mu}\text{-a.e.},$$

hence $(\bar{L}_1, D(\bar{L}_1))$ is an extension of $(\bar{L}, C_c^\infty((0, \infty) \times \mathbb{R}^d))$ on $L^1([0, \infty) \times \mathbb{R}^d; \bar{\mu})$.

Step OP1.4: We claim that there exists a resolvent of Markov kernels $\bar{U}^0 := (\bar{U}_\alpha^0)_{\alpha>0}$ on $[0, \infty) \times \mathbb{R}^d$ such that $\mathcal{E}(\mathcal{U}_\beta^0)$ is min stable, contains the constant functions, and generates $\mathcal{B}([0, \infty) \times \mathbb{R}^d)$ for one (hence all) $\beta > 0$; that is, the equivalent statements in Proposition 6.8 from Elements of right processes hold.

This claim follows mostly from [17]. More precisely, by [17, Theorem 2.2 and Remark 2.3], we can directly deduce that a resolvent of sub-Markov kernels \tilde{U} exists having all the other claimed properties satisfied. Thus, we only need to prove the Markov property pointwise on $[0, \infty) \times \mathbb{R}^d$. To this end we use Lemma 6.53 which is also taken from [17], as follows: Let

$$\begin{aligned} E_0 &:= \left\{ x \in [0, \infty) \times \mathbb{R}^d : \alpha \tilde{U}_\alpha 1(x) = 1, \alpha \in \mathbb{Q}_+ \right\} \\ &= \left\{ x \in [0, \infty) \times \mathbb{R}^d : \alpha \tilde{U}_\alpha 1(x) = 1, \alpha > 0 \right\}, \end{aligned}$$

where the (second) equality holds due to the resolvent equation and the fact that \tilde{U} is already sub-Markovian. Since $\alpha\tilde{U}_\alpha 1 = 1$ $\bar{\mu}$ -a.e. for every $\alpha > 0$, we get that $\bar{\mu}([0, \infty) \times \mathbb{R}^d \setminus E_0) = 0$. But now we can apply Lemma 6.53 simultaneously for the kernels $\alpha\tilde{U}_\alpha$, $0 < \alpha \in \mathbb{Q}_+$ to find $F \in \mathcal{B}([0, \infty) \times \mathbb{R}^d)$ such that $F \subset E_0$,

$$\bar{\mu}([0, \infty) \times \mathbb{R}^d \setminus F) = 0, \quad \text{and} \quad \tilde{U}_\alpha 1_{[0, \infty) \times \mathbb{R}^d \setminus F} = 0 \text{ on } F, \quad 0 < \alpha \in \mathbb{Q}_+, \text{ hence for all } \alpha > 0.$$

Therefore, by Definition 6.30 and Definition 6.29 from Elements of right processes, we can consider the resolvent of Markovian kernels $\bar{U}^0 := (\bar{U}_\alpha^0)_{\alpha > 0}$ on $[0, \infty) \times \mathbb{R}^d$ obtained by restricting first \tilde{U} to F , and then by taking its trivial extension to $[0, \infty) \times \mathbb{R}^d$. \bar{U}^0 enjoys all the properties required in the claim.

Step OP1.5: *There exists $\mathcal{N} \in \mathcal{B}([0, \infty) \times \mathbb{R}^d)$ with*

$$\bar{\mu}(\mathcal{N}) = 0 \quad \text{and} \quad \bar{U}_\alpha^0 1_{\mathcal{N}} = \bar{U}_\alpha^N 1_{\mathcal{N} \cap ([0, N] \times \mathbb{R}^d)} = 0 \text{ on } ([0, N] \times \mathbb{R}^d) \setminus \mathcal{N} \text{ for all } \alpha, N > 0,$$

such that if we set

$$E := ([0, \infty) \times \mathbb{R}^d) \setminus \mathcal{N}, \quad E^N := ([0, N] \times \mathbb{R}^d) \setminus \mathcal{N}, \quad N > 0,$$

then for every $N > 0$ and $f \in b\mathcal{B}([0, \infty))$ we have

$$\bar{U}_\alpha^0((s, x), \cdot)|_{E^N} = \bar{U}_\alpha^N((s, x), \cdot), (s, x) \in E^N, \quad \text{and} \quad \bar{U}_\alpha^0 f(s, x) = \int_0^\infty e^{-\alpha t} f(s+t) dt, (s, x) \in E. \quad (2.33)$$

To prove the above assertion, let $(f_k)_{k \geq 1} \in b\mathcal{B}([0, \infty) \times \mathbb{R}^d)$ and $(g_k)_{k \geq 1} \in b\mathcal{B}([0, \infty) \times \mathbb{R}^d)$ be such that the two families of functions separate the space of finite measures on the corresponding spaces. Thus, using Lemma 2.19,

$$\bar{U}_\alpha^0(f_k 1_{[0, N] \times \mathbb{R}^d}) = \bar{U}_\alpha^N(f_k|_{[0, N] \times \mathbb{R}^d}) 1_{[0, N] \times \mathbb{R}^d} \quad \bar{\mu}\text{-a.e. for all } \alpha, N > 0, k \geq 1,$$

so if we set

$$\mathcal{N}^0 := \bigcup_{k \geq 1, N, \alpha \in \mathbb{Q}_+^*} \left[\bar{U}_\alpha^0(f_k 1_{[0, N] \times \mathbb{R}^d}) \neq \bar{U}_\alpha^N(f_k|_{[0, N] \times \mathbb{R}^d}) 1_{[0, N] \times \mathbb{R}^d} \right],$$

then $\bar{\mu}(\mathcal{N}^0) = 0$; throughout, for $f \in \mathcal{B}([0, N] \times \mathbb{R}^d)$, by $f 1_{[0, N] \times \mathbb{R}^d}$ we denote the extension of f to $[0, \infty) \times \mathbb{R}^d$ by setting it to be zero on $[N, \infty) \times \mathbb{R}^d$. By the same Lemma 2.19, we get $\bar{\mu}$ -a.e. for all $\alpha > 0, k \geq 1$ that

$$\begin{aligned} \bar{U}_\alpha^0(g_k) &= \sup_N \bar{U}_\alpha^N(g_k|_{[0, N]}) 1_{[0, N] \times \mathbb{R}^d} = \sup_N \int_0^\infty e^{-\alpha t} g_k|_{[0, N]}(s+t) dt 1_{[0, N] \times \mathbb{R}^d} \\ &= \int_0^\infty e^{-\alpha t} g_k(s+t) dt. \end{aligned}$$

Thus, if we set

$$\mathcal{N}^{00} := \bigcup_{k \geq 1, \alpha \in \mathbb{Q}_+^*} \left[\bar{U}_\alpha^0 g_k(s, x) \neq \int_0^\infty e^{-\alpha t} g_k(s+t) dt \right],$$

then $\bar{\mu}(\mathcal{N}^{00}) = 0$. Now we set $\mathcal{N}' := \mathcal{N}^0 \cup \mathcal{N}^{00}$ which satisfies $\bar{\mu}(\mathcal{N}') = 0$, so we can apply Lemma 6.53 simultaneously for the kernels $\bar{U}_\alpha^0, \alpha \in \mathbb{Q}_+^*$ to find $\mathcal{N} \in \mathcal{B}([0, \infty) \times \mathbb{R}^d)$ such that

$$\mathcal{N}' \subset \mathcal{N} \quad \text{and} \quad \bar{U}_\alpha^0 1_{\mathcal{N}} = 0 \text{ on } E \text{ for all } \alpha \in \mathbb{Q}_+^*, \text{ hence for all } \alpha > 0.$$

Moreover, for every $(s, x) \in E^N, k \geq 1, \alpha \in \mathbb{Q}_+^*$ we have

$$\bar{U}_\alpha^0(f_k 1_{[0, N] \times \mathbb{R}^d})(s, x) = \bar{U}_\alpha^N(f_k|_{[0, N] \times \mathbb{R}^d})(s, x) 1_{[0, N] \times \mathbb{R}^d}(s, x),$$

and since $f_k, k \geq 1$ is measure separating the above equality holds with f_k replaced by any $f \in b\mathcal{B}([0, \infty) \times \mathbb{R}^d)$. Furthermore, by the resolvent equation, it is easy to see that the equality holds for all $\alpha > 0$. Thus, we obtain

$$\bar{U}_\alpha^0((s, x), \cdot)|_{E^N} = \bar{U}_\alpha^N((s, x), \cdot), \quad (s, x) \in E^N, \alpha > 0.$$

By a similar reasoning we also have

$$\bar{U}_\alpha^0 f(s, x) = \int_0^\infty e^{-\alpha t} f(s+t) dt, \quad f \in b\mathcal{B}([0, \infty)), (s, x) \in E, \alpha > 0.$$

Step OP1.6: Keeping the notations from above, let $\tilde{U}^0 := (\tilde{U}_\alpha^0)_{\alpha>0}$ and $\tilde{U}^N := (\tilde{U}_\alpha^N)_{\alpha>0}$ be the restrictions of \bar{U}^0 and \bar{U}^N to E and E^N , respectively, given by Definition 6.30. Furthermore, set

$$V : E \rightarrow [0, 1], \quad V(s, x) = e^{-s}, (s, x) \in E.$$

Then V is \tilde{U}^0 -excessive.

Indeed, for all $(s, x) \in E$, by (2.33) we have

$$\alpha \tilde{U}_\alpha^0 V(s, x) = \alpha \int_0^\infty e^{-\alpha t} e^{-(s+t)} dt = \frac{\alpha}{\alpha+1} e^{-s} \leq e^{-s}, \quad \alpha > 0,$$

and we clearly also have $\lim_{\alpha \rightarrow \infty} \alpha \tilde{U}_\alpha^0 V(s, x) = V(s, x)$.

Step OP1.7: Let (E_1, \tilde{U}^1) be the saturation of (E, \tilde{U}^0) (with respect to \tilde{U}_β^0 for some $\beta > 0$) given by Definition 6.25. Consequently, by Remark 6.26 the function V has a unique extension V_1 by fine continuity from E to E_1 ; V_1 is \tilde{U}^1 -excessive. Furthermore, by Theorem 6.27, there exists a right process

$$\tilde{X} = (\tilde{\Omega}, \tilde{\mathcal{F}}, \tilde{\mathcal{F}}_t, \tilde{X}(t), \theta(t), \tilde{\mathbb{P}}_z, z \in E_1) \quad (2.34)$$

on E_1 with resolvent \tilde{U}^1 .

Let $\tilde{U}^{1,N}$ be the resolvent of the process \tilde{X} killed upon leaving the finely open set $[V_1 > e^{-N}]$, $N > 0$. That is

$$\tilde{U}^{1,N} = \tilde{U}^1 - B_\alpha^{[V_1 \leq e^{-N}]} \tilde{U}^1 \quad \text{on } [V_1 > e^{-N}], \alpha > 0, \quad (2.35)$$

where recall that the balayage operator B_α^A is defined in Definition 6.17; see also (6.2).

We claim that for every $N > 0$, the resolvent $\tilde{U}^{1,N}$ is an extension of \tilde{U}^N from E^N to $[V_1 > e^{-N}]$, in the sense of Definition 6.49.

To prove the claim, we need to check conditions (i.1) – (i.2) from Definition 6.49, namely that

$$\tilde{U}_\alpha^{1,N}(z, [V_1 > e^{-N}] \setminus E^N) = 0, \quad z \in [V_1 > e^{-N}], \text{ and that}$$

$$\tilde{U}_\alpha^{1,N}((s, x), \cdot)|_{E^N} = \tilde{U}_\alpha^N((s, x), \cdot) \text{ as measures on } E^N, \text{ for every } (s, x) \in E^N \text{ and } N > 0.$$

To check the first condition, note that

$$E \cap ([V_1 > e^{-N}] \setminus E^N) = \mathcal{N} \cap [0, N] \times \mathbb{R}^d, \quad \text{where } \mathcal{N} \text{ is given in Step OP1.5.}$$

Thus, $\tilde{U}_\alpha^{1,N}((s, x), E \cap ([V_1 > e^{-N}] \setminus E^N)) \leq \tilde{U}_\alpha^1((s, x), \mathcal{N}) = 0$, $(s, x) \in E^N$, hence

$$\tilde{U}_\alpha^{1,N}([V_1 > e^{-N}] \setminus E^N) \leq \tilde{U}_\alpha^{1,N}(E_1 \setminus E) \leq \tilde{U}_\alpha^1(E_1 \setminus E) = 0.$$

Checking the second condition is a little bit more involved: Let $N > 0$ and $(s, x) \in E^N$, so that $s < N$. Note that

$$\tilde{U}_\alpha^1((s, x), \cdot)|_{E^N} = \tilde{U}_\alpha^0((s, x), \cdot)|_{E^N} = \tilde{U}_\alpha^N((s, x), \cdot), \quad \alpha > 0,$$

where the last equality has been obtained in (2.27). Moreover, since

$$\tilde{U}_\alpha^{1,N} = \tilde{U}_\alpha^1 - B_\alpha^{[V_1 \leq e^{-N}]} \tilde{U}_\alpha^1,$$

it is sufficient to show that

$$B_\alpha^{[V_1 \leq e^{-N}]} \tilde{U}_\alpha^1((s, x), \cdot)|_{E^N} = 0, \quad N > 0. \quad (2.36)$$

To this end, let $F \in b\mathcal{B}(E_N)$ with compact support in E_N , in particular

$$\text{supp}(F) \subset [0, N - \varepsilon) \times \mathbb{R}^d \quad \text{for some } \varepsilon > 0.$$

We claim that

$$\tilde{U}_\alpha^1 F = 0 \quad \text{on } [V_1 \leq e^{-N}], \quad N > 0, \quad (2.37)$$

hence $B_\alpha^{[V_1 \leq e^{-N}]} \tilde{U}_\alpha^1 F(s, x) \leq \tilde{U}_\alpha^1 F(s, x) = 0$; from this, by a straightforward density argument, we get (2.36). To prove (2.37), let $z \in [V_1 \leq e^{-N}]$, so that by the fine density of E in E_1 , there exists a sequence $(z_k = (s_k, x_k))_{k \geq 1} \subset E$ such that

$$e^{-s_k} = V(s_k, x_k) = V_1(s_k, x_k) \rightarrow_k V_1(z) \leq e^{-N}.$$

Also,

$$\tilde{U}_\alpha^1 F(z) = \lim_k \tilde{U}_\alpha^1 F(z_k) = \lim_k \tilde{U}_\alpha^0 F(z_k).$$

Thus, if we set

$$f(s) := \sup_{x \in \mathbb{R}^d} F(s, x), \quad s \in [0, N),$$

and take into account that $\lim_k s_k \geq N$ and that $f(r) = 0, r \geq N - \varepsilon$, we get

$$\tilde{U}_\alpha^1 F(z) \leq \limsup_k \tilde{U}_\alpha^0 f(s_k, x_k) = \limsup_k \int_0^\infty e^{-\alpha t} f(s_k + t) dt = 0,$$

which concludes the claim.

Step Q1.8: Let $(\tilde{U}^{N,1}, E^{N,1})$ be the saturation of (\tilde{U}^N, E^N) , $N \geq 1$. Then, by **Step OP1.7** and Remark 6.52, we have that $(\tilde{U}^{N,1}, E^{N,1})$ is an extension of $(\tilde{U}^{1,N}, [V_1 > e^{-N}])$, and that the latter is an extension of (\tilde{U}^N, E^N) , in the sense of Definition 6.49. Or, using the notation introduced in Definition 6.49, we have

$$(\tilde{U}^N, E^N) \subset (\tilde{U}^{1,N}, [V_1 > e^{-N}]) \subset (\tilde{U}^{N,1}, E^{N,1}). \quad (2.38)$$

Further, set

$$\hat{E}^N := \left\{ (s, x) \in [0, N) \times \mathbb{R}^d : \bar{U}_0^N((s, x), [0, N) \times \mathbb{R}^d \setminus E^N) = 0 \right\}.$$

Note that $\hat{E}^N \supset E^N$ and $\hat{E}^N \in \mathcal{A}(\bar{U}^N)$, $N \geq 1$; see Section 6.3.2. Let

$$\hat{U}^N = (\hat{U}_\alpha^N)_{\alpha > 0} := \bar{U}^N|_{\hat{E}^N}, \quad N \geq 1,$$

the restriction of \bar{U}^N to \hat{E}^N given by Definition 6.30. Thus, we also have

$$(\tilde{U}^N, E^N) \subset (\hat{U}^N, \hat{E}^N) \subset (\tilde{U}^{N,1}, E^{N,1}). \quad (2.39)$$

Further, set

$$\check{E}^N := \hat{E}^N \cap [V_1 > e^{-N}] \quad \text{in } E^{N,1}, \quad \check{U}^N := \tilde{U}^{N,1}|_{\check{E}^N}, \quad N > 0.$$

Thus, we now have

$$(\tilde{U}^N, E^N) \subset (\check{U}^N, \check{E}^N) \subset (\tilde{U}^{1,N}, [V_1 > e^{-N}]) \subset (\tilde{U}^{N,1}, E^{N,1}), \quad (2.40)$$

$$(\tilde{U}^N, E^N) \subset (\check{U}^N, \check{E}^N) \subset (\hat{U}^N, \hat{E}^N) \subset (\tilde{U}^{N,1}, E^{N,1}). \quad (2.41)$$

Moreover, note that $E^{N,1} \setminus [V_1 > e^{-N}]$ and $E^{N,1} \setminus \hat{E}^N$ are polar sets in $E^{N,1}$; this follows by Theorem 6.27, since both resolvents \hat{U}^N and $\check{U}^{1,N}$ are the resolvents of two right processes on the corresponding spaces. Consequently, for every $N \geq 1$,

$$E^{N,1} \setminus \check{E}^N \text{ is polar for } \tilde{U}^{N,1}, \quad [V_1 > e^{-N}] \setminus \hat{E}^N \text{ is polar for } \tilde{U}^{1,N}, \quad \text{and} \quad \hat{E}^N \setminus \check{E}^N \text{ is polar for } \hat{U}^N.$$

Step OP1.9: Recall that X^N is a right process associated to \bar{U}^N on $[0, N) \times \mathbb{R}^d$, as in Theorem 2.10, Lemma 2.19, and Lemma 2.20. Also, since $\check{E}^N \subset [0, N) \times \mathbb{R}^d$ and $\check{E}^N \in \mathcal{A}(\bar{U}^N)$, we can consider

$$\check{X}^N := X^N|_{\check{E}^N}, \quad (2.42)$$

the restriction of X^N from $[0, N) \times \mathbb{R}^d$ to \check{E}^N , which is a right process on \check{E}^N with lifetime ζ^N and (norm-)continuous paths, as it follows from Proposition 6.35 and Definition 6.36. Clearly, the resolvent of \check{X}^N is \check{U}^N , $N \geq 1$.

Also, let $\check{X}^{1,N}$ be the right process on \check{E}^N obtained by first killing the right process \check{X} upon leaving $[V_1 > e^{-N}]$ and then taking the restriction to \check{E}^N , $N \geq 1$; this is done following the procedures in Section 6.3.2. Thus, the resolvent of $\check{X}^{1,N}$ is precisely $\check{U}^{1,N}|_{\check{E}^N} = \hat{U}^N$, $N \geq 1$. Now, since the norm topology on \check{E}^N is a natural topology with respect to \check{U}^N , the latter being the resolvent of both processes \check{X}^N and $\check{X}^{1,N}$, by Proposition 6.16 we get

$$\check{X}^{1,N} \text{ and } \check{X}^N \text{ have the same laws on the path space } C([0, \infty); \check{E}^N), \quad (2.43)$$

with \check{E}^N endowed with the norm-topology, and for all starting points $(s, x) \in \check{E}^N$, $N \geq 1$.

Consequently, relying also on Lemma 2.19, *the right process \check{X} given by (2.34) satisfies for all $N > 0$*

$$\tilde{\mathbb{P}}_{\delta_s \otimes \mu_s}(T_{[V_1 \leq e^{-N}]} = N - s) = \mathbb{P}_{\delta_s \otimes \mu_s}(\zeta^N = N - s) = 1, \quad s \in (0, N), \quad (2.44)$$

$$\tilde{\mathbb{P}}_{\delta_s \otimes \mu_s}([0, N - s) \ni t \mapsto \check{X}(t) \ni \check{E}^N \text{ is well defined and norm-continuous}) = 1, \quad s \in (0, N). \quad (2.45)$$

$$\tilde{\mathbb{P}}_{(s,x)} \circ (\check{X}(t))^{-1} = \mathbb{P}_{(s,x)} \circ (X^N(t))^{-1}, \quad s \in (0, N), t \in (0, N - s), \mu_s\text{-a.e. } (s, x) \in \check{E}^N. \quad (2.46)$$

Step OP1.10: Define

$$\check{E} := \bigcup_{N \geq 1} \check{E}^N, \quad \text{so that we clearly have } \bar{\mu}([0, \infty) \times \mathbb{R}^d \setminus \check{E}) = 0.$$

By **Step OP1.9**, we have

$$\tilde{\mathbb{P}}_{\bar{\mu}}\left(\sup_N T_{[V_1 \leq e^{-N}]} = \infty\right) = 1, \quad (2.47)$$

$$\tilde{\mathbb{P}}_{\bar{\mu}}([0, \infty) \ni t \mapsto \check{X}(t) \ni \check{E} \text{ is well defined and norm-continuous}) = 1. \quad (2.48)$$

In particular, if we extend $\bar{\mu}$ from E to E_1 by setting $\bar{\mu}(E_1 \setminus E) := 0$, then $E_1 \setminus \check{E}$ is $\bar{\mu}$ -polar and $\bar{\mu}$ -negligible (with respect to \mathcal{U}^1), hence by Remark 6.22, (iii), there exists a further $\mathcal{B}([0, \infty) \times \mathbb{R}^d)$ -measurable set $\check{E}_0 \subset \check{E}$ such that $E_1 \setminus \check{E}_0$ is $\bar{\mu}$ -inessential (see Definition 6.18).

By Proposition 6.35, \check{X} can be restricted to \check{E}_0 , the restricted process being denoted by \check{X}^0 .

Furthermore, note that each \check{E}^N is finely open in E_1 , since it is finely open in $[V_1 > e^{-N}]$, and the latter is finely open in E_1 . Consequently, by (2.43), if $f \in C_b([0, \infty) \times \mathbb{R}^d)$ then

$$[0, \infty) \ni t \mapsto f(\check{X}^0(t)) \in \mathbb{R} \text{ is (right) continuous at } t = 0, \quad \tilde{\mathbb{P}}_{(s,x)}\text{-a.s. for all } (s, x) \in \check{E}_0,$$

hence by Remark 6.14 we easily deduce that the norm topology on \check{E}^0 is a natural topology with respect to the restriction $\mathcal{U}^1|_{\check{E}^0}$. We can now apply (2.48) and Proposition 6.37 to deduce that there exists another subset $E^\dagger \subset \check{E}_0$ set with $\check{E}_0 \setminus E^\dagger$ $\bar{\mu}$ -inessential such that if we restrict the process \check{X}^0 to E^\dagger , then it has $\tilde{\mathbb{P}}_{(s,x)}$ -a.s. norm-continuous paths in E^\dagger for every $(s, x) \in E^\dagger$. We denote this last right process with by

$$X^\dagger = \left(\mathbb{P}_{(s,x)}^\dagger, (s, x) \in E^\dagger, X^\dagger(t) = (X_1^\dagger(t), X_2^\dagger(t), t \geq 0)\right), \quad \text{with resolvent } \mathcal{U}^\dagger. \quad (2.49)$$

Note that X^\dagger has $\mathbb{P}_{(s,x)}^\dagger$ -a.s. norm-continuous paths for every $(s, x) \in E^\dagger$, and it is conservative since its resolvent \mathcal{U}^\dagger is the restriction of $\tilde{\mathcal{U}}^1$ from E^1 to E^\dagger , hence

$$\alpha U_\alpha^\dagger 1(x) = \alpha \tilde{U}_\alpha^1 1_{E^\dagger}(x) = \alpha \tilde{U}_\alpha^1 1(x) = 1, \quad x \in E^\dagger. \quad (2.50)$$

Step OP1.11: *The right process X^\dagger satisfies*

$$\mathbb{P}_{(s,x)}^\dagger \left(X_1^\dagger(t) = t + s, t \geq 0 \right) = 1, \quad (s, x) \in E^\dagger. \quad (2.51)$$

Indeed, we can argue as follows: First, note that by (2.33) and the fact that $\tilde{\mathcal{U}}^0$ is the restriction of $\bar{\mathcal{U}}^0$ from $[0, \infty) \times \mathbb{R}^d$ to E , we have

$$\tilde{\mathcal{U}}_\alpha^0 f|_E(s, x) = \int_0^\infty e^{-\alpha t} f(s+t) dt, \quad (s, x) \in E, f \in b\mathcal{B}([0, \infty)).$$

Since E is finely dense in E_1 , it is straightforward that, for all $(s, x) \in E^1$,

$$\tilde{\mathcal{U}}_\alpha^1 \tilde{f}(s, x) = \int_0^\infty e^{-\alpha t} f(s+t) dt, \quad f \in b\mathcal{B}([0, \infty)),$$

where \tilde{f} is any bounded and \mathcal{B}_1 -measurable extension of $f|_E$ to E_1 . Consequently, since \mathcal{U}^\dagger is the restriction of \mathcal{U}^1 to E^\dagger , and $E^\dagger \in \mathcal{B}([0, \infty) \times \mathbb{R}^d)$, we get

$$U_\alpha^\dagger f|_{E^\dagger}(s, x) = \int_0^\infty e^{-\alpha t} f(s+t) dt, \quad f \in b\mathcal{B}([0, \infty)), (s, x) \in E^\dagger.$$

It is now straightforward to deduce (2.51).

Step OP1.12: *For every $N \geq 1$, $s \in (0, N)$ and $t \in [0, N - s)$ we have the identification*

$$\mathbb{P}_{(s,x)}^\dagger \circ (X^\dagger(t))^{-1} = \mathbb{P}_{(s,x)}^N \circ (X^N(t))^{-1}, \quad \mu_s\text{-a.e. } (s, x) \in E^\dagger \cap \check{E}^N. \quad (2.52)$$

Indeed, note first that by the previous steps,

$$E^\dagger \subset \check{E} = \bigcup_{N \geq 1} \check{E}^N \subset [0, \infty) \times \mathbb{R}^d, \quad E^\dagger \in \mathcal{A}(\mathcal{U}^1), \quad X^\dagger = \tilde{X}|_{E^\dagger}.$$

Hence, using also (2.46) we get the identification

$$\mathbb{P}_{(s,x)}^\dagger \circ (X^\dagger(t))^{-1} = \tilde{\mathbb{P}}_{(s,x)} \circ (\tilde{X}(t))^{-1} = \mathbb{P}_{(s,x)}^N \circ (X^N(t))^{-1}, \quad \mu_s\text{-a.e. } (s, x) \in E^\dagger \cap \check{E}^N.$$

Step OP1.13: *We have $\bar{\mu}([0, \infty) \times \mathbb{R}^d \setminus E^\dagger) = 0$. Furthermore, for every $s > 0$ let us set*

$$E^\dagger(s) := \{x \in \mathbb{R}^d : (s, x) \in E^\dagger\}. \quad (2.53)$$

Then

$$\mu_s(E^\dagger(s)) = 1 \text{ for all } s \in (0, \infty) \quad \text{and} \quad \mathbb{P}_{\delta_s \otimes \mu_s}^\dagger \circ (X^\dagger(t))^{-1} = \delta_{t+s} \otimes \mu_{t+s} \quad \text{for every } s > 0 \text{ and } t \geq 0. \quad (2.54)$$

Indeed, using **Step OP1.10**,

$$\bar{\mu}([0, \infty) \times \mathbb{R}^d \setminus E^\dagger) = \bar{\mu}([0, \infty) \times \mathbb{R}^d \setminus \check{E}) = 0.$$

Moreover, from (2.26) and (2.52) we obtain that

$$\mathbb{P}_{\delta_s \otimes \mu_s}^\dagger \circ (X^\dagger(t))^{-1} = \delta_{t+s} \otimes \mu_{t+s} \quad \text{for every } s > 0, t \geq 0.$$

Step OP1.14: By construction, it is clear that $\bar{\mu}$ is sub-invariant for the transition function $(P_t^\dagger)_{t \geq 0}$ of the process X^\dagger , hence $(P_t^\dagger)_{t \geq 0}$ can be extended to a C_0 -semigroup $(T_t^\dagger)_{t \geq 0}$ on every $L^p([0, \infty) \times \mathbb{R}^d; \bar{\mu})$, $1 \leq p < \infty$. Moreover, the resolvent of $(T_t^\dagger)_{t \geq 0}$ is precisely $\bar{\mathcal{U}}$ constructed at **Step OP1.2**. Consequently,

$$P_t^\dagger f = f + \int_0^t P_s^\dagger \bar{\Gamma} f ds \quad \bar{\mu}\text{-a.e. for all } t \geq 0 \text{ and } f \in C_c^\infty([0, \infty) \times \mathbb{R}^d).$$

Step OP1.15: Let $s > 0$ and $\mathcal{P}(\mathbb{R}^d) \ni \nu_s \leq c\mu_s$ for some constant $c \in (0, \infty)$, and set

$$\nu_{t+s} := \mathbb{P}_{\delta_s \otimes \nu_s}^\dagger \circ (X_2^\dagger(t))^{-1}, \quad t \geq 0.$$

We claim that for all $N > s$, the family $(\nu_{t+s})_{t \in [0, N-s]}$ is the unique weakly continuous solution to the LFPE (2.3) on $[s, N)$ which consists of probability measures, belongs to $\mathcal{A}_{\lesssim u}^{s, N}$, and whose initial condition (at time s) is ν_s .

Indeed, the fact that $(\nu_{t+s})_{t \in [0, N-s]}$ is a weakly continuous curve in $\mathcal{P}(\mathbb{R}^d)$ is now clear since $(X_2^\dagger(t))_{t \geq 0}$ is conservative and has a.s. continuous paths.

The fact that $(\nu_{t+s})_{t \in [0, N-s]}$ is a solution to the LFPE (2.3) on $[s, N)$ follows e.g. similarly to the proof of Lemma 2.17, (iii). Finally,

$$\nu_{t+s} = \int_{\mathbb{R}^d} \mathbb{P}_{s,x}^\dagger \circ (X_2^\dagger(t))^{-1} \frac{d\nu_s}{d\mu_s}(dx) \leq c \mathbb{P}_{\delta_s \otimes \nu_s}^\dagger \circ (X_2^\dagger(t))^{-1} = c\mu_{t+s}, \quad t \geq 0.$$

Consequently, the uniqueness part also follows since condition $\mathbf{H}_{\lesssim u}$ is in force.

Step OP1.16: Proof of Theorem 2.12, (ii). Let E^\dagger and X^\dagger defined in (2.49). Note that X^\dagger is a conservative norm-continuous right process on E , so it satisfies (ii.1) from Theorem 2.12 due to (2.54).

Concerning assertion (ii.2), note that by Theorem 6.13, (i), we have that the (trace of the) norm topology on E^\dagger is a natural topology, hence for any $f \in C_b([0, \infty) \times \mathbb{R}^d)$ we have that $f|_{E^\dagger}$ is finely continuous, hence $P_t^\dagger f|_{E^\dagger}$ is finely continuous due to Remark 6.15. Furthermore, recall that \mathcal{U}^\dagger is a $\bar{\mu}$ -version of \bar{U} , hence by **Step OP1.3**, (ii.2) is as well satisfied by X^\dagger .

Assertion (ii.3) is also fulfilled by X^\dagger due to the conclusion in **Step OP1.10**.

Assertion (ii.4) with X^\dagger instead of X follows by (2.51) and (2.54).

Regarding assertion (ii.5), let us notice first that by a straightforward density argument we have that $C_c^2([0, \infty) \times \mathbb{R}^d) \in \mathbf{D}(\bar{\mathbf{L}}_1)$. Let $f \in C_c^2([0, \infty) \times \mathbb{R}^d)$ and set

$$M_f(t) := f(X^\dagger(t)) - f(X^\dagger(0)) - \int_0^t \bar{\mathbf{L}}f(X^\dagger(r)) dr, \quad t \geq 0. \quad (2.55)$$

Let us define $\bar{\mu}_1 \in \mathcal{P}([0, \infty) \times \mathbb{R}^d)$ by

$$\bar{\mu}_1 := e^{-t} u(t, x) dt \otimes dx \quad \text{on } [0, \infty) \times \mathbb{R}^d.$$

Notice that by Theorem 2.12, (i)-(ii.2), proved above, we have

$$\bar{\mu}_1 (U_\alpha^\dagger |\bar{\mathbf{L}}f|) < \infty \quad \text{for every } \alpha > 0, \quad (2.56)$$

and also that $(P_t^\dagger f)_{t \geq 0}$ solves the Cauchy problem for $(\bar{\mathbf{L}}_1, \mathbf{D}(\bar{\mathbf{L}}_1))$ on $L^p([0, \infty) \times \mathbb{R}^d; \bar{\mu})$. Consequently, using also the continuity of $t \mapsto M(t)$ we easily get

$$\mathbb{E}_{s,x}^\dagger \{M_f(t)\} = 0, \quad t \geq 0, \quad (s, x) \in E^\dagger, \quad \mathbb{P}_{\bar{\mu}_1}^\dagger \text{-a.e.},$$

and thus,

$$\mathbb{E}_{\tilde{\nu}}^\dagger \{M_f(t)\} = 0, \quad t \geq 0, \quad \text{for every } \tilde{\nu} \in \mathcal{P}(E^\dagger) \text{ for which there exists } c > 0 \text{ such that } \tilde{\nu} \leq c\bar{\mu}_1.$$

Now, by Proposition 6.56, there exists a subset $E_f \subset E^\dagger$ which is μ -inessential and such that

$$(M_f(t))_{t \geq 0} \quad \text{is an } (\mathcal{F}_t^\dagger)\text{-martingale under } \mathbb{P}_{\delta_{(s,x)}}^\dagger \text{ for every } (s, x) \in E_f. \quad (2.57)$$

Furthermore, for a countable subset $\mathcal{T} \subset C_c^\infty((0, \infty) \times \mathbb{R}^d)$ which is dense with respect to the uniform convergence up to (and including) the second derivatives, set

$$E := \left(\bigcap_{f \in \mathcal{T}} E_f \right) \cap \left(\bigcap_n [U_1^\dagger |\bar{\mathbf{L}}\psi_n| < \infty] \right),$$

where $(\psi_n)_{n \geq 1} \subset C_c^\infty((0, \infty) \times \mathbb{R}^d)$ is such that

$$\psi_n \geq 0 \quad \text{and} \quad \psi_n = 1 \text{ on } [1/n, n] \times B(0, n), \quad n \geq 1.$$

Note that by (2.56), Remark 6.33, and Remark 6.34, we have that $[U_1^\dagger |\bar{L}\psi_n| = \infty]$ is $\tilde{\mu}$ -inessential for every $f \in \mathcal{T}$, as well as $E^\dagger \setminus E$.

Now, we take X to be the restriction of X^\dagger from E^\dagger to E in the sense of Definition 6.36. It is clear that (ii.1)-(ii.4) are still fulfilled, whilst (ii.5) follows easily from (2.57) and the above mentioned density of \mathcal{T} in $C_c^\infty((0, \infty) \times \mathbb{R}^d)$.

(ii.6) Let $s > 0$ and $\nu_s \in \mathcal{P}(\mathbb{R}^d)$ such that $\nu_s \leq c\mu_s$ for some finite $c > 0$. By (2.12) and (2.54) we easily get that

$$\mathbb{E}_{\delta_s \otimes \mu_s}^\dagger \{|M(t)|\} < \infty, \quad t \geq 0.$$

Now, by assertion (ii.5) and the notation M_f introduced in (2.57), we deduce that

$$M_f \text{ is an } (\mathcal{F}(t))\text{-martingale under } \mathbb{P}_{\delta_s \otimes \nu_s}, \quad f \in C_c^\infty((0, \infty) \times \mathbb{R}^d).$$

Now, assertion (ii.6) follows from **Step OP1.15**.

2.4.3 Proof of Theorem 2.15

Because the process ensured by Theorem 2.12 requires some modifications, let us use the notation

$$X^\dagger = (\Omega^\dagger, \mathcal{F}^\dagger, \mathcal{F}^\dagger(t), X^\dagger(t), \mathbb{P}_{s,x}^\dagger, (s, x) \in E^\dagger, t \geq 0)$$

for the process constructed in Theorem 2.12, with transition function denoted by $(P_t^\dagger)_{t \geq 0}$ and resolvent $\mathcal{U}^\dagger = (U_\alpha^\dagger)_{\alpha > 0}$. Note that we have also used the notation X^\dagger to define the process from (2.49) in **Step OP1.10**; it is useful to note that the process we now denote by X^\dagger , namely the one furnished by Theorem 2.12, is in fact the restriction (in the sense of Definition 6.36) of X^\dagger from (2.49), so it inherits all the properties deduced for X^\dagger in all steps that follows after (including) **Step OP1.10**.

Proving that in the conclusions (ii.1), (ii.4) and (ii.7) from Theorem 2.12 we can allow $s = 0$, if condition $\mathbf{H}_6^{\dagger V}$ is in force, is highly nontrivial. We show this systematically, in what follows.

Step OP2.1: First of all, since the uniqueness assumption $\mathbf{H}_{\lesssim u}$ holds, one can apply [64, Lemma 2.12] to deduce that the family $(\eta^N)_{N > 0}$ given by Theorem 2.5 is a consistent family (see [57] for the notion), hence by the general theory (see again [57, Theorem 4.1]) admits a unique projective limit $\eta \in \mathcal{P}(C([0, \infty); \mathbb{R}^d))$, with one dimensional marginals $(\mu_t)_{t \in [0, \infty)}$. Let

$$(\eta^x)_{x \in \mathbb{R}^d} \subset \mathcal{P}(C([0, \infty); \mathbb{R}^d))$$

be the Borel family of probability measures obtained by disintegrating the law $\eta \in \mathcal{P}(C([0, \infty); \mathbb{R}^d))$ with respect to μ_0 , in the same spirit of Theorem 2.5. Also, let

$$(\eta_t^x)_{t \in [0, \infty)} \text{ be the one dimensional marginals of } \eta^x \text{ for every } x \in \mathbb{R}^d.$$

Step OP2.2: We claim that there exists $M \in \mathcal{B}(\mathbb{R}^d)$ such that $\mu_0(M) = 1$ and for all $x \in M$ and $f \in b\mathcal{B}([0, \infty) \times \mathbb{R}^d)$

$$\eta_s^x \left(P_t^\dagger f |_{E^\dagger}(s, \cdot) \right) = \eta_{s+t}^x (f(s+t, \cdot)), \quad ds \otimes dt\text{-a.e. } (s, t) \in (0, \infty)^2. \quad (2.58)$$

To prove the claim, note first that by a monotone class argument, it is sufficient to fix $f \in b\mathcal{B}([0, \infty) \times \mathbb{R}^d)$ and prove (2.58) with M possibly depending on f . Further, let $\nu_0 \in \mathcal{P}(\mathbb{R}^d)$ such that $\nu_0 \leq c\mu_0$ for some finite constant c , and define $\eta^\nu \in \mathcal{P}(C([0, \infty]; \mathbb{R}^d))$ by

$$\eta^{\nu_0}(A) := \int_A \eta^x(A) \nu_0(dx), \quad A \in \mathcal{B}(C([0, \infty]; \mathbb{R}^d)). \quad (2.59)$$

Consequently, for every $s > 0$ we have $\eta_s^{\nu_0} \leq c\mu_s$, and hence by Theorem 2.12, (ii.7) we have that

$$\nu_{t+s} := \mathbb{P}_{\delta_s \otimes \eta_s^{\nu_0}}^\dagger \circ (X_2^\dagger(t))^{-1}, \quad t \geq 0,$$

is the unique weakly continuous solution to the LFPE (2.3) on $[s, N]$ which consists of probability measures, belongs to $\mathcal{A}_{\leq u}^{s, N}$, and whose initial condition (at time s) is $\eta_s^{\nu_0}$, for all $N > s$. But $(\eta_{t+s}^{\nu_0})_{t \geq 0}$ is also such a solution (see Theorem 2.5, in fact [64, Proposition 2.7]), so we must have

$$\nu_{s+t} = \eta_{s+t}^{\nu_0}, \quad t \geq 0, \quad (2.60)$$

hence, using also Theorem 2.12, (ii.4),

$$\eta_s^{\nu_0} (P_t^\dagger f |_{E^\dagger}(s, \cdot)) = \mathbb{E}^{\dagger, \delta_s \otimes \eta_s^{\nu_0}} \{f |_{E^\dagger}(X^\dagger(t))\} = \eta_{s+t}^{\nu_0} (f(s+t, \cdot)), \quad s > 0, t \geq 0.$$

Since ν_0 is arbitrarily chosen such that $\nu_0 \leq c\mu_0$, the claim follows.

Remark 2.21. *It is worth to point out that relation (2.58) is not obtained for every $s, t > 0$ because we do not know whether the left hand side of the equality from (2.58) is continuous in both arguments (s, t) , the challenging continuity being the one in the s variable.*

Step OP2.3: *We claim that there exists $M' \in \mathcal{B}(\mathbb{R}^d)$ such that*

$$\mu_0(M') = 1, \quad \eta_0^x = \delta_x, \quad \text{and} \quad \eta_t^x(E^\dagger(t)) = 1, \quad x \in M', \quad dt\text{-a.e. } t > 0,$$

where recall that $E^\dagger(t)$ is given by (2.53).

Indeed, using Theorem 2.12, (ii.4), we get

$$\mu_0(\cdot) = \int_{\mathbb{R}^d} \eta_0^x(\cdot) \mu_0(dx), \quad \text{and} \quad \int_{\mathbb{R}^d} \eta_t^x(E^\dagger(t)) \mu_0(dx) = \eta_t(E^\dagger(t)) = \mu_t(E^\dagger(t)) = 1, \quad t > 0,$$

from which the claim can be easily deduced.

Step Q2.4: *Let us set*

$$M_0 := M \cap M', \quad E^{\dagger, 0} := (\{0\} \times M_0) \cup (E^\dagger \setminus \{0\} \times E^\dagger(0)),$$

and

$$P_t^{\dagger, 0} f(s, x) := \begin{cases} \eta_t^x(f(t, \cdot)), & \text{if } s = 0 \\ P_t^\dagger(f |_{E^\dagger})(s, x), & \text{if } s > 0 \end{cases}, \quad f \in b\mathcal{B}(E_0^\dagger), \quad t \geq 0, \quad (s, x) \in E^{\dagger, 0}. \quad (2.61)$$

We claim that

- (1) $P_0^{\dagger, 0} = \text{Id}$ and $P_t^{\dagger, 0} P_{t'}^{\dagger, 0} = P_{t+t'}^{\dagger, 0}$, as kernels on $E^{\dagger, 0}$, $dt \otimes dt'$ -a.e. $(t, t') \in (0, \infty)^2$.
- (2) $\lim_{t \rightarrow 0} P_t^{\dagger, 0} f = f$ pointwise on E_0^\dagger for all $f \in C_b(E^{\dagger, 0})$.
- (3) If we denote by $\mathcal{U}^{\dagger, 0} := \left(U_\alpha^{\dagger, 0} = \int_0^\infty e^{-\alpha t} P_t^{\dagger, 0} dt \right)_{\alpha > 0}$ the resolvent of $(P_t^{\dagger, 0})_{t \geq 0}$, then $\mathcal{U}^{\dagger, 0}$ is a resolvent family of kernels on $E^{\dagger, 0}$ such that

$$\alpha U_\alpha^{\dagger, 0} 1 = 1, \quad \alpha > 0, \quad \text{and} \quad \lim_{\alpha \rightarrow \infty} \alpha U_\alpha^{\dagger, 0} f = f, \quad f \in C_b(E^{\dagger, 0}). \quad (2.62)$$

To prove the claim, first of all notice that

the mapping $[0, \infty) \times E^{\dagger, 0} \ni (t, s, x) \mapsto P_t^{\dagger, 0} f(s, x) \in \mathbb{R}$ is measurable for every $f \in b\mathcal{B}(E^{\dagger, 0})$.

Let us prove (1), by noticing first that $P_0^{\dagger,0} = \text{Id}$ follows from the fact that $\eta_0^x = \delta_x$, $x \in M_0$ and $\delta_{(s,x)} \circ P_0^{\dagger,0} = \delta_{(s,x)}$, $(s,x) \in E^{\dagger,0}$. To show the semigroup property for a.e. $t, t' > 0$, note that taking first $s = 0$ we have by (2.58) that

$$\begin{aligned} P_t^{\dagger,0} P_{t'}^{\dagger,0} f(0, x) &= \eta_t^x \left(P_{t'}^{\dagger,0} f(t, \cdot) \right) = \eta_t^x \left(P_{t'}^{\dagger} (f|_{E^\dagger}) (t, \cdot) \right) = \eta_{t+t'}^x (f(t+t', \cdot)) \\ &= P_{t+t'}^{\dagger,0} f(0, x), \quad dt \otimes dt' \text{-a.e.} \end{aligned}$$

The case $s > 0$ is immediate since $P_t^{\dagger} P_{t'}^{\dagger} = P_{t+t'}^{\dagger}$ on E^\dagger .

To prove (2), it is sufficient to show the desired convergence for every f which is bounded and Lipschitz on $E^{\dagger,0}$. With this reduction, if $(s, x) \in E^\dagger$ then by using the fact that $\left(P_t^{\dagger} \right)_{t \geq 0}$ is the transition function of a right process with a.s. continuous paths in E^\dagger , we get

$$\lim_{t \rightarrow 0} P_t^{\dagger,0} f(s, x) = \lim_{t \rightarrow 0} P_t^{\dagger} (f|_{E^\dagger}) (s, x) = f(s, x).$$

If $s = 0$ and $x \in M_0$, then

$$\begin{aligned} \lim_{t \rightarrow 0} \left| P_t^{\dagger,0} f(0, x) - f(0, x) \right| &= \lim_{t \rightarrow 0} \left| \eta_t^x (f(t, \cdot)) - f(0, x) \right| \\ &= \limsup_{t \rightarrow 0} \left[\eta_t^x (\min(\|\cdot - x\|, 2\|f\|_\infty)) + |f|_{\text{Lip}} t \right] = 0, \end{aligned}$$

where $|f|_{\text{Lip}}$ denotes the Lipschitz norm of f .

The proof of (3) is now straightforward, using (1) and (2) from above, as well as **Step OP2.3** and the fact that $P_t^\dagger 1 = 1$ on E^\dagger .

Step OP2.5: We claim that there exists $M_{00} \in \mathcal{B}(\mathbb{R}^d)$, $M_{00} \subset M_0$ such that $\mu_0(M_{00}) = 1$, and if we set

$$E^{\dagger,00} := (\{0\} \times M_{00}) \cup E^\dagger,$$

then

$$U_1^{\dagger,0} ((s, x), E^{\dagger,0} \setminus E^{\dagger,00}) = 0, \quad (s, x) \in E^{\dagger,00},$$

and the restricted resolvent $\mathcal{U}^{\dagger,00} := \mathcal{U}^{\dagger,0}|_{E^{\dagger,00}}$ is a Markovian resolvent of kernels with no branch points, that is for all $(s, x) \in E^{\dagger,00}$, $f, g \in b\mathcal{B}(E^{\dagger,00})$ we have

$$\lim_{\alpha \rightarrow \infty} \alpha U_{\alpha+1}^{\dagger,00} \left(U_1^{\dagger,00} f \wedge U_1^{\dagger,00} g \right) (s, x) = U_1^{\dagger,00} f(s, x) \wedge U_1^{\dagger,00} g(s, x). \quad (2.63)$$

The proof of this step is involved, so let us split it in several sub-steps:

(1) Let us set

$$E^{\dagger,00} := \left\{ (s, x) \in E^{\dagger,0} : \lim_{\alpha \rightarrow \infty} \alpha U_{\alpha+1}^{\dagger,0} \left(U_1^{\dagger,0} f \wedge U_1^{\dagger,0} g \right) (s, x) = U_1^{\dagger,0} f(s, x) \wedge U_1^{\dagger,0} g(s, x), f, g \in b\mathcal{B}(E^{\dagger,0}) \right\}$$

Then by [14, Theorem 1.2.11 and assertion 1 of Remark at page 23] we have that $E^{\dagger,00} \in \mathcal{B}(E^{\dagger,0})$ and $U_1^{\dagger,0}(\cdot, E^{\dagger,0} \setminus E^{\dagger,00}) = 0$ on $E^{\dagger,0}$. Moreover, notice that if we set $E' := E^\dagger \setminus \{0\} \times E^\dagger(0)$ then $E' \in \mathcal{A}(\mathcal{U}^\dagger)$ (use e.g. Theorem 2.12, (ii.4)), $\mathcal{U}^{\dagger,0}|_{E'} = \mathcal{U}^\dagger|_{E'}$, and \mathcal{U}^\dagger has no branch points since it is the resolvent of a right process on E^\dagger ; see Proposition 6.8 and the remark after it. Thus,

$$E^{\dagger,0} \setminus E^{\dagger,00} \subset \{0\} \times M_0 \subset \{0\} \times E^{\dagger,0}(0),$$

and it remains to show that

$$\mu_0(E^{\dagger,00}(0)) = 1.$$

(2) The next step is to notice that by the monotone class argument [14, Lemma 1.2.10], it is sufficient to fix $f, g \in b\text{Lip}([0, \infty) \times \mathbb{R}^d)$ with compact support, where $b\text{Lip}([0, \infty) \times \mathbb{R}^d)$ denotes the space of bounded and Lipschitz functions on $[0, \infty) \times \mathbb{R}^d$, and show that

$$\mu_0 \left(\left\{ x \in E^{\dagger,0}(0) : \lim_{\alpha \rightarrow \infty} \alpha U_{\alpha+1}^{\dagger,0} \left(U_1^{\dagger,0} f \wedge U_1^{\dagger,0} g \right) (0, x) = U_1^{\dagger,0} f(0, x) \wedge U_1^{\dagger,0} g(0, x) \right\} \right) = 1.$$

(3) In light of (2) from above, let us fix $f, g \in b\text{Lip}([0, \infty) \times \mathbb{R}^d)$ and define for $x \in M_0$ and $\alpha > 0$

$$\begin{aligned}\psi_\alpha(x) &:= \alpha U_{\alpha+1}^{\dagger,0} \left(U_1^{\dagger,0} f \wedge U_1^{\dagger,0} g \right) (0, x) = \int_0^\infty \alpha e^{-(\alpha+1)t} \eta_t^x \left(\left(U_1^{\dagger} f \wedge U_1^{\dagger} g \right) (t, \cdot) \right) dt, \\ \psi(x) &:= U_1^{\dagger,0} f(0, x) \wedge U_1^{\dagger,0} g(0, x)\end{aligned}$$

Note that $U_1^{\dagger,0} f \wedge U_1^{\dagger,0} g$ is $\mathcal{U}_1^{\dagger,0}$ -supermedian, hence $\psi_\alpha \leq \psi_\beta \leq \psi$ for every $0 < \alpha \leq \beta$, and we can define

$$\psi_\infty := \lim_{\alpha \rightarrow \infty} \psi_\alpha = \sup_{\alpha > 0} \psi_\alpha \leq \psi, \quad \text{pointwise on } M_0.$$

Consequently, to achieve **Step OP2.5** it is sufficient to show that

$$(\mu_0(\psi_\infty) =) \lim_{\alpha \rightarrow \infty} \mu_0(\psi_\alpha) = \mu_0(\psi). \quad (2.64)$$

(4) Proceeding to prove (2.64), let us notice first that using e.g. Theorem 2.12, (ii.4), and the fact that $P_t^\dagger \left(U_1^\dagger f \wedge U_1^\dagger g \right) \leq e^t \left(U_1^\dagger f \wedge U_1^\dagger g \right)$, $t \geq 0$, we get

$$\begin{aligned}\mu_0(\psi_\alpha) &= \int_0^\infty \alpha e^{-(\alpha+1)t} \mu_t \left(\left(U_1^\dagger f \wedge U_1^\dagger g \right) (t, \cdot) \right) dt \\ e^{-t} \mu_t \left(\left(U_1^\dagger f \wedge U_1^\dagger g \right) (t, \cdot) \right) &= e^{-t} \mu_s \left(\left(P_{t-s}^\dagger \left(U_1^\dagger f \wedge U_1^\dagger g \right) \right) (s, \cdot) \right) \\ &\leq e^{-s} \mu_s \left(\left(U_1^\dagger f \wedge U_1^\dagger g \right) (s, \cdot) \right) dt, \quad s \leq t,\end{aligned}$$

we obtain that $t \mapsto e^{-t} e^{-t} \mu_t \left(\left(U_1^\dagger f \wedge U_1^\dagger g \right) (t, \cdot) \right)$ is non-decreasing. This easily implies that (2.64) is satisfied if we show that

$$\exists (t_n)_n \searrow 0 : \quad \lim_n \mu_{t_n} \left(\left(U_1^\dagger f \wedge U_1^\dagger g \right) (t_n, \cdot) \right) = \mu_0(\psi). \quad (2.65)$$

(5) We focus now on proving (2.65). To this end, we write for $t > 0$

$$\begin{aligned}\mu_t \left(\left(U_1^\dagger f \wedge U_1^\dagger g \right) (t, \cdot) \right) &= \mu_t (1_{E^{\dagger,0}(0)} \psi) + e_t, \\ |e_t| &\leq \mu_t \left(\left| 1_{E^{\dagger,0}(t)}(\cdot) U_1^{\dagger,0} f(t, \cdot) - 1_{E^{\dagger,0}(0)}(\cdot) U_1^{\dagger,0} f(0, \cdot) \right| \right. \\ &\quad \left. + \left| 1_{E^{\dagger,0}(t)}(\cdot) U_1^{\dagger,0} g(t, \cdot) - 1_{E^{\dagger,0}(0)}(\cdot) U_1^{\dagger,0} g(0, \cdot) \right| \right).\end{aligned}$$

Hence, using \mathbf{H}_0^{TV} we have

$$\limsup_{t \rightarrow 0} \left| \mu_t (1_{E^{\dagger,0}(0)} \psi) - \mu_0(\psi) \right| \leq \|\psi\|_\infty \limsup_{t \rightarrow 0} \int_{\mathbb{R}^d} |u(t, x) - u(0, x)| dx = 0,$$

hence it remains to prove that $\exists (t_n)_n \searrow 0 : \quad \lim_n e_{t_n} = 0$, which boils down to show that for every $f \in b\text{Lip}([0, \infty) \times \mathbb{R}^d)$ with compact support we have

$$\exists (t_n)_n \searrow 0 : \quad \lim_n \mu_{t_n} \left(\left| 1_{E^{\dagger,0}(t_n)}(\cdot) U_1^{\dagger,0} f(t_n, \cdot) - 1_{E^{\dagger,0}(0)}(\cdot) U_1^{\dagger,0} f(0, \cdot) \right| \right) = 0.$$

Remark 2.22. Here and in general, if $v : F \rightarrow \mathbb{R}$ and $w : A \rightarrow \mathbb{R}$ with $A \subset F$, then by vw we understand either $v|_A w$ or $v(1_A w)$, where recall that $1_A w$ stands for the extension of w from A to F by setting $w = 0$ on $F \setminus A$.

Furthermore, by \mathbf{H}_0^{TV} , we deduce that **Step OP2.5** is achieved if we show that for every $f \in b\text{Lip}([0, \infty) \times \mathbb{R}^d)$ with compact support,

$$\exists (t_n)_n \searrow 0 : \quad \lim_n \int_{\mathbb{R}^d} \left| u(t_n, x) U_1^{\dagger,0} f(t_n, x) - u(0, x) U_1^{\dagger,0} f(0, x) \right| dx = 0. \quad (2.66)$$

(6) To prove (2.66) we need first a compactness result. More precisely, fixing $f \in b\text{Lip}([0, \infty) \times \mathbb{R}^d)$ with support (for simplicity) in $[0, 1] \times \mathbb{R}^d$ and $h \in C_c^\infty(\mathbb{R}^d)$, we claim that by setting

$$w(t, x) := u(t, x)U_1^\dagger f(t, x)h(x), \quad x \in \mathbb{R}^d, t \in (0, 1],$$

then there exists a $dt \otimes dx$ -version \tilde{w} such that

$$\tilde{w} \in C([0, 1]; L^2(\mathbb{R}^d)).$$

In fact, the claim follows by Lions-Magenes interpolation lemma [51, Chapter 3, Theorem 3.1], once we show that

$$w \in L^2([0, 1]; H^1(\mathbb{R}^d)) \quad \text{and} \quad \frac{dw}{dt} \in L^2([0, 1]; H^{-1}(\mathbb{R}^d)). \quad (2.67)$$

To this end, following [61] let us consider $(\mathcal{A}_0, \mathbf{D}(\mathcal{A}_0))$ the closure in $L^2([0, 1] \times \mathbb{R}^d; \bar{\mu}_1)$ of the symmetric bilinear form

$$\mathcal{A}_0(\varphi, \psi) := \int_{[0, 1] \times \mathbb{R}^d} \langle A(t, x) \nabla \varphi(t, x), \nabla \psi(t, x) \rangle \bar{\mu}_1(dt, dx), \quad \varphi, \psi \in C_0^\infty([0, 1] \times \mathbb{R}^d).$$

Also, note that by Theorem 2.12, (i), on $L^p([0, 1] \times \mathbb{R}^d; \bar{\mu}_1)$ we can identify U_1^\dagger with \bar{U}_1^1 , where \bar{U}_1^1 is given in Theorem 2.10. Consequently, by [61, Theorem 1.7], b) and c), we have that

$$U_1^\dagger f \text{ and } q(\cdot, \cdot) := h(\cdot)U_1^\dagger f(\cdot, \cdot) \text{ belong to } b\mathbf{D}(\bar{\mathbf{L}}_1^1) \subset \mathbf{D}(\mathcal{A}_0), \quad (2.68)$$

where $\mathbf{D}(\bar{\mathbf{L}}_1^1)$ is given in Theorem 2.10, and also that

$$- \int_{[0, 1] \times \mathbb{R}^d} \frac{d\varphi}{dt} q \, d\bar{\mu}_1 = \mathcal{A}_0(q, \varphi) + \int_{[0, 1] \times \mathbb{R}^d} \langle b - b^0, \nabla \varphi \rangle q \, d\bar{\mu}_1 + \int_{[0, 1] \times \mathbb{R}^d} \varphi \bar{\mathbf{L}}_1^1 q \, d\bar{\mu}_1, \quad (2.69)$$

where

$$b_i^0 := \sum_{j=1}^d \left(\partial_{x_j} a_{ij} + 2a_{ij} \frac{\partial_{x_j} \sqrt{u}}{\sqrt{u}} \right), \quad 1 \leq i \leq d.$$

Now, let us show that $w \in L^2([0, 1]; H^1(\mathbb{R}^d))$. To this end, note that it is not a priori clear that we can consider the classical weak gradient in the x -variable, ∇w . Thus, to make the computations rigorous, let us consider

$$(v_k)_k \subset C_0^\infty([0, 1] \times \mathbb{R}^d) \subset \mathbf{D}(\mathcal{A}_0) \text{ uniformly bounded,}$$

and such that

$$\lim_k \mathcal{A}_0(v_k - U_1^\dagger f, v_k - U_1^\dagger f) = 0, \quad q_k := hv_k, \quad w_k := uq_k, \quad k \geq 1.$$

For later use, note that it is straightforward to see that $\lim_k \mathcal{A}_0(q_k - q, q_k - q) = 0$. Furthermore, it is easy to see that for our current goal it is sufficient to show that

$$(w_k)_k \text{ is bounded in } L^2([0, 1]; H^1(\mathbb{R}^d)). \quad (2.70)$$

By the product rule,

$$\nabla w_k = q_k \nabla u + u \nabla q_k = 2q_k \sqrt{u} \nabla \sqrt{u} + u \nabla q_k,$$

so that by $\mathbf{H}^{\sqrt{u}}$ and using that q_k is uniformly bounded with respect to k, t, x , and has compact support uniformly with respect to k , the first summand $2q_k \sqrt{u} \nabla \sqrt{u}$ can be easily bounded in $L^2([0, 1]; L^2(\mathbb{R}^d))$. Hence, it remains to show that

$$(u \nabla q_k)_k \text{ is bounded in } L^2([0, 1] \times \mathbb{R}^d; dt dx).$$

But, according to (2.68),

$$\int_{[0,1] \times \mathbb{R}^d} |u \nabla q_k|^2 dt dx = \int_{[0,1] \times \mathbb{R}^d} u |\nabla q_k|^2 d\bar{\mu}_1 \leq c \mathcal{A}_0(q_k, q_k),$$

where c is a constant depending on the ellipticity of A and the L^∞ -norm of u on $[0, 1] \times \text{supp}(h)$, where $\text{supp}(h)$ is the support of h in \mathbb{R}^d . Now, (2.70) follows since it is easy to see that $\sup_k \mathcal{A}_0(q_k, q_k) < \infty$ due to the fact that $\sup_k \mathcal{A}_0(v_k, v_k) < \infty$.

Finally, let us show that $\frac{dw}{dt} \in L^2([0, 1]; H^{-1}(\mathbb{R}^d))$, namely that there exists $\dot{w} \in L^2([0, 1]; H^{-1}(\mathbb{R}^d))$ and a constant C such that

$$\left| \int_{[0,1] \times \mathbb{R}^d} \frac{d\varphi}{dt} w dt dx \right| \leq C \|\varphi\|_{L^2((0,1) \times H^1(\mathbb{R}^d))}, \quad \text{for all } \varphi \in C_0^\infty((0,1) \times \mathbb{R}^d).$$

To do so, note that by (2.69)

$$\begin{aligned} \left| \int_{[0,1] \times \mathbb{R}^d} \frac{d\varphi}{dt} w dt dx \right| &= \left| \int_{[0,1] \times \mathbb{R}^d} \frac{d\varphi}{dt} q u dt dx \right| \\ &\leq |\mathcal{A}_0(q, \varphi)| + \left| \int_{[0,1] \times \mathbb{R}^d} \langle b - b^0, \nabla \varphi \rangle q d\bar{\mu}_1 \right| + \left| \int_{[0,1] \times \mathbb{R}^d} \varphi \bar{\mathbb{L}}_1^{-1} q d\bar{\mu}_1 \right| \end{aligned}$$

We treat now each of the three terms in the last sum:

$$\begin{aligned} |\mathcal{A}_0(q, \varphi)| &= \lim_k |\mathcal{A}_0(q_k, \varphi)| \leq \sup_k \int_{[0,1] \times \text{supp}(h)} \|A^{1/2} \nabla q_k\| \|A^{1/2} \varphi\| d\bar{\mu}_1 \\ &\leq \sup_k \mathcal{A}_0(q_k, q_k)^{1/2} \left(\int_{[0,1] \times \text{supp}(h)} \|A^{1/2} \nabla \varphi\|^2 d\bar{\mu}_1 \right)^{1/2} \\ &\leq C \sup_k \mathcal{A}_0(q_k, q_k)^{1/2} \|\varphi\|_{L^2([0,1]; H^1(\mathbb{R}^d))}, \end{aligned}$$

where the constant C is independent of k and φ , and its existence is ensured by \mathbf{H}_a and $\mathbf{H}^{\sqrt{u}}$. also, note that $\sup_k \mathcal{A}_0(q_k, q_k) < \infty$.

Let us now treat the second term, for which we have

$$\begin{aligned} \left| \int_{[0,1] \times \mathbb{R}^d} \langle b - b^0, \nabla \varphi \rangle q d\bar{\mu}_1 \right| &\leq \|\varphi\|_{L^2((0,1) \times H^1(\mathbb{R}^d))} \left(\int_{[0,1] \times \mathbb{R}^d} \|b - b^0\|^2 q^2 u d\bar{\mu}_1 \right)^{1/2} \\ &\leq C \|\varphi\|_{L^2([0,1]; H^1(\mathbb{R}^d))}, \end{aligned}$$

where the constant C is independent of φ , and its existence is ensured by $\mathbf{H}_{a,b}^{\sqrt{u}}$.

We now treat the last term, as follows:

$$\left| \int_{[0,1] \times \mathbb{R}^d} \varphi \bar{\mathbb{L}}_1^{-1} q d\bar{\mu}_1 \right| \leq \|\varphi\|_{L^2((0,1) \times \mathbb{R}^d)} \|u \bar{\mathbb{L}}_1^{-1} q\|_{L^2((0,1) \times \mathbb{R}^d; \bar{\mu}_1)},$$

whilst $\|u \bar{\mathbb{L}}_1^{-1} q\|_{L^2((0,1) \times \mathbb{R}^d; \bar{\mu}_1)}$ is finite since $\bar{\mathbb{L}}_1^{-1} q$ is bounded with compact support by [61, Theorem 1.7, (c)], whilst u is bounded on compact sets by $\mathbf{H}^{\sqrt{u}}$.

Thus, the claim in item (6) is proved.

- (7) *We now claim that there exists $F : [0, 1] \times \mathbb{R}^d \rightarrow \mathbb{R}$ measurable, which is a $dt \otimes dx$ -version of $u(\cdot, \cdot) U_1^\dagger f(\cdot, \cdot)$ on $[0, 1] \times \mathbb{R}^d$ such that*

$$hF \in C([0, 1]; L^2(\mathbb{R}^d)), \quad \text{for every } h \in C_c^\infty(\mathbb{R}^d).$$

To prove this, for each $n \geq 0$ let $h_n \in C_c^\infty(\mathbb{R}^d)$ such that $h_n = 1$ on $B(0, n)$. Then, by (6), let $\tilde{w}_n \in C([0, 1]; L^2(\mathbb{R}^d))$ be a $dt \otimes dx$ -version of $h_n(\cdot)u(\cdot, \cdot)U_1^\dagger f(\cdot, \cdot)$, and for every $t \in [0, 1]$ define

$$F(t, \cdot) := \sum_{n \geq 0} \tilde{w}_{n+1}(t, \cdot) 1_{B(0, n+1) \setminus B(0, n)}(\cdot).$$

It is now easy to see that F satisfies the requirements from (7).

(8) Let F be as in item (7). We claim that

$$\forall (t_n)_n \searrow 0 \text{ there exists a subsequence } (t_{n_k})_k \text{ such that } \lim_k F(t_{n_k}, x) = F(0, x) \quad dx\text{-a.e.}$$

To show this, let $V(x) := e^{-\|x\|}$, $x \in \mathbb{R}^d$, and $(t_k)_k \subset (0, 1]$ such that $\lim_k t_k = 0$. Then

$$\begin{aligned} \|V(\cdot)F(t_k, \cdot) - V(\cdot)F(0, \cdot)\|_{L^1(\mathbb{R}^d)} &\leq \sum_{n \geq 0} \int_{B(0, n+1) \setminus B(0, n)} V(x) |\tilde{w}_{n+1}(t_k, x) - \tilde{w}_{n+1}(0, x)| dx \\ &\leq \sum_{n \geq 0} e^{-n} \int_{B(0, n+1) \setminus B(0, n)} |\tilde{w}_{n+1}(t_k, x) - \tilde{w}_{n+1}(0, x)| dx \\ &= \sum_{n \geq 0} e^{-n} a_{k, n}, \end{aligned}$$

where

$$a_{k, n} := \int_{B(0, n+1) \setminus B(0, n)} |\tilde{w}_{n+1}(t_k, x) - \tilde{w}_{n+1}(0, x)| dx, \quad k \geq 1, n \geq 0,$$

whilst $\tilde{w}_n, n \geq 0$ are given in (7). Now, on the one hand we have that $\tilde{w}_{n+1} \in C([0, 1]; L^2(\mathbb{R}^d))$, hence $\lim_k a_{k, n} = 0$. On the other hand,

$$\begin{aligned} \sup_k a_{k, n} &\leq 2 \sup_{t \in [0, 1]} \int_{B(0, n+1)} |\tilde{w}_{n+1}(t, x)| dx = 2 \|\tilde{w}_{n+1}(t, x)\|_{L^\infty([0, 1]; L^1(B(0, n+1)))} \\ &\leq 2 \|f\|_{L^\infty([0, 1] \times \mathbb{R}^d)} < \infty. \end{aligned}$$

Therefore, by dominated convergence, we obtain

$$\lim_k V(\cdot)F(t_k, \cdot) = V(\cdot)F(0, \cdot) \quad \text{in } L^1(\mathbb{R}^d),$$

hence, taking also into account that $V > 0$, there exists a subsequence $(t_{k_n})_{n \geq 1}$ such that

$$\lim_n F(t_{k_n}, \cdot) = F(0, \cdot) \quad dx\text{-a.e.}$$

(9) We claim that

$$\exists F_0 : \mathbb{R}^d \rightarrow \mathbb{R} \text{ and } (t_n)_n \searrow 0 \text{ such that } \lim_n 1_{E^{\dagger, 0}(t_n)}(\cdot) U_1^{\dagger, 0} f(t_n, \cdot) = F_0(\cdot) \quad dx\text{-a.e.} \quad (2.71)$$

Indeed, first of all, since F given in (7) is a $dt \otimes dx$ -version of $u(\cdot, \cdot)U_1^\dagger f(\cdot, \cdot)$ on $[0, 1] \times \mathbb{R}^d$, there exists $(t_n)_{n \geq 1} \searrow 0$ such that

$$F(t_n, \cdot) = u(t_n, \cdot)U_1^\dagger(t_n, \cdot), \quad n \geq 1, dx\text{-a.e.}$$

Therefore, using also that $U_1^{\dagger, 0}(t, \cdot) = U_1^\dagger(t, \cdot)$ for $t > 0$, we have

$$\exists (t_{n_k})_k : \lim_k u(t_{n_k}, \cdot)U_1^{\dagger, 0}(t_{n_k}, \cdot) = F(0, \cdot) \quad dx\text{-a.e.}$$

Now, since $\lim_k u(t_{n_k}, \cdot) = u(0, \cdot)$ in $L^1(\mathbb{R}^d)$ and $\int_{E^{\dagger, 0}(t)} u(t, x) dx = 1, t \in [0, 1]$, we can easily deduce that, by passing to a further subsequence,

$$\exists (t_{k_l})_l : \lim_l 1_{E^{\dagger, 0}(t_{k_l})}(\cdot) U_1^{\dagger, 0}(t_{k_l}, \cdot) = 1_{E^{\dagger, 0}(0)}(\cdot) \frac{F(0, \cdot)}{u(0, \cdot)} \quad dx\text{-a.e.},$$

which proves the claim.

(10) Let us fix $s \geq 0$, $f \in b\text{Lip}([0, \infty) \times \mathbb{R}^d)$ with compact support, and define

$$\varphi_t(x) = \varphi_t^s(x) = 1_{E^{\dagger,0}(t)}(x) P_s^{\dagger,0} f(t, x), \quad t \in [0, 1], x \in \mathbb{R}^d,$$

where $1_{E^{\dagger,0}(t)}(x)g(x)$, $x \in \mathbb{R}^d$ denotes the extension by zero (from $E^{\dagger,0}(t)$ to \mathbb{R}^d) of any given function $g : E^{\dagger,0}(t) \rightarrow \mathbb{R}$.

We claim that

$$\lim_{t \rightarrow 0} \mu_0(v_0 \varphi_t) = \mu_0(v_0 \varphi_0) \quad \text{for all } v_0 \in \mathcal{A}_0, \quad (2.72)$$

where \mathcal{A}_0 is the class of probability densities employed in \mathbf{H}_0^{TV} .

To prove the above assertion, let $v_0 \in \mathcal{A}_0$,

$$v_0 := v_0(x)u(0, x)dx \leq \|v_0\|_{\infty}\mu_0, \quad \text{and } \eta_t^{\nu_0}, t \geq 0 \text{ be given by (2.59).}$$

Then,

$$\begin{aligned} \mu_0(v_0 \varphi_t) &= \int_{\mathbb{R}^d} v_0(x) \varphi_t(x) \mu_0(dx) = \int_{\mathbb{R}^d} v_0(x) u(0, x) \varphi_t(x) dx \\ &= \int_{\mathbb{R}^d} \left(v_0(x) u(0, x) - \frac{d\eta_t^{\nu_0}}{dx}(x) \right) \varphi_t(x) dx + \int_{\mathbb{R}^d} \varphi_t(x) \eta_t^{\nu_0}(dx) \\ &= \int_{\mathbb{R}^d} \left(v_0(x) u(0, x) - \frac{d\eta_t^{\nu_0}}{dx}(x) \right) \varphi_t(x) dx + \eta_t^{\nu_0}(P_s^{\dagger,0} f(t, \cdot)) \end{aligned}$$

hence by (2.58) we can continue with

$$= \int_{\mathbb{R}^d} \left(v_0(x) u(0, x) - \frac{d\eta_t^{\nu_0}}{dx}(x) \right) \varphi_t(x) dx + \eta_{s+t}^{\nu_0}(f(s+t, \cdot)).$$

Now, on the one hand, since \mathbf{H}_0^{TV} is in force we have that the first term in the right-hand side of the last equality converges to zero when $t \rightarrow 0$. On the other hand, by the weak continuity of $(\eta_t^{\nu_0})_{t \geq 0}$ and the fact that $f \in b\text{Lip}([0, \infty) \times \mathbb{R}^d)$, we easily get

$$\lim_{t \rightarrow 0} \eta_{s+t}^{\nu_0}(f(s+t, \cdot)) = \eta_s^{\nu_0}(f(s, \cdot)) = \int_{\mathbb{R}^d} v_0(x) \varphi_0(x) \mu_0(dx).$$

(11) Let us now finish the proof of the claim in **Step OP2.5**. First, notice that

$$\int_0^{\infty} e^{-s} \varphi_t^s(x) ds = 1_{E^{\dagger,0}(t)}(x) U_1^{\dagger,0} f(t, x), \quad t \geq 0, x \in \mathbb{R}^d.$$

Consequently, on the one hand we have

$$\lim_{t \rightarrow 0} \mu_0(v_0(\cdot) 1_{E^{\dagger,0}(t)}(x) U_1^{\dagger,0} f(t, \cdot)) = \mu_0(v_0(\cdot) 1_{E^{\dagger,0}(0)}(\cdot) U_1^{\dagger,0} f(0, \cdot)) \quad \text{for all } v_0 \in \mathcal{A}_0,$$

On the other hand, by item (9) and using dominated convergence, we deduce that

$$\exists (t_n)_n \searrow 0 : \lim_{t \rightarrow 0} \mu_0(v_0(\cdot) 1_{E^{\dagger,0}(t)}(x) U_1^{\dagger,0} f(t, \cdot)) = \mu_0(v_0(\cdot) F_0(\cdot)) \quad \text{for all } v_0 \in \mathcal{A}_0,$$

where F_0 is the function entailed in (9). Since \mathcal{A}_0 is assumed to have the separation property from \mathbf{H}_0^{TV} , we deduce that

$$F_0 = 1_{E^{\dagger,0}(0)}(\cdot) U_1^{\dagger,0} f(0, \cdot) \quad \mu_0\text{-a.e.} \quad (2.73)$$

Finally, by (2.71) and (2.73) we can ensure (2.66), hence **Step OP2.5** is accomplished.

Step OP2.6: Let $\mathcal{U}^{\dagger,00} = \mathcal{U}^{\dagger,0}|_{E^{\dagger,00}}$ be the resolvent of Markov kernels constructed in **Step OP2.5**. Also, recall that \mathcal{U}^{\dagger} is the resolvent of the right process X^{\dagger} on E^{\dagger} , as in (2.49), and note that using **Step OP1.10** we have

$$U_1^{\dagger}(\cdot, \{0\} \times E^{\dagger}(0)) = 0 \quad \text{and} \quad E^{\dagger} \setminus \{0\} \times E^{\dagger}(0) \in \mathcal{A}(\mathcal{U}^{\dagger}).$$

Hence, we can consider the restricted resolvent and process (X', \mathcal{U}') of $(X^{\dagger}, \mathcal{U}^{\dagger})$ from E^{\dagger} to $E' := E^{\dagger} \setminus \{0\} \times E^{\dagger}(0)$. Now, let us notice that we have the natural extension (see Definition 6.49)

$$(E', \mathcal{U}') \subset (E^{\dagger,00}, \mathcal{U}^{\dagger,00}).$$

Moreover, by **Step OP2.4** and **Step OP2.5**, one can easily see that $\mathcal{E}(\mathcal{U}_\beta^{\dagger,00})$ is min stable, contains the constant functions, and generates $\mathcal{B}(E^{\dagger,00})$ for one (hence all) $\beta > 0$; see Section 6 for details. Consequently, by Proposition 6.8 and Theorem 6.51, we deduce that there exists a right process on $E^{\dagger,00}$ denoted by $X^{\dagger,00}$

$$X^{\dagger,00} := \left(\Omega^{\dagger,00}, \mathcal{F}^{\dagger,00}, \mathcal{F}_t^{\dagger,00}, X^{\dagger,00}(t), \theta^{\dagger,00}(t), \mathbb{P}_{(s,x)}^{\dagger,00}, (s, x) \in E^{\dagger,00} \right),$$

which is a natural extension of X' from E' to $E^{\dagger,00}$; using the notation in Definition 6.49, this means that

$$(E', X') \subset (E^{\dagger,00}, X^{\dagger,00}). \quad (2.74)$$

In particular,

$$\{0\} \times E^{\dagger,00}(0) \text{ is polar with respect to } \mathcal{U}^{\dagger,00}. \quad (2.75)$$

Note that $X^{\dagger,00}$ is also conservative, namely

$$\mathbb{P}_{s,x}^{\dagger,00} (X^{\dagger,00}(t) \in E^{\dagger,00} \text{ for all } t \geq 0) = 1, \quad \text{for all } (s, x) \in E^{\dagger,00}.$$

We claim that for all $(s, x) \in E^{\dagger,00}$, $s > 0$, we have

$$\mathbb{P}_{s,x}^{\dagger,00} (X^{\dagger,00} \text{ has norm-continuous paths in } E^{\dagger,00} \subset [0, \infty) \times \mathbb{R}^d) = 1. \quad (2.76)$$

Indeed, if $(s, x) \in E^{\dagger,00}$ with $s > 0$ then

$$\mathbb{P}_{s,x}^{\dagger,00} \circ (X^{\dagger,00})^{-1} = \mathbb{P}_{s,x}^{\dagger} \circ (X^{\dagger})^{-1},$$

and in this case the claim is clear since X^{\dagger} has norm-continuous paths a.s. for all starting points.

Step OP2.7: We claim that

$$\mathbb{P}_{\delta_0 \otimes \mu_0}^{\dagger,00} ([0, \infty) \ni t \mapsto X^{\dagger,00}(t) \text{ is norm-continuous in } E^{\dagger,00} \subset [0, \infty) \times \mathbb{R}^d) = 1. \quad (2.77)$$

To show this end, note first that by (2.74)-(2.75), and the fact that X' has a.s. (for every starting point) norm-continuous paths in E' , we deduce that it is enough to show

$$\mathbb{P}_{\delta_0 \otimes \mu_0}^{\dagger,00} ([0, N) \ni t \mapsto X^{\dagger,00}(t) \text{ is norm-continuous in } E^{\dagger,00}) = 1, \text{ for some } N > 0. \quad (2.78)$$

As a matter of fact, by the same (2.74)-(2.75) we clearly have

$$\mathbb{P}_{\delta_0 \otimes \mu_0}^{\dagger,00} ((0, N) \ni t \mapsto X^{\dagger,00}(t) \text{ is norm-continuous in } E^{\dagger,00}) = 1, \text{ for every } N > 0,$$

so that showing (2.78) boils down to showing

$$\mathbb{P}_{\delta_0 \otimes \mu_0}^{\dagger,00} \left(\lim_{t \rightarrow 0} X^{\dagger,00}(t) = X^{\dagger,00}(0) \text{ with respect to the norm-topology} \right) = 1. \quad (2.79)$$

To show this, recall that by the above and **Step OP1.9** we have

$$X^{\dagger,00}|_{E'} = X', \quad E^{\dagger,00} \setminus E' \text{ is polar}, \quad X' = X^{\dagger}|_{E'}, \quad X^{\dagger} = \check{X}^0|_{E^{\dagger}}, \quad \check{X}^0 = \tilde{X}|_{\tilde{E}_0},$$

whilst by **Step OP1.8** we have

$$\tilde{\mathbb{P}}_{\delta_s \otimes \mu_s} \circ \left[(\tilde{X}(t))_{t \in [0, N-s]} \right]^{-1} = \mathbb{P}_{\delta_s \otimes \mu_s}^N \circ \left[(X^N(t))_{t \in [0, N-s]} \right]^{-1}, \quad 0 < s < N.$$

Consequently, we obtain

$$\mathbb{P}_{\delta_s \otimes \mu_s}^{\dagger, 00} \circ \left[(X^{\dagger, 00}(t))_{t \in [0, N-s]} \right]^{-1} = \mathbb{P}_{\delta_s \otimes \mu_s}^N \circ \left[(X^N(t))_{t \in [0, N-s]} \right]^{-1}, \quad 0 < s < N,$$

hence, by Lemma 2.20, there exists $\eta^N \in \mathcal{P}(C([0, N]; \mathbb{R}^d))$ such that

$$\eta^N|_{C([\varepsilon, N]; \mathbb{R}^d)} = \mathbb{P}_{\delta_\varepsilon \otimes \mu_\varepsilon}^{\dagger, 00} \circ \left[(X^{\dagger, 00}(t - \varepsilon))_{t \in [\varepsilon, N]} \right]^{-1} \quad \text{as probabilities on } C([\varepsilon, N]; \mathbb{R}^d), \quad \text{for every } \varepsilon \in (0, N).$$

Furthermore, using (2.61) we get

$$U_\alpha^{\dagger, 00} f(0, x) = \int_0^\infty e^{-\alpha t} \eta_t^x(f(t, \cdot)) dt, \quad \alpha > 0, f \in b\mathcal{B}(E^{\dagger, 00}). \quad (2.80)$$

Consequently,

$$\delta_0 \otimes \mu_0 (U_\alpha^{\dagger, 00} f) = \int_0^\infty e^{-\alpha t} \delta_t \otimes \mu_t(f) dt, \quad \alpha > 0,$$

which means that

$$\mathbb{P}_{\delta_0 \otimes \mu_0}^{\dagger, 00} \circ [X^{\dagger, 00}(t)]^{-1} = \delta_t \otimes \mu_t, \quad t \geq 0.$$

Consequently, using also that $X^{\dagger, 00}$ is a (right) Markov process, we get

$$\begin{aligned} \mathbb{P}_{\delta_0 \otimes \mu_0}^{\dagger, 00} \circ \left[(X_2^{\dagger, 00}(t))_{t \in [\varepsilon, N]} \right]^{-1} &= \mathbb{P}_{\delta_\varepsilon \otimes \mu_\varepsilon}^{\dagger, 00} \circ \left[(X_2^{\dagger, 00}(t - \varepsilon))_{t \in [\varepsilon, N]} \right]^{-1} \\ &= \eta^N|_{C([\varepsilon, N]; \mathbb{R}^d)} \quad \text{on } C([\varepsilon, N]; \mathbb{R}^d) \quad 0 < \varepsilon < N, \end{aligned}$$

Getting rid of ε , the above can be rewritten as

$$\mathbb{P}_{\delta_0 \otimes \mu_0}^{\dagger, 00} \circ \left[(X_2^{\dagger, 00}(t))_{t \in (0, N)} \right]^{-1} = \eta^N|_{C((0, N); \mathbb{R}^d)}, \quad \text{for some (in fact all) } N > 0. \quad (2.81)$$

Since $\eta^N \in \mathcal{P}(C([0, N]; \mathbb{R}^d))$, the above (2.81) shows that

$$\exists Z : \Omega \rightarrow \mathbb{R}^d \text{ which is } \mathcal{F}^{\dagger, 00}(0)\text{-measurable such that } \lim_{t \rightarrow 0} \|X_2^{\dagger, 00}(t) - Z\| = 0 \quad \mathbb{P}_{\delta_0 \otimes \mu_0}^{\dagger, 00}\text{-a.e.} \quad (2.82)$$

Moreover, by (2.80) it is easy to deduce that

$$X_1^{\dagger, 00}(t) = t, \quad t \geq 0, \quad \mathbb{P}_{(0, x)}^{\dagger, 00}\text{-a.e. for every } (0, x) \in E^{\dagger, 00}.$$

Hence, in order to finish **Step OP2.7**, we need to prove the identification

$$X_2^{\dagger, 00}(0) = Z \quad \mathbb{P}_{\delta_0 \otimes \mu_0}^{\dagger, 00}\text{-a.e.} \quad (2.83)$$

To prove this, let us set for each $z \in \mathbb{R}^d$

$$F_z \in C_b(E^{\dagger, 00}), \quad F_z(t, x) := \|z - x\| \wedge 1, \quad (t, x) \in \mathbb{R}^d,$$

and note that by (2.62) we have

$$\lim_{\alpha \rightarrow \infty} \alpha U_\alpha^{\dagger, 00} F_z(0, x) = F_z(0, x), \quad (0, x) \in E^{\dagger, 00}.$$

Hence, using that $X^{\dagger,00}$ is a Markov process with transition function $(P_t^{\dagger,00})_{t \geq 0}$ and resolvent $\mathcal{U}^{\dagger,00}$, we have

$$\lim_{\alpha \rightarrow \infty} \int_0^\infty e^{-\alpha t} \mathbb{E}_{\delta_0 \otimes \mu_0}^{\dagger,00} \left\{ \left\| X_2^{\dagger,00}(t) - X_2^{\dagger,00}(0) \right\| \wedge 1 \right\} dt \quad (2.84)$$

$$= \lim_{\alpha \rightarrow \infty} \int_0^\infty e^{-\alpha t} \mathbb{E}_{\delta_0 \otimes \mu_0}^{\dagger,00} \left\{ F_{X_2^{\dagger,00}(0)}(X^{\dagger,00}(t)) \right\} dt \quad (2.85)$$

$$= \lim_{\alpha \rightarrow \infty} \int_0^\infty e^{-\alpha t} \int_{E^{\dagger,00}(0)} [P_t^{\dagger,00} F_x]((0, x)) \mu_0(dx) \quad (2.86)$$

$$= \lim_{\alpha \rightarrow \infty} \int_{E^{\dagger,00}(0)} \alpha [U_\alpha^{\dagger,00} F_x]((0, x)) \mu_0(dx) \quad (2.87)$$

$$= \int_{E^{\dagger,00}(0)} \alpha F_x((0, x)) \mu_0(dx) = 0. \quad (2.88)$$

The above convergence entails that

$$\exists (t_n)_n \searrow 0 \quad \text{such that} \quad \lim_n \mathbb{E}_{\delta_0 \otimes \mu_0}^{\dagger,00} \left\{ \left\| X_2^{\dagger,00}(t_n) - X_2^{\dagger,00}(0) \right\| \wedge 1 \right\} = 0,$$

hence by passing to a subsequence, we get

$$\exists (t_n)_n \searrow 0 \quad \text{such that} \quad \lim_n \left\| X_2^{\dagger,00}(t_n) - X_2^{\dagger,00}(0) \right\| \wedge 1 = 0, \quad \mathbb{P}_{\delta_0 \otimes \mu_0}^{\dagger,00}\text{-a.e.}$$

Thus, using (2.82) and (2.4.3), we get (2.83), hence **Step OP2.7** is finished.

Step OP2.8: Let now

$$E_0^{\dagger,00} := \left\{ (s, x) \in E^{\dagger,00} : \mathbb{P}_{(s,x)}^{\dagger,00}([0, \infty) \ni t \mapsto X^{\dagger,00}(t) \in E^{\dagger,00} \text{ is norm-continuous}) = 1 \right\},$$

and notice that from the previous steps we have that

$$E_0^{\dagger,00}(t) = E^\dagger(t), \quad \text{for every } t > 0.$$

Using now (2.76) and (2.77), we can apply Proposition 6.39 to deduce that $E_0^{\dagger,00}$ is $\delta_0 \otimes \mu_0$ -inessential and satisfies

$$E_0^{\dagger,00}(0) \in \mathcal{B}(\mathbb{R}^d), \quad E_0^{\dagger,00}(0) \subset E^{\dagger,00}(0), \quad \mu_0(E_0^{\dagger,00}(0)) = 1.$$

Step OP2.9: Finalizing the proof of Theorem 2.15. Based on the previous step, let us consider the right process

$$X = (\Omega, \mathcal{F}, \mathcal{F}(t), X(t), \theta(t), \mathbb{P}_{(s,x)}, (s, x) \in E),$$

with transition semigroup $(P_t)_{t \geq 0}$ and resolvent $\mathcal{U} = (U_\alpha)_{\alpha > 0}$, obtained by taking the restriction of $X^{\dagger,00}$ from $E^{\dagger,00}$ to $E := E_0^{\dagger,00}$ as in Proposition 6.35-Definition 6.36. Then clearly X has all the desired properties in Theorem 2.12, with the mention that assertion (ii.5) from Theorem 2.12 is at the moment satisfied for every $(s, x) \in E$ with $s > 0$. To extend (ii.5) to $s = 0$ we first notice that

$$\delta_0 \otimes \mu_0 (U_\alpha |\bar{L}f|) < \infty, \quad f \in C_c^\infty((0, \infty) \times \mathbb{R}^d), \quad \alpha > 0. \quad (2.89)$$

Then, by the same arguments as in **Step OP1.3** that follow after (2.57), we can find a further set $\tilde{E} \subset E$ such that $E \setminus E_0$ is $\delta_0 \otimes \mu_0$ -inessential, and M_f given by (2.55) with X^\dagger replaced by X , is well defined and integrable under $\mathbb{P}_{(0,x)}$ for every $(0, x) \in \tilde{E}$ and $f \in C_c^\infty((0, \infty) \times \mathbb{R}^d)$. Moreover, $[0, 1] \ni t \mapsto M_f(t)$ is right continuous in $L^1(\mathbb{P}_{\delta_{(0,x)}})$. Consequently, using the fact that $X(s) \in E^{\dagger,00}$, $s > 0$, for every $t > s > 0$

$$\begin{aligned} \mathbb{E}_{(0,x)}(M_f(t)) &= \mathbb{E}_{(0,x)} \{ M_f(s) + M_f(t) \circ \theta(s) \} = \mathbb{E}_{(0,x)} \{ M_f(s) \} + \mathbb{E}_{(0,x)} \{ \mathbb{E}_{X(s)} \{ M_f(t) \} \} \\ &= \mathbb{E}_{(0,x)} \{ M_f(s) \}. \end{aligned}$$

By the above mentioned right continuity in $L^1(\mathbb{P}_{\delta_{(0,x)}})$, and using that we clearly have $\mathbb{E}_{(0,x)} \{ M_f(0) \} = 0$, it follows by Lemma 6.54 that M_f is a martingale under $\mathbb{P}_{\delta_{(0,x)}}$ for every $(0, x) \in \tilde{E}$. Taking the restriction of the process X from E to \tilde{E} we can thus ensure (ii.5) from Theorem 2.12.

To conclude, (ii.1') and (ii.4') in Theorem 2.15 are clearly satisfied, whilst (ii.6') follows similarly to **Step OP1.16**, (ii.6).

3 Consequences of the main results

In this part we use the right processes constructed in Section 2, Theorem 2.12 and Theorem 2.15, in order to deduce new results about time-dependent linear PDEs. More precisely: We construct fundamental flow solutions; we solve parabolic Dirichlet problems by probabilistic means; by means of Lyapunov functions we derive maximal tail estimates for the path-space law obtained by the superposition of the fundamental flow solution to the time-dependent linear Fokker-Planck equation (2.13); we solve Fokker-Planck equations perturbed by nonlocal operators.

3.1 Existence of fundamental flow solutions to linear Fokker-Planck equations

Let $(L_t)_{t \geq 0}$ be given by (2.1)-(2.2) and $(\mu_t)_{t \geq 0} \subset \mathcal{P}(\mathbb{R}^d)$ be a weakly continuous solution to the linear Fokker-Planck equation (2.13) on $[0, \infty)$.

Definition 3.1. A pair (E, Γ_μ) , where $E \in \mathcal{B}([0, \infty) \times \mathbb{R}^d)$ and Γ_μ is a mapping

$$[0, \infty) \times E \times \mathcal{B}(\mathbb{R}^d) \ni (t, s, x, A) \mapsto \Gamma_{\mu, t}^{s, x}(A) \in [0, 1],$$

is called a μ -restricted fundamental solution flow to the linear Fokker-Planck equation (2.3) if the following properties are fulfilled:

(i) We have $\mu_s(E(s)) = 1 = \Gamma_{\mu, t}^{s, x}(E(s+t))$ for every $t \geq 0, (s, x) \in E$.

(ii) Γ_μ is a probability kernel from $[0, \infty) \times E$ to \mathbb{R}^d

(iii) For every $t \geq 0, (s, x) \in E$ it holds that $(\Gamma_{\mu, t}^{s, x})_{t \geq 0}$ is a weakly continuous solution to the linear Fokker-Planck equation (2.3) on $[s, \infty)$, with initial condition $\Gamma_{\mu, 0}^{s, \nu_s} = \delta_x$,

(iv) Let us set

$$\Gamma_{\mu, t}^{s, \xi} := \int_{E(s)} \Gamma_{\mu, t}^{s, x} \xi(dx), \quad t, s \geq 0, \xi \in \mathcal{P}(E(s)).$$

For every $s \geq 0$ and $\nu_s \in \mathcal{P}(E(s))$, the following flow property holds:

$$\Gamma_{\mu, t+r}^{s, \nu_s} = \Gamma_{\mu, t+r}^{s+r, \Gamma_{\mu, r}^{s, \nu_s}}, \quad r, t \geq 0.$$

(v) For every $s \geq 0$ and $\nu_s \in \mathcal{P}(E(s))$ for which

$$\int_s^N \int_{B(0, N)} (\|b(t, x)\| + \|a(t, x)\|) \Gamma_{\mu, t}^{s, \nu_s}(dx) dt < \infty, \quad N > 0,$$

it holds that $(\Gamma_{\mu, t}^{s, \nu_s})_{t \geq 0}$ is a weakly continuous solution to the linear Fokker-Planck equation (2.3) on $[s, \infty)$.

Corollary 3.2. Let $(\mu_t := u(t, x)dx)_{t \geq 0}$ be a given solution to the linear Fokker-Planck equation (2.3), and assume that $\mathbf{H}_{\mathbf{a}, \mathbf{b}}^{\sqrt{\mathbf{u}}}$, $\mathbf{H}_{\lesssim \mathbf{u}}$, and \mathbf{H}_0^{TV} hold. Let

$$X = (\Omega, \mathcal{F}, \mathcal{F}(t), X(t) = (X_1(t), X_2(t)), \mathbb{P}_{s, x}, (s, x) \in E, t \geq 0)$$

be the right process provided by both Theorem 2.12 and Theorem 2.15. Further, set

$$\Gamma_\mu := (\Gamma_{\mu, t}^{s, x})_{t \geq 0, (s, x) \in E}, \quad \Gamma_{\mu, t}^{s, x}(A) := \mathbb{P}_{s, x}(X_2(t) \in A \cap E(t+s)), \quad A \in \mathcal{B}(\mathbb{R}^d), s, t \geq 0.$$

Then (E, Γ_μ) is a μ -restricted fundamental solution flow to the linear Fokker-Planck equation (2.3).

Proof. Note that (i) is automatically satisfied by Theorem 2.15, whilst (ii) follows immediately from the properties of $(P_t)_{t \geq 0}$.

Assertion (iii) can be deduced directly from Theorem 2.12-Theorem 2.15, (ii.5), and a similar observation as in Remark 2.4.

Assertion (iv) follows from the semigroup property of $(P_t)_{t \geq 0}$ and the fact that X_1 is the uniform motion to the right.

Finally, (v) follows by Theorem 2.15, (ii.6'). \square

Remark 3.3. *Under the hypotheses of Corollary 3.2, it follows that for every $s \geq 0$ and $\nu_s \in \mathcal{P}(E(s))$ such that there exists a constant $c \in (0, \infty)$ with $\nu_s \leq \mu_s$, it holds that it holds that $(\Gamma_{\mu, t}^{s, \nu_s})_{t \geq 0}$ is the unique weakly continuous solution to the linear Fokker-Planck equation (2.13) on $[s, \infty)$ which satisfies $\Gamma_{\mu, t}^{s, \nu_s} \leq c\mu_{t+s}$ for every $t \geq 0$.*

3.2 Backward Kolmogorov equations and the parabolic Dirichlet problem

Let $(\mu_t := u(t, x)dx)_{t \geq 0}$ be a given solution to the linear Fokker-Planck equation (2.3), and assume that $\mathbf{H}_{\mathbf{a}, \mathbf{b}}^{\sqrt{u}}$, $\mathbf{H}_{\lesssim u}$, and \mathbf{H}_0^{TV} hold. Also, recall that \mathbf{L}_t is given by (2.1).

In this part, given an open subset $D \subset \mathbb{R}^d$ and a time horizon $0 < T < \infty$, we are interested in giving a meaning of and constructing a solution to the following final value problem

$$\left\{ \begin{array}{ll} \frac{d}{dt} \rho(t, x) + \mathbf{L}_t \rho(t, x) = 0, & (t, x) \in (0, T) \times D, \text{ dt} \otimes \mu_t(dx)\text{-a.e.} \\ \rho(t, x) = F(t, x), & (t, x) \in (0, T) \times \partial D, \\ \rho(T, x) = F(T, x), & x \in D, \end{array} \right. \quad (3.1)$$

where $F \in C_c^2([0, \infty) \times \mathbb{R}^d)$ is given.

In order to solve (3.1), we adopt the strategy developed by E.B. Dynkin in [42] in order to treat the so called *stochastic Dirichlet problem*; see also [56, Section 9], or [10]. Crucially, we shall rely on the regularity properties of the balayage operator, as derived in Lemma 6.23 and Proposition 6.24 in Section 6, which essentially allow us to rigorously employ Doob's maximal inequality for continuous-time martingales, by ensuring the right-continuity of the trajectories. Let us point out that such a regularity argument seems to have been overlooked in the proof of [56, Theorem 9.2.5] in order to guarantee that, using the notation from [56], the martingale $(N_t)_t$ is right-continuous.

Proposition 3.4. *Let $(\mu_t := u(t, x)dx)_{t \geq 0}$ be a given solution to the linear Fokker-Planck equation (2.3), and assume that $\mathbf{H}_{\mathbf{a}, \mathbf{b}}^{\sqrt{u}}$, $\mathbf{H}_{\lesssim u}$, and \mathbf{H}_0^{TV} hold. Further, let*

$$X = (\Omega, \mathcal{F}, \mathcal{F}(t), X(t) = (X_1(t), X_2(t)), \mathbb{P}_{s, x}, (s, x) \in E, t \geq 0)$$

be the right process provided by both Theorem 2.12 and Theorem 2.15, $D \subset \mathbb{R}^d$ be open, $T \in (0, \infty)$, and $F \in b\mathcal{B}([0, \infty) \times \mathbb{R}^d)$ for which there exists $N > 0$ such that $F(s, \cdot) = 0, s \geq N$. Furthermore, let us set

$$\tau := \inf \{ t > 0 : X(t) \in [0, \infty) \times \mathbb{R}^d \setminus (0, T) \times D \}, \text{ and note that } \tau \leq T - s \text{ } \mathbb{P}_{s, x}\text{-a.s. for } s \in [0, T],$$

$$\rho(t, x) := \mathbb{E}_{t, x} \{ F(X(\tau)) \}, \quad t \in [0, \infty), x \in E(t).$$

The following assertions hold.

(i) ρ is a solution to problem (3.1) in the following sense:

(i.1) ρ is $\mathcal{B}(E)$ -measurable, finely continuous on $[[0, T) \times D] \cap E$ (it is finely continuous on E if F is finely continuous), and under $\mathbb{P}_{s, x}, (s, x) \in [[0, T) \times D] \cap E$,

$$(\rho(X(t \wedge \tau))_{t \geq 0}) \text{ is a càdlàg } (\mathcal{F}(t \wedge \tau))\text{-martingale} \quad \text{and} \quad \lim_{t \nearrow \tau} \rho(X(t)) = F(X(\tau)) \text{ a.s.}$$

Moreover, we have $\rho(T, x) = F(T, x), \quad x \in D \text{ } \mu_T\text{-a.e.}$

(i.2) If F is finely continuous and $F \in \mathbf{D}(\bar{\mathbf{L}}_p)$ for some $p \in (1, \infty)$ then $\rho \in \mathbf{D}(\bar{\mathbf{L}}_q)$ and

$$\mathbf{L}_q \rho = 0 \quad \text{on } [0, T] \times D \text{ } \bar{\mu}\text{-a.e. for every } 1 \leq q \leq p.$$

(ii) ρ is the unique solution to problem (3.1) in both of the following meanings:

(ii.1) Let $v \in \mathbf{b}\mathcal{B}(E)$ be such that under $\mathbb{P}_{s,x}$, $(s, x) \in [[0, T] \times D] \cap E$,

$$(v(X(t \wedge \tau)))_{t \geq 0} \text{ is a càdlàg } (\mathcal{F}(t \wedge \tau))\text{-martingale and } \lim_{t \nearrow \tau} v(X(t)) = F(X(\tau)) \text{ a.s.}$$

Then $v(s, x) = \rho(s, x)$ for every $(s, x) \in [[0, T] \times D] \cap E$.

(ii.2) Let $v \in \mathbf{b}\mathcal{B}(E)$ be finely continuous such that

$$v \in \mathbf{D}(\bar{\mathbf{L}}_1), \quad \mathbf{L}_1 v = 0 \quad \text{on } [0, T] \times D \text{ } \bar{\mu}\text{-a.e., and } \lim_{t \nearrow \tau} v(X(t)) = F(X(\tau)) \quad \mathbb{P}_{\bar{\mu}|_{(0, T) \times D}}\text{-a.e.} \quad (3.2)$$

Then there exists a set $E' \in \mathcal{B}(E)$ such that $E \setminus E'$ is $\bar{\mu}$ -inessential and $v(t, x) = \rho(t, x)$ for every $(t, x) \in [[0, T] \times D] \cap E'$.

Proof. (i.1). Note first that in terms of the balayage operator B^A given by Definition 6.17 and Definition 6.21, we have the identification

$$\rho(t, x) = (B^A F)(t, x), \quad t \in [0, T], x \in E(t), \quad \text{where } A = E \setminus (0, T) \times D.$$

The fact that ρ is \mathcal{B} -measurable and finely continuous on $[[0, T] \times D] \cap E$ (and finely continuous on E if F is finely continuous) follow from Proposition 6.44 and Lemma 6.23. Further, note that X is clearly a Hunt process (being a conservative right process with continuous paths), F is finely continuous (since the norm topology is a natural topology), whilst τ is a predictable stopping time under $\mathbb{P}_{s,x}$ for $(s, x) \in [[0, T] \times D] \cap E$ since if we define

$$\tau_k := \inf \{ t > 0 : X(t) \in [0, \infty) \times \mathbb{R}^d \setminus D_k \}, \quad \text{with } \bar{D}_k \subset D_{k+1}, D_k \nearrow [0, T] \times D, k \geq 1$$

then, under $\mathbb{P}_{s,x}$ for $(s, x) \in [[0, T] \times D] \cap E$, and using the path-continuity of X , there exists k_0 such that $\tau_k < \tau$ for $k \geq k_0$ and $\tau_k \nearrow \tau$ a.s. Consequently, the first part of assertion (i.1) follows from Proposition 6.24, (v).

To prove the second part, recall that $\mu_T(E(T)) = 1$. By the path-continuity of X and the fact that X_1 is the uniform motion to the right, we get

$$\tau = 0 \quad \mathbb{P}_{T,x}\text{-a.s. } x \in E(T) \cap D.$$

(i.1) is now completely proved.

(i.2). First of all, since there exists $N > 0$ such that $F(s, \cdot) = 0, s \geq N$, by the Markov property we easily get

$$\begin{aligned} \bar{\mu}(|\rho|^q) &= \int_0^\infty \delta_0 \otimes \mu_0 (P_t |\rho|^q) dt \leq \int_0^\infty \delta_0 \otimes \mu_0 (\mathbb{E}_{(\cdot, \cdot)} \{|F|^q(t + \tau \circ \theta(t), X_2(t + \tau \circ \theta(t)))\}) dt \\ &= \int_0^N \delta_0 \otimes \mu_0 (\mathbb{E}_{(\cdot, \cdot)} \{|F|^q(t + \tau \circ \theta(t), X_2(t + \tau \circ \theta(t)))\}) dt \\ &\leq N \|F\|_\infty^q, \quad 1 \leq q < \infty. \end{aligned}$$

Thus, $\rho \in L^q(\bar{\mu})$ for every $q \geq 1$.

Furthermore, by the strong Markov property and using that τ is a *terminal* time ([60, p.65]), hence $\tau = t + \tau \circ \theta(t)$ on $\tau > t$ a.s. for every $t \geq 0$, we have

$$\begin{aligned} P_t \rho(s, x) &= \mathbb{E}_{s,x} \{ \mathbb{E}_{X(t)} \{ F(X(\tau)) \} \} = \mathbb{E}_{s,x} \{ F(X(t + \tau \circ \theta(t))) \} \\ &= \mathbb{E}_{s,x} \{ F(X(\tau)); \tau > t \} + \mathbb{E}_{s,x} \{ F(X(t + \tau \circ \theta(t))); \tau \leq t \} \\ &= \rho(s, x) + \mathbb{E}_{s,x} \{ F(X(t + \tau \circ \theta(t))) - F(X(\tau)); \tau \leq t \}, \quad (s, x) \in E. \end{aligned}$$

Since $t + (\tau \wedge t) \circ \theta(t)$ is also a stopping time (see e.g. [33, Theorem (8.6)]), and using that X solves the martingale problem (see Theorem 2.12 and Theorem 2.15), we deduce that

$$\begin{aligned}
& \mathbb{E}_{s,x} \{F(X(t + \tau \circ \theta(t))) - F(X(\tau)); \tau \leq t\} \\
&= \mathbb{E}_{s,x} \{F(X(t + (\tau \wedge t) \circ \theta(t))) - F(X(\tau \wedge t)); \tau \leq t\} \\
&= \mathbb{E}_{s,x} \{ \mathbb{E}_{s,x} \{F(X(t + (\tau \wedge t) \circ \theta(t))) - F(X(\tau \wedge t)) | \mathcal{F}(\tau \wedge t)\}; \tau \leq t\} \\
&= \mathbb{E}_{s,x} \left\{ \mathbb{E}_{s,x} \left\{ \int_{\tau \wedge t}^{t + (\tau \wedge t) \circ \theta(t)} \bar{\Gamma} F(X(r)) dr \mid \mathcal{F}(\tau \wedge t) \right\}; \tau \leq t \right\} \\
&= \mathbb{E}_{s,x} \left\{ \int_{\tau \wedge t}^{t + (\tau \wedge t) \circ \theta(t)} \bar{\Gamma} F(X(r)) dr; \tau \leq t \right\}, \quad (s, x) \in E.
\end{aligned}$$

Thus, corroborating the above we get

$$\begin{aligned}
\left| \frac{P_t \rho(s, x) - \rho(s, x)}{t} \right| &\leq \mathbb{E}_{s,x} \left\{ \frac{1}{t} \int_0^{2t} |\bar{\Gamma} F|(X(r)) dr; \tau \leq t \right\} \\
&\leq \frac{1}{t} \int_0^{2t} P_s(|\bar{\Gamma} F|) dr,
\end{aligned}$$

hence, by e.g. Jensen's inequality,

$$\begin{aligned}
\bar{\mu} \left(\left| \frac{P_t \rho(s, x) - \rho(s, x)}{t} \right|^p \right) &\leq 2^{p-1} \frac{1}{t} \int_0^{2t} \bar{\mu} \left(P_r(|\bar{\Gamma} F|^p) \right) dr \\
&\leq 2^p \bar{\mu} \left(|\bar{\Gamma} F|^p \right).
\end{aligned}$$

Since by Theorem 2.12 we have $\bar{\Gamma} F \in L^p(\bar{\mu})$, we deduce that

$$\sup_{t \in [0,1]} \left\| \frac{P_t \rho - \rho}{t} \right\|_{L^p(\bar{\mu})} < \infty.$$

But now it is a standard fact in the theory of strongly continuous semigroups that $\rho \in \mathbf{D}(\bar{\Gamma}_p)$; see e.g. the discussion before Proposition 4.4 in [20]. Furthermore, note that $\bar{\Gamma}_p \rho = 0$ on $(N, \infty) \times \mathbb{R}^d$ $\bar{\mu}$ -a.e., where N has been in the beginning of the proof of (iii). Thus, since $\bar{\mu}([0, N] \times \mathbb{R}^d) = N < \infty$, we deduce that $\bar{\Gamma}_p \rho \in L^1(\bar{\mu})$. Since we also have $\rho \in L^1(\bar{\mu})$, we get that $\rho \in \mathbf{D}(\bar{\Gamma}_1)$ as well.

Moreover, using Hölder's inequality,

$$\begin{aligned}
\left| \frac{P_t \rho(s, x) - \rho(s, x)}{t} \right| &\leq \mathbb{E}_{s,x} \left\{ \frac{1}{t} \int_0^{2t} |\bar{\Gamma} F|(X(r)) dr; \tau \leq t \right\} \\
&\leq 2^{1/p} \left[\mathbb{E}_{s,x} \left\{ \frac{1}{t} \int_0^{2t} |\bar{\Gamma} F|^p(X(r)) dr \right\} \right]^{1/p} [\mathbb{P}_{s,x}(\tau \leq t)]^{(p-1)/p} \\
&= 2^{1/p} \left[\frac{1}{t} \int_0^{2t} P_s(|\bar{\Gamma} F|^p) dr \right]^{1/p} [\mathbb{P}_{s,x}(\tau \leq t)]^{(p-1)/p}.
\end{aligned}$$

Consequently,

$$\begin{aligned}
& \int_{(0,T) \times D} \left| \frac{P_t \rho(s, x) - \rho(s, x)}{t} \right| \bar{\mu}(dt, dx) \\
&\leq 2^{1/p} \bar{\mu} \left(\left[\frac{1}{t} \int_0^{2t} P_s(|\bar{\Gamma} F|^p) dr \right]^{1/p} \left(\int_{(0,T) \times D} \mathbb{P}_{s,x}(\tau \leq t) \bar{\mu}(dt, dx) \right)^{(p-1)/p} \right) \\
&\leq 2^{2/p} \bar{\mu} \left(|\bar{\Gamma} F|^p \right)^{1/p} \left(\int_{[(0,T) \times D] \cap E} \mathbb{P}_{s,x}(\tau \leq t) \bar{\mu}(dt, dx) \right)^{(p-1)/p}
\end{aligned}$$

which converges to 0 by dominated convergence since due to Remark 6.20 we have

$$\mathbb{P}_{s,x}(\tau > 0) = 1, \quad (s, x) \in [(0, T) \times D] \cap E.$$

This clearly completes the proof of (i.2).

(ii.1). Let $(\tau_k)_{k \geq 1}$ be as in the proof of (i.1) and $\tau = 0$, so that from the hypothesis of (ii.2) we deduce that, under $\mathbb{P}_{s,x}$, $(s, x) \in [(0, T) \times D] \cap E$, $(v(X(\tau_k)))_{k \geq 0}$ is a $(\mathcal{F}(\tau_k))$ -martingale which is closed by $F(X(\tau))$. Therefore,

$$v(s, x) = \mathbb{E}_{s,x} \{v(X(0))\} = \mathbb{E}_{s,x} \{F(X(\tau))\}, \quad (s, x) \in [(0, T) \times D] \cap E,$$

which proves the claim.

(ii.2). Let $v : E \rightarrow \mathbb{R}$ be $\mathcal{B}(E)$ -measurable, bounded, and finely continuous such that (3.2) is satisfied, and consider the process

$$M(t) := v(X(t)) - v(X(0)) - \int_0^t \bar{\mathbf{L}}_1 v(X(s)) ds, \quad t \geq 0,$$

which, under $\mathbb{P}_{\bar{\mu}}$, is well defined; in particular it does not depend on the $\bar{\mu}$ -version of $\bar{\mathbf{L}}_1 v$. Let us take $g : E \rightarrow \mathbb{R}$ $\mathcal{B}(E)$ -measurable such that $g(x) = 0, x \in (0, T) \times D$ and $g = \bar{\mathbf{L}}_1 v$ $\bar{\mu}$ -a.e. and consider the process

$$\tilde{M}(t) := v(X(t)) - v(X(0)) - \int_0^t g(X(s)) ds, \quad t \geq 0.$$

Since $P_t v = v + \int_0^t P_s g ds$ $\bar{\mu}$ -a.e., it follows by Lemma 6.54 and Proposition 6.56 that there exists a set $E' \in \mathcal{B}(E)$ such that $E \setminus E'$ is $\bar{\mu}$ -inessential and $(\tilde{M}(t))_{t \geq 0}$ is a well-defined right-continuous $\mathcal{F}(t)$ -martingale under $\mathbb{P}_{s,x}$ for every $(s, x) \in E'$. Moreover, since $g(x) = 0, x \in (0, T) \times D$, we have $\mathbb{P}_{s,x}$ -a.s. for $(s, x) \in [(0, T) \times D] \cap E'$ that

$$v(s, x) = v(X(t)) - \tilde{M}(t), \quad t \in [0, \tau],$$

hence by (3.2) we get

$$v(s, x) = F(X(\tau)) - \tilde{M}(\tau -).$$

Note that since X is in fact a Hunt process, it follows from [60, Theorem (47.6) and Exercise (47.7)], that $(\mathcal{F}(t))_{t \geq 0}$ is *quasi-left continuous* (see [60] for details), hence by e.g. [58, pp. 191-192], no $(\mathcal{F}(t))$ -martingale can jump at predictable times. Thus, $\mathbb{P}_{s,x}$ -a.s. for $(s, x) \in [(0, T) \times D] \cap E'$ we have

$$\tilde{M}(\tau -) = \tilde{M}(\tau) \quad \text{and} \quad F(X(\tau)) = v(X(\tau)).$$

Since τ is bounded a.s. we get $\mathbb{P}_{s,x}$ -a.s. for $(s, x) \in [(0, T) \times D] \cap E'$ that $\mathbb{E}_{s,x} \{M(\tau)\} = 0$, hence finally

$$v(s, x) = \mathbb{E}_{s,x} \{F(X(\tau))\} (= \rho(s, x)).$$

□

Remark 3.5. Let $(\mu_t := u(t, x)dx)_{t \geq 0}$ be a given solution to the linear Fokker-Planck equation (2.3), and assume that $\mathbf{H}_{\mathbf{a}, \mathbf{b}}^{\vee \bar{\mathbf{u}}}$, $\mathbf{H}_{\lesssim \mathbf{u}}$, and $\mathbf{H}_{\mathbf{0}}^{\text{TV}}$ hold. Further, let $T > 0$, $F \in C_c^2(\mathbb{R}^d)$, and ρ be the solution from Proposition 3.4 for $D = \mathbb{R}^d$. Then

$$\int_{\mathbb{R}^d} F(x) \mu_T(dx) = \int_{\mathbb{R}^d} \rho(0, x) \mu_0(dx). \quad (3.3)$$

To see this, first of all note that since $D = \mathbb{R}^d$ whilst $(X_1(t))_{t \geq 0}$ is the uniform motion to the right as ensured by Theorem 2.12, it follows that τ given in Proposition 3.4 satisfies

$$\tau = T - t \quad \mathbb{P}_{t,x}\text{-a.s.} \quad t \in [0, T], x \in E(t)$$

Thus,

$$\rho(t, x) := \mathbb{E}_{t,x} \{F(X_2(T - t))\}, \quad t \in [0, T], x \in E(t),$$

and the claim follows by Theorem 2.15, (ii.4').

From the numerical standpoint, (3.3) indicates that in order to compute the integral $\mu_T(F)$, given μ_0, T , and F , one could attempt to solve backward in time equation (3.1) in order to compute $\rho(0, \cdot)$, and then use (3.3). However, solving (3.1) should be carefully understood in the μ -restricted sense given by Proposition 3.4; from a probabilistic numerical perspective, one should be able to reverse in time the diffusion $(X_2(t))_{t \in [0, T]}$.

3.3 Lyapunov functions

Let $(\mu_t := u(t, x)dx)_{t \geq 0}$ be a given solution to the linear Fokker-Planck equation (2.3), and assume that $\mathbf{H}_{\mathbf{a}, \mathbf{b}}^{\sqrt{u}}$, $\mathbf{H}_{\leq u}$, and \mathbf{H}_0^{TV} hold. Further, let X be the process provided by Theorem 2.12 and Theorem 2.15, with transition function $(P_t)_{t \geq 0}$ and resolvent $\mathcal{U} = (U_\alpha)_{\alpha > 0}$. In this part we are interested in deriving general maximal tail estimates of the type

$$\mathbb{P}_{(s,x)} \left(\sup_{t \in [0, T]} V(X_2(t)) \geq \varepsilon \right) \leq \frac{e^{cT}}{\varepsilon} V(x),$$

where V is a suitable Lyapunov function. Such estimates have been obtained in [35] with $\mathbb{P}_{(s,x)}$ replaced by $\mathbb{P}_{\delta_0 \otimes \mu_0}$. Here, we are interested in obtaining such estimates for every $(s, x) \in E$; see Proposition 3.9 below.

Throughout, for a function $V \in p\mathcal{B}(\mathbb{R} \times \mathbb{R}^d)$, by $P_t V$ and $U_\alpha V$ we mean $P_t(V|_E)$ and $U_\alpha(V|_E)$, respectively.

We also introduce the following measures

$$\bar{\mu}_\alpha(dx, dt) := e^{-\alpha t} \mu_t(dx) dt, \quad \alpha > 0.$$

Proposition 3.6. *Assume that there exists a Borel measurable function $V : [0, \infty) \times \mathbb{R}^d \rightarrow [0, \infty]$ and a constant $\delta_V \in (0, \infty)$ such that*

$$\mu_0(V(0, \cdot)) < \infty \quad \text{and} \quad \mu_t(V(t, \cdot)) \leq e^{\delta_V t} \mu_0(V(0, \cdot)), \quad t \geq 0. \quad (3.4)$$

Then the following hold:

- (i) $\bar{\mu}_\beta(U_{\alpha+\beta} F) \leq \bar{\mu}(F) \wedge \frac{\bar{\mu}_\beta(F)}{\alpha}$, $\beta, \alpha > 0$, $F \in p\mathcal{B}([0, \infty) \times \mathbb{R}^d)$.
- (ii) $\delta_0 \otimes \mu_0(U_\alpha F) \leq \bar{\mu}_\alpha(F)$ $F \in p\mathcal{B}([0, \infty) \times \mathbb{R}^d)$.
- (iii) $\bar{\mu}_\beta(U_\alpha V) \leq \frac{\mu_0(V(0, \cdot))}{(\alpha - \delta_V)(\beta - \delta_V)}$, $\alpha, \beta > \delta_V$.
- (iv) $\delta_0 \otimes \mu_0(U_\alpha V) \leq \frac{\mu_0(V(0, \cdot))}{\alpha - \delta_V}$, $\alpha > \delta_V$.

Proof.

$$\begin{aligned}
\bar{\mu}_\beta(U_{\alpha+\beta}F) &= \int_0^\infty e^{-(\alpha+\beta)t} \bar{\mu}_\beta(P_t F(\cdot, \cdot)) dt = \int_0^\infty e^{-(\alpha+\beta)t} \int_0^\infty e^{-\beta s} \mu_s(P_t F(s, \cdot)) ds dt \\
&= \int_0^\infty e^{-(\alpha+\beta)t} \int_0^\infty e^{-\beta s} \mu_{t+s}(F(t+s, \cdot)) ds dt \\
&= \int_0^\infty e^{-\alpha t} \int_0^\infty e^{-\beta(t+s)} \mu_{t+s}(F(t+s, \cdot)) ds dt \\
&\leq \bar{\mu}(F) \wedge \left(\bar{\mu}_\beta(F) \int_0^\infty e^{-\alpha t} dt \right) \\
&= \bar{\mu}(F) \wedge \frac{\bar{\mu}_\beta(F)}{\alpha}, \quad \beta, \alpha > 0, F \in p\mathcal{B}([0, \infty) \times \mathbb{R}^d). \\
\delta_0 \otimes \mu_0(U_\alpha F) &= \int_0^\infty e^{-\alpha t} \mu_0(P_t F(0, \cdot)) dt = \int_0^\infty e^{-\alpha t} \mu_t(F(t, \cdot)) dt = \bar{\mu}_\alpha(F), \quad F \in p\mathcal{B}([0, \infty) \times \mathbb{R}^d) \\
\bar{\mu}_\beta(U_\alpha V) &= \int_0^\infty e^{-\alpha t} \bar{\mu}_\beta(P_t V(\cdot, \cdot)) dt = \int_0^\infty e^{-\alpha t} \int_0^\infty e^{-\beta s} \mu_s(P_t V(s, \cdot)) ds dt \\
&= \int_0^\infty e^{-\alpha t} \int_0^\infty e^{-\beta s} \mu_{t+s}(V(t+s, \cdot)) ds dt \\
&\leq \mu_0(V(0, \cdot)) \int_0^\infty e^{-(\alpha-\delta_V)t} \int_0^\infty e^{-(\beta-\delta_V)s} ds dt \\
&= \frac{\mu_0(V(0, \cdot))}{(\alpha-\delta_V)(\beta-\delta_V)}, \quad \alpha, \beta > \delta_V. \\
\delta_0 \otimes \mu_0(U_\alpha V) &= \int_0^\infty e^{-\alpha t} \mu_0(P_t V(0, \cdot)) dt = \int_0^\infty e^{-\alpha t} \mu_t(V(t, \cdot)) dt \\
&\leq \mu_0(V(0, \cdot)) \int_0^\infty e^{-(\alpha-\delta_V)t} dt \\
&= \frac{\mu_0(V(0, \cdot))}{\alpha-\delta_V}, \quad \alpha > \delta_V.
\end{aligned}$$

□

Proposition 3.7. *Assume that there exists $0 \leq V_n \in C_b^2([0, \infty) \times \mathbb{R}^d)$, $n \geq 1$, and $\delta_V \in (0, \infty)$ such that*

$$\bar{L}V_n(s, x) \leq \delta_V V(s, x), \quad V(s, x) := \liminf_n V_n(s, x), \quad (s, x) \in E. \quad (3.5)$$

Then

$$P_t V \leq e^{\delta_V t} V, \quad t \geq 0. \quad (3.6)$$

Furthermore, if $\liminf_n (\bar{L}V_n)^- \geq \bar{L}V$, then

$$\int_0^t P_r |\bar{L}V| dr \leq V(2e^{\delta_V t} - 1), \quad t \geq 0. \quad (3.7)$$

Proof. Let us consider

$$\tau_k := \inf\{t > 0 : X(t) \in E \setminus [0, k] \times B(0, k)\}, \quad k \geq 1.$$

Since X has continuous paths and it is conservative, we deduce that

$$\tau_k \nearrow_k \infty \quad \mathbb{P}_{s,x}\text{-a.s.}, \quad (s, x) \in E. \quad (3.8)$$

Indeed, if we set $\tau_\infty := \sup_k \tau_k$, then, since X_1 is the uniform motion to the right,

$$\lim_k \|X_2(\tau_k)\| = \infty \quad \text{on } \tau_\infty < \infty \quad \text{a.s.}$$

But $\sup_{t \in [0, \tau_\infty]} \|X_2(t)\| < \infty$ on $\tau_\infty < \infty$ a.s., hence we must have (3.8).

Further, for each $n, k \geq 1$ there exists $0 \leq V_{n,k} \in C_c^2([0, \infty) \times \mathbb{R}^d)$ such that

$$V_{n,k-1} = V_n \quad \text{on } (-1/k, k) \times B(0, k). \quad (3.9)$$

On the one hand, by Theorem 2.12 and Theorem 2.15, for every $(s, x) \in E$ and $k \geq 1$ we have that the stopped process

$$V_{n,k}(X(t \wedge \tau_k)) - V_{n,k}(X(0)) - \int_0^{t \wedge \tau_k} \bar{\mathbb{L}}V_{n,k}(X(r)) dr, \quad t \geq 0 \quad (3.10)$$

is a continuous (\mathcal{F}_t) -martingale under $\mathbb{P}_{(s,x)}$. Consequently, using (3.9), taking expectations in (3.10), and using also (3.5), we have for every $(s, x) \in E$

$$\begin{aligned} \mathbb{E}_{(s,x)} \{V_{n,k}(X(t \wedge \tau_k))\} &= V_{n,k}(s, x) + \mathbb{E} \left\{ \int_0^{t \wedge \tau_k} \bar{\mathbb{L}}V_{n,k}(X(r)) dr \right\} \\ &= V_{n,k}(s, x) + \mathbb{E} \left\{ \int_0^{t \wedge \tau_k} \bar{\mathbb{L}}V_n(X(r)) dr \right\} \\ &\leq V_{n,k}(s, x) + \delta_V \mathbb{E} \left\{ \int_0^{t \wedge \tau_k} V(X(r)) dr \right\} \\ &\leq V_{n,k}(s, x) + \delta_V \int_0^t \mathbb{E} \{V(X(r \wedge \tau_k))\} dr, \quad t \geq 0, n, k \geq 1. \end{aligned}$$

On the other hand, note that

$$\liminf_n V_{n,k} = V \quad \text{on } [0, k) \times B(0, k), \quad k \geq 1.$$

Therefore, by Fatou's lemma we get for $k \geq 1$ and $(s, x) \in E \cap [0, k) \times B(0, k)$ that

$$\mathbb{E}_{(s,x)} \{V(X(t \wedge \tau_k))\} \leq V(s, x) + \delta_V \int_0^t \mathbb{E} \{V(X(r \wedge \tau_k))\} dr, \quad t \geq 0.$$

By Gronwall's lemma we thus have for $k \geq 1$ and $(s, x) \in E \cap [0, k) \times B(0, k)$ that

$$\mathbb{E}_{(s,x)} \{V(X(t \wedge \tau_k))\} \leq e^{\delta_V t} V(s, x), \quad t \geq 0.$$

Now, (3.6) follows by letting k go to infinity.

To prove (3.7), we use the second equality after (3.10) to get

$$\begin{aligned} \mathbb{E} \left\{ \int_0^{t \wedge \tau_k} (\bar{\mathbb{L}}V_n)^-(X(r)) dr \right\} &\leq V_{n,k}(s, x) + \mathbb{E} \left\{ \int_0^{t \wedge \tau_k} (\bar{\mathbb{L}}V_n)^+(X(r)) dr \right\} \\ &\leq V_{n,k}(s, x) + \delta_V \mathbb{E} \left\{ \int_0^{t \wedge \tau_k} V(X(r)) dr \right\}, \quad t \geq 0, n, k \geq 1. \end{aligned}$$

Letting k go to infinity we thus get

$$\begin{aligned} \mathbb{E}_{s,x} \left\{ \int_0^t (\bar{\mathbb{L}}V_n)^-(X(r)) dr \right\} &\leq V_n(s, x) + \delta_V \mathbb{E} \left\{ \int_0^t V(X(r)) dr \right\} \\ &\leq V_n(s, x) + \delta_V V \int_0^t e^{\delta_V s} ds \\ &= V_n(s, x) + V(e^{\delta_V t} - 1), \quad t \geq 0, n \geq 1. \end{aligned}$$

Letting now n go to infinity and using Fattou's lemma we get

$$\int_0^t P_r (\bar{\mathbb{L}}V)^- dr \leq e^{\delta_V t} V, \quad t \geq 0.$$

Since we clearly have

$$\int_0^t P_r (\bar{\mathbb{L}}V)^+ dr \leq (e^{\delta_V t} - 1)V, \quad t \geq 0,$$

relation (3.7) follows. \square

Proposition 3.8. *Let $V \in p\mathcal{B}([0, \infty) \times \mathbb{R}^d)$ and $\delta_V \in (0, \infty)$ such that*

$$\mu_0(V(0, \cdot)) < \infty, \quad P_t V \leq e^{\delta_V t} V, \quad t \geq 0.$$

Then $\mu_t(V(t, \cdot)) \leq e^{\delta_V t} \mu_0(V(0, \cdot))$, $t \geq 0$, and hence assertions (i)-(iii) from Proposition 3.6 hold.

Proof. First of all notice that

$$e^{-\delta_V t} P_t(V \wedge n) \leq V \wedge n, \quad t \geq 0, n \geq 1.$$

Consequently,

$$e^{-\delta_V t} \mu_0(P_t(V \wedge n)(0, \cdot)) \leq \mu_0(V(0, \cdot) \wedge n),$$

hence

$$e^{-\delta_V t} \mu_t(V(t, \cdot) \wedge n) \leq \mu_0(V(0, \cdot) \wedge n), \quad n \geq 1,$$

from which the statement follows. □

The main result of this subsection is the following.

Proposition 3.9. *Let V be such that*

$$0 \leq V \in C^2(\mathbb{R}^d), \quad \bar{\mathbb{L}}V \leq \delta_V V \quad \text{for some } \delta_V \in (0, \infty).$$

Then the following assertions hold:

(i) $P_t V \leq e^{\delta_V t} V, t \geq 0.$

(ii) $U_{\alpha + \delta_V} |\bar{\mathbb{L}}V| \leq 2 \frac{(\alpha + \delta_V)}{\alpha} V, \quad \alpha > 0.$

(iii) *For every $\varepsilon > 0$ we have*

$$\mathbb{P}_{(s,x)} \left(\sup_{t \in [0, T]} V(X_2(t)) \geq \varepsilon \right) \leq \frac{e^{\delta_V T}}{\varepsilon} V(x), \quad (s, x) \in E. \quad (3.11)$$

Proof. (i). For each $N \geq 1$ let $\psi_N \in C_b^\infty(\mathbb{R})$ such that

$$\psi_N(t) = t, t \leq N - 1, \quad \psi_N(t) = N, t > N + 1, \quad 0 \leq \psi'_N \leq 1, \quad \psi''_N \leq 0.$$

Now, as in [35, Proof of Proposition 2.2], we have

$$\bar{\mathbb{L}}\psi_N(V) = \psi'_N(V) \bar{\mathbb{L}}V + \psi''_N(V) \|\sqrt{a} \nabla V\|^2 \leq \psi'_N(V) \bar{\mathbb{L}}V,$$

from which we deduce, by setting $V_N := \psi_N(V) \in C_b^2(\mathbb{R}^d)$, $N \geq 1$, that

$$\bar{\mathbb{L}}V_N \leq \delta_V V, \quad N \geq 1.$$

By applying Proposition 3.7 we conclude that $P_t V \leq e^{\delta_V t} V, t \geq 0.$

(ii). We first note that

$$\lim_N \bar{\mathbb{L}}V_N = \bar{\mathbb{L}}V \quad \text{pointwise on } [0, \infty) \times \mathbb{R}^d.$$

Therefore, we can apply Proposition 3.7 to get (3.7), namely

$$\int_0^t P_r |\bar{\mathbb{L}}V| dr \leq V(2e^{\delta_V t} - 1), \quad t \geq 0.$$

Consequently, integrating by parts

$$\begin{aligned} \int_0^t e^{-(\alpha+\delta_V)r} P_r |\bar{\mathbb{L}}V| dr &= e^{-(\alpha+\delta_V)t} \int_0^t P_r |\bar{\mathbb{L}}V| dr + (\alpha + \delta_V) \int_0^t e^{-(\alpha+\delta_V)s} \int_0^s P_r |\bar{\mathbb{L}}V| dr \\ &\leq e^{-\alpha t} V + 2(\alpha + \delta_V) V \int_0^t e^{-\alpha s} ds, \\ &= e^{-\alpha t} V + 2 \frac{(\alpha + \delta_V)}{\alpha} V (1 - e^{-\alpha t}), \quad t \geq 0. \end{aligned}$$

Letting t go to infinity we get (3.7).

(iii). To prove (3.11), note that by the first part we have that V is δ_V -excessive, hence (see e.g. Theorem 6.13) $(e^{-\delta_V t} V(X_2(t)))_{t \geq 0}$ is a right-continuous $\mathcal{F}(t)$ -supermartingale under $\mathbb{P}_{(s,x)}$ for every $(s,x) \in E$. \square

Thus, by Doob's maximal inequality we get

$$\begin{aligned} \mathbb{P}_{(s,x)} \left(\sup_{t \in [0,T]} V(X_2(t)) \geq \varepsilon \right) &\leq \mathbb{P}_{(s,x)} \left(\sup_{t \in [0,T]} e^{-\delta_V t} V(X_2(t)) \geq e^{-\delta_V T} \varepsilon \right) \\ &\leq \frac{e^{\delta_V T}}{\varepsilon} V(x), \quad (s,x) \in E. \end{aligned}$$

Example. We follow [35, Example 2.3] and take

$$V(s,x) := V(x) = \log(1 + \|x\|^2), \quad (s,x) \in \mathbb{R} \times \mathbb{R}^d.$$

Assume further that there exists $C \in (0, \infty)$ such that

$$\|a(t,x)\| \leq C \|x\|^2 V(x) + C, \quad \langle b(t,x), x \rangle \leq C \|x\|^2 V(x) + C.$$

It follows that there exists some constant $\delta_V \in (0, \infty)$ such that

$$\bar{\mathbb{L}}V \leq \delta_V V + \delta_V, \quad \text{or equivalently, } \bar{\mathbb{L}}(V+1) \leq \delta_V(V+1)$$

Thus, by considering $V+1$ instead of V , the assumptions (and hence the conclusion) in Proposition 3.9 are verified.

3.4 Adding jumps to time-dependent linear Fokker-Planck equations

Let $(\mu_t := u(t,x)dx)_{t \geq 0}$ be a given solution to the linear Fokker-Planck equation (2.3), and assume that $\mathbf{H}_{\mathbf{a},\mathbf{b}}^{\sqrt{u}}$, $\mathbf{H}_{\lesssim u}$, and $\mathbf{H}_0^{\mathbb{T}V}$ hold. Also, let X be the process on E provided by Theorem 2.12 and Theorem 2.15, with transition function $(P_t)_{t \geq 0}$ and resolvent $\mathcal{U} = (U_\alpha)_{\alpha > 0}$.

Let $(\mathcal{L}, \mathcal{D}_b(\mathcal{L}))$ be the generator of X defined through the resolvent as in (6.19).

Remark 3.10. *It is not hard to deduce from Theorem 2.12, (ii.5), (see e.g. Lemma 3.12 from below), that*

$$\bar{\mathbb{L}}f(t,x) = \mathcal{L}f(t,x), \quad (t,x) \in E, \quad f \in C_c^2([0, \infty) \times \mathbb{R}^d) \cap \mathcal{D}_b(\mathcal{L}).$$

However, despite that both spaces $C_c^2([0, \infty) \times \mathbb{R}^d)$ and $\mathcal{D}_b(\mathcal{L})$ are rich enough in the sense that they both separate the set of finite measures on E , their intersection can be small if b merely satisfies \mathbf{H}_b from the beginning of Section 2.2. For this reason, in what follows we pay special attention to the two operators, distinctly.

Let $K : [0, \infty) \times \mathbb{R}^d \times \mathcal{B}(\mathbb{R}^d) \rightarrow [0, \infty)$ be a bounded kernel from $[0, \infty) \times \mathbb{R}^d$ to \mathbb{R}^d , that is K is measurable in the first two arguments, and a measure in the third argument, and

$$\sup_{t \geq 0, x \in \mathbb{R}^d} K(t,x, \mathbb{R}^d) < \infty.$$

Let us consider the perturbed operators

$$\begin{aligned}\mathcal{L}^K f(t, x) &:= \mathcal{L}f(t, x) + \int_{\mathbb{R}^d} [f(t, y) - f(t, x)] K(t, x, dy), \quad f \in \mathcal{D}_b(\mathcal{L}), \quad (t, x) \in E \\ \bar{\mathcal{L}}^K f(t, x) &:= \bar{\mathcal{L}}f(t, x) + \int_{\mathbb{R}^d} [f(t, y) - f(t, x)] K(t, x, dy), \quad f \in C_b^2([0, \infty) \times \mathbb{R}^d), \quad (t, x) \in [0, \infty) \times \mathbb{R}^d.\end{aligned}\tag{3.12}$$

Proposition 3.11. *Let us keep the context fixed in the beginning of this section. Then there exists a right process X^K on E ,*

$$X^K = (\Omega^K, \mathcal{F}^K, \mathcal{F}^K(t), X^K(t) = (X_1^K(t), X_2^K(t)), \mathbb{P}_{s,x}^K, (s, x) \in E, t \geq 0)$$

which has a.s. càdlàg trajectories with respect to the (trace of the) norm topology on E , and such that the following properties hold:

(i) For $(s, x) \in E$ we have

$$\mathbb{P}_{s,x}^K(X_1^K(t) = t + s; 0 \leq t < \infty) = 1.$$

(ii) For every $\psi \in pb(\mathcal{B}(\mathbb{R}^d) \otimes \mathcal{B}(\mathbb{R}^d))$, $\psi(y, y) = 0, y \in \mathbb{R}^d$, we have

$$\mathbb{E}_{s,x}^K \left\{ \sum_{r \leq t} \psi(X_2^K(r-), X_2^K(r)) \right\} = \mathbb{E}_{s,x}^K \left\{ \int_0^t \int_{\mathbb{R}^d} \psi(X_2^K(r), y) K(s+r, X_2^K(r), dy) dr \right\}, \quad t \geq 0.$$

(iii) If there exists $\delta_V \in (0, \infty)$ and $0 \leq V \in C^2(\mathbb{R}^d)$ such that

$$\bar{\mathcal{L}}V \leq \delta_V V \quad \text{and} \quad K(t, x, V(\cdot)) \leq \delta_V V(x), \quad t \geq 0, x \in \mathbb{R}^d,$$

then the transition semigroup $(P_t^K)_{t \geq 0}$ of X^K satisfies

$$P_t^K V \leq e^{2\delta_V t} V, \quad t \geq 0.$$

Moreover, the bound (3.11) holds for X_2^K instead of X_2 and $2\delta_V$ instead of δ_V .

(iv) We have that \mathcal{L}^K coincides with the generator on $b\mathcal{B}(E^K)$ associated to X^K in the sense given by (6.19), $\mathcal{D}_b(\mathcal{L}) = \mathcal{D}_b(\mathcal{L}^K)$, and for every $f \in \mathcal{D}_b(\mathcal{L}) = \mathcal{D}_b(\mathcal{L}^K)$ and every $(s, x) \in E$, the process

$$f(X^K(t)) - f(X^K(0)) - \int_0^t \mathcal{L}^K f(X^K(s)) ds, \quad t \geq 0,\tag{3.13}$$

is a càdlàg $(\mathcal{F}^K(t))$ -martingale under $\mathbb{P}_{s,x}^K$.

Proof. The idea is to regard K as a bounded kernel on E ,

$$Kf(t, x) := K(t, x, f(t, \cdot)), \quad f \in b\mathcal{B}(E), \quad (t, x) \in E,$$

where $f(t, \cdot)$ is extended by zero on the complement of $E(t) = \{x \in \mathbb{R}^d : (t, x) \in E\}$. Thus, by Proposition 6.59 there exists a càdlàg right Markov process

$$X^K = (\Omega^K, \mathcal{F}^K, \mathcal{F}^K(t), X^K(t) = (X_1^K(t), X_2^K(t)), \mathbb{P}_{s,x}^K, (s, x) \in E, t \geq 0)$$

with transition semigroup $(P_t^K)_{t \geq 0}$ and resolvent $U^K := (U_\alpha^K)_{\alpha > 0}$ such that

$$P_t^K = P_t' + \int_0^t P_s' K P_s^K ds \quad \text{and} \quad U_\alpha^K = U_\alpha' + U_\alpha' K U_\alpha^K, \quad t, \alpha > 0,\tag{3.14}$$

where $(P_t')_{t \geq 0}$ and $(U_\alpha')_{\alpha > 0}$ denote the transition semigroup and the resolvent of the process X killed by the multiplicative functional induced by the function $K1$, as in Section 6.5.

Now, since K in fact acts only on the spatial variable, it is straightforward to check, using (3.14), that X_1^K is the uniform motion to the right, i.e. (i) holds.

Assertions (ii) and (iv) are now just consequences of Proposition 6.59 and Proposition 6.60.

Let us prove (iii). First of all, notice that from Proposition 3.9 and the fact that \mathcal{U}' is obtained from \mathcal{U} by killing, we have

$$U'_{\alpha+\delta_V} V \leq U_{\alpha+\delta_V} V \leq \frac{1}{\alpha} V, \quad \alpha > 0.$$

Thus, by (3.14) we get

$$\begin{aligned} U_{\alpha+2\delta_V}^K V &= U'_{\alpha+2\delta_V} V + U'_{\alpha+2\delta_V} \sum_{n \leq 1} (K U'_{\alpha+2\delta_V})^n V \\ &\leq \frac{1}{\alpha + \delta_V} V + \frac{1}{\alpha + \delta_V} V \sum_{n \geq 1} \left(\frac{\delta_V}{\alpha + \delta_V} \right)^n \\ &= \frac{1}{\alpha} V, \quad \alpha > 0. \end{aligned}$$

Now (iii) follows from by the inverse Laplace transform the proof of Proposition 3.9, (iii). \square

We now go further and present two cases when in (3.13) we can take $f \in C_c^2([0, \infty) \times \mathbb{R}^d)$ and replace \mathcal{L}^K by $\bar{\mathcal{L}}^K$.

In the first case we consider locally bounded coefficients, and we start with the following lemma:

Lemma 3.12. *If b is locally bounded then*

$$C_c^2([0, \infty) \times \mathbb{R}^d) \subset \mathcal{D}(\mathcal{L}) \quad \text{and} \quad \bar{\mathcal{L}}f = \mathcal{L}f, \quad f \in C_c^2([0, \infty) \times \mathbb{R}^d).$$

Proof. Recall that by Theorem 2.12, (ii.5) we have that for every $(s, x) \in E$ and every $f \in C_c^2([0, \infty) \times \mathbb{R}^d)$ we have that

$$f(X(t)) - f(X(0)) - \int_0^t \bar{\mathcal{L}}f(X(r)) dr, \quad t \geq 0$$

is a continuous (\mathcal{F}_t) -martingale under $\mathbb{P}_{(s,x)}$. Moreover, since b is locally bounded by assumption, whilst a is locally bounded since \mathbf{H}_a is fulfilled, it follows that

$$\bar{\mathcal{L}}f \in b\mathcal{B}(E), \quad f \in C_c^2([0, \infty) \times \mathbb{R}^d).$$

Now the claim follows from Proposition 6.57. \square

Proposition 3.13. *If b is locally bounded then*

$$C_c^2([0, \infty) \times \mathbb{R}^d) \subset \mathcal{D}(\mathcal{L}^K) \quad \text{and} \quad \bar{\mathcal{L}}^K f = \mathcal{L}^K f, \quad f \in C_c^2([0, \infty) \times \mathbb{R}^d). \quad (3.15)$$

Consequently, if X^K is the right process constructed in Proposition 3.11, then for every $(s, x) \in E$ and every $f \in C_c^2([0, \infty) \times \mathbb{R}^d)$ we have that

$$\mathbb{E}_{(s,x)}^K \left\{ \int_0^t \left| \bar{\mathcal{L}}^K f(X^K(r)) \right| dr \right\} < \infty, \quad t \geq 0,$$

and

$$\left(f(X^K(t)) - f(X^K(0)) - \int_0^t \bar{\mathcal{L}}^K f(X^K(r)) dr \right)_{t \geq 0} \quad \text{is a càdlàg } (\mathcal{F}_t^K) \text{-martingale under } \mathbb{P}_{(s,x)}^K.$$

Proof. Clearly, if we prove (3.15), then the second part of the statement is just a direct consequence of Proposition 3.11, (iv).

Let us prove (3.15). Recall first that by Proposition 6.59, (iv), it follows that the generator $(\mathcal{L}^K, \mathcal{D}_b(\mathcal{L}^K))$ associated to X^K in the sense of (6.19) satisfies $\mathcal{D}_b(\mathcal{L}) = \mathcal{D}_b(\mathcal{L}^K)$, and \mathcal{L}^K is given by (3.12). Now, (3.15) follows by Lemma 3.12. \square

In the second case the idea is to regard $\bar{\mathbb{L}}^K$ as a bounded perturbation of $\bar{\mathbb{L}}$ in $L^1(\bar{\mu})$, where recall that $\bar{\mu}$ denotes the measure $\mu_t(dx)dt$ on $[0, \infty) \times \mathbb{R}^d$. To this end we introduce the following condition:

H_u(K) There exists $\delta \in (0, \infty)$

$$\int_0^\infty \int_{\mathbb{R}^d} K(t, x, f(t, \cdot)) \mu_t(dx) dt \leq \delta \int_0^\infty \int_{\mathbb{R}^d} f(t, x) \mu_t(dx) dt, \quad f \in L^1(\bar{\mu}). \quad (3.16)$$

Remark 3.14. If the kernel K admits a density $K(t, x, dy) = k(t, x, y)dy$ then (3.16) is equivalent to

$$\int_{\mathbb{R}^d} k(t, x, y)u(t, x) dx \leq \delta u(t, y) \quad dy \otimes dt\text{-a.e.}$$

Lemma 3.15. Let X^K be the right process given by Proposition 3.11, with resolvent $\mathcal{U}^K := (U_\alpha^K)_{\alpha > 0}$. If **H_u(K)** is fulfilled, then

$$\|U_\alpha^K f\|_{L^1(E; \delta_0 \otimes \mu_0)} \leq \left(1 + \frac{\delta}{\alpha}\right) \|f\|_{L^1(E; \bar{\mu})}, \quad f \in L^1(E; \bar{\mu}).$$

Proof. First of all note that by Proposition 3.6 we have

$$\|U_\alpha f\|_{L^1(E; \delta_0 \otimes \mu_0)} \leq \|f\|_{L^1(E; \bar{\mu}_\alpha)}, \quad f \in L^1(E; \bar{\mu}_\alpha), \alpha > 0.$$

By the fact that $U_\alpha \leq U'_\alpha$ as kernels, the fact that $\bar{\mu}$ is sub-invariant for \mathcal{U} , and using also **H_u(K)**, we have

$$\bar{\mu}(KU'_\alpha f) \leq \delta \bar{\mu}(U'_\alpha f) \leq \delta \bar{\mu}(U_\alpha f) \leq \frac{\delta}{\alpha} \bar{\mu}(f), \quad f \in p\mathcal{B}(E).$$

Now we note that by employing the resolvent formula (3.14) we get

$$U_\alpha^K f = U'_\alpha f + U'_\alpha \sum_{n \geq 1} (KU'_\alpha)^n f, \quad f \in p\mathcal{B}(E), \alpha > 0.$$

Consequently, for $f \in p\mathcal{B}(E)$ we have

$$\begin{aligned} \delta_0 \otimes \mu_0 (U_{\alpha+\delta}^K f) &\leq \bar{\mu}(f) + \bar{\mu} \left(\sum_{n \geq 1} (KU'_{\alpha+\delta})^n f \right) \\ &\leq \bar{\mu}(f) + \bar{\mu}(f) \sum_{n \geq 1} \left(\frac{\delta}{\alpha + \delta} \right)^n \\ &= \left(1 + \frac{\delta}{\alpha} \right) \bar{\mu}(f), \quad f \in p\mathcal{B}(E). \end{aligned}$$

□

Proposition 3.16. Assume that **H_u(K)** hold. Furthermore, let $(\bar{\mathbb{L}}_1, D(\bar{\mathbb{L}}_1))$ be the generator on $L^1([0, \infty) \times \mathbb{R}^d; \bar{\mu})$ provided by Theorem 2.12.

Then there exists a maximal operator $\bar{\mathbb{L}}_1^K$ on $L^1([0, \infty) \times \mathbb{R}^d; \bar{\mu})$ which extends $\bar{\mathbb{L}}^K$ given in (3.12), namely

$$\begin{aligned} D(\bar{\mathbb{L}}_1) \ni f &\mapsto \bar{\mathbb{L}}_1^K f \in L^1([0, \infty) \times \mathbb{R}^d; \bar{\mu}) \\ \bar{\mathbb{L}}_1^K f(t, x) &:= \bar{\mathbb{L}}_1 f(t, x) + \int_{\mathbb{R}^d} [f(t, y) - f(t, x)] K(t, x, dy), \quad (t, x) \in E \bar{\mu}\text{-a.e.}, \quad f \in D(\bar{\mathbb{L}}_1), \end{aligned}$$

which is well defined and generates a C_0 -semigroup $(T_t^K)_{t \geq 0}$ on $L^1([0, \infty) \times \mathbb{R}^d; \bar{\mu})$. Moreover, if X^K is the right process constructed in Proposition 3.11, then the following properties hold:

- (i) The transition semigroup $(P_t^K)_{t \geq 0}$ of X^K coincides with $(T_t^K)_{t \geq 0}$ as family of operators on $L^1([0, \infty) \times \mathbb{R}^d; \bar{\mu})$

(ii) There exists $E^K \in \mathcal{B}(E)$ such that $E \setminus E^K$ is $\bar{\mu}$ -inessential with respect to \mathcal{U}^K , $\mu_t(E^K(t)) = 1$, $t \geq 0$, where recall that for a set $A \in [0, \infty) \times \mathbb{R}^d$ we define $A(t) := \{x \in \mathbb{R}^d : (t, x) \in A\}$, and for every $(s, x) \in E^K$ and every $f \in C_c^\infty((0, \infty) \times \mathbb{R}^d)$ we have that

$$\mathbb{E}_{(s,x)}^K \left\{ \int_0^t \left| \bar{\Gamma}^K f(X^K(r)) \right| dr \right\} < \infty, \quad t \geq 0, \quad (3.17)$$

and

$$\left(f(X^K(t)) - f(X^K(0)) - \int_0^t \bar{\Gamma}^K f(X^K(r)) dr \right)_{t \geq 0} \text{ is a càdlàg } (\mathcal{F}_t^K)\text{-martingale under } \mathbb{P}_{(s,x)}^K. \quad (3.18)$$

Proof. From $\mathbf{H}_u(\mathbf{K})$ we deduce that K induces a bounded operator on $L^1([0, \infty) \times \mathbb{R}^d; \bar{\mu})$. Therefore, the existence of the operator $\bar{\Gamma}_1^K$ with domain $D(\bar{\Gamma}_1^K)$ and generating a C_0 -semigroup on $L^1([0, \infty) \times \mathbb{R}^d; \bar{\mu})$ follows by Proposition 6.64.

Moreover, assertion (i) follows from the same Proposition 6.64.

Let us prove (ii). To this end, first of all note that by (2.12) we deduce that

$$1_{[0,N] \times B(0,N)} (\|a\| + \|b\|) \in L^1([0, \infty) \times \mathbb{R}^d; \bar{\mu}), \quad N \geq 1,$$

hence, by (i) and Lemma 3.15,

$$U_1^K (1_{[0,N] \times B(0,N)} (\|a\| + \|b\|)) \in L^1([0, \infty) \times \mathbb{R}^d; \bar{\mu} + \delta_0 \otimes \mu_0), \quad N \geq 1.$$

Consequently, by setting

$$\tilde{E}^K := \bigcap_{N \geq 1} [U_1^K (1_{[0,N] \times B(0,N)} (\|a\| + \|b\|)) < \infty]$$

we get $\tilde{E}^K \in \mathcal{A}(\mathcal{U}^K)$, see Remark 6.34, and $\bar{\mu}(E \setminus \tilde{E}^K) + \delta_0 \otimes \mu_0(E \setminus \tilde{E}^K) = 0$, hence $E \setminus \tilde{E}^K$ is $(\bar{\mu} + \delta_0 \otimes \mu_0)$ -inessential with respect to \mathcal{U}^K .

Note that (3.17) holds for every $(s, x) \in \tilde{E}^K$.

Let us fix $f \in C_c^\infty((0, \infty) \times \mathbb{R}^d)$, and apply once again Proposition 6.64 to deduce that there exists a set $E_f \in \mathcal{B}(E)$ such that $E \setminus E_f$ is $\bar{\mu}$ -inessential and for every $(s, x) \in E_f$ we have that

$$f(X^K(t)) - f(X^K(0)) - \int_0^t \bar{\Gamma}_1^K f(X^K(r)) dr, \quad t \geq 0$$

is a càdlàg (\mathcal{F}_t^K) -martingale under $\mathbb{P}_{(s,x)}^K$. Now let us note that by (i) we have

$$\bar{\mu} \left(U_1^K \left| \bar{\Gamma}_1^K f - \bar{\Gamma}^K f \right| \right) = 0,$$

hence, by Remark 6.33, by passing to a further subset, we can assume that E_f is such that

$$U_1^K \left| \bar{\Gamma}_1^K f - \bar{\Gamma}^K f \right| (s, x) = 0, \quad (s, x) \in E_f,$$

and thus, for every $(s, x) \in E_f$,

$$\mathbb{E}_{(s,x)}^K \left\{ \int_0^t \left| \bar{\Gamma}_1^K f - \bar{\Gamma}^K f \right| (X^K(r)) dr \right\} = 0, \quad t \geq 0.$$

Let now $\mathcal{C} \subset C_c^2((0, \infty) \times \mathbb{R}^d)$ be dense (with respect to the uniform convergence up to the second derivatives) and also countable, and consider

$$E_0^K := \tilde{E}^K \cap \bigcap_{f \in \mathcal{C}} E_f.$$

Then $E_0^K \in \mathcal{A}(\mathcal{U}^K)$, see Remark 6.34, and $\bar{\mu}(E \setminus E_0^K) = 0$, hence $E \setminus E_0^K$ is $\bar{\mu}$ -inessential with respect to \mathcal{U}^K . Moreover, for every $f \in \mathcal{C}$ we have that (3.18) is fulfilled. Finally, by dominated convergence, it is easy to see that (3.18) holds for every $f \in C_c^\infty((0, \infty) \times \mathbb{R}^d)$.

Let us show that $\mu_t(E_0^K(t)) = 1$ for $t > 0$ as well. To this end, note that $\bar{\mu}(E \setminus E_0^K) = 0$, there exists $t_n \searrow 0$ such that $\mu_{t_n}(E_0^K(t_n)) = 1, n \geq 1$. Moreover, from (3.14) we have that for every $(s, x) \in E_0^K$

$$P_t^K(A)(s, x) = 0 \quad \text{implies} \quad P_t(A)(s, x) = 0, \quad A \in \mathcal{B}(E).$$

Since X_1^K is the uniform motion to the right and $E_0^K \in \mathcal{A}(\mathcal{U}^K)$ (see Remark 6.32), we have

$$P_{t-t_n}^K([0, \infty) \times (E(t) \setminus E_0^K(t)))(t_n, x) = 0, \quad t > t_n > 0, (t_n, x) \in E_0^K.$$

Consequently,

$$P_{t-t_n}([0, \infty) \times (E(t) \setminus E_0^K(t)))(t_n, x) = 0, \quad t > t_n > 0, (t_n, x) \in E_0^K,$$

and thus

$$\mu_t(E(t) \setminus E_0^K(t)) = \mu_{t_n}(P_{t-t_n}([0, \infty) \times (E(t) \setminus E_0^K(t)))(t_n, \cdot)) = 0, \quad t > t_n > 0.$$

This means that

$$\mu_t(E_0^K(t)) = 1, \quad t > 0.$$

Now, let us define

$$E_{00}^K := E_0^K \setminus (\{0\} \times E_0^K(0)),$$

and note that since X_1^K is the uniform motion to the right, we have $E_{00}^K \in \mathcal{A}(\mathcal{U}^K)$ and in fact that $E^K \setminus E_{00}^K$ is $\bar{\mu}$ -inessential.

By Lemma 3.15 we have that

$$\delta_0 \otimes \mu_0(U_1^K 1_{E \setminus E_{00}^K}) \leq \bar{\mu}(E \setminus E_{00}^K) = 0,$$

from which we deduce that there exist $0 < t_n \searrow 0$ such that

$$\delta_0 \otimes \mu_0(P_{t_n}^K(1_{E_{00}^K})) = 1, \quad n \geq 1. \quad (3.19)$$

Let us consider

$$Z_0 := \{x \in E^K(0) : P_{t_n}^K(1_{E_{00}^K})(0, x) = 1, n \geq 1\},$$

so that, using (3.19), we have $\mu_0(Z_0) = 1$. Now let us set

$$E^K := [(\{0\} \times Z_0) \cup E_{00}^K] \cap \tilde{E}^K.$$

Note that by (3.19), the fact that X_1^K is the uniform motion to the right, and using also Remark 6.34, we have that $E^K \in \mathcal{A}(\mathcal{U}^K)$, and in fact that $E \setminus E^K$ is $\bar{\mu}$ -inessential. Moreover, from the above discussion we also have

$$\mu_t(E^K(t)) = 1, \quad t \geq 0.$$

Furthermore, notice that if $(s, x) \in E^K$ with $s > 0$, we have already obtained that (3.18) holds for every $f \in C_c^\infty((0, \infty) \times \mathbb{R}^d)$. It thus remains to show that (3.18) holds for every $f \in C_c^\infty((0, \infty) \times \mathbb{R}^d)$ and for every $(0, x) \in E^K$. To this end, note that $X^K(s) \in E_{00}^K, s > 0$, hence if we denote by $(M_f(t))_{t \geq 0}$ the process in (3.18), then for every $t > s > 0$

$$\begin{aligned} \mathbb{E}_{(0,x)}^K(M_f(t)) &= \mathbb{E}_{(0,x)}^K \{M_f(s) + M_f(t) \circ \theta^K(s)\} = \mathbb{E}_{(0,x)}^K \{M_f(s)\} + \mathbb{E}_{(0,x)}^K \left\{ \mathbb{E}_{X^K(s)}^K \{M_f(t)\} \right\} \\ &= \mathbb{E}_{(0,x)}^K \{M_f(s)\}. \end{aligned}$$

Since $[0, \infty) \ni t \mapsto M_f(t)$ is right continuous in $L^1(\mathbb{P}_{\delta_{(0,x)}^K}^K)$, and using that we clearly have $\mathbb{E}_{(0,x)}^K \{M_f(0)\} = 0$, it follows by Lemma 6.54 that M_f is a martingale under $\mathbb{P}_{\delta_{(0,x)}^K}^K$ for every $(0, x) \in E^K$. \square

Definition 3.17. Let $0 \leq s < T < \infty$. Let $(\nu_t)_{t \in (s, T)} \subset \mathcal{M}(\mathbb{R}^d)$ be a Borel curve in $\mathcal{M}(\mathbb{R}^d)$.

(i) $(\nu_t)_{t \in (s, T)}$ is called a (weak, or distributional) solution on (s, T) to the linear Fokker-Planck equation

$$\frac{d}{dt} \nu_t = (\mathcal{L}_t^K)^* \nu_t, \quad t \in (s, T), \quad (3.20)$$

if

$$\int_s^T \int_{B(0, R)} \|b(t, x)\| + \|a(t, x)\| \nu_t(dx) dt < \infty, \quad R > 0,$$

and

$$\int_s^T \int_{\mathbb{R}^d} \bar{\Gamma}^K f(t, x) \nu_t(dx) dt = 0, \quad f \in C_c^\infty((s, T) \times \mathbb{R}^d).$$

(ii) $(\nu_t)_{t \in [s, T]}$ is called a (weak, or distributional) solution on (s, T) to the linear Fokker-Planck equation

$$\frac{d}{dt} \nu_t = (\mathcal{L}_t^K)^* \nu_t, \quad t \in (s, T), \quad (3.21)$$

if it is weakly continuous and

$$\nu_T(f(T, \cdot)) = \nu_s(f(s, \cdot)) + \int_s^T \int_{\mathbb{R}^d} \mathcal{L}^K f(t, x) \nu_t(dx) dt, \quad f \in \mathcal{D}_b(\mathcal{L}),$$

where recall that $\mathcal{D}_b(\mathcal{L})$ is given by (6.19), whilst \mathcal{L}^K is given by (3.12); see also Proposition 3.11, (iv).

If $(\nu_t)_{t \in [s, T]} \subset \mathcal{M}(\mathbb{R}^d)$ is weakly continuous and is a solution to (3.21) or (3.20) from above, then ν_s is considered to be the initial condition of the corresponding Fokker-Planck equation.

Definition 3.18. A pair (E, Γ_μ) , where $E \in \mathcal{B}([0, \infty) \times \mathbb{R}^d)$ and Γ_μ is a mapping

$$[0, \infty) \times E \times \mathcal{B}(\mathbb{R}^d) \ni (t, s, x, A) \xrightarrow{\Gamma_\mu} \Gamma_{\mu, t}^{s, x}(A) \in [0, 1],$$

is called a μ -restricted fundamental solution flow to the linear Fokker-Planck equation (3.21) (respectively (3.20)) associated to \mathcal{L}^K (respectively $\bar{\Gamma}^K$) if (i)-(v) from Definition 3.1 are fulfilled with the linear Fokker-Planck equation (2.13) replaced by (3.21) (respectively (3.20)).

Corollary 3.19. Let X^K be the right process provided by Proposition 3.11. Further, for every $A \in \mathcal{B}(\mathbb{R}^d)$, $s, t \geq 0$, $(s, x) \in E$, set

$$\Gamma_\mu^K := \left(\Gamma_{\mu, t}^{K, s, x} \right)_{t \geq 0, (s, x) \in E}, \quad \Gamma_{\mu, t}^{K, s, x}(A) := \mathbb{P}_{s, x}^K(X_2^K(t) \in A \cap E(t + s)).$$

Then the following assertions hold

(i) (E, Γ_μ) is a μ -restricted fundamental solution flow to the linear Fokker-Planck equation (3.21). Moreover, if b is locally bounded then (E, Γ_μ) is a μ -restricted fundamental solution flow to the linear Fokker-Planck equation (3.20).

(ii) If $\mathbf{H}_u(\mathbf{K})$ is fulfilled and E^K is given by Proposition 3.16, then $\left(E^K, \left(\Gamma_{\mu, t}^{K, s, x} \right)_{t \geq 0, (s, x) \in E^K} \right)$ is a μ -restricted fundamental solution flow to the linear Fokker-Planck equation (3.20).

Proof. The proof is very similar to that of Corollary 3.2, so we omit it. \square

4 Construction of reference densities satisfying $\mathbf{H}_{a, b}^{\sqrt{u}}$, $\mathbf{H}_{\lesssim u}$ and \mathbf{H}_0^{TV}

In this section we first show that condition $\mathbf{H}^{\sqrt{u}}$ from the beginning of Section 2 is valid under very mild additional assumptions. We mention that when the diffusion coefficients a satisfy a Lipschitz condition in the space variable and a 1/2-Hölder in the time variable, then $\mathbf{H}^{\sqrt{u}}$ can be ensured through [36, Theorem 7.4.1]. However, it turns out (see Proposition 4.4) that $\mathbf{H}^{\sqrt{u}}$ is valid under the much weaker conditions (4.2), (4.3), and (4.5) from below.

Then, we show that the Trevisan's well-posedness conditions from [64] in the nondegenerate case are explicit sufficient conditions for the main general assumptions $\mathbf{H}_{a, b}^{\sqrt{u}}$, $\mathbf{H}_{\lesssim u}$, and \mathbf{H}_0^{TV} from Section 2 to be fulfilled.

4.1 The square root H^1 -regularity for Fokker-Planck equations

Let $N > 0$, $(u(t, x)dx)_{t \in [0, N]} \subset \mathcal{P}(\mathbb{R}^d)$ be a weakly-continuous solution on $[0, N]$ to the linear Fokker-Planck equation

$$\int_{(0, N) \times \mathbb{R}^d} \left[\frac{d}{dt} f(t, x) + \mathbf{L}_t f(t, x) \right] u(t, x) dx dt = 0, \quad f \in C_c^\infty((0, N) \times \mathbb{R}^d), \quad (4.1)$$

where

$$b : [0, \infty) \times \mathbb{R}^d \rightarrow \mathbb{R}^d \quad \text{and} \quad a : [0, \infty) \times \mathbb{R}^d \rightarrow \mathbb{R}^{d \times d}$$

are Borel measurable coefficients, such that $A(t, x) := (a_{i,j}(t, x))_{1 \leq i, j \leq d}$ is a symmetric and non-negative definite real matrix for all $(t, x) \in [0, \infty) \times \mathbb{R}^d$.

Furthermore, we assume that for every $r > 0$ we have

$$b_i \text{ and } \partial_{x_j} a_{i,j} \text{ belong to } L^2([0, N] \times B(0, r)), \quad 1 \leq i, j \leq d, \quad (4.2)$$

and there exists a constant $C = C(r, N) > 0$ such that

$$C^{-1} \|\xi\|^2 \leq \langle A(t, x)\xi, \xi \rangle \leq C \|\xi\|^2, \quad \xi \in \mathbb{R}^d, (t, x) \in [0, N] \times B(0, r). \quad (4.3)$$

Let

$$\tilde{b}_i(t, x) := \frac{1}{2} \sum_{i=1}^d \partial_{x_j} a_{i,j}(t, x), \quad (t, x) \in [0, N] \times \mathbb{R}^d, \quad (4.4)$$

so that the operator L_t can be rewritten as

$$L_t f = \langle b - \tilde{b}, \nabla f \rangle + \frac{1}{2} \operatorname{div}(A \nabla f), \quad f \in C_c^\infty((0, N) \times \mathbb{R}^d).$$

Further, we assume that

$$hu \in L^\infty([0, N] \times \mathbb{R}^d; \mathbb{R}) \quad \text{and} \quad h \nabla u \in L^2([0, N] \times \mathbb{R}^d; \mathbb{R}^d), \quad \text{for all } h \in C_c^\infty(\mathbb{R}^d). \quad (4.5)$$

Lemma 4.1. *If (4.2), (4.3), and (4.5) hold, then*

$$hu \in W^{1,2}([0, N]; H^{-1}(\mathbb{R}^d)), \text{ hence } hu \in C([0, T]; L^2(\mathbb{R}^d)), \quad h \in C_c^\infty(\mathbb{R}^d).$$

Proof. Let $h \in C_c^\infty(\mathbb{R}^d)$ and $f \in C_c^\infty((0, N) \times \mathbb{R}^d)$, so that

$$\int_{(0, N) \times \mathbb{R}^d} \left[\frac{d}{dt} (hf) + \langle b - \tilde{b}, \nabla(hf) \rangle + \frac{1}{2} \operatorname{div}(A \nabla(hf)) \right] u dx dt = 0.$$

Consequently,

$$\begin{aligned} \int_{(0, N) \times \mathbb{R}^d} hu \frac{d}{dt} f dx dt &= - \int_{(0, N) \times \mathbb{R}^d} u \langle b - \tilde{b}, \nabla(hf) \rangle dx dt - \int_{(0, N) \times \mathbb{R}^d} u \operatorname{div}(A \nabla(hf)) dx dt \\ &= - \int_{(0, N) \times \mathbb{R}^d} u \langle b - \tilde{b}, \nabla(hf) \rangle dx dt - \int_{(0, N) \times \mathbb{R}^d} \langle A \nabla u, \nabla(hf) \rangle dx dt. \end{aligned}$$

Thus, by the Cauch-Schwartz inequality

$$\begin{aligned} \left| \int_{(0, N) \times \mathbb{R}^d} hu \frac{d}{dt} f dx dt \right| &\leq \int_{(0, N) \times \mathbb{R}^d} \left(\|u(b - \tilde{b})\| + \|A \nabla u\| \right) (\|\nabla h\| + |h|) (|f| + \|\nabla f\|) dx dt \\ &\leq \left(\int_{(0, N) \times \mathbb{R}^d} \left(\|u(b - \tilde{b})\| + \|A \nabla u\| \right)^2 (\|\nabla h\| + |h|)^2 dx dt \right)^{1/2} \\ &\quad \times \left(\int_{(0, N) \times \mathbb{R}^d} (|f| + \|\nabla f\|)^2 dx dt \right)^{1/2}. \end{aligned}$$

Now, using (4.2), (4.3), and (4.5), it is straightforward to get that there exists a constant $C < \infty$, independent of f , such that

$$\left| \int_{(0,N) \times \mathbb{R}^d} hu \frac{d}{dt} f \, dxdt \right| \leq C \|f\|_{L^2([0,N]; H^1(\mathbb{R}^d))},$$

which proves the first part of the statement. The second part of the statement follows by the first one, (4.5), and Lions-Magenes lemma. \square

Remark 4.2. Note that by Lemma 4.1 and the fact that $(\mu_t)_{t \in [0,N]}$ is weakly continuous, hence tight, it is straightforward to deduce that

$$\lim_{t \rightarrow 0} u(t, \cdot) = u(0, \cdot) \quad \text{in } L^1(\mathbb{R}^d).$$

The following regularity result is the first main step towards the main result of this subsection, namely Proposition 4.4 from below. For its proof we got inspired from [34], where the uniform ellipticity is employed to prove the regularity of invariant measures for time-independent operators.

Proposition 4.3. If (4.2), (4.3), and (4.5) hold, then for every $\gamma \in (0, 1)$ we have

$$\int_0^N \int_{\mathbb{R}^d} h \frac{\|\nabla(hu)\|^2}{(hu)^\gamma} \, dxdt < \infty, \quad 0 \leq h \in C_c^\infty((0, N) \times \mathbb{R}^d).$$

Proof. Let $0 \leq h \in C_c^\infty((0, N) \times \mathbb{R}^d)$ and set

$$u_c := hu + c, \quad v_c := (hu + c)^{1-\gamma}, \quad c > 0.$$

By Fatou's lemma we have

$$\int_0^N \int_{\mathbb{R}^d} h \frac{\|\nabla(hu)\|^2}{(hu)^\gamma} \, dxdt \leq \liminf_{c \rightarrow 0} \int_0^N \int_{\mathbb{R}^d} h \frac{\|\nabla u_c\|^2}{u_c^\gamma} \, dxdt.$$

We proceed with the key fact that

$$\begin{aligned} \int_0^N \int_{\mathbb{R}^d} h \frac{\|\nabla u_c\|^2}{u_c^\gamma} \, dxdt &\leq C \int_0^N \int_{\mathbb{R}^d} \frac{h}{u_c^\gamma} \langle A \nabla u_c, \nabla u_c \rangle \, dxdt \\ &= \frac{C}{1-\gamma} \int_0^N \int_{\mathbb{R}^d} h \langle \nabla v_c, A \nabla u_c \rangle \, dxdt \\ &= \frac{C}{1-\gamma} \int_0^N \int_{\mathbb{R}^d} \langle \nabla(hv_c), A \nabla u_c \rangle \, dxdt - \frac{C}{1-\gamma} \int_0^N \int_{\mathbb{R}^d} \langle v_c \nabla h, A \nabla u_c \rangle \, dxdt. \end{aligned}$$

Further, note that for $0 < c \leq 1$

$$\begin{aligned} \left| \int_0^N \int_{\mathbb{R}^d} \langle v_c \nabla h, A \nabla u_c \rangle \, dxdt \right| &= \left| \int_0^N \int_{\mathbb{R}^d} \langle v_c A \nabla h, \nabla(hu) \rangle \, dxdt \right| \\ &\leq \int_0^N \int_{\mathbb{R}^d} v_c \|A \nabla h\| \|\nabla(hu)\| \, dxdt \\ &\leq \left(\int_0^N \int_{\mathbb{R}^d} (hu + 1)^{2(1-\gamma)} \|A \nabla h\|^2 \, dxdt \right)^{1/2} \left(\int_0^N \int_{\mathbb{R}^d} \|\nabla(hu)\|^2 \, dxdt \right)^{1/2}. \end{aligned}$$

The last two integrals are finite by (4.3), (4.5), and the fact that $h \in C_c^\infty((0, N) \times \mathbb{R}^d)$.

It remains to show that

$$\liminf_{c \rightarrow 0} \int_0^N \int_{\mathbb{R}^d} \langle \nabla(hv_c), A \nabla u_c \rangle \, dxdt < \infty. \quad (4.6)$$

To this end, note that the equation (4.1) satisfied by u can be rewritten as

$$0 = \int_0^N \int_{\mathbb{R}^d} \left(\frac{d}{dt} f \right) u + f \langle b - \tilde{b}, \nabla u \rangle - \langle \nabla f, A \nabla u \rangle \, dxdt, \quad (4.7)$$

hence, for every $f \in C_c^\infty((0, \infty) \times \mathbb{R}^d)$ we have

$$\int_0^N \int_{\mathbb{R}^d} \langle \nabla f, A \nabla(hu) \rangle dxdt = \int_0^N \int_{\mathbb{R}^d} \langle h \nabla f, A \nabla u \rangle dxdt + \langle \nabla f, u A \nabla h \rangle dxdt \quad (4.8)$$

$$= \int_0^N \int_{\mathbb{R}^d} \langle \nabla(hf), A \nabla u \rangle dxdt - \langle f \nabla h, A \nabla u \rangle + \langle \nabla f, u A \nabla h \rangle dxdt \quad (4.9)$$

$$= \int_0^N \int_{\mathbb{R}^d} u \frac{d}{dt}(hf) + hf \langle b - \tilde{b}, \nabla u \rangle - \langle f \nabla h, A \nabla u \rangle + \langle \nabla f, u A \nabla h \rangle dxdt, \quad (4.10)$$

where for the last equality we have employed (4.7). Also, note that by a density argument, (4.10) remains valid for

$$f \in W_0^{1,2}((0, N); L^2(\mathbb{R}^d)) \cap L^2((0, N); H^1(\mathbb{R}^d)) \cap L^\infty((0, N) \times \mathbb{R}^d).$$

Let us now consider the following (smoothing) bounded linear operators

$$\Gamma_\alpha := \alpha(\alpha - \Delta)^{-1} : L^2(\mathbb{R}^d) \rightarrow H^2(\mathbb{R}^d), \quad \alpha > 0,$$

and recall that Γ_α extend to linear bounded operators

$$\Gamma_\alpha : H^{-1}(\mathbb{R}^d) \rightarrow H^1(\mathbb{R}^d), \quad \Gamma_\alpha : L^p(\mathbb{R}^d) \rightarrow L^p(\mathbb{R}^d), \quad p \in [1, \infty].$$

Moreover, we have that Γ_α is a Markov operator on $L^p(\mathbb{R}^d)$ and

$$\lim_{\alpha \rightarrow \infty} \Gamma_\alpha u = u \text{ for } u \in H^1(\mathbb{R}^d), \quad \lim_{\alpha \rightarrow \infty} \Gamma_\alpha u = u \text{ for } u \in L^p(\mathbb{R}^d), p \in [1, \infty) \quad (4.11)$$

Let us define

$$u_{\alpha,c} := \Gamma_\alpha u_c, \quad v_{\alpha,c} := u_{\alpha,c}^{1-\gamma}, \quad \alpha, c > 0,$$

so that, by Lemma 4.1 and (4.5),

$$\Gamma_\alpha(hv_{\alpha,c}) \in W_0^{1,2}((0, N); L^2(\mathbb{R}^d)) \cap L^2((0, N); H^1(\mathbb{R}^d)) \cap L^\infty((0, N) \times \mathbb{R}^d), \quad \alpha, c > 0.$$

Plugging $\Gamma_\alpha(hv_{\alpha,c})$ instead of f in (4.10) we obtain

$$\begin{aligned} & \int_0^N \int_{\mathbb{R}^d} \langle \nabla \Gamma_\alpha(hv_{\alpha,c}), A \nabla(hu) \rangle dxdt \\ &= \int_0^N \int_{\mathbb{R}^d} u \frac{d}{dt}(h \Gamma_\alpha(hv_{\alpha,c})) + h \Gamma_\alpha(hv_{\alpha,c}) \langle b - \tilde{b}, \nabla u \rangle \\ & \quad - \langle \Gamma_\alpha(hv_{\alpha,c}) \nabla h, A \nabla u \rangle + \langle \nabla \Gamma_\alpha(hv_{\alpha,c}), u A \nabla h \rangle dxdt. \end{aligned}$$

Note that by (4.11), (4.5), and the fact that ∇ and Γ_α commute, and Γ_α is $L^2(\mathbb{R}^d)$ -symmetric, we have

$$\begin{aligned} \lim_{\alpha \rightarrow \infty} \int_0^N \int_{\mathbb{R}^d} \langle \nabla(\Gamma_\alpha(hv_{\alpha,c})), A \nabla(hu) \rangle dxdt &= \lim_{\alpha \rightarrow \infty} \int_0^N \int_{\mathbb{R}^d} \langle \nabla(hv_{\alpha,c}), \Gamma_\alpha(A \nabla(hu)) \rangle dxdt \\ &= \int_0^N \int_{\mathbb{R}^d} \langle \nabla(hv_c), A \nabla(hu) \rangle dxdt, \end{aligned}$$

hence showing (4.6) boils down to showing

$$\liminf_{c \rightarrow 0} \lim_{\alpha \rightarrow \infty} \left| \int_0^N \int_{\mathbb{R}^d} u \frac{d}{dt}(h \Gamma_\alpha(hv_{\alpha,c})) + h(\Gamma_\alpha(hv_{\alpha,c})) \langle b - \tilde{b}, \nabla u \rangle \right. \quad (4.12)$$

$$\left. - \langle (\Gamma_\alpha(hv_{\alpha,c})) \nabla h, A \nabla u \rangle + \langle \nabla(\Gamma_\alpha(hv_{\alpha,c})), u A \nabla h \rangle dxdt \right| < \infty. \quad (4.13)$$

In what follows we treat all the four terms in the last integral, separately: For the first term, since $\frac{d}{dt}$ and Γ_α commute, and Γ_α is $L^2(\mathbb{R}^d)$ -symmetric, we have

$$\begin{aligned}
& \int_0^N \int_{\mathbb{R}^d} u \frac{d}{dt} (h\Gamma_\alpha(hv_{\alpha,c})) \, dxdt \\
&= \int_0^N \int_{\mathbb{R}^d} u\Gamma_\alpha(hv_{\alpha,c}) \frac{d}{dt} h + u_{\alpha,c} \frac{d}{dt} (hv_{\alpha,c}) \, dxdt \\
&= \int_0^N \int_{\mathbb{R}^d} hv_{\alpha,c} \Gamma_\alpha \left(u \frac{d}{dt} h \right) + \left(\frac{d}{dt} h \right) u_{\alpha,c}^{2-\gamma} + h(1-\gamma) u_{\alpha,c}^{1-\gamma} \frac{d}{dt} u_{\alpha,c} \, dxdt \\
&= \int_0^N \int_{\mathbb{R}^d} hv_{\alpha,c} \Gamma_\alpha \left(u \frac{d}{dt} h \right) + \left(\frac{d}{dt} h \right) u_{\alpha,c}^{2-\gamma} + h \frac{1-\gamma}{2-\gamma} \frac{d}{dt} u_{\alpha,c}^{2-\gamma} \, dxdt \\
&= \int_0^N \int_{\mathbb{R}^d} hv_{\alpha,c} \Gamma_\alpha \left(u \frac{d}{dt} h \right) + \left(\frac{d}{dt} h \right) u_{\alpha,c}^{2-\gamma} - \left(\frac{d}{dt} h \right) \frac{1-\gamma}{2-\gamma} u_{\alpha,c}^{2-\gamma} \, dxdt \\
&\xrightarrow{\alpha \rightarrow \infty} \int_0^N \int_{\mathbb{R}^d} hv_c u \frac{d}{dt} h + \left(\frac{d}{dt} h \right) u_c^{2-\gamma} - \left(\frac{d}{dt} h \right) \frac{1-\gamma}{2-\gamma} u_c^{2-\gamma} \, dxdt \\
&\xrightarrow{c \rightarrow 0} \int_0^N \int_{\mathbb{R}^d} (hu)^{2-\gamma} \frac{d}{dt} h + \left(\frac{d}{dt} h \right) u^{2-\gamma} - \left(\frac{d}{dt} h \right) \frac{1-\gamma}{2-\gamma} u^{2-\gamma} \, dxdt.
\end{aligned}$$

Proceeding to the second term in (4.12), we have

$$\begin{aligned}
\lim_{\alpha \rightarrow \infty} \int_0^N \int_{\mathbb{R}^d} h\Gamma_\alpha(hv_{\alpha,c}) \langle b - \tilde{b}, \nabla u \rangle \, dxdt &= \lim_{\alpha \rightarrow \infty} \int_0^N \int_{\mathbb{R}^d} hv_{\alpha,c} \Gamma_\alpha(h \langle b - \tilde{b}, \nabla u \rangle) \, dxdt \\
&= \int_0^N \int_{\mathbb{R}^d} h^2 (hu + c)^{1-\gamma} \langle b - \tilde{b}, \nabla u \rangle \, dxdt \\
&\xrightarrow{c \rightarrow 0} \int_0^N \int_{\mathbb{R}^d} h^2 (hu)^{1-\gamma} \langle b - \tilde{b}, \nabla u \rangle \, dxdt < \infty \text{ by (4.2), (4.5)}.
\end{aligned}$$

For the third term in (4.12), we have

$$\begin{aligned}
\int_0^N \int_{\mathbb{R}^d} \langle \Gamma_\alpha(hv_{\alpha,c}) \nabla h, A \nabla u \rangle \, dxdt &= \int_0^N \int_{\mathbb{R}^d} \Gamma_\alpha(hv_{\alpha,c}) \langle \nabla h, A \nabla u \rangle \, dxdt \\
&= \int_0^N \int_{\mathbb{R}^d} hv_{\alpha,c} \Gamma_\alpha(\langle \nabla h, A \nabla u \rangle) \, dxdt \\
&\xrightarrow{\alpha \rightarrow \infty} \int_0^N \int_{\mathbb{R}^d} h(hu + c)^{1-\gamma} \langle \nabla h, A \nabla u \rangle \, dxdt \\
&\xrightarrow{c \rightarrow 0} \int_0^N \int_{\mathbb{R}^d} h(hu)^{1-\gamma} \langle \nabla h, A \nabla u \rangle \, dxdt.
\end{aligned}$$

For the last term in (4.12), we have

$$\begin{aligned}
\int_0^N \int_{\mathbb{R}^d} \langle \nabla \Gamma_\alpha(hv_{\alpha,c}), u A \nabla h \rangle \, dxdt &= \int_0^N \int_{\mathbb{R}^d} \langle \nabla(hv_{\alpha,c}), \Gamma_\alpha(u A \nabla h) \rangle \, dxdt \\
&= - \int_0^N \int_{\mathbb{R}^d} \langle hv_{\alpha,c}, \nabla \Gamma_\alpha(u A \nabla h) \rangle \, dxdt \\
&\xrightarrow{\alpha \rightarrow \infty} - \int_0^N \int_{\mathbb{R}^d} \langle h(hu + c)^{1-\gamma}, \nabla(u A \nabla h) \rangle \, dxdt \\
&\xrightarrow{c \rightarrow 0} - \int_0^N \int_{\mathbb{R}^d} \langle h(hu)^{1-\gamma}, \nabla(u A \nabla h) \rangle \, dxdt \\
&< \infty \text{ by (4.2), (4.3), (4.5)}.
\end{aligned}$$

Thus, the proof is finished. \square

We are now able to improve Proposition 4.3 by allowing $\gamma = 1$:

Proposition 4.4. *If (4.2), (4.3), and (4.5) hold, then condition $\mathbf{H}^{\sqrt{u}}$ is fulfilled.*

Proof. Let $0 \leq h \in C_c^\infty((0, N) \times \mathbb{R}^d)$ and set

$$u_c := hu + c, \quad v_c := \log(hu + c), \quad c > 0.$$

By Fatou's lemma we have

$$\int_0^N \int_{\mathbb{R}^d} h^2 \frac{\|\nabla(hu)\|^2}{hu} dxdt \leq \liminf_{c \rightarrow 0} \int_0^N \int_{\mathbb{R}^d} h^2 \frac{\|\nabla u_c\|^2}{u_c} dxdt.$$

We proceed with the key fact that

$$\begin{aligned} \int_0^N \int_{\mathbb{R}^d} h^2 \frac{\|\nabla u_c\|^2}{u_c} dxdt &\leq C \int_0^N \int_{\mathbb{R}^d} \frac{h^2}{u_c} \langle A\nabla u_c, \nabla u_c \rangle dxdt = \frac{C}{1-\gamma} \int_0^N \int_{\mathbb{R}^d} h^2 \langle \nabla v_c, A\nabla u_c \rangle dxdt \\ &= \frac{C}{2(1-\gamma)} \int_0^N \int_{\mathbb{R}^d} \langle \nabla(h^2 v_c), A\nabla u_c \rangle dxdt - \frac{C}{2(1-\gamma)} \int_0^N \int_{\mathbb{R}^d} \langle hv_c \nabla h, A\nabla u_c \rangle dxdt. \end{aligned}$$

Further, note that for $0 < c \leq 1$ Note that

$$\begin{aligned} &\left| \int_0^N \int_{\mathbb{R}^d} \langle hv_c \nabla h, A\nabla u_c \rangle dxdt \right| \\ &= \left| \int_0^N \int_{\mathbb{R}^d} \langle hv_c A\nabla h, \nabla(hu) \rangle dxdt \right| = \left| \int_0^N \int_{\mathbb{R}^d} \operatorname{div}(hv_c A\nabla h) hu dxdt \right| \\ &= \left| \int_0^N \int_{\mathbb{R}^d} \langle \nabla(hu), hA\nabla h \rangle \frac{hu}{hu+c} dxdt + \int_0^N \int_{\mathbb{R}^d} \operatorname{div}(hA\nabla h) hu \log(hu+c) dxdt \right| \\ &\leq \int_0^N \int_{\mathbb{R}^d} |\langle h\nabla u, hA\nabla h \rangle| dxdt + \int_0^N \int_{\mathbb{R}^d} |\operatorname{div}(hA\nabla h)| |hu \log(hu+c)| dxdt \\ &\leq \int_0^N \int_{\mathbb{R}^d} |\langle \nabla(hu), hA\nabla h \rangle| dxdt + \int_0^N \int_{\mathbb{R}^d} |\operatorname{div}(hA\nabla h)| [hu \log(hu+1) + e^{-1} + 1] dxdt \\ &< \infty \text{ by (4.2), (4.3), and (4.5),} \end{aligned}$$

where we have also used

$$|x \log(x+c)| \leq |x \log(x)| + c, \quad |x \log(x)| \leq e^{-1} + x \log(x+1), \quad x \geq 0, c > 0. \quad (4.14)$$

It remains to show that

$$\liminf_{c \rightarrow 0} \int_0^N \int_{\mathbb{R}^d} \langle \nabla(h^2 v_c), A\nabla u_c \rangle dxdt < \infty. \quad (4.15)$$

Let us define

$$u_{\alpha,c} := \Gamma_\alpha u_c, \quad v_{\alpha,c} := \log(u_{\alpha,c}), \quad \alpha, c > 0,$$

so that, by Lemma 4.1 and (4.5),

$$\Gamma_\alpha(hv_{\alpha,c}) \in W_0^{1,2}((0, N); L^2(\mathbb{R}^d)) \cap L^2((0, N); H^1(\mathbb{R}^d)) \cap L^\infty((0, N) \times \mathbb{R}^d), \quad \alpha, c > 0.$$

Plugging $\Gamma_\alpha(hv_{\alpha,c})$ instead of f in (4.10) we obtain

$$\begin{aligned} &\int_0^N \int_{\mathbb{R}^d} \langle \nabla \Gamma_\alpha(h^2 v_{\alpha,c}), A\nabla(hu) \rangle dxdt \\ &= \int_0^N \int_{\mathbb{R}^d} u \frac{d}{dt} (h\Gamma_\alpha(h^2 v_{\alpha,c})) + h\Gamma_\alpha(h^2 v_{\alpha,c}) \langle b - \tilde{b}, \nabla u \rangle \\ &\quad - \langle \Gamma_\alpha(h^2 v_{\alpha,c}) \nabla h, A\nabla u \rangle + \langle \nabla \Gamma_\alpha(h^2 v_{\alpha,c}), uA\nabla h \rangle dxdt. \end{aligned}$$

Note that by (4.11), (4.5), and the fact that ∇ and Γ_α commute, and Γ_α is $L^2(\mathbb{R}^d)$ -symmetric, we have

$$\begin{aligned} \lim_{\alpha \rightarrow \infty} \int_0^N \int_{\mathbb{R}^d} \langle \nabla(\Gamma_\alpha(h^2 v_{\alpha,c})), A\nabla(hu) \rangle dxdt &= \lim_{\alpha \rightarrow \infty} \int_0^N \int_{\mathbb{R}^d} \langle \nabla(hv_{\alpha,c}), \Gamma_\alpha(A\nabla(hu)) \rangle dxdt \\ &= \int_0^N \int_{\mathbb{R}^d} \langle \nabla(h^2 v_c), A\nabla(hu) \rangle dxdt, \end{aligned}$$

hence showing (4.15) boils down to showing

$$\liminf_{c \rightarrow 0} \lim_{\alpha \rightarrow \infty} \left| \int_0^N \int_{\mathbb{R}^d} u \frac{d}{dt} (h\Gamma_\alpha(h^2 v_{\alpha,c})) + h(\Gamma_\alpha(h^2 v_{\alpha,c})) \langle b - \tilde{b}, \nabla u \rangle \right. \quad (4.16)$$

$$\left. - \langle (\Gamma_\alpha(h^2 v_{\alpha,c})) \nabla h, A\nabla u \rangle + \langle \nabla(\Gamma_\alpha(h^2 v_{\alpha,c})), u A\nabla h \rangle dxdt \right| < \infty. \quad (4.17)$$

We treat all the four terms in the last integral, separately: For the first term, since $\frac{d}{dt}$ and Γ_α commute, and Γ_α is $L^2(\mathbb{R}^d)$ -symmetric, we have

$$\begin{aligned} &\int_0^N \int_{\mathbb{R}^d} u \frac{d}{dt} (h\Gamma_\alpha(h^2 v_{\alpha,c})) dxdt \\ &= \int_0^N \int_{\mathbb{R}^d} u \Gamma_\alpha(h^2 v_{\alpha,c}) \frac{d}{dt} h + u_{\alpha,c} \frac{d}{dt} (h^2 v_{\alpha,c}) dxdt \\ &= \int_0^N \int_{\mathbb{R}^d} h^2 v_{\alpha,c} \Gamma_\alpha \left(u \frac{d}{dt} h \right) + \left(\frac{d}{dt} h^2 \right) u_{\alpha,c} \log(u_{\alpha,c}) + h^2 \frac{d}{dt} u_{\alpha,c} dxdt \\ &= \int_0^N \int_{\mathbb{R}^d} h^2 v_{\alpha,c} \Gamma_\alpha \left(u \frac{d}{dt} h \right) + \left(\frac{d}{dt} h^2 \right) u_{\alpha,c} \log(u_{\alpha,c}) - \left(\frac{d}{dt} h^2 \right) u_{\alpha,c} dxdt \\ &\xrightarrow{\alpha \rightarrow \infty} \int_0^N \int_{\mathbb{R}^d} h^2 v_c u \frac{d}{dt} h + \left(\frac{d}{dt} h^2 \right) u_c \log(u_c) - \left(\frac{d}{dt} h^2 \right) u_c dxdt \\ &= \int_0^N \int_{\mathbb{R}^d} h^2 u \log(hu + c) \frac{d}{dt} h + \left(\frac{d}{dt} h^2 \right) (hu + c) \log(hu + c) - \left(\frac{d}{dt} h^2 \right) (hu + c) dxdt \\ &\xrightarrow{c \rightarrow 0} \int_0^N \int_{\mathbb{R}^d} h^2 u \log(hu) \frac{d}{dt} h + \left(\frac{d}{dt} h^2 \right) hu \log(hu) - \left(\frac{d}{dt} h^2 \right) hu dxdt. \end{aligned}$$

Proceeding to the second term in (4.16), we have

$$\begin{aligned} &\lim_{\alpha \rightarrow \infty} \left| \int_0^N \int_{\mathbb{R}^d} h\Gamma_\alpha(h^2 v_{\alpha,c}) \langle b - \tilde{b}, \nabla u \rangle dxdt \right| \\ &= \left| \lim_{\alpha \rightarrow \infty} \int_0^N \int_{\mathbb{R}^d} h^2 v_{\alpha,c} \Gamma_\alpha(h \langle b - \tilde{b}, \nabla u \rangle) dxdt \right| = \left| \int_0^N \int_{\mathbb{R}^d} h^3 \log(hu + c) \langle b - \tilde{b}, \nabla u \rangle dxdt \right| \quad (4.18) \\ &\leq \int_0^N \int_{\mathbb{R}^d} |h^2 \log(hu + c) \langle b - \tilde{b}, \nabla(hu) \rangle| + |h^2 u \log(hu + c) \langle b - \tilde{b}, \nabla h \rangle| dxdt \\ &\leq \int_0^N \int_{\mathbb{R}^d} \|h \log(hu + c) \nabla(hu)\| \|h(b - \tilde{b})\| + |h^2 u \log(hu + c)| \|b - \tilde{b}\| \|\nabla h\| dxdt \\ &\leq \left(\int_0^N \int_{\mathbb{R}^d} h^2 (\|\log(hu + c) \nabla(hu)\|)^2 dxdt \right)^{1/2} \left(\int_0^N \int_{\mathbb{R}^d} h^2 \|(b - \tilde{b})\|^2 dxdt \right)^{1/2} \\ &\quad + \int_0^N \int_{\mathbb{R}^d} |h^2 u \log(hu + c) \langle b - \tilde{b}, \nabla h \rangle| dxdt. \end{aligned}$$

Now, let us note that

$$\begin{aligned} \int_0^N \int_{\mathbb{R}^d} h^2 \|b - \tilde{b}\|^2 dxdt &< \infty, \quad \text{by (4.2)} \\ \int_0^N \int_{\mathbb{R}^d} |h^2 u \log(hu + c) \langle b - \tilde{b}, \nabla h \rangle| dxdt &\xrightarrow{c \rightarrow 0} \int_0^N \int_{\mathbb{R}^d} |h^2 u \log(hu) \langle b - \tilde{b}, \nabla h \rangle| dxdt < \infty \quad \text{by (4.2), (4.5)}. \end{aligned}$$

Furthermore, by fixing $\gamma \in (0, 1)$,

$$\begin{aligned} \int_0^N \int_{\mathbb{R}^d} h^2 |\log(hu + c)|^2 \|\nabla(hu)\|^2 dxdt &= \int_0^N \int_{\mathbb{R}^d} h^2 (hu)^\gamma |\log(hu + c)|^2 \frac{\|\nabla(hu)\|^2}{(hu)^\gamma} dxdt \\ &\leq \|h(hu)^\gamma |\log(hu + c)|^2\|_{L^\infty((0, N) \times \mathbb{R}^d)} \int_0^N \int_{\mathbb{R}^d} h \frac{\|\nabla(hu)\|^2}{(hu)^\gamma} dxdt. \end{aligned}$$

Finally,

$$\sup_{c \in (0, 1)} \|h(hu)^\gamma |\log(hu + c)|^2\|_{L^\infty((0, N) \times \mathbb{R}^d)} < \infty \quad \text{by using (4.5),}$$

whilst

$$\int_0^N \int_{\mathbb{R}^d} h \frac{\|\nabla(hu)\|^2}{(hu)^\gamma} dxdt < \infty \quad \text{by Proposition 4.3.}$$

For the third term in (4.16), we have

$$\begin{aligned} &\int_0^N \int_{\mathbb{R}^d} \langle \Gamma_\alpha(h^2 v_{\alpha, c}) \nabla h, A \nabla u \rangle dxdt \\ &= \int_0^N \int_{\mathbb{R}^d} \Gamma_\alpha(h^2 v_{\alpha, c}) \langle \nabla h, A \nabla u \rangle dxdt = \int_0^N \int_{\mathbb{R}^d} h^2 v_{\alpha, c} \Gamma_\alpha(\langle \nabla h, A \nabla u \rangle) dxdt \\ &\xrightarrow{\alpha \rightarrow \infty} \int_0^N \int_{\mathbb{R}^d} h^2 \log(hu + c) \langle A \nabla h, \nabla u \rangle dxdt, \end{aligned}$$

and this last term can be further treated precisely like $\int_0^N \int_{\mathbb{R}^d} h^3 \log(hu + c) \langle b - \tilde{b}, \nabla u \rangle dxdt$ in (4.18).

For the last term in (4.16), we have

$$\begin{aligned} &\int_0^N \int_{\mathbb{R}^d} \langle \nabla \Gamma_\alpha(h^2 v_{\alpha, c}), u A \nabla h \rangle dxdt \\ &= \int_0^N \int_{\mathbb{R}^d} \langle \nabla(h^2 v_{\alpha, c}), \Gamma_\alpha(u A \nabla h) \rangle dxdt \xrightarrow{\alpha \rightarrow \infty} \int_0^N \int_{\mathbb{R}^d} \langle \nabla(h^2 \log(hu + c)), u A \nabla h \rangle dxdt \\ &= - \int_0^N \int_{\mathbb{R}^d} h^2 \log(hu + c) \operatorname{div}(u A \nabla h) dxdt = - \int_0^N \int_{\mathbb{R}^d} h^2 \log(hu + c) [\langle \nabla u, A \nabla h \rangle + u \operatorname{div}(A \nabla h)] dxdt. \end{aligned}$$

Now, the integral

$$\int_0^N \int_{\mathbb{R}^d} h^2 \log(hu + c) \langle \nabla u, A \nabla h \rangle dxdt$$

can be further treated precisely like $\int_0^N \int_{\mathbb{R}^d} h^3 \log(hu + c) \langle b - \tilde{b}, \nabla u \rangle dxdt$ in (4.18), whilst

$$\begin{aligned} &\int_0^N \int_{\mathbb{R}^d} h^2 \log(hu + c) u \operatorname{div}(A \nabla h) dxdt \\ &\leq \sup_{c \in (0, 1)} \|h^2 u \log(hu + c)\|_{L^\infty((0, N) \times \mathbb{R}^d)} \int_0^N \int_{\mathbb{R}^d} |\operatorname{div}(A \nabla h)| dxdt < \infty \quad \text{by (4.2), (4.5)}. \end{aligned}$$

Thus, the proof is finished. \square

If conditions (4.2), (4.3), and (4.5) are satisfied globally in space, then, by following essentially the same steps as above, the results from Lemma 4.1 and Proposition 4.4 are satisfied globally in space as well. More precisely, we have the following.

Proposition 4.5. *Assume the following global in space versions of (4.2), (4.3), and (4.5):*

- b_i and $\partial_{x_j} a_{i,j}$ belong to $L^2([0, N] \times B(0, r))$, $1 \leq i, j \leq d$,
- $C^{-1} \|\xi\|^2 \leq \langle A(t, x)\xi, \xi \rangle \leq C \|\xi\|^2$, $\xi \in \mathbb{R}^d, (t, x) \in [0, N] \times B(0, r)$
- $u \in L^\infty([0, N] \times \mathbb{R}^d; \mathbb{R})$ and $\nabla u \in L^2([0, N] \times \mathbb{R}^d; \mathbb{R}^d)$.

Then

$$u \in W^{1,2}([0, N]; H^{-1}(\mathbb{R}^d)), \text{ hence } u \in C([0, T]; L^2(\mathbb{R}^d)), \text{ and } \sqrt{u} \in L^2([0, T]; H^1(\mathbb{R}^d)).$$

4.2 Trevisan's conditions

Let us consider the diffusion operator $(L_t)_{t \in [0, \infty)}$ with coefficients a, b given by (2.1)-(2.4), such that the well-posedness conditions in [64] for the elliptic case are satisfied for every $N > 0$, namely:

$$(T1) \quad \exists C = C(N) > 0 : \langle A(t, x)\xi, \xi \rangle \geq C \|\xi\|^2, \quad t \in [0, N], x, \xi \in \mathbb{R}^d, \quad (4.19)$$

$$(T2) \quad a \in L^\infty([0, N] \times \mathbb{R}^d), \quad \frac{d}{dt} a \in L^\infty([0, N] \times \mathbb{R}^d), \quad (4.20)$$

$$(T3) \quad a \in L^1([0, N]; W^{1,p}(\mathbb{R}^d)) \text{ for some } p \in [2, \infty], \quad \left(\sum_{i,j} \frac{\partial^2}{\partial x_i \partial x_j} a_{i,j} \right)^+ \in L^1([0, N]; L^\infty(\mathbb{R}^d)) \quad (4.21)$$

$$(T4) \quad b \in L^1([0, N]; L^\infty(\mathbb{R}^d)). \quad (4.22)$$

Then, by [64, Theorem 33]:

For every $\mu_0 = u_0 dx \in \mathcal{P}(\mathbb{R}^d)$ with $u_0 \in L^r(\mathbb{R}^d)$, $r \geq 2p/(p-2)$, and for every $N > 0$, there exists a unique weakly continuous solution $\mu := (\mu(t) = u(t, x) dx)_{t \in [0, N]} \subset \mathcal{P}(\mathbb{R}^d)$ of the linear Fokker-Planck equation (2.3), with initial condition μ_0 and such that $u \in L^\infty([0, N]; L^r(\mathbb{R}^d))$. Moreover, it satisfies

$$u \in L^2([0, N]; W^{1,2}(\mathbb{R}^d)), \quad \text{for every } N > 0. \quad (4.23)$$

Remark 4.6. (i) Note that in [64, Theorem 33], the uniqueness part is apparently claimed in the class of weakly continuous solutions which belong to $L^\infty([0, N]; L^r(\mathbb{R}^d))$ and are necessarily probability solutions. However, by inspecting the proof of [64, Theorem 33], one can see that uniqueness holds among all weakly continuous solutions to (2.3) whose densities belongs to $L^\infty([0, N]; L^r(\mathbb{R}^d))$.

(ii) We emphasize that the uniqueness is guaranteed in the space $u \in L^\infty([0, N]; L^r(\mathbb{R}^d))$, whilst (4.23) is a byproduct of the construction of the solution; see [64, Subsection 3.3] and the discussion on page 20.

We also need the following slight modifications of (T3) and (T4):

$$(T3') \quad \partial_{x_j} a_{i,j} \in L^2_{\text{loc}}([0, \infty) \times \mathbb{R}^d), \quad 1 \leq i, j \leq d$$

$$(T4') \quad b_i \in L^2_{\text{loc}}([0, \infty) \times \mathbb{R}^d), \quad 1 \leq i \leq d$$

Corollary 4.7. *Assume that conditions (T1)-(T4), and (T3')-(T4') are satisfied. Further, let $u_0(x)dx \in \mathcal{P}(\mathbb{R}^d)$ be such that $u_0 \in L^\infty(\mathbb{R}^d)$, and consider $(\mu(t) = u(t, x)dx)_{t \in [0, T]} \subset \mathcal{P}(\mathbb{R}^d)$ the unique weakly continuous solution to the linear Fokker-Planck equation (2.3), with initial condition $u_0(x)dx$, and such that $u \in L^\infty([0, N] \times \mathbb{R}^d)$. Then the conditions $\mathbf{H}_{\mathbf{a}, \mathbf{b}}^{\sqrt{u}}$, $\mathbf{H}_{\lesssim u}$, and \mathbf{H}_0^{TV} are satisfied, hence the conclusion of Theorem 2.12 & Theorem 2.15 hold.*

Proof. First, by (T1),(T2), (T3'), and the fact that u is bounded, we get that $\mathbf{H}_{\mathbf{a}}$ is fulfilled.

Further, $\mathbf{H}_{\mathbf{b}}$ follows from (T4') and the fact that u is bounded.

Condition $\mathbf{H}^{\sqrt{u}}$ is ensured by Proposition 4.4 since u is bounded and (4.23) holds.

Further, $\mathbf{H}_{\lesssim u}$ follows from [64, Theorem 33] presented above, and Remark 4.6, (i).

Finally, \mathbf{H}_0^{TV} follows by Remark 4.2 and Remark 2.14. \square

5 Regularization of the superposition principle for nonlinear time-dependent Fokker-Planck equations

In this part we extend to nonlinear Fokker-Planck equations the results presented above for the linear case. We construct Choquet capacities for such nonlinear equations by probabilistic means and we investigate explicit sufficient conditions on the coefficients for a class of generalized porous media equations to which our theory applies.

5.1 The strong Markov property and a Choquet capacity for the solution of the corresponding McKean-Vlasov equation

In this section we start from a general nonlinear Fokker-Planck equation, namely we consider the (t, μ) -dependent diffusion operator

$$\mathbf{L}_{t,\nu} f(t, x) := b(t, x, \nu) \nabla_x f(t, x) + \frac{1}{2} \sum_{i,j=1}^d a_{i,j}(t, x, \nu) \frac{\partial^2}{\partial x_i \partial x_j} f(t, x), \quad (t, x, \nu) \in (0, \infty) \times \mathbb{R}^d \times \mathcal{P}(\mathbb{R}^d),$$

that acts on test functions $f \in C_c^\infty((0, \infty) \times \mathbb{R}^d)$ which are twice continuously differentiable in the second argument, whilst b and a are functions between the following spaces

$$b : [0, \infty) \times \mathbb{R}^d \times \mathcal{P}(\mathbb{R}^d) \rightarrow \mathbb{R}^d \quad \text{and} \quad a : [0, \infty) \times \mathbb{R}^d \times \mathcal{P}(\mathbb{R}^d) \rightarrow \mathbb{R}_{\text{sym}^+}^{d \times d}.$$

Remark 5.1. *Note that above we did not assume any kind of regularity/measurability for the coefficients b and a , the reason being that our aim is to cover irregular coefficients, like those from below which depend on point-wise evaluations of the density of ν assuming it exists, also called Nemytskii-type operators. This is in fact the setting adopted in [4], to which we refer for more details and examples.*

Let us now extend the notion of solution to the Fokker-Planck equation from the linear case (see Definition 2.1) to the nonlinear one:

Definition 5.2. *Let $0 \leq s < T < \infty$.*

(i) *A family $(\mu_t)_{t \in (s, T)} \subset \mathcal{M}(\mathbb{R}^d)$ is called a (weak, or distributional) solution on (s, T) to the nonlinear Fokker-Planck equation*

$$\frac{d}{dt} \mu_t = \mathbf{L}_{t, \mu_t}^* \mu_t, \quad t \in (s, T), \quad (5.1)$$

if $(\mu_t)_{t \in (s, T)}$ is a Borel curve in $\mathcal{M}(\mathbb{R}^d)$,

- $[0, \infty) \times \mathbb{R}^d \ni (t, x) \mapsto b(t, x, \mu_t) \in \mathbb{R}^d$ and $(t, x) \mapsto a(t, x, \mu_t) \in \mathbb{R}^{d \times d}$ are Borel measurable,

- $\int_{(s, T) \times B(0, R)} \|b(t, x, \mu_t)\| + \|A(t, x, \mu_t)\| \mu_t(dx) dt < \infty, \quad R > 0 \quad (5.2)$

- $\int_{(s, T) \times \mathbb{R}^d} \left[\frac{d}{dt} f(t, x) + \mathbf{L}_{t, \mu_t} f(t, x) \right] \mu_t(dx) dt = 0, \quad f \in C_c^\infty((s, T) \times \mathbb{R}^d). \quad (5.3)$

(ii) *We say that $(\mu_t)_{t \in [s, T)} \subset \mathcal{M}(\mathbb{R}^d)$ is a solution to the nonlinear Fokker-Planck equation (5.1) with initial condition μ_s if $(\mu_t)_{t \in (s, T)}$ is a solution on (s, T) as in (i), and $\lim_{t \searrow s} \mu_t = \mu_s$ weakly.*

(iii) *We say that $(\mu_t)_{t \in [s, T)} \subset \mathcal{M}(\mathbb{R}^d)$ is a weakly continuous solution to the nonlinear Fokker-Planck equation (5.1) if it is a solution on (s, T) in the sense of (i), and the curve of measures $[s, T) \ni t \mapsto \mu_t \in \mathcal{M}(\mathbb{R}^d)$ is weakly continuous; this definition is obviously extended when $(\mu_t)_{t \in [s, T]}$ is merely weakly right-continuous.*

Further, if $\mu = (\mu_t = u(t, x)dx)_{t \geq 0} \subset \mathcal{P}(\mathbb{R}^d)$ is a Borel curve, we consider the following "linearized" objects:

$$\begin{aligned} b^u(t, x) &:= b(t, x, \mu_t), \quad a_{ij}^u(t, x) := a_{ij}(t, x, \mu_t), \quad (t, x) \in [0, \infty) \times \mathbb{R}^d \\ \mathbb{L}^u f(t, x) &:= b^u(t, x) \nabla_x f(t, x) + \frac{1}{2} \sum_{i, j=1}^d a_{ij}^u(t, x) \frac{\partial^2}{\partial x_i \partial x_j} f(t, x) \\ &= \mathbb{L}_{t, \mu_t} f(t, x), \quad f \in C_c^\infty((0, \infty) \times \mathbb{R}^d), \quad (t, x) \in [0, \infty) \times \mathbb{R}^d. \end{aligned}$$

Let us now recall the counterpart of the nonlinear Fokker-Planck equations considered above to the theory of McKean-Vlasov SDEs, also called distribution dependent SDEs. First, assume that

$$\sigma : [0, \infty) \times \mathbb{R}^d \times \mathcal{P}(\mathbb{R}^d) \rightarrow \mathbb{R}^{d \times d} \quad \text{is such that} \quad a = \sigma \sigma^*.$$

Definition 5.3. Let $(\Omega, \mathcal{F}, (\mathcal{F}(t))_{t \geq 0}, \mathbb{P})$ be a filtered probability space. An $\mathcal{F}(t)$ -adapted and \mathbb{P} -a.s. path-continuous stochastic process $X = (X(t))_{t \geq 0}$ with values in \mathbb{R}^d is called a weak solution to the McKean-Vlasov SDE

$$dX(t) = b(t, x, \mathcal{L}_{X(t)}) dt + \sigma(t, x, \mathcal{L}_{X(t)}) dW(t), \quad X(0) \sim \mu_0 \in \mathcal{P}(\mathbb{R}^d) \quad (5.4)$$

where W is a d -dimensional standard Brownian motion defined on $(\Omega, \mathcal{F}, (\mathcal{F}(t))_{t \geq 0}, \mathbb{P})$, whilst $\mathcal{L}_{X(t)} := \mathbb{P} \circ (X(t))^{-1}$, $t \geq 0$, if

$$\int_0^N \int_{B(0, R)} (\|b(t, x, \mathcal{L}_{X(s)})\| + \|a(t, x, \mathcal{L}_{X(s)})\|) \mathcal{L}_{X(s)}(dx) dt < \infty, \quad N > 0, R > 0,$$

and

$$M(t) := f(t, X(t)) - f(0, X(0)) - \int_0^t \left(\frac{d}{ds} f + \mathbb{L}_{s, \mathcal{L}_{X(s)}} \right) (s, X(s)) ds, \quad t \geq 0$$

is an $\mathcal{F}(t)$ -martingale for every $f \in C_c^\infty([0, \infty) \times \mathbb{R}^d)$.

Theorem 5.4. Let $(\mu_t)_{t \geq 0} \subset \mathcal{P}(\mathbb{R}^d)$ be a weakly continuous solution to the nonlinear Fokker-Planck equation (5.1) on $(0, \infty)$, such that $\mu_t := u(t, x)dx$, $t \geq 0$ for a jointly measurable density $u : [0, \infty) \times \mathbb{R}^d \rightarrow \mathbb{R}_+$. Furthermore, assume that the "linearized" conditions

$$\mathbf{H}_{\mathbf{a}^u, \mathbf{b}^u}^{\sqrt{\mathbf{a}}}, \quad \mathbf{H}_{\lesssim u}(\mathbb{L}^u), \quad \mathbf{H}_0^{\text{TV}}(\mathbb{L}^u) \quad \text{are fulfilled.}$$

Then there exists a set $E \in \mathcal{B}([0, \infty) \times \mathbb{R}^d)$ and a conservative right (hence strong Markov) process on E

$$X = (\Omega, \mathcal{F}, \mathcal{F}(t), X(t), \mathbb{P}_{s, x}, (s, x) \in E, t \geq 0), \quad \text{with transition function } (P_t)_{t \geq 0},$$

such that the following assertions hold:

(i) $\bar{\mu}([0, \infty) \times \mathbb{R}^d \setminus E) = 0$ and $\mu_s(\{x \in \mathbb{R}^d, (s, x) \in E\}) = 1$, $s \geq 0$.

(ii) For $(s, x) \in E$ we have

$$\mathbb{P}_{s, x}([0, \infty) \ni t \mapsto X(t) \in E \subset [0, \infty) \times \mathbb{R}^d \text{ is continuous}) = 1.$$

(iii) If we set $X(t) = (X_1(t), X_2(t))$, $t \geq 0$, then for $(s, x) \in E$ we have

$$\mathbb{P}_{s, x}(X_1(t) = t + s; 0 \leq t < \infty) = 1.$$

(iv) The process $(X_2(t))_{t \geq 0}$ has the following properties:

(iv.1) For every $s, t \geq 0$

$$\mathbb{P}_{s, \mu_s} \circ (X_2(t))^{-1} = \mu_{t+s}.$$

(iv.2) Regarded on the filtered probability space $(\Omega, \mathcal{F}, \mathcal{F}(t), \mathbb{P}_{0, \mu_0})$, the process $(X_2(t))_{t \geq 0}$ is the unique weak solution to the McKean-Vlasov SDE (5.4), with initial condition $\mathcal{L}_{X_2(0)} = \mu_0$ and such that $\mathcal{L}_{X_2(t)} \leq c\mu_t$, $t \geq 0$ for some constant $c \in (0, \infty)$.

(iv.3) Regarded on the filtered probability space $(\Omega, \mathcal{F}, \mathcal{F}(t), \mathbb{P}_{0, \mu_0})$ the process $(X_2(t))_{t \geq 0}$ is strong Markov, in the sense that for every finite $\mathcal{F}(t)$ -stopping times τ we have that

$$\mathbb{E}_{\delta_0 \otimes \mu_0} \{f(t + \tau, X_2(t + \tau)) \mid \mathcal{F}(\tau)\} = Pf(\tau, X_2(\tau)), \quad t \geq 0, f \in b\mathcal{B}([0, \infty) \times \mathbb{R}^d).$$

Proof. The existence of X with all the desired properties follows by the main results Theorem 2.12 and Theorem 2.15, with a, b replaced by a^u, b^u . \square

Remark 5.5. Let $(X_2(t))_{t \geq 0}$ be the weak solution constructed in Theorem 5.4, and define for every $A \in \mathcal{B}(\mathbb{R}^d)$

$$\tau_A^{(2)} := \inf \{t > 0 : X_2(t) \in A\} \quad (\text{hitting time}), \quad D_A^{(2)} := \inf \{t \geq 0 : X_2(t) \in A\} \quad (\text{entry time}).$$

Then both $\tau_A^{(2)}$ and $D_A^{(2)}$ are $(\mathcal{F}(t))$ -stopping times.

Corollary 5.6. Let $\mu := (\mu_t)_{t \geq 0} \subset \mathcal{P}(\mathbb{R}^d)$ be a weakly continuous solution to the nonlinear Fokker-Planck equation (5.1), such that $\mu_t := u(t, x)dx$, $t \geq 0$ for a jointly measurable density $u : [0, \infty) \times \mathbb{R}^d \rightarrow \mathbb{R}_+$. Furthermore, assume that the "linearized" conditions

$$\mathbf{H}_{\mathbf{a}^u, \mathbf{b}^u}^{\sqrt{u}}, \quad \mathbf{H}_{\lesssim u}(\mathbf{L}^u), \quad \mathbf{H}_0^{\text{TV}}(\mathbf{L}^u) \quad \text{are fulfilled.}$$

Let X be the process provided by Theorem 5.4, $\alpha > 0$, and consider the mapping

$$\mathbb{R}^d \supset A \longmapsto \text{Cap}_{(\mu_t)}^\alpha(A) := \inf \left\{ \mathbb{E}_{\delta_0 \otimes \mu_0} \left\{ e^{-\alpha D_G^{(2)}} \right\} : A \subset G, G \subset \mathbb{R}^d \text{ open} \right\}.$$

Then $\text{Cap}_{(\mu_t)}^\alpha$ is a Choquet capacity on \mathbb{R}^d endowed with the norm topology. Moreover, $\text{Cap}_{(\mu_t)}^\alpha$ is tight.

Proof. Recall that X is a path-continuous right process on E . Therefore, we can consider the (tight) Choquet capacity $\text{Cap}_{\delta_0 \otimes \mu_0}^\alpha$ given by (6.10), with \mathcal{T} being the norm topology on E , $\nu = \delta_0 \otimes \mu_0$, and $f = 1/\alpha$. Then it is straightforward to check that

$$\text{Cap}_{(\mu_t)}^\alpha(A) = \text{Cap}_{\delta_0 \otimes \mu_0}^\alpha([0, \infty) \times A), \quad A \subset \mathbb{R}^d.$$

is a Choquet capacity on \mathbb{R}^d endowed with the norm topology, which is also tight. \square

Remark 5.7. The capacity $\text{Cap}_{\mu_0}^\alpha$ constructed in Corollary 5.6 depends only on the given solution $\mu = (\mu_t)_{t \geq 0}$ of the nonlinear Fokker-Planck equation (5.1), through the corresponding superposition on the path space $C([0, \infty); \mathbb{R}^d)$ provided by Theorem 2.5.

5.2 Example: Generalized porous media equations

We place in the framework of [5] and consider the nonlinear Fokker-Planck equation

$$\begin{aligned} \frac{d}{dt} u(t, x; u_0) &= \Delta \beta(x, u(t, x; u_0)) - \text{div}(D(x)b(u(t, x; u_0))u(t, x; u_0)), \quad t > 0, x \in \mathbb{R}^d \\ u(0, \cdot) &= u_0(\cdot). \end{aligned} \quad (5.5)$$

As we will see, like in Section 4, in order to apply our results to (5.5), we have to construct a reference density satisfying the linearized conditions $\mathbf{H}_{\mathbf{a}^u, \mathbf{b}^u}^{\sqrt{u}}$, $\mathbf{H}_{\lesssim u}(\mathbf{L}^u)$, $\mathbf{H}_0^{\text{TV}}(\mathbf{L}^u)$.

To this end, we consider the following set of assumptions on (β, D, b) , most of them being taken over from [5]:

$$\beta : \mathbb{R}^2 \times \mathbb{R} \rightarrow \mathbb{R} \quad \text{and} \quad D : \mathbb{R}^d \times \mathbb{R} \rightarrow \mathbb{R}, \quad b : \mathbb{R} \rightarrow \mathbb{R}$$

satisfy the following:

(H_β) $\beta \in C^2(\mathbb{R}^d \times \mathbb{R}; \mathbb{R})$ is such that:

$$(H_{\beta,1}) \quad \beta_r(x, r) := \frac{d}{dr}\beta(x, r) > 0, \quad \beta_r \in L^\infty(\mathbb{R}^d \times \mathbb{R}), \quad \nabla_x \beta_r, \beta_{rr} \in L^\infty(\mathbb{R}^d \times [0, N]), \quad N > 0,$$

$$\beta(x, 0) = 0, \quad x \in \mathbb{R}^d.$$

$$(H_{\beta,2}) \quad \int_{\mathbb{R}^d} \sup_{|r| \leq N} |\Delta_x \beta(x, r)| dx < \infty, \quad N > 0.$$

(H_D) For $m > d/2$ if $d \geq 2$, and $m = 1$ if $d = 1$,

$$D \in L^\infty(\mathbb{R}^d; \mathbb{R}^d), \quad \operatorname{div} D \in L^m_{\text{loc}}(\mathbb{R}^d), \quad (\operatorname{div} D)^- \in L^\infty(\mathbb{R}^d).$$

(H_b) $b \in C^1(\mathbb{R}) \cap C_b(\mathbb{R})$ and $b \geq 0$.

$(H_{b,\beta})$ $|b(r)r - b(s)s| \leq c|\beta(x, r) - \beta(x, s)|$, $x \in \mathbb{R}^d, s, t \in \mathbb{R}$, for some $c \in (0, \infty)$.

The following well posedness result follows from [5, Theorem 6.1]: *For every $u_0 \in L^1(\mathbb{R}^d)$ there exists a unique mild solution $u(\cdot, \cdot; u_0) \in C([0, \infty); L^1(\mathbb{R}^d))$ to the nonlinear Fokker-Planck equation (5.5); see [5, Definition 1.1] for details on the notion of solution.*

In fact, the above mentioned well-posedness result is valid under more general assumptions on the coefficients, see [5, Theorem 2.1 & Theorem 6.1]. The role of the assumptions $(H_\beta), (H_D), (H_b), (H_{b,\beta})$ from above, besides the well-posedness part, is that they also guarantee the following key properties of the solution, according to [5].

(1) Cf. [5, (1.16) & (2.4)]:

$$u(t + s, \cdot; u_0) = u(t, \cdot; u(s, \cdot; u_0)), \quad t, s \geq 0.$$

(2) Cf. [5, (1.17) & (2.4)]:

$$\|u(t, \cdot; u_0) - u(t, \cdot; v_0)\|_{L^1(\mathbb{R}^d)} \leq \|u_0 - v_0\|_{L^1(\mathbb{R}^d)}, \quad t \geq 0, u_0, v_0 \in L^1(\mathbb{R}^d).$$

(3) Cf. [5, Theorem 6.1]:

If $u_0(x)dx \in \mathcal{P}(\mathbb{R}^d)$ then $u_t(x)dx \in \mathcal{P}(\mathbb{R}^d)$ for every $t > 0$.

(4) Cf. [5, Theorem 6.1]:

If $u_0 \in L^1(\mathbb{R}^d) \cap L^\infty(\mathbb{R}^d)$ then $u(\cdot, \cdot; u_0) \in L^\infty((0, N) \times \mathbb{R}^d)$, for every $N > 0$.

(5) Cf. [5, Theorems 2.2 & 5.2 & Section 6]:

If $u_0 \in L^1(\mathbb{R}^d) \cap L^2(\mathbb{R}^d)$ then

$$(5.1) \quad \beta(\cdot, u(\cdot, \cdot; u_0)) \in L^2((0, N); H^1(\mathbb{R}^d)) \cap L^\infty((0, N); L^2(\mathbb{R}^d))$$

$$(5.2) \quad \frac{d}{dt} u(t, \cdot; u_0) \in L^2((0, N); H^{-1}(\mathbb{R}^d)), \quad N > 0.$$

Lemma 5.8. *Under (H_β) there exists a function $\hat{\beta}$*

$$\hat{\beta} : \mathbb{R}^d \times \mathbb{R} \rightarrow \mathbb{R}, \quad \beta(x, \hat{\beta}(x, r)) = r = \hat{\beta}(x, \beta(x, r)), \quad (x, r) \in \mathbb{R}^d \times \mathbb{R},$$

such that the following properties hold

$$\hat{\beta} \in C^1(\mathbb{R}^d \times \mathbb{R}), \quad \hat{\beta}_r, \nabla_x \hat{\beta}_r \in L^\infty(\mathbb{R}^d \times [0, N]), \quad N > 0 \tag{5.6}$$

Proof. First of all, note that by the implicit function theorem, for every $r \in \mathbb{R}$ there exists a unique $C^1(\mathbb{R}^d)$ -function $\hat{\beta}(\cdot, r)$ such that

$$\beta(x, \hat{\beta}(x, r)) = r, \quad x \in \mathbb{R}^d.$$

Now, by the inverse function theorem we have that for every $x \in \mathbb{R}^d$ there is a $C^1(\mathbb{R})$ -function $\bar{\beta}(x, \cdot)$ such that

$$\beta(x, \bar{\beta}(x, r)) = r = \bar{\beta}(x, \beta(x, r)), \quad r \in \mathbb{R}.$$

Clearly, from the injectivity of $\beta(x, \cdot)$ for each $x \in \mathbb{R}^d$ we deduce that $\hat{\beta} = \bar{\beta}$ and that $\hat{\beta}$ is C^1 in each variable. Moreover

$$\nabla_x \hat{\beta}(x, r) = \frac{-\nabla_x \beta(x, t)|_{t=\hat{\beta}(x, r)}}{\beta_r(x, t)|_{t=\hat{\beta}(x, r)}}, \quad \hat{\beta}_r(x, r) = \frac{1}{\beta_r(x, t)|_{t=\hat{\beta}(x, r)}}, \quad (x, r) \in \mathbb{R}^d \times \mathbb{R}.$$

Now, the properties in (5.6) follow immediately from (H_β) . \square

Proposition 5.9. *Assume that β, D, b satisfy conditions $(H_\beta), (H_D), (H_b), (H_{b,\beta})$, let $u_0 \in L^1(\mathbb{R}^d) \cap L^\infty(\mathbb{R}^d)$ and consider $u(\cdot, \cdot) := u(\cdot, \cdot; u_0)$ to be the solution to (5.5). Then*

$$u(\cdot, \cdot; u_0) \in L^2([0, T]; H^1(\mathbb{R}^d)), \quad T > 0.$$

Proof. Let $\hat{\beta}$ be the function provided by Lemma 5.8, so that

$$u(t, x) = \hat{\beta}(x, \beta(x, u(t, x))), \quad x, t \in \mathbb{R}^d \times \mathbb{R} \text{ a.e.}$$

Consequently, the statement follows if we show that

$$(t, x) \mapsto \nabla_x [\hat{\beta}(x, \beta(x, u(t, x)))] \in L^2([0, T]; L^2(\mathbb{R}^d)).$$

To this end, by (5.6) and the fact that we already know that $\beta(u) \in L^2([0, T]; H^1(\mathbb{R}^d))$, we have

$$\begin{aligned} \nabla_x [\hat{\beta}(x, \beta(x, u(t, x)))] &= \underbrace{\nabla_x \hat{\beta}(x, r)|_{r=\beta(x, u(t, x))}}_{\in L^\infty([0, T] \times \mathbb{R}^d)} + \underbrace{\hat{\beta}_r(x, \beta(x, u(t, x)))}_{\in L^\infty([0, T] \times \mathbb{R}^d)} \underbrace{\nabla_x [\beta(x, u(t, x))]}_{\in L^2([0, T]; L^2(\mathbb{R}^d))} \\ &\in L^2([0, T]; L^2(\mathbb{R}^d)). \end{aligned}$$

\square

Let $0 \leq u_0 \in L^\infty(\mathbb{R}^d)$ such that $\int_{\mathbb{R}^d} u_0(x) dx = 1$, and let $u(\cdot, \cdot; u_0)$ be the corresponding solution presented above. If we set $\mu_t(dx) := u(t, x; u_0) dx$, $t \geq 0$, then (5.5) recasts as:

For every $f \in C_c^\infty((0, \infty) \times \mathbb{R}^d)$

$$\int_0^\infty \int_{\mathbb{R}^d} \frac{d}{dt} f(t, x) + \underbrace{\frac{\beta(x, u(t, x))}{u(t, x)}}_{=: a(x, \mu_t)} \Delta f(t, x) + \underbrace{\langle D(x)b(u(t, x)), \nabla f(t, x) \rangle}_{=: b(x, \mu_t)} \mu_t(dx) dt = 0, \quad (5.7)$$

$$u(0, \cdot) = u_0(\cdot). \quad (5.8)$$

Since $\mu := (\mu_t)_{t \geq 0} \subset \mathcal{P}(\mathbb{R}^d)$ is also continuous in the total variation norm, we thus have that μ is a weakly continuous solution (in the sense of Definition 5.2) to the non-linear Fokker-Planck equation

$$\begin{aligned} \frac{d}{dt} \mu_t &= \mathbf{L}_{t, \mu_t}^* \mu_t, \quad (5.9) \\ \underbrace{\mathbf{L}_{t, \mu_t}^u}_{\mathbf{L}_t^u} f(t, x) &:= \underbrace{b(x, \mu_t)}_{=: b^u(t, x)} \nabla_x f(t, x) + \frac{1}{2} \underbrace{2a(x, \mu_t)}_{=: a^u(t, x)} \Delta_x f(t, x), \quad (t, x) \in (0, \infty) \times \mathbb{R}^d, \end{aligned}$$

where $b(x, \mu_t)$ and $a(x, \mu_t)$ are given in (5.7).

Now, we need a convenient well-posedness result for the *linearized* Fokker-Planck equation

$$\frac{d}{dt} \nu_t = (\mathbf{L}_t^u)^* \nu_t, \quad \text{where } \mathbf{L}_t^u \text{ is defined in (5.9) with } u \text{ being fixed in the coefficients.} \quad (5.10)$$

Proposition 5.10. *Assume that β, D, b satisfy conditions $(H_\beta), (H_D), (H_b), (H_{b,\beta})$, let $0 \leq u_0 \in L^\infty(\mathbb{R}^d)$ such that $u_0(x) dx \in \mathcal{P}(\mathbb{R}^d)$, and let $u(\cdot, \cdot; u_0)$, $\mu_t(dx) := u(t, x; u_0) dx$, $t \geq 0$ be the solution to (5.5)/(5.7). Let $0 \leq s < T$ and $\nu_s = \nu_s(x) dx$ such that $0 \leq \nu_s \in L^1(\mathbb{R}^d) \cap L^\infty(\mathbb{R}^d)$. Then the following assertions hold*

(i) Then there exists a solution $(\nu_t = v(t, x)dx)_{t \in [s, T]}$ to the linearized Fokker-Planck equation (5.10) such that

$$v \in L^\infty([0, T]; L^1 \cap L^\infty(\mathbb{R}^d)) \quad \text{and} \quad \sqrt{a^u} \nabla v \in L^2([0, T]; L^2(\mathbb{R}^d)).$$

Moreover, it follows that

$$v \in C([0, T]; L^2(\mathbb{R}^d)), \quad T > 0. \quad (5.11)$$

(ii) If $\nu^{(1)} := (\nu_t^{(1)} = v^{(1)}(t, x)dx)_{t \in [s, T]}$ and $\nu_t^{(2)} := (\nu_t^{(2)} = v^{(2)}(t, x)dx)_{t \in [s, T]}$ are two right weakly-continuous solutions to the linearized Fokker-Planck equation (5.10), such that $v^{(1)}, v^{(2)} \in L^\infty([s, T] \times \mathbb{R}^d)$, sharing the same initial condition $v_s(x)dx$, then

$$\nu^{(1)} = \nu^{(2)}.$$

Proof. (i). The main idea is to rely on the results [38, Proposition 2 & Proposition 3], so let v as in the statement and note that (5.10) written in divergence form [38, (5.8)] becomes

$$\frac{d}{dt} v + \operatorname{div} \left(\left(b^u - \frac{1}{2} \nabla a^u \right) v \right) - \frac{1}{2} \operatorname{div}(a^u \nabla v) = 0$$

In fact, it turns out that the Girsanov-type result [38, Proposition 3] is much more suitable to our context, so we go on and rewrite the above equation in the special form

$$\frac{d}{dt} v + \operatorname{div}(\sqrt{a^u} \theta v) - \frac{1}{2} \operatorname{div}(a^u \nabla v) = 0, \quad \text{where } \theta := \frac{b^u}{\sqrt{a^u}} - \frac{\nabla a^u}{2\sqrt{a^u}}.$$

Now, assertion (i) follows by [38, Proposition 4] as soon as we check the following conditions as soon as we check the following conditions (see [38, (6.12)-(6.13)]),

$$(i.1) \quad \sqrt{a^u} \in L^2([0, T]; W_{\text{loc}}^{1,2}(\mathbb{R}^d)), \quad \frac{\sqrt{a^u}}{1+\|x\|} \in L^2([0, T]; L^2 + L^\infty(\mathbb{R}^d))$$

$$(i.2) \quad \theta \in L^2([0, T]; L^2 + L^\infty(\mathbb{R}^d)).$$

Let us check (i.1) first. Note that by condition (H_β) on β we have

$$\exists c \in (0, \infty) : \quad a^u = \beta(u)/u \leq c, \quad \frac{|\beta_r(u)u - \beta(u)|}{u^2} \leq c, \quad \frac{|\nabla_x \beta(x, u)|}{u} \leq c. \quad (5.12)$$

Now, by the first bound in (5.12), the second part of (i.1) is immediate. Regarding the first part, note that by (H_β) and the fact that $u \in L^\infty([0, T]; \mathbb{R}^d)$ we have

$$\exists \delta > 0 : \quad a^u = \frac{\beta(u)}{u} \geq \delta. \quad (5.13)$$

Moreover, since $\beta \in C^2$ whilst $u \in L^\infty$, by $(H_{\beta,1})$ we have

$$\begin{aligned} \frac{|\beta_r(x, u)u - \beta(x, u)|}{u^2} &\leq \frac{|\sup_{\xi \in [0, \|u\|_\infty]} \beta_{rr}(x, \xi)u^2|}{2u^2} = \frac{|\sup_{\xi \in [0, \|u\|_\infty]} \beta_{rr}(x, \xi)|}{2} \\ &\in L^\infty([0, T] \times \mathbb{R}^d). \end{aligned}$$

Consequently, using also (5.13), (5.12), as well as the fact that $u \in L^2([0, T]; H^1(\mathbb{R}^d))$ due to Proposition 5.9,

$$\begin{aligned} |\nabla \sqrt{a^u}| &= \frac{|\nabla a^u|}{\sqrt{a^u}} \leq \frac{|\nabla a^u|}{\sqrt{\delta}} \leq \frac{1}{\sqrt{\delta}} |\nabla a^u| \\ &\leq \frac{1}{\sqrt{\delta}} |\nabla u| \frac{|\nabla_x \beta(x, u)u + \beta_r(u)u - \beta(u)|}{u^2} \\ &\leq \frac{1}{\sqrt{\delta}} |\nabla u| \left(\frac{|\nabla_x \beta(x, u)|}{u} + \frac{|\beta_r(u)u - \beta(u)|}{u^2} \right) \\ &\leq \frac{2c}{\sqrt{\delta}} |\nabla u| \in L^2([0, T]; L^2(\mathbb{R}^d)). \end{aligned}$$

Thus, assertion (i.1) is completely proved. Assertion (i.2) is also clear now, by (5.13), the fact that $b \in C_b(\mathbb{R})$ according to (H_b) , and the property $\nabla\sqrt{a^u} \in L^2([0, T]; L^2(\mathbb{R}^d))$ which was proved above.

In order to completely prove assertion (i), we need to justify that $v \in C([0, T]; L^2(\mathbb{R}^d))$. To this end, note that by the first part of (i) proved above, using also (5.13), it follows that $v \in L^2([0, T]; H^1(\mathbb{R}^d))$. To conclude, we can now apply Proposition 4.5 for L_t^u instead of L_t , for which the conditions from Proposition 4.5 are clearly satisfied.

(ii) This follows directly from [5, Corollary 3.5] as well as the last line on page 34 from [5]. \square

Theorem 5.11. *Assume that β, D, b satisfy conditions (i)-(iv), let $0 \leq u_0 \in L^\infty(\mathbb{R}^d)$ such that $\int_{\mathbb{R}^d} u_0(x) dx = 1$, and let $u(\cdot, \cdot; u_0)$, $\mu_t(dx) := u(t, x; u_0)dx, t \geq 0$ be the solution to (5.5)/(5.7). Then the "linearized" conditions*

$$\mathbf{H}_{\mathbf{a}^u, \mathbf{b}^u}^{\sqrt{u}}, \quad \mathbf{H}_{\lesssim u}(\mathbf{L}^u), \quad \mathbf{H}_0^{\text{TV}}(\mathbf{L}^u)$$

are fulfilled, where a^u, b^u , and \mathbf{L}^u are given in (5.9). Consequently, the conclusions of Theorem 5.4 and Corollary 5.6 are valid in this situation.

Proof. Condition $\mathbf{H}_{\mathbf{a}^u}$ is fulfilled by e.g. (i.1) from the proof of Proposition 5.10, (5.12), (5.13), and the fact that $u \in L^\infty$. Condition $\mathbf{H}_{\mathbf{b}^u}$ follows from (H_b) and (H_D) . Condition $\mathbf{H}^{\sqrt{u}}$ follow from Proposition 4.4. Thus, condition $\mathbf{H}_{\mathbf{a}^u, \mathbf{b}^u}^{\sqrt{u}}$ is fulfilled.

Finally, it is straightforward to see that condition $\mathbf{H}_{\lesssim u}(\mathbf{L}^u)$ is fulfilled by Proposition 5.10.

Finally, condition $\mathbf{H}_0^{\text{TV}}(\mathbf{L}^u)$ is also satisfied due to Proposition 5.10, (i), and the dominated convergence, using also that $u \in C([0, \infty); L^1(\mathbb{R}^d))$. \square

6 Elements of right processes

Throughout this section we follow mainly the terminology of [33], [60], and [14]. The aim here is two-fold:

- We present a concise introduction to the theory of probabilistic potential theory, with special emphasis on those tools that are fundamentally used in order to derive the main results from Section 2 and Section 3.
- We prove new results on right processes that are of general nature and which are, as well, needed in Section 2 and Section 3; see e.g. Lemma 6.23, Proposition 6.24, Proposition 6.39, Proposition 6.44, Proposition 6.56, as well as those from Section 6.5.

6.1 Fundamentals on right processes and their potential theory

Let (E, \mathcal{B}) be a Lusin measurable space, i.e. it is measurably isomorphic to a Borel subset of a compact metric space endowed with the corresponding Borel σ -algebra. We denote by $(b)p\mathcal{B}$ the set of all numerical, (bounded) positive \mathcal{B} -measurable functions on E . Throughout, by $\mathcal{U} = (U_\alpha)_{\alpha>0}$ we denote a resolvent family of (sub-)Markovian kernels on (E, \mathcal{B}) ; see e.g. [14, Subsection 1.1]. If $\beta > 0$, we set $\mathcal{U}_\beta := (U_{\beta+\alpha})_{\alpha>0}$.

$A \subset E$ is called *universally \mathcal{B} -measurable* if for every (positive) finite measure μ on (E, \mathcal{B}) there exist $A_1, A_2 \in \mathcal{B}$ such that $A_1 \subset A \subset A_2$ and $\mu(A_2 \setminus A_1) = 0$. By \mathcal{B}^u we denote the σ -algebra of all universally \mathcal{B} -measurable sets in E .

Definition 6.1. *A universally \mathcal{B} -measurable function $v : E \rightarrow \overline{\mathbb{R}}_+$ is called \mathcal{U} -excessive provided that $\alpha U_\alpha v \leq v$ for all $\alpha > 0$ and $\sup_\alpha \alpha U_\alpha v = v$ point-wise; by $\mathcal{E}(\mathcal{U})$ we denote the convex cone of all \mathcal{B} -measurable \mathcal{U} -excessive functions.*

Remark 6.2. *Recall that if $f \in p\mathcal{B}$, then $U_\alpha f$ is \mathcal{U}_α -excessive.*

If a universally \mathcal{B} -measurable function $w : E \rightarrow \overline{\mathbb{R}}_+$ is merely \mathcal{U}_β -supermedian (i.e. $\alpha U_{\beta+\alpha} w \leq w$ for all $\alpha > 0$), then its \mathcal{U}_β -excessive regularization is defined as $\widehat{w} := \sup_\alpha \alpha U_{\beta+\alpha} w$. The function \widehat{w} is \mathcal{U}_β -excessive and $\widehat{w} \in \mathcal{E}(\mathcal{U}_\beta)$ provided that w is in addition \mathcal{B} -measurable.

Definition 6.3. *The fine topology on E (associated with \mathcal{U}) is the coarsest topology on E such that every \mathcal{U}_α -excessive function is continuous for some (hence all) $\alpha > 0$.*

Proposition 6.4. *If u is \mathcal{U}_β -excessive for some $\beta > 0$, then the set $[u = 0]$ is finely open.*

Definition 6.5. *A topology τ on E is called natural if it is a Lusin topology (i.e. (E, τ) is homeomorphic to a Borel subset of a compact metrizable space) which is coarser than the fine topology, and whose Borel σ -algebra is \mathcal{B} .*

Remark 6.6. (i) *The necessity of considering natural topologies comes from the fact that, in general, the fine topology is neither metrizable, nor countably generated.*

(ii) *Let (E, \mathcal{B}) be a Lusin measurable space, and $(f_n)_{n \geq 1} \subset \mathcal{B}$ be such that it separates the points of E . Then the topology on E generated by $(f_n)_{n \geq 1}$ is a Lusin topology; see e.g. [60, Theorem A2.12] for the similar Radonian case.*

There is a convenient class of natural topologies to work with (as we do in Section 2), especially when the aim is to construct a right process associated with \mathcal{U} (see Definition 6.10). These topologies are called Ray topologies, and are defined as follows.

Definition 6.7. (i) *If $\beta > 0$ then a Ray cone associated with \mathcal{U}_β is a convex cone \mathcal{R} of bounded, \mathcal{B} -measurable \mathcal{U}_β -excessive functions which is separable in the supremum norm, min-stable, contains the constant function 1, generates \mathcal{B} , $U_\alpha(\mathcal{R}) \subset \mathcal{R}$ for all $\alpha > 0$, and $U_\beta((\mathcal{R} - \mathcal{R})_+) \subset \mathcal{R}$.*

(ii) *A Ray topology on E is a topology generated by a Ray cone \mathcal{R} , and it is denoted by $\tau_{\mathcal{R}}$.*

In order to ensure the existence of a Ray topology, as well as many other useful properties related to the fine topology and excessive functions, we need the following condition:

(H_C). \mathcal{C} is a min-stable convex cone of non-negative \mathcal{B} -measurable functions such that

- (i) $1 \in \mathcal{C}$ and there exists a countable subset of \mathcal{C} which separates the points of E ,
- (ii) $U_\alpha(\mathcal{C}) \subset \mathcal{C}$ for all $f \in \mathcal{C}$ and $\alpha > 0$,
- (iii) $\lim_{\alpha \rightarrow \infty} \alpha U_\alpha f = f$ pointwise on E for all $f \in \mathcal{C}$.

The following result is basically Corollary 2.3 from [29]:

Proposition 6.8. *The following assertions are equivalent.*

- (i) *There exists a cone \mathcal{C} such that (H_C) is fulfilled.*
- (ii) *$\mathcal{E}(\mathcal{U}_\beta)$ is min stable, contains the constant functions, and generates \mathcal{B} for one (hence all) $\beta > 0$.*
- (iii) *For any $\beta > 0$ there exists a Ray cone associated with \mathcal{U}_β .*

For the rest of this section we assume that one (hence all) of the assertions from 6.8 is valid. This is always satisfied when there is a right Markov process with resolvent \mathcal{U} , whose definition is given below.

Remark 6.9. (i) *It is clear that any Ray topology is a natural topology. A useful converse is true: For any natural topology there exists a finer Ray topology; see Proposition 2.1 from [16]. In fact, a key ingredient here is that given a countable family of finite \mathcal{U}_β -excessive functions, one can always find a Ray cone that contains it; we shall use this fact later on.*

(ii) *As a matter of fact, if (H_C) holds, then one can explicitly construct plenty of Ray topologies, following e.g. [14], Proposition 1.5.1, or [29], Proposition 2.2. As we shall need such a construction later on, let us detail it here. Let $\mathcal{A}_0 \subset \text{bp}\mathcal{B}$ be countable and such that it separates the set of finite measures on (E, \mathcal{B}) . The Ray cone generated by \mathcal{A}_0 and associated with \mathcal{U}_β is defined inductively as follows:*

$$\begin{aligned} \mathcal{R}_0 &:= U_\beta(\mathcal{A}_0) \cup \mathbb{Q}_+, \\ \mathcal{R}_{n+1} &:= \mathbb{Q}_+ \cdot \mathcal{R}_n \cup \left(\sum_f \mathcal{R}_n \right) \cup \left(\bigwedge_f \mathcal{R}_n \right) \cup \left(\bigcup_{\alpha \in \mathbb{Q}_+} U_\alpha(\mathcal{R}_n) \right) \cup U_\beta((\mathcal{R}_n - \mathcal{R}_n)_+), \end{aligned}$$

where by $\sum_f \mathcal{R}_n$ resp. $\bigwedge_f \mathcal{R}_n$ we denote the space of all finite sums (resp. infima) of elements from \mathcal{R}_n . Then, the Ray cone $\mathcal{R}_{\mathcal{A}_0}$ generated by \mathcal{A}_0 is obtained by taking the closure of $\bigcup_n \mathcal{R}_n$ w.r.t. the supremum norm.

Right processes. Let now $X = (\Omega, \mathcal{F}, \mathcal{F}(t), X(t), \theta(t), \mathbb{P}_x)$ be a normal Markov process with state space E , shift operators $\theta(t) : \Omega \rightarrow \Omega$, $t \geq 0$, and lifetime ζ ; if Δ denotes the cemetery point attached to E , then any numerical function f on E shall be extended to Δ by setting $f(\Delta) = 0$.

We assume that X has the resolvent \mathcal{U} fixed above, i.e. for all $f \in b\mathcal{B}$ and $\alpha > 0$

$$U_\alpha f(x) = \mathbb{E}_x \int_0^\infty e^{-\alpha t} f(X(t)) dt, \quad x \in E.$$

To each probability measure μ on (E, \mathcal{B}) we associate the probability

$$\mathbb{P}^\mu(A) := \int \mathbb{P}^x(A) \mu(dx)$$

for all $A \in \mathcal{F}$, and we consider the following enlarged filtration

$$\widetilde{\mathcal{F}}(t) := \bigcap_\mu \mathcal{F}^\mu(t), \quad \widetilde{\mathcal{F}} := \bigcap_\mu \mathcal{F}^\mu,$$

where \mathcal{F}^μ is the completion of \mathcal{F} under \mathbb{P}^μ , and $\mathcal{F}^\mu(t)$ is the completion of $\mathcal{F}(t)$ in \mathcal{F}^μ w.r.t. \mathbb{P}^μ ; in particular, $(x, A) \mapsto \mathbb{P}^x(A)$ is assumed to be a kernel from (E, \mathcal{B}^u) to (Ω, \mathcal{F}) , where \mathcal{B}^u denotes the σ -algebra of all universally measurable subsets of E .

Definition 6.10. *The Markov process X is called a right (Markov) process if the following additional hypotheses are satisfied:*

- (i) *The filtration $(\mathcal{F}(t))_{t \geq 0}$ is right continuous and $\mathcal{F}(t) = \widetilde{\mathcal{F}}(t)$, $t \geq 0$.*
- (ii) *For one (hence all) $\alpha > 0$ and for each $f \in \mathcal{E}(\mathcal{U}_\alpha)$ the process $f(X)$ has right continuous paths \mathbb{P}^x -a.s. for all $x \in E$.*
- (iii) *There exists a natural topology on E with respect to which the paths of X are \mathbb{P}^x -a.s. right continuous for all $x \in E$.*

Remark 6.11 (Strong Markov property). *Recall that if X is a right process on (E, \mathcal{B}) with transition function $(P_t)_{t \geq 0}$, then it is a strong Markov process, i.e. for any $\mu \in \mathcal{P}(E)$, $t \geq 0$, and τ an $(\mathcal{F}(t))$ -stopping time we have*

$$\mathbb{E}_\mu \{f(X(t + \tau)) \mid \mathcal{F}_\tau\} = P_t f(X(\tau)) \quad \mathbb{P}_\mu\text{-a.s.}, \quad f \in b\mathcal{B}(E); \text{ recall that } f(\Delta) := 0.$$

Moreover, for any μ, τ as above and any $Y \in b\mathcal{F}$, the strong Markov property with shifts holds, namely

$$\mathbb{E}_\mu \{Y \circ \theta(\tau) \mid \mathcal{F}_\tau\} = \mathbb{E}_{X(\tau)} \{Y\} \quad \mathbb{P}_\mu\text{-a.s.}, \quad Y \in b\mathcal{F}; \text{ here, } Y \circ \theta(\infty) := 0.$$

Sometimes, for shortness and if there is no risk of confusion, we use $(X, \mathbb{P}_x, x \in E)$ to denote a right process X , instead of specifying the full tuple $X = (\Omega, \mathcal{F}, \mathcal{F}_t, X(t), \theta(t), \mathbb{P}^x)$.

Remark 6.12. *It is worth to mention that if $X = (\Omega, \mathcal{F}, \mathcal{F}(t), X(t), \theta(t), \mathbb{P}_x)$ is a right process as defined above, and τ is a given natural topology on E , then then one can construct another right process $\tilde{X} = (\tilde{\Omega}, \tilde{\mathcal{F}}, \tilde{\mathcal{F}}(t), \tilde{X}(t), \tilde{\theta}(t), \tilde{\mathbb{P}}_x)$ sharing the same resolvent (or transition function) as X , i.e. X and \tilde{X} are equivalent in the sense of [33, Definition 4.1], such that assertion (iii) from Definition 6.10 can be replaced by a stronger one, namely:*

- (iii') *The path $[0, \infty) \ni t \rightarrow \tilde{X}(t, \omega) \in E$ is right-continuous with respect to τ for every $\omega \in \tilde{\Omega}$.*

Indeed, one can construct such \tilde{X} by considering first a second process equivalent to X which is of function space type, according to [33, Definition 4.2 and Theorem 4.3], and then simply deleting from the path space of this second process the set Λ of those paths that are not right-continuous with respect to τ . Such a removal is allowed since Λ is universally measurable according to [40, Ch. IV, Theorem 34].

The probabilistic description of the fine topology is given by the following key result, according to [33, Chapter II, Theorem 4.8], or [60, Proposition 10.8 and Exercise 10.18], and [22, Corollary A.10]. The second assertion states the relation between excessive functions and right-continuous supermartingales (cf. e.g., [19, Proposition 1], see also [20]).

Theorem 6.13. *Let X be a right process and f a universally \mathcal{B} -measurable function. Then the following assertions hold.*

(i) *The function f is finely continuous if and only if $(f(X(t)))_{t \geq 0}$ has \mathbb{P}^x -a.s. right continuous paths for all $x \in E$. In particular, X has a.s. right continuous paths in any natural topology on E .*

(ii) *If $\beta > 0$ then the function f is \mathcal{U}_β -excessive if and only if $(e^{-\beta t} f(X(t)))_{t \geq 0}$ is a right continuous \mathcal{F}_t -supermartingale w.r.t. \mathbb{P}^x for all $x \in E$.*

Remark 6.14. *Inspecting the proof of [33, Chapter II, Theorem 4.8], we see that if f is universally \mathcal{B} -measurable and the mapping $[0, \infty) \ni t \mapsto f(X(t)) \in \mathbb{R}$ is right continuous merely at $t = 0$ \mathbb{P}^x -a.s. for all $x \in E$, then f is finely continuous.*

Remark 6.15 (P_t is finely Feller). *A remarkable property which has been relatively overlooked in the literature is that given a right process X on a general Lusin measurable space (E, \mathcal{B}) , the operators P_t have the property that $P_t f$ is finely continuous whenever f is finely continuous, bounded, and \mathcal{B}^u -measurable; see e.g. [60, Exercise 10.24], [33, Exercise 4.14], or [20, Lemma 2.8] for a proof of this result. In other words, P_t is Feller with respect to the fine topology.*

Proposition 6.16. *Let $X = (\Omega, \mathcal{F}, \mathcal{F}(t), X(t), \theta(t), \mathbb{P}_x), x \in E$ and $\tilde{X} = (\tilde{\Omega}, \tilde{\mathcal{F}}, \tilde{\mathcal{F}}(t), \tilde{X}(t), \tilde{\theta}(t), \tilde{\mathbb{P}}_x, x \in E)$ be two right processes on a Lusin measurable space E sharing the same resolvent \mathcal{U} . Further, let τ be a natural (hence Lusin) topology on E (with respect to \mathcal{U}), such that X has continuous paths with respect to τ \mathbb{P}_μ -a.s. for some $\mu \in \mathcal{P}$. Then*

$$\mathbb{P}_\mu \circ (X)^{-1} = \tilde{\mathbb{P}}_\mu \circ (\tilde{X})^{-1} \quad \text{as laws on the path-space } C([0, \infty); E, \tau), \quad (6.1)$$

where $C([0, \infty); E, \tau)$ denotes the space of continuous functions defined on $[0, \infty)$ with values in E endowed with the topology τ .

Proof. Let D be the countable set of dyadics in $[0, \infty)$, and denote by W^E the space of the restrictions to D of all paths from $[0, \infty)$ to E which are τ -continuous. Also, consider the product E^D endowed with the canonical σ -algebra, and consider the laws

$$\nu := \mathbb{P}_\nu \circ (X)^{-1} \quad \text{and} \quad \tilde{\nu} := \tilde{\mathbb{P}}_\mu \circ (\tilde{X})^{-1}$$

of the two processes X and \tilde{X} on E^D . Now, on the one hand any distribution on E^D is uniquely determined by its finite dimensional marginals, hence, by the Markov property, ν and $\tilde{\nu}$ are uniquely determined by their one-dimensional marginals. The latter marginals are in turn uniquely determined by their Laplace transforms, and since X and \tilde{X} share the same resolvent \mathcal{U} , we deduce that

$$\nu = \tilde{\nu}.$$

Next, by [40, Chapter IV, Theorem 34; see also pages 91-92], we have that W^E is a universally measurable subset of E^D . Consequently, we have

$$\nu(W^E) = \tilde{\nu}(W^E) = 1,$$

so, on the one hand, the paths of \tilde{X} are restrictions to D of τ -continuous paths in E $\tilde{\mathbb{P}}_\mu$ -a.s. On the other hand, by Theorem 6.13 it follows that \tilde{X} has right-continuous paths with respect to τ \mathbb{P}_x -a.s. for every $x \in E$. This means that the right process \tilde{X} has E and has τ -continuous paths in E $\tilde{\mathbb{P}}_\mu$ -a.s., and clearly (6.1) is satisfied since a probability measure on $C([0, \infty); E, \tau)$ is uniquely determined by its finite dimensional time-marginals. \square

We end this paragraph by recalling the analytic and probabilistic descriptions of *polar* (and other "small") sets.

Definition 6.17. Let $\alpha \geq 0$

(i) If $u \in \mathcal{E}(\mathcal{U}_\alpha)$ and $A \in \mathcal{B}$, then the α -order reduced function of u on A is given by

$$R_\alpha^A u = \inf\{v \in \mathcal{E}(\mathcal{U}_\alpha) : v \geq u \text{ on } A\}.$$

$R_\alpha^A u$ is merely \mathcal{U}_α -supermedian and we denote by $B_\alpha^A u := \widehat{R_\alpha^A u}$ its \mathcal{U}_α -excessive regularization, called the balayage of u on A .

One has $B_\alpha^A u = R_\alpha^A u$ on $E \setminus A$. If in addition A is finely open then $B_\alpha^A u = R_\alpha^A u \in \mathcal{E}(\mathcal{U}_\alpha)$.

(ii) A set $A \in \mathcal{B}$ is called polar if $B_\alpha^A 1 \equiv 0$.

When $\alpha = 0$ we drop the index from notation and we simply write R^A and B^A , respectively.

Definition 6.18. Let m be a σ -finite measure on E .

(i) A set $A \in \mathcal{B}$ is called

- " \mathcal{U} -negligible" if $U_q(1_A) \equiv 0$ for one (hence all) $q > 0$.
- " α -polar" if $B_\alpha^A 1 \equiv 0$ for some $\alpha > 0$.
- " m -polar" if $B_\alpha^A 1 = 0$ m -a.e. for some $\alpha > 0$.
- " m -inessential" provided that it is m -negligible and $R_\alpha^A 1(x) = 0$ for all $x \in E \setminus A$, for some $\alpha > 0$.

(ii) A property is said to hold m -quasi-everywhere w.r.t. \mathcal{U} (resp. \mathcal{U} -a.e.), if there exists an m -inessential (resp. a \mathcal{U} -negligible) set N such that the property holds for all $x \in E \setminus N$; on short, we write m -q.e. instead of m -quasi-everywhere.

Remark 6.19. (i) It is well known that if $V \in \mathcal{E}(\mathcal{U}_\alpha)$ for some $\alpha \geq 0$ such that $U_1(1_{[V=\infty]}) \equiv 0$, then the set $[V = \infty]$ is polar.

(ii) We have the following probabilistic characterization due to G.A. Hunt holds (see e.g. [40]): If X is a right process with resolvent \mathcal{U} , then for all $\alpha > 0$, $u \in \mathcal{E}(\mathcal{U}_\alpha)$, $A \in \mathcal{B}$, and $x \in E$

$$R_\alpha^A u(x) = \mathbb{E}_x\{e^{-\alpha D_A} u(X(D_A))\}, \quad B_\alpha^A u(x) = \mathbb{E}_x\{e^{-\alpha T_A} u(X(T_A))\}, \quad (6.2)$$

where $D_A := \inf\{t \geq 0 : X(t) \in A\}$ is the entry time of A and $T_A := \inf\{t > 0 : X(t) \in A\}$ is the hitting time of A . In particular, A is polar if and only if $\mathbb{P}_x(T_A < \infty) = 0$ for all $x \in E$. Furthermore, A is m -polar if and only if $\mathbb{P}_x(T_A < \infty) = 0$ m -a.e. $x \in E$.

A similar interpretation holds for m -inessential sets: $A \in \mathcal{B}$ is m -inessential if and only if $m(A) = 0$ and $\mathbb{P}_x(T_A < \infty) = 0$ for every $x \in E \setminus A$.

Remark 6.20 (Probabilistic description of finely open sets). Recall that a set $D \in \mathcal{B}$ is finely open if and only if

$$\mathbb{P}_x(T_{E \setminus D} > 0) = 1, \quad x \in D;$$

see, e.g. [60, p. 52-53].

Definition 6.21 (The kernels R_α^A and B_α^A). Let $A \in \mathcal{B}$. We consider the extension of the reduced function, $\mathcal{E}(\mathcal{U}_\alpha) \ni u \mapsto R_\alpha^A u$ to a sub-Markovian kernel on \mathcal{B}^u and clearly, by (6.2) we have

$$R_\alpha^A f(x) = \mathbb{E}_x\{e^{-\alpha D_A} f(X(D_A))\}, \quad f \in p\mathcal{B}^u.$$

Analogously, B_α^A is extended to a sub-Markovian kernel on \mathcal{B}^u .

Remark 6.22. For the reader convenience, let us recall several potential theoretic facts; cf. [14], [29], and [25].

(i) If $A \in \mathcal{B}$ is finely open and U -negligible, then $A = \emptyset$.

(ii) If m is a σ -finite measure on E such that $m(A) = 0$ implies $U_1(1_A) = 0$ m -a.e. for all $A \in \mathcal{B}$, then any finely open and m -negligible set $A \in \mathcal{B}$ is m -polar.

(iii) Any m -inessential set is m -polar and m -negligible. The following converse holds (see e.g. [14], page 168): A set which is m -polar and m -negligible is the subset of an m -inessential set. For the reader's convenience we present here a sketch of the proof of this sentence. Let $A \in \mathcal{B}$. Assume that A is m -polar and m -negligible, or equivalently, $m(R_\alpha^A 1) = 0$, since $R_\alpha^A 1 = B_\alpha^A 1$ on $E \setminus A$. We may suppose that m is a finite measure and further we argue as in the first part of the proof of Theorem 1.7.29 from [14]. Let $(v_n)_n$ be a decreasing sequence of bounded U_α -excessive functions such that $v_n \geq R_\alpha^A 1$ for all n and $\inf_n v_n = R_\alpha^A 1$ m -a.e. Let $A_o := \{\inf_n v_n > 0\}$. Then $A \subset A_o$ and $R_\alpha^{A_o} 1 = 0$ on $E \setminus A_o$. Because $m(\inf_n v_n) = 0$ it follows that $m(A_o) = 0$, hence A_o is m -inessential.

(iv) If $u \in b\mathcal{B}$ and v is \mathcal{B} -measurable such that $v = u$ q.e., then $B_1^A v$ is well defined and equal to $B_1^A u$ m -a.e.

The following two results generalize several classical tools due to Dynkin [42, Chapter XII, Sections 4–5] that concern the balayage operator and are key ingredients in solving the *stochastic parabolic Dirichlet problem* in Section 3.2.

Lemma 6.23. *Let X be a right process on E , $u \in b\mathcal{B}$, $A \in \mathcal{B}$ and $\alpha \geq 0$. The following assertions hold:*

(i) *The function $B_\alpha^A u$ is finely continuous on $(E \setminus A)^{\circ f}$, where $B^{\circ f}$ denotes the interior of a set $B \subset E$ with respect to the fine topology.*

(ii) *If u is finely continuous then $B_\alpha^A u$ is finely continuous.*

Proof. Let $A \in \mathcal{B}$ and set $f := B_\alpha^A u$. In order to show that f is finely continuous (resp. on $(E \setminus A)^{\circ f}$) it is sufficient to prove that if $\varepsilon > 0$ then x is irregular for $V := f^{-1}([f(x) + \varepsilon, \infty))$ and for $W := f^{-1}((-\infty, f(x) - \varepsilon])$, for every $x \in E$ (resp. $x \in (E \setminus A)^{\circ f}$). We treat only the case of V , the one for W being similar. Fix $x \in E$. Let $(V_n)_n$ be an increasing sequence of finely closed sets such that $T_{V_n} \searrow T_V$ \mathbb{P}_x -a.s., which exists by e.g. [60, Chapter I, Exercise (10.30)] [33, Theorem 10.19]. By the zero-one law ([60, Chapter I, Corollary (3.11)] or [33, Proposition 5.17]), x is either regular for V , i.e. $\mathbb{P}_x([T_V = 0]) = 1$, or irregular for V , i.e. $\mathbb{P}_x([T_V = 0]) = 0$.

Assume that x is regular. By the strong Markov property and recalling that $u(\Delta) = 0$ we have

$$\begin{aligned} f(x) + \varepsilon &\leq \mathbb{E}_x\{f(X(T_{V_n}))\} \\ &= \mathbb{E}_x\{\mathbb{E}_{X(T_{V_n})}\{e^{-\alpha T_A} u(X_{T_A}); T_A < \infty\}\} \\ &= \mathbb{E}_x\{\mathbb{E}_x\{e^{-\alpha T_A \circ \theta_{T_{V_n}}} u(X(T_{V_n} + T_A \circ \theta_{T_{V_n}})); T_A \circ \theta_{T_{V_n}} < \infty \mid \mathcal{F}_{T_{V_n}}\}\} \\ &= \mathbb{E}_x\{\mathbb{E}_x\{e^{-\alpha T_A \circ \theta_{T_{V_n}}} u(X(T_{V_n} + T_A \circ \theta_{T_{V_n}})); T_A \circ \theta_{T_{V_n}} < \infty \mid \mathcal{F}_{T_{V_n}}\}\} \\ &= \mathbb{E}_x\{e^{-\alpha T_A \circ \theta_{T_{V_n}}} u(X(T_{V_n} + T_A \circ \theta_{T_{V_n}})); T_A \circ \theta_{T_{V_n}} < \infty\}. \end{aligned}$$

Also, note that for every $x \in E$,

$$\lim_n T_{V_n} = T_V = 0, \quad \mathbb{P}_x\text{-a.s.}, \quad (6.3)$$

and also that $T_{V_n} + T_A \circ \theta_{T_{V_n}} \searrow_n T_A$ \mathbb{P}_x -a.s., see e.g. [60, p. 65] or [33, (10.3)].

Now, if u is finely continuous, by Theorem 6.13 it follows that the trajectories $t \mapsto u(X_t)$ are right continuous \mathbb{P}_x -a.s for all $x \in E$. Thus, by dominated convergence, $\lim_n \mathbb{E}_x\{f(X(T_{V_n}))\} = f(x)$, which is a contradiction with the inequality from above, hence (ii) is proved.

Let us now prove (i). To this end, let $x \in (E \setminus A)^{\circ f}$, so that by Remark 6.20 we have $\mathbb{P}_x(T_A > 0) = 1$. Using again that T_A is a terminal time ([60, p. 65]) and (6.3), we get

$$\lim_n u(X(T_{V_n} + T_A \circ \theta_{T_{V_n}})) = u(X(T_A)), \quad \mathbb{P}_x\text{-a.s.}$$

Now, by the same argument as in the proof of (ii) we arrive immediately at a contradiction. \square

Let (E, \mathcal{B}) be endowed with a topology τ such that $\sigma(\tau) = \mathcal{B}$. Recall that a right process X on E is a *Hunt process* if for $\mu \in \mathcal{P}(E)$, X has left τ -limits in E \mathbb{P}_μ -a.s., and for every every sequence of $\mathcal{F}(t)$ -stopping times $(T_n)_{n \geq 1}$ such that $T_n \nearrow_n T$ we have

$$\tau - \lim_n X(T_n) = X(T) \text{ on } [T < \infty] \quad \mathbb{P}_\mu\text{-a.s.}$$

Proposition 6.24. *Let $A \in \mathcal{B}$, $u : E \rightarrow \mathbb{R}$ be \mathcal{B} -measurable, $\alpha \geq 0$, and $x \in E$ such that*

$$\mathbb{E}_x \{e^{-\alpha T_A} |u|(X(T_A))\} < \infty. \quad (6.4)$$

Further, consider the process

$$M = (M(t))_{t \geq 0} := \left(e^{-\alpha(t \wedge T_A)} B_\alpha^A u(X(t \wedge T_A)) \right)_{t \geq 0}.$$

The following assertions hold:

- (i) M is an $(\mathcal{F}(t \wedge T_A))$ -martingale under \mathbb{P}_x which is closed by $e^{-\alpha T_A} u(X(T_A))$.
- (ii) If $E \setminus A$ is finely open and $x \in E \setminus A$ then M is càdlàg \mathbb{P}_x -a.s.
- (iii) If u is finely continuous, then M is càdlàg \mathbb{P}_x -a.s.
- (iv) If X is a Hunt process with respect to a given topology τ on E such that $\sigma(\tau) = \mathcal{B}$, then for every increasing sequence of $(\mathcal{F}(t))$ -stopping times $(\tau_k)_{k \geq 1}$ such that $\lim_k \tau_k = T_A$ \mathbb{P}_x -a.s. we have

$$\lim_{k \rightarrow \infty} M(\tau_k) = e^{-\alpha T_A} u(X(T_A)) \quad \mathbb{P}_x\text{-a.s. and in } L^1(\mathbb{P}_x).$$

- (v) If X is a Hunt process, in addition to (6.4) we have that either u is finely continuous or $E \setminus A$ is finely open and $x \in E \setminus A$, whilst T_A is predictable under \mathbb{P}_x , then

$$\lim_{t \nearrow T_A} M(t) = e^{-\alpha T_A} B_\alpha^A u(X(T_A)) \quad \mathbb{P}_x\text{-a.s.} \quad (6.5)$$

Proof. (i). First of all, recall that by Remark 6.11, or by [33, Corollary (8.6)], for all $x \in E$, $Y \in b\mathcal{F}$, and τ and $(\mathcal{F}(t))$ -stopping time, we have

$$\mathbb{E}_x \{Y \circ \theta(\tau) | \mathcal{F}(\tau)\} = E^{X(\tau)} \{Y\} \quad \mathbb{P}_x\text{-a.s.} \quad (6.6)$$

As it easily follows by dominated convergence, (6.6) also holds for $x \in E$ and $Y \in \mathcal{F}$ for which

$$\mathbb{E}_x \{|Y| \circ \theta(\tau)\} < \infty.$$

Now, let us notice that since (6.4) is in force and using the fact that T_A is a terminal stopping time, we get

$$\begin{aligned} \mathbb{E}_x \{ [e^{-\alpha T_A} |u|(X(T_A))] \circ \theta(t \wedge T_A) \} &= \mathbb{E}_x \{ e^{-\alpha(T_A \circ \theta(t \wedge T_A))} |u|(X(t \wedge T_A + T_A \circ \theta(t \wedge T_A))) \circ \theta(t \wedge T_A) \} \\ &= \mathbb{E}_x \{ e^{-\alpha(T_A - t \wedge T_A)} |u|(X(T_A)) \} \\ &\leq e^{\alpha t} \mathbb{E}_x \{ e^{-\alpha T_A} |u|(X(T_A)) \} < \infty. \end{aligned}$$

Consequently, by choosing $Y := e^{-\alpha T_A} u(X(T_A))$ and $\tau := t \wedge T_A$ we can apply (6.6) to deduce that

$$\begin{aligned} \mathbb{E}_x \{ [e^{-\alpha T_A} u(X(T_A))] \circ \theta(t \wedge T_A) | \mathcal{F}(t \wedge T_A) \} &= E^{X(t \wedge T_A)} \{ e^{-\alpha T_A} u(X(T_A)) \} \\ &= B_\alpha^A u(X(t \wedge T_A)) \quad \mathbb{P}_x\text{-a.s.} \end{aligned}$$

But, using again that T_A is a terminal time and by similar computations as above, the left hand side of the above inequality also satisfies

$$\mathbb{E}_x \{ [e^{-\alpha T_A} u(X(T_A))] \circ \theta(t \wedge T_A) | \mathcal{F}(t \wedge T_A) \} = e^{\alpha(t \wedge T_A)} \mathbb{E}_x \{ e^{-\alpha T_A} u(X(T_A)) | \mathcal{F}(t \wedge T_A) \}.$$

Thus,

$$M(t) = e^{-\alpha(t \wedge T_A)} B_\alpha^A u(X(t \wedge T_A)) = \mathbb{E}_x \{ e^{-\alpha T_A} u(X(T_A)) | \mathcal{F}(t \wedge T_A) \}, \quad t \geq 0, \quad (6.7)$$

which proves that M is a closed (hence uniformly integrable) $(\mathcal{F}(t \wedge T_A))$ -martingale under \mathbb{P}_x .

(ii)&(iii). Clearly, it is enough to consider the case $u \in \mathcal{B}$ satisfies $u \geq 0$. Moreover, if we consider $u_n := u \wedge n, n \geq 1$, by the first part of the proof we get that the processes

$$M_n = (M_n(t))_{t \geq 0} := \left(e^{-\alpha(t \wedge T_A)} B_\alpha^A u_n(X(t \wedge T_A)) \right)_{t \geq 0}, \quad n \geq 1$$

are $(\mathcal{F}(t \wedge T_A))$ -martingales under \mathbb{P}_x . Moreover, under the hypotheses in (ii) or (iii), by Lemma 6.23 and Theorem 6.13-Remark 6.48, they are all right-continuous, hence càdlàg \mathbb{P}_x -a.s. Since $\lim_n u_n = u$ increasingly, we get that by monotone convergence that $B_\alpha^A u_n \nearrow_n B_\alpha^A u$ and thus

$$M_n(t) \nearrow_n M(t) \quad \text{for every } t \geq 0 \quad \mathbb{P}_x\text{-a.s.}$$

The fact that M is càdlàg \mathbb{P}_x -a.s. now follows from a standard result of convergence for increasing families of càdlàg supermartingales, see e.g. [40, Chapter VI, Theorem 18] or [41, Part II, Chapter IV, Section 4].

(iv). As in the proof of (i), by the strong Markov property we deduce that $(M(\tau_k))_{k \geq 1}$ is a (discrete time) $(\mathcal{F}(\tau_k))$ -martingale closed by $e^{-\alpha T_A} u(X(T_A))$ under \mathbb{P}_x . Since X is a Hunt process and $\tau_k \nearrow_k T_A$ \mathbb{P}_x -a.s., by arguing as in [33, (4.1)-(4.2)], we get that \mathbb{P}_x -a.s.

$$X(T_A) = \lim_k X(\tau_k) 1_{[T_A < \infty]} + \Delta 1_{[T_A = \infty]}, \quad [T_A < \infty] = \bigcup_{k \geq 1} \bigcap_{n \geq k} [\tau_n \leq k],$$

and thus

$$X(T_A) \in \sigma \left(\bigcup_{k \geq 1} \mathcal{F}(\tau_k) \right),$$

so now (iv) follows by Lévy's convergence theorem.

(v). Let the assumptions in (v) be satisfied. Then by (i), (ii), and (iii), we have that M is a càdlàg $(\mathcal{F}(t \wedge T_A))$ -martingale under \mathbb{P}_x which is closed by $M(\infty) := e^{-\alpha T_A} u(X(T_A))$. Let $(\tau_k)_{k \geq 1}$ be an increasing sequence of finite $(\mathcal{F}(t))$ -stopping times such that $\tau_k < T_A, k \geq 1$ on $[T_A > 0]$ and $\lim \tau_k = T_A$ \mathbb{P}_x -a.s. To prove (iv) it is clearly enough to show

$$\lim_k \sup_{t \in [\tau_k, \infty)} |M(t) - M(\infty)| = 0 \quad \mathbb{P}_x\text{-a.s.}$$

Note that the above limit is decreasing with respect to k , so the required a.s. convergence is in fact equivalent with the convergence in probability under \mathbb{P}_x . Thus, by introducing the processes

$$M_k(t) := M(\tau_k + t) - M(\tau_k), \quad t \geq 0, \quad k \geq 1,$$

the previous convergence boils down to showing

$$\limsup_k \sup_{t \geq 0} |M_k(t) - M_k(\infty)| = 0 \quad \text{in probability with respect to } \mathbb{P}_x.$$

By (iv) we have that $\lim_k M_k(\infty) = 0$ \mathbb{P}_x -a.s. and in $L^1(\mathbb{P}_x)$, hence the above convergence, thus it remains to show that

$$\limsup_k \sup_{t \geq 0} |M_k(t)| = 0 \quad \text{in probability with respect to } \mathbb{P}_x.$$

Now, let us note that by Doob's stopping theorem we have that for every $k \geq 1$ the process $(M_k(t))_{t \geq 0}$ is a càdlàg $(\mathcal{F}((\tau_k + t) \wedge T_A))$ -martingale under \mathbb{P}_x which is closed by $M_k(\infty)$. Consequently, by Doob's maximal inequality for (continuous-time) right-continuous martingales, and using again (iv), we get

$$\mathbb{P}_x \left(\sup_{t \geq 0} |M_k(t)| \geq \varepsilon \right) \leq \frac{1}{\varepsilon} \mathbb{E}_x (|M_k(\infty)|) \xrightarrow[k]{} 0,$$

which finishes the proof. \square

Existence of a right process on a larger space. Condition (\mathbf{H}_C) , although necessary, is not sufficient to guarantee the existence of a right process with the given resolvent \mathcal{U} , even if the resolvent is strong Feller; see [23]. However, under (\mathbf{H}_C) , there is always a larger space on which an associated right process exists.

We denote by $Exc(\mathcal{U}_\beta)$ the set of all \mathcal{U}_β -excessive measures: $\xi \in Exc(\mathcal{U}_\beta)$ if and only if ξ is a σ -finite measure on E and $\xi \circ \alpha U_{\beta+\alpha} \leq \xi$ for all $\alpha > 0$.

Definition 6.25 (The saturation of (E, \mathcal{U}) w.r.t. \mathcal{U}_β). *Let $\beta > 0$.*

(i) *The energy functional associated with \mathcal{U}_β is $L^\beta : Exc(\mathcal{U}_\beta) \times \mathcal{E}(\mathcal{U}_\beta) \rightarrow \overline{\mathbb{R}}_+$ given by*

$$\begin{aligned} L^\beta(\xi, v) &:= \sup\{\mu(v) : \mu \text{ is a } \sigma\text{-finite measure, } \mu \circ U_\beta \leq \xi\} \\ &= \sup\{\xi(f) : f \geq 0 \text{ is a } \mathcal{B}\text{-measurable function such that } U_\beta f \leq v\}. \end{aligned}$$

(ii) *The saturation of E (with respect to \mathcal{U}_β) is the set E_1 of all extreme points of the set $\{\xi \in Exc(\mathcal{U}_\beta) : L^\beta(\xi, 1) = 1\}$.*

(iii) *The map $E \ni x \mapsto \delta_x \circ U_\beta \in Exc(\mathcal{U}_\beta)$ is a measurable embedding of E into E_1 and every \mathcal{U}_β -excessive function v has an extension v_1 to E_1 , defined as $v_1(\xi) := L^\beta(\xi, v)$. The set E_1 is endowed with the σ -algebra \mathcal{B}_1 generated by the family $\{v_1 : v \in \mathcal{E}(\mathcal{U}_\beta)\}$. In addition, as in [17, Sections 1.1 and 1.2], there exists a unique resolvent of kernels $\mathcal{U}^1 = (U_\alpha^1)_{\alpha>0}$ on (E_1, \mathcal{B}_1) which is an extension of \mathcal{U} in the sense that*

(iii.1) $U^1 1_{E_1 \setminus E} = 0$ on E_1 ,

(iii.2) $(U^1 f)|_E = U(f|_E)$ for all $f \in b\mathcal{B}_1$.

More precisely, it is given by

$$\begin{aligned} U_\alpha^1 f(\xi) &= L^\beta(\xi, U_\alpha(f|_E)) \\ &= L^\beta(\xi, U_\beta(f|_E + (\beta - \alpha)U_\alpha)) \\ &= \xi(f + (\beta - \alpha)U_\alpha f) \text{ for all } f \in b\mathcal{B}_1, \xi \in E_1, \alpha > 0. \end{aligned} \tag{6.8}$$

Remark 6.26. *Note that (E_1, \mathcal{B}_1) is a Lusin measurable space, the map $x \mapsto \delta_x \circ U_\beta$ identifies E with a subset of E_1 , $E \in \mathcal{B}_1$, and $\mathcal{B} = \mathcal{B}_1|_E$. Furthermore, E is dense in E_1 with respect to the fine topology on E_1 associated with \mathcal{U}^1 . Hence, if $\alpha > 0$ and v is a \mathcal{U}_α -excessive function then its \mathcal{U}_α^1 -excessive extension v_1 given by Definition 6.25 is the unique extension of v from E to E_1 by fine continuity.*

The existence of a right process on a larger space, more precisely on E_1 , is given by the following result, for which we refer to [29, (2.3)], in [14, Sections 1.7 and 1.8], [17, Theorem 1.3], and [18, Section 3]; see also [63].

Theorem 6.27. *There is always a right process X^1 on the saturation (E_1, \mathcal{B}_1) , associated with \mathcal{U}^1 . Moreover, the following assertions are equivalent:*

(i) *There exists a right process on E associated with \mathcal{U} .*

(ii) *The set $E_1 \setminus E$ is polar (w.r.t. \mathcal{U}^1).*

6.2 Choquet capacity

Let $X = (\Omega, \mathcal{F}, \mathcal{F}(t), X(t), \theta(t), \mathbb{P}_x), x \in E$ be a right process on a Lusin measurable space (E, \mathcal{B}) , and \mathcal{T} be a natural topology on E . It is well known that for every $\nu \in \mathcal{P}(E)$, $\alpha > 0$ and $f \in b\mathcal{B}, f > 0$, the mapping

$$\mathbb{R}^d \supset A \mapsto \text{Cap}_\nu^\alpha(A) := \inf \{\nu(R_\alpha^G(U_\alpha f)) : A \subset G, G \in \mathcal{T}\}, \tag{6.10}$$

is a Choquet capacity on (E, \mathcal{T}) . Moreover, if X has left limits in (E, \mathcal{T}) \mathbb{P}_ν -a.s., then the capacity Cap_ν^α is tight; cf. [52], [16], and [30].

For convenience, let us briefly recall the notion of capacity (for details see e.g. [40, Ch. III] and [14, Appendix A.1]):

Definition 6.28 (Choquet capacity). Let \mathcal{K} be the family of all \mathcal{T} -compact subsets of E . A map

$$\text{Cap} : \mathcal{P}(E) \rightarrow \overline{\mathbb{R}}_+, \quad \mathcal{P}(E) := \{A : A \subset E\},$$

is called Choquet capacity on E provided that:

- (i) $\text{Cap}(K) < \infty$ for every $K \in \mathcal{K}$;
- (ii) $\text{Cap}(A_1) \leq \text{Cap}(A_2)$ if $A_1 \subset A_2 \in \mathcal{P}(E)$;
- (iii) If $(A_n)_{n \geq 1} \subset \mathcal{P}(E)$ is such that $A_n \nearrow_n A \in \mathcal{P}(E)$, then $\text{Cap}(A_n) \nearrow \text{Cap}(A)$;
- (iv) If $(K_n)_{n \geq 1} \subset \mathcal{K}$ is such that $K_n \searrow_n A$, then $\text{Cap}(K_n) \searrow \text{Cap}(A)$.

Also, recall that a Choquet capacity Cap is tight provided that there exists an increasing sequence $(K_n)_n \subset \mathcal{K}$ such that

$$\inf_n \text{Cap}(E \setminus K_n) = 0.$$

6.3 Basic operations on right processes

6.3.1 Trivial modification and restriction of the resolvent

Let $M \in \mathcal{B}$ be such that $U_\alpha(1_M) = 0$ on $E \setminus M$ for one (and therefore for all) $\alpha > 0$. Following [18] and [32] we can define the trivial modification of \mathcal{U} on M and the restriction of \mathcal{U} to $E \setminus M$ as follows.

Definition 6.29 (Trivial modification). For all $\alpha > 0$ we define the kernel

$$U'_\alpha f = 1_{E \setminus M} U_\alpha f + \frac{1}{\alpha} 1_M f, \quad f \in p\mathcal{B}. \quad (6.11)$$

The family $\mathcal{U}' = (U'_\alpha)_{\alpha > 0}$ is also a sub-Markovian resolvent of kernels on (E, \mathcal{B}) and there exists a cone \mathcal{C} such that $(\mathbf{H}_\mathcal{C})$ is fulfilled. \mathcal{U}' is called the trivial modification of the resolvent \mathcal{U} on M .

Definition 6.30 (Restriction). The family $\mathcal{U}|_{E \setminus M} = (U_\alpha|_{E \setminus M})_{\alpha > 0}$ is a sub-Markovian resolvent of kernels on $(E \setminus M, \mathcal{B}|_{E \setminus M})$, called the restriction of \mathcal{U} to $E \setminus M$. Moreover, $\mathcal{U}|_{E \setminus M}$ satisfies condition $(\mathbf{H}_\mathcal{C})$ on the measurable space $(E \setminus M, \mathcal{B}|_{E \setminus M})$.

A function $u \in p\mathcal{B}|_{E \setminus M}$ is $\mathcal{U}_\beta|_{E \setminus M}$ -excessive if and only if there exists a function \bar{u} which is \mathcal{U}_β -excessive, such that $u = \bar{u}|_{E \setminus M}$. Consequently, the fine topology on $E \setminus M$ with respect to the restriction of \mathcal{U} to $E \setminus M$ is the trace on $E \setminus M$ of the fine topology on E with respect to \mathcal{U} . A function $v \in p\mathcal{B}$ will be \mathcal{U}'_β -excessive if and only if $v|_{E \setminus M}$ is $\mathcal{U}_\beta|_{E \setminus M}$ -excessive.

Remark 6.31. If $M \in \mathcal{B}$ is a set such that $U_\alpha(1_M) = 0$ on $E \setminus M$ then the following assertions hold.

- (i) Let \mathcal{R} be a Ray cone associated with \mathcal{U}_β and \mathcal{U}' be the trivial modification of \mathcal{U} on M . Then there exists a Ray cone \mathcal{R}' with respect to \mathcal{U}'_β such that $\mathcal{R} \subset \mathcal{R}'$.
- (ii) If \mathcal{R}° is a Ray cone associated with $\mathcal{U}_\beta|_{E \setminus M}$ then there exists a Ray cone \mathcal{R} with respect to \mathcal{U}_β such that $\mathcal{R}|_{E \setminus M} = \mathcal{R}^\circ$.

6.3.2 Restriction, trivial extension, and trivial modification of the process

Strongly supermedian functions. Recall first that a positive numerical function f on E is called *nearly Borel* provided that for any finite measure λ on E there exist two positive numerical Borelian functions g and h on E such that $g \leq f \leq h$ and the set $[g < h]$ is λ -polar and λ -negligible. It is known that any U_α -excessive function on E , $\alpha > 0$, is nearly Borel. Let \mathcal{B}^n denote the σ -algebra of all nearly Borel sets.

A positive numerical function f on E is called *strongly supermedian* (with respect to \mathcal{U}_α) if it is nearly Borel and for any two finite measures μ and ν on E we have

$$\mu \circ U_\alpha \leq \nu \circ U_\alpha \implies \mu(f) \leq \nu(f).$$

We denote by \mathcal{S}_α the set of all strongly supermedian functions on E (with respect to \mathcal{U}_α). Further, we present several basic properties of the strongly supermedian functions, cf. e.g. [13], [14], [15], and [26].

\mathcal{S}_α is a convex cone. By [13], assertion (1) of the Remark at page 369, it follows that a function $f \in p\mathcal{B}^n$ is strongly supermedian (with respect to \mathcal{U}_α) if and only if $R_\alpha^M f \leq f$ for any $M \in \mathcal{B}$. Let $(v_n)_n$ be a sequence in \mathcal{S}_α , then $\inf_n v_n \in \mathcal{S}_\alpha$. If in addition $(v_n)_n$ is increasing then $\sup_n v_n \in \mathcal{S}_\alpha$. For all $v \in \mathcal{S}_\alpha$ we have $v = \inf\{u : u \in \mathcal{E}(\mathcal{U}_\alpha), u \geq v\}$. In particular, every strongly supermedian function is finely upper semicontinuous.

If $u \in \mathcal{E}(\mathcal{U}_\alpha)$ and $A \in \mathcal{B}$, then $R_\alpha^A u$, the α -order reduced function of u on A , is a strongly supermedian function, particularly, it is nearly Borel measurable. More generally, the following assertions are equivalent for a positive nearly Borel measurable function v on E :

- (i) The function v is strongly supermedian with respect to \mathcal{U}_α .
- (ii) There exists a family \mathcal{G} of \mathcal{U}_α -excessive functions such that $v = \inf \mathcal{G}$.
- (iii) We have $v = \inf\{u \in \mathcal{E}(\mathcal{U}_\alpha) : u \geq v\}$.

Following [29], we set

$$\mathcal{A}(\mathcal{U}) := \left\{ A \in \mathcal{B} : R_\alpha^{E \setminus A} 1 = 0 \text{ on } A \text{ for some } \alpha > 0 \right\} \quad (6.12)$$

An element $A \in \mathcal{A}(\mathcal{U})$ is called an *absorbing set*. An absorbing set is finely open and in what follows we shall recall several key facts concerning the restriction of a right process to an absorbing set.

Remark 6.32. By [26, A.1.2] the following assertions are equivalent for a set $A \in \mathcal{B}$.

- (i) The set A belongs to $\mathcal{A}(\mathcal{U})$.
- (ii) We have \mathbb{P}_x -a.s. $D_{E \setminus A} = \infty$ for every $x \in A$.
- (iii) There exists a strongly supermedian function v with respect to \mathcal{U}_α such that $A = [v = 0]$.

Remark 6.33. If $v \in \mathcal{S}_\alpha \cap p\mathcal{B}$, then

$$[v < \infty], [v = 0] \in \mathcal{A}(\mathcal{U}).$$

The fact that the set $[v = 0]$ belongs to $\mathcal{A}(\mathcal{U})$ follows from the above considerations. To show that $[v < \infty] \in \mathcal{A}(\mathcal{U})$, we complete the arguments from [29] and A.1.4 from [26], where the function v was supposed to be \mathcal{U}_α -excessive. If $n \geq 1$ then we have $1 \leq \frac{1}{n} v$ on $[v = \infty]$. Let $u \in \mathcal{E}(\mathcal{U}_\alpha)$, $u \geq v$. Then $R_\alpha^{[v = \infty]} 1 \leq \frac{1}{n} u$ and consequently

$$R_\alpha^{[v = \infty]} 1 \leq v/n = \inf\{u \in \mathcal{E}(\mathcal{U}_\alpha) : u \geq v\}/n.$$

We conclude that $R_\alpha^{[v = \infty]} 1 = 0$ on $[v < \infty]$, hence $[v < \infty] \in \mathcal{A}(\mathcal{U})$.

Remark 6.34. (i) If $(A_n)_{n \geq 1} \subset \mathcal{A}(\mathcal{U})$, then $\bigcap_n A_n \in \mathcal{A}(\mathcal{U})$.

- (ii) A set $A \in \mathcal{B}$ is μ -inessential (see Definition 6.18), if and only if $A \in \mathcal{A}(\mathcal{U})$ and $\mu(A) = 0$. In particular, a countable intersection of μ -inessential sets is again μ -inessential.

If $A \in \mathcal{A}(\mathcal{U})$ then $U_\beta(1_{E \setminus A}) = 0$ on A and therefore we may consider the restriction $\mathcal{U}|_A$ of \mathcal{U} to A . In particular, the resolvent $\mathcal{U}|_A = (U_\alpha|_A)_{\alpha > 0}$ satisfies condition (\mathbf{H}_C) on the measurable space $(A, \mathcal{B}|_A)$ and we already observed that the fine topology on A with respect to the restriction of \mathcal{U} to A is the trace on A of the fine topology on E with respect to \mathcal{U} . If \mathcal{R} is a Ray cone associated with \mathcal{U}_β then $\mathcal{R}|_A$ is a Ray cone associated with $\mathcal{U}_\beta|_A$.

Proposition 6.35 (cf. [29, Corollary 3.2]). *If \mathcal{U} is the resolvent of a right process with state space E and $A \in \mathcal{A}(\mathcal{U})$, then the restriction of \mathcal{U} to A is the resolvent of a right process with state space A .*

Definition 6.36 (Restriction of the process). *If X is a right process with resolvent \mathcal{U} , and $A \in \mathcal{A}(\mathcal{U})$, then one can construct a right process \widetilde{X} associated to the restriction of \mathcal{U} to A as in Proposition 6.35, by setting*

$$\begin{aligned}\widetilde{\Omega} &:= \{\omega \in \Omega : X_t(\omega) \in A, \text{ for all } 0 \leq t < \zeta(\omega)\} \\ \widetilde{\mathbb{P}}^x &:= \mathbb{P}_x|_{\widetilde{\Omega}}, x \in A, \quad \widetilde{X}_t(\omega) := X_t(\omega), \omega \in \widetilde{\Omega}.\end{aligned}$$

The right process \widetilde{X} defined above is called the restriction of X to A ; for more details, see [60].

Proposition 6.37. *Let $(X, \mathbb{P}_x, x \in E)$ be a right process associated to the resolvent \mathcal{U} on the Lusin measurable space (E, \mathcal{B}) . Further, let m be a σ -finite measure on E , and τ be a natural topology on E . Furthermore, suppose that*

$$\mathbb{P}_m([0, \infty) \ni t \mapsto X(t) \in E \text{ is } \tau\text{-continuous}) = 1.$$

Then there exists a set $E_0 \subset E$ such that $E \setminus E_0$ is m -inessential and

$$\mathbb{P}_x([0, \infty) \ni t \mapsto X(t) \in E_0 \text{ is } \tau\text{-continuous}) = 1 \quad \text{for all } x \in E_0. \quad (6.13)$$

Consequently, one can consider the restriction \widetilde{X} of X from E to E_0 , so that \widetilde{X} is a right process with a.s. τ -continuous paths in E_0 ; in particular, τ is a natural topology on E_0 . Furthermore, the above assertion remains true if instead of " τ -continuous" we use " τ -right-continuous" or " τ -càdlàg".

Proof. The proof of (6.13) follows by the same arguments as in [29, Proposition 5.1], so we skip it. The fact that τ is a natural topology on E_0 follows immediately from Theorem 6.13, (i). \square

Remark 6.38. *We would like to stress that in Proposition 6.37, the assumption that τ is a natural topology is crucial. Indeed, the proof of [29, Proposition 5.1], relies heavily on the fact that X is a priori a.s. right-continuous at $t = 0$ with respect to τ , and this is automatically satisfied if τ is a natural topology, by Theorem 6.13. In the next result, which is needed in the main body of the paper, we drop the assumption that τ is a natural topology; however, we need to assume (i) below, instead.*

Proposition 6.39. *Let $(X, \mathbb{P}_x, x \in E)$ be a right process associated to the resolvent \mathcal{U} on the Lusin measurable space (E, \mathcal{B}) . Further, let τ be a Lusin topology on E whose Borel σ -algebra is precisely \mathcal{B} . Furthermore, suppose that*

(i) *For every $x \in E$ we have*

$$\mathbb{P}_x((0, \infty) \ni t \mapsto X(t) \in E \text{ is } \tau\text{-continuous}) = 1.$$

(ii) *m is a σ -finite measure on E such that*

$$\mathbb{P}_m([0, \infty) \ni t \mapsto X(t) \in E \text{ is } \tau\text{-continuous}) = 1.$$

Further, set

$$E_0 := \{x \in E : \mathbb{P}_x([0, \infty) \ni t \mapsto X(t) \in E \text{ is } \tau\text{-continuous}) = 1\}.$$

Then $E \setminus E_0$ is m -inessential.

Consequently, one can consider the restriction \widetilde{X} of X from E to E_0 , so that \widetilde{X} is a right process with a.s. τ -continuous paths in E_0 ; in particular, τ is a natural topology on E_0 . Furthermore, the above assertion remains true if instead of " τ -continuous" we use " τ -right-continuous" or " τ -càdlàg".

Proof. Let

$$\Omega_o := \{\omega \in \Omega : (0, \infty) \ni t \mapsto X(t) \in E \text{ is } \tau\text{-continuous}\}.$$

By assumption (i) we have $\mathbb{P}(\Omega_o) = 1$. Let

$$\Lambda := \{\omega \in \Omega : [0, \infty) \ni t \mapsto X(t, \omega) \in E \text{ is not } \tau\text{-continuous}\}$$

and

$$\Lambda_o := \{\omega \in \Omega : [0, \infty) \ni t \mapsto X(t, \omega) \in E \text{ is not } \tau\text{-continuous in zero}\}$$

Clearly we have $\Omega_o \cap \Lambda = \Omega_o \cap \Lambda_o$ and

$$\begin{aligned}\Lambda_o^c &= \{\omega \in \Omega : X(t, \omega) \rightarrow X(0, \omega) \text{ when } t \rightarrow 0\} \\ &= \bigcap_{\varepsilon \in \mathbb{Q}_+} \bigcup_{t \in \mathbb{Q}_+} \{w \in \Omega : \sup_{0 < s \leq t} |X(s, \omega) - X(0, \omega)| \leq \varepsilon\} \\ &= \bigcap_{\varepsilon \in \mathbb{Q}_+} \bigcup_{t \in \mathbb{Q}_+} \bigcap_n \{w \in \Omega : \sup_{\frac{1}{n} < s \leq t} |X(s, \omega) - X(0, \omega)| \leq \varepsilon\}.\end{aligned}$$

Further, let

$$\Lambda_1 := \bigcap_{\varepsilon \in \mathbb{Q}_+} \bigcup_{t \in \mathbb{Q}_+} \bigcap_n \{w \in \Omega : \sup_{\frac{1}{n} < s \leq t, s \in \mathbb{Q}} |X(s, \omega) - X(0, \omega)| \leq \varepsilon\},$$

and $\mathcal{F}^0 := \sigma(X(t) : t \geq 0)$. Then $\Lambda_1 \in \mathcal{F}^0$ and $\Omega_o \cap \Lambda_o^c = \Omega_o \cap \Lambda_1$ since the trajectories from Ω_o are τ -continuous on $(0, \infty)$.

Define the function V on E as

$$V(x) := \mathbb{P}_x(\Lambda), \quad x \in E.$$

Let us show that V is a Borel function on E . Indeed, recall that by [33, Theorem (3.6)] we have that $E \ni x \mapsto \mathbb{P}_x(\Gamma)$ is a Borel function on E for every Γ in \mathcal{F}^0 . Because $\mathbb{P}_x(\Omega_o) = 1$ we get that $V(x) = \mathbb{P}_x(\Lambda_o) = 1 - \mathbb{P}_x(\Omega_o \cap \Lambda_o^c) = 1 - \mathbb{P}_x(\Lambda_1)$.

We show that V is a strongly supermedian function w.r.t. \mathcal{U}_α for some $\alpha > 0$. It is sufficient to prove that $R_\alpha^K V \leq V$ for any Ray compact subset K of E . Using the probabilistic characterization of the reduced function due to G.A. Hunt, see (6.2), we have for every $x \in E$:

$$\begin{aligned}R_\alpha^K V(x) &= \mathbb{E}_x(e^{-\alpha D_K} V(X(D_K))) = \mathbb{E}_x(e^{-\alpha D_K} \mathbb{P}_{X(D_K)}(\Lambda)) \\ &= \mathbb{E}_x(e^{-\alpha D_K} 1_\Lambda \circ \theta(D_K)) = \mathbb{E}_x \{e^{-\alpha D_K} 1_{[\theta(D_K)]^{-1}(\Lambda)}\} \leq \mathbb{P}_x(\Lambda) \\ &= V(x),\end{aligned}$$

where the last inequality holds because $\theta_T^{-1}(\Lambda) \subset \Lambda$. By Remark 6.32 it follows that the set $E_0 = \{x \in E : \mathbb{P}_x(\Lambda) = 0\}$ belongs to $\mathcal{A}(\mathcal{U})$. By assumption (ii) we have $m(V) = 0$ and therefore $m(E \setminus E_0) = m([V > 0]) = 0$. We conclude that $E \setminus E_0$ is m -inessential, completing the proof. \square

Let $M \in \mathcal{B}$, $F := E \setminus M$ and assume that $\mathcal{U} = (U_\alpha)_{\alpha > 0}$ is a sub-Markovian resolvent of kernels on $(F, \mathcal{B}|_F)$, the resolvent of a right process with state space F . For all $\alpha > 0$ define the kernel

$$U'_\alpha f = 1_F U_\alpha f + \frac{1}{\alpha} 1_M f, \quad f \in p\mathcal{B}.$$

Then, as in Definition 6.29, the family $\mathcal{U}' = (U'_\alpha)_{\alpha > 0}$ is a sub-Markovian resolvent of kernels on (E, \mathcal{B}) . We can expose now the trivial extension of a right process to a larger set.

Proposition 6.40. *Assume that $\mathcal{U} = (U_\alpha)_{\alpha > 0}$ is the resolvent on F of a right process X with state space F . Then $\mathcal{U}' = (U'_\alpha)_{\alpha > 0}$ is the resolvent of a right process $X' = (\Omega', \mathcal{F}', X'_t, \mathbb{P}^{t,x}, x \in E)$ with state space E and each point of M is a trap for X' , that is, $\mathbb{P}^{t,x}[X'_t = x \text{ for all } t \geq 0] = 1$ for all $x \in M$. Moreover, X is the restriction of X' from E to F in the sense of Definition 6.36.*

For the construction of the process X' , the extension of X from F to E , we refer to [53, p. 118].

Definition 6.41 (Trivial extension of the process). *The process X' from Proposition 6.40 is called the trivial extension of the process X from F to E .*

Proposition 6.42. *Let X be a right process with state space E and resolvent family $\mathcal{U} = (U_\alpha)_{\alpha > 0}$ and let $M \in \mathcal{B}$ be such that $R_\beta^M 1 = 0$ on $E \setminus M$. Consider the resolvent $\mathcal{U}' = (U'_\alpha)_{\alpha > 0}$ on E defined by (6.11). Then $\mathcal{U}' = (U'_\alpha)_{\alpha > 0}$ is the resolvent of a right process X' with state space E and each point of M is a trap for X' .*

Proof. Because $E \setminus M \in \mathcal{A}(\mathcal{U})$, we can apply Proposition 6.35, to the restrict of X to $E \setminus M$. The requested process X' on E is obtained now applying Proposition 6.40 to this restriction. \square

Definition 6.43 (Trivial modification of the process). *The process X' from Proposition 6.42 is called the trivial modification of the process X on M .*

6.3.3 The killed process

In this part we assume that X is a right process on (E, \mathcal{B}) with transition semigroup $(P_t)_{t \geq 0}$ and resolvent $\mathcal{U} = (U_\alpha)_{\alpha > 0}$.

Proposition 6.44. *Let F be a closed set in a natural topology. Then R_α^F and B_α^F are Borel kernels on E , that is, if $g \in p\mathcal{B}$ then $R_\alpha^F g$ and $B_\alpha^F g$ are Borel functions on E .*

Proof. Let $\mathcal{F}^0 := \sigma(X(t), t \geq 0)$. The main observation is that $\Omega \ni \omega \mapsto D_F(\omega)$ is \mathcal{F}^0 -measurable, which follows arguing as in (6.16) from [33], page 35. Indeed, if $s \geq 0$ then we have

$$\begin{aligned} [D_F > s] &= \{\omega \in \Omega : \text{there exists } \varepsilon \in \mathbb{Q}_+^* \text{ such that } X(t)(\omega) \in E \setminus F \text{ for all } t \in [s, s + \varepsilon]\} \\ &= \bigcup_{\varepsilon \in \mathbb{Q}_+^*} \bigcap_{t \in [s, s + \varepsilon] \cap \mathbb{Q}} X(t)^{-1}(E \setminus F), \end{aligned}$$

where for the second equality we used the right continuity of the paths of X (in the natural topology) and the fact that $E \setminus F$ is an open set. Since clearly $X(t)^{-1}(E \setminus F) \in \mathcal{F}^0$, we conclude that also $[D_F > s] \in \mathcal{F}^0$. The process X is progressively measurable and consequently $(t, \omega) \mapsto g(X(t)(\omega))$ is $\mathcal{B}(\mathbb{R}_+) \otimes \mathcal{F}^0$ -measurable. It follows that

$$\Omega \ni \omega \mapsto Y(\omega) := e^{-\alpha D_F(\omega)} g(X(D_F(\omega))(\omega)) \text{ is } \mathcal{F}^0\text{-measurable.}$$

Hence $R_\alpha^F g(x) = \mathbb{E}_x(Y)$, with $Y \in b\mathcal{F}^0$. On the other hand, by [33, Theorem (3.6)], the function $E \ni x \mapsto \mathbb{E}_x(Y)$ is \mathcal{B} -measurable for all $Y \in b\mathcal{F}^0$, since the transition semigroup consists of Borel kernels on E . Consequently, $R_\alpha^F g \in p\mathcal{B}$. If $g \in b\mathcal{E}(\mathcal{U}_\alpha)$ then $B_\alpha^F g = \widehat{R_\alpha^F g} \in \mathcal{E}(\mathcal{U}_\alpha)$, so, it is also a Borel function on E . By a monotone class argument it results that $B_\alpha^F g \in p\mathcal{B}$ for all $g \in p\mathcal{B}$. \square

Let $D \in \mathcal{B}$ be a finely open subset of E . For every $\alpha > 0$ we denote by U_α^D the kernel on (E, \mathcal{B}^n) given by

$$U_\alpha^D f := U_\alpha f - B_\alpha^{E \setminus D} U_\alpha f, \quad f \in bp\mathcal{B}^n.$$

Using (6.2) one can see that

$$U_\alpha^D f(x) = \mathbb{E}_x \left\{ \int_0^{T_D} e^{-\alpha t} f(X(t)) dt \right\} \text{ for } f \in bp\mathcal{B}^n, x \in E, \text{ and } \alpha > 0.$$

Then the family $\mathcal{U}^D = (U_\alpha^D)_{\alpha > 0}$ is a sub-Markovian resolvent of kernels on (E, \mathcal{B}^n) , $U_\alpha^D \leq U_\alpha$, and $U_\alpha^D(1_{E \setminus D}) = 0$ for all $\alpha > 0$; cf. e.g. [12] and [14]. In particular, we may consider the restriction of \mathcal{U}^D to D , it is a sub-Markovian resolvent of kernels on (D, \mathcal{B}^n) , also denoted by \mathcal{U}^D .

Let

$$\widetilde{D} := \{x \in E : \exists u \in \mathcal{E}(\mathcal{U}_\alpha) \text{ such that } B_\alpha^{E \setminus D} u(x) < u(x)\}.$$

Then by [12] and Section 3.6 in [14] the set \widetilde{D} is finely open, $\widetilde{D} \in \mathcal{B}^n$, and $D \subset \widetilde{D}$.

Proposition 6.45. *Let $D \in \mathcal{B}$ be a finely open subset of E . Then the following assertions hold.*

(i) *There exists σ -algebra $\widetilde{\mathcal{B}}$ on \widetilde{D} such that $(\widetilde{D}, \widetilde{\mathcal{B}})$ is a Radon measurable space, $\mathcal{B}|_{\widetilde{D}} \subset \widetilde{\mathcal{B}} \subset \mathcal{B}^n|_{\widetilde{D}}$, and the family $\mathcal{U}^D = (\mathcal{U}_\alpha^D)_{\alpha > 0}$ is a sub-Markovian resolvent of kernels on $(\widetilde{D}, \widetilde{\mathcal{B}})$ such that $\mathcal{E}(\mathcal{U}_\alpha^D)$ generates $\widetilde{\mathcal{B}}$.*

(ii) *Let X^D denote the process X killed at time $T_{E \setminus D}$, that is,*

$$X^D(t)(\omega) := \begin{cases} X(t)(\omega), & t < T_{E \setminus D}(\omega) \\ \Delta, & t \geq T_{E \setminus D}(\omega) \end{cases}, \quad \omega \in \Omega.$$

Then X^D is a right process with state space \widetilde{D} and life time $T_{E \setminus D}$, and \mathcal{U}^D is its associated resolvent. In addition, the set $\widetilde{D} \setminus D$ is polar with respect to X^D .

(iii) If D is an open set in a natural topology then \mathcal{U}^D is a resolvent of kernels on $(D, \mathcal{B}|_D)$ and X^D is a right process with state space D and resolvent \mathcal{U}^D .

Proof. Assertion (i) and the fact that the set $\widetilde{D} \setminus D$ is polar follow from Proposition 3.6.7 from [14].

Assertion (ii) is a consequence of Corollary 12.24 from [60], since $T_{E \setminus D}$ is an exact terminal time.

To prove assertion (iii), observe first that using Proposition 6.44, U_α^D is a kernel on $(D, \mathcal{B}|_D)$ for each $\alpha > 0$. It follows that X^D is a right process with state space D . \square

Remark 6.46. (i) Because the set $\widetilde{D} \setminus D$ is polar (according to assertion (ii) of Proposition 6.45), we can always restrict the process \overline{X} from \widetilde{D} to D , not only in the case when D is open in a natural topology, as in assertion (iii) of the above proposition. So, we can regard X^D as a right process on D , but in a wider sense, having the transition semigroup only universally $\mathcal{B}|_D$ -measurable.

(ii) The following properties are equivalent for a subset D of E .

(ii.1) D is open in a natural topology.

(ii.2) D is open in a Ray topology.

(ii.3) D is open in the topology generated by a countable family of α -excessive functions for some $\alpha > 0$.

Definition 6.47 (Killing). We say that the process X^D from Proposition 6.45 is obtained by killing X at the first entry time in $E \setminus D$.

Remark 6.48. Let $f : E \rightarrow \mathbb{R}$ be universally \mathcal{B} -measurable and $D \in \mathcal{B}$ be finely open. Then, it follows from Proposition 6.45 and Theorem 6.13 that f is finely continuous on D if and only if $(f(X(t)))_{0 \leq t < T_{E \setminus D}}$ has \mathbb{P}^x -a.s. right continuous paths for all $x \in D$.

6.3.4 The natural extension

In this paragraph we recall one of the key tools (see Theorem 6.51 below) that we employ in the proof of Theorem 2.15 for solving the open problem OP2.

Definition 6.49 (Natural extension). Let (E, \mathcal{B}) and $(\overline{E}, \overline{\mathcal{B}})$ be two Lusin measurable spaces such that $E \in \overline{\mathcal{B}}$ and $\mathcal{B} = \overline{\mathcal{B}}|_{\mathcal{B}}$.

(i) Let \mathcal{U} be a sub-Markovian resolvent of kernels on E . A second sub-Markovian resolvent of kernels $\overline{\mathcal{U}} := (\overline{U}_\alpha)_{\alpha > 0}$ on $(\overline{E}, \overline{\mathcal{B}})$ is called an extension of \mathcal{U} if:

(i.1) $\overline{U}_\alpha(1_{\overline{E} \setminus E}) = 0$.

(i.2) $(\overline{U}_\alpha f)|_E = U_\alpha(f|_E)$ (on E) for all $\alpha > 0$ and $f \in b\overline{\mathcal{B}}$.

In this case, for shortness, we write

$$(E, \mathcal{U}) \subset (\overline{E}, \overline{\mathcal{U}}) \tag{6.14}$$

(ii) Given a right process X on E , we say that a second Markov process $\overline{X} = (\overline{\Omega}, \overline{\mathcal{F}}, \overline{\mathcal{F}}_t, \overline{X}(t), \overline{\theta}(t), \overline{\mathbb{P}}^x)$, with state space \overline{E} , is a natural extension of X if the following conditions are fulfilled:

(ii.1) \overline{X} is a right process.

(ii.2) The processes $((X(t))_{t \geq 0}, \mathbb{P}_x)$ and $((\overline{X}(t))_{t \geq 0}, \overline{\mathbb{P}}^x)$ are equal in distribution for all $x \in E$;

(ii.3) The set $\overline{E} \setminus E$ is polar with respect to (the resolvent of) \overline{X} .

In this case, again for shortness, we write

$$(E, X) \subset (\overline{E}, \overline{X}). \tag{6.15}$$

Proposition 6.50. If \overline{X} is a natural extension of X , then its resolvent denoted by $\overline{\mathcal{U}}$ is an extension of \mathcal{U} .

Theorem 6.51 (cf [22, Theorem 1.10]). Let $\overline{\mathcal{U}}$ be an extension of \mathcal{U} . Then there exists a natural extension \overline{X} of X , with resolvent $\overline{\mathcal{U}}$, if and only if (\mathbf{H}_C) is satisfied.

Remark 6.52. As already stated in [24], the following assertion can be easily deduced from Theorem 6.51: Let X be a right process on E . The right process X^1 on the saturation $\widetilde{E^1}$ given by Theorem 6.27 is maximal, in the sense that $\widetilde{X^1}$ is a natural extension of any natural extension \widetilde{X} of X . In particular, for any such natural extension \widetilde{X} on \widetilde{E} , we have that $\widetilde{E} \setminus E$ is polar.

Lemma 6.53 (cf [17, Lemma 2.1]). Let K be a kernel and ν be a σ -finite measure on a measurable space (E, \mathcal{B}) , such that if $B \in \mathcal{B}$ and $\nu(B) = 0$, then $K1_B = 0$ ν -a.e. If $E_0 \in \mathcal{B}$ satisfies $\nu(E \setminus E_0) = 0$, then there exists $F \in \mathcal{B}$, $F \subset E_0$, such that $\nu(E \setminus F) = 0$ and $K1_{E \setminus F} = 0$ on F .

6.4 On martingale additive functionals

Let $X = (\Omega, \mathcal{F}, \mathcal{F}(t), X(t), \theta(t), \mathbb{P}_x, t \geq 0, x \in E)$ be a right process on a Lusin measurable space (E, \mathcal{B}) , with transition $(P_t)_{t \geq 0}$ and resolvent $\mathcal{U} = (U_\alpha)_{\alpha > 0}$. Further, let $\mu \in \mathcal{P}(E)$ and $f, g : E \rightarrow \mathbb{R}$ be \mathcal{B} -measurable such that

$$f \text{ is bounded and finely continuous} \quad \text{and} \quad \mu(U_{\alpha_0}|g|) < \infty \quad \text{for some } \alpha_0 > 0. \quad (6.16)$$

Consequently, the process $M = (M(t))_{t \geq 0}$ given by

$$M(t) := f(X(t)) - f(X(0)) - \int_0^t g(X(r))dt, \quad t \geq 0, \quad (6.17)$$

is well defined \mathbb{P}_μ -a.e. and from $L^1(\mathbb{P}_\mu)$. Furthermore, it is easily seen that M is a right-continuous *additive functional* under \mathbb{P}_μ , namely

$$[0, \infty) \ni r \mapsto M(r) \text{ is right-continuous and } M(t+s) = M(s) + M(t) \circ \theta(s) \quad \mathbb{P}_\mu\text{-a.s.}, \quad t, s \geq 0.$$

As a consequence, we deduce:

Lemma 6.54. The process M given by (6.17) is a right-continuous (\mathcal{F}_t) -martingale under \mathbb{P}_μ if and only if $\mathbb{E}_\mu \{M(t)\} = 0$ for every $t > 0$.

Another straightforward fact is the following:

Lemma 6.55. If the process M given by (6.17) is a well-defined and right-continuous (\mathcal{F}_t) -martingale under \mathbb{P}_x for $x \in E$ μ -a.e., then M is a well-defined and right-continuous (\mathcal{F}_t) -martingale under \mathbb{P}_ν for every $\nu \in \mathcal{P}$ for which there exists a constant $c > 0$ such that $\nu \leq c\mu$.

The following fine converse of Lemma 6.55 is useful:

Proposition 6.56. Let X be a right process on a Lusin measurable space (E, \mathcal{B}) , with transition $(P_t)_{t \geq 0}$ and resolvent $\mathcal{U} = (U_\alpha)_{\alpha > 0}$. Let $\mu \in \mathcal{P}(E)$ such that $\mu \circ U_\alpha \ll \mu$ for one (hence all) $\alpha > 0$, $f, g : E \rightarrow \mathbb{R}$ be \mathcal{B} -measurable such that (6.16) holds, and moreover, assume that the process M given by (6.17) is an $(\mathcal{F}(t))$ -martingale under \mathbb{P}_ν for all $\nu \in \mathcal{P}(E)$ for which there exists a constant $c > 0$ such that $\nu \leq c\mu$. Then there exists a set $E' \in \mathcal{B}$ such that $E \setminus E'$ is μ -inessential and M is a well-defined right-continuous (\mathcal{F}_t) -martingale under \mathbb{P}_x for every $x \in E'$.

Proof. Let us consider $E_0 := [U_{\alpha_0}|g| < \infty]$, so that, by Remark 6.33, $E_0 \in \mathcal{A}(\mathcal{U})$. Since $\mu(E \setminus E_0) = 0$ we deduce that $E \setminus E_0$ is μ -inessential.

Consequently, we can restrict the process X (see Definition 6.36), the measure μ , as well as the functions f and g , from E to E_0 . Note that by (6.16) we have that $M(t)$ is well defined \mathbb{P}_x -a.s. for every $x \in E_0$ and $t \geq 0$. Moreover,

$$M(t) \in L^1(\mathbb{P}_x) \cap L^1(\mathbb{P}_\mu), \quad t \geq 0, x \in E_0,$$

and it is easy to see that M is a right-continuous *additive functional* on E_0 , namely

$$[0, \infty) \ni r \mapsto M(r) \text{ is right-continuous and } M(t+s) = M(s) + M(t) \circ \theta(s) \quad \mathbb{P}_x\text{-a.s.},$$

for all $t, s \geq 0, x \in E_0$. As a consequence, by Lemma 6.54 we deduce: If $x \in E_0$, then the process M given by (6.17) is a right-continuous (\mathcal{F}_t) -martingale under \mathbb{P}_x if and only if $\mathbb{E}_x \{M(t)\} = 0$ for every $t \geq 0$.

Note that $\mathbb{E}_\nu \{M(t)\} = 0$ for every $t \geq 0$ and ν as in the statement. Thus,

$$\mathbb{E}_x \{M(t)\} = 0, \quad x \in E_0 \text{ } \mu\text{-a.e.}, \quad t \geq 0.$$

Furthermore, since for $x \in E_0$ we have that M is right-continuous \mathbb{P}_x -a.s., we deduce on the one hand that

$$\mathbb{E}_x \{M(t)\} = 0, \quad t \geq 0, \quad x \in E_0 \text{ } \mu\text{-a.e.} \quad (6.18)$$

We also have

$$\mathbb{E}_x \{M(t)\} = P_t f(x) - f(x) - \int_0^t P_r g(x) dr, \quad t \geq 0, x \in E_0.$$

Noticing that

$$\begin{aligned} e^{-\alpha t} \int_0^t P_r g(x) dr &= e^{-(\alpha-\alpha_0)t} \int_0^t e^{-\alpha_0 r} P_r g(x) dr \leq e^{-(\alpha-\alpha_0)t} \int_0^t e^{-\alpha_0 r} P_r g(x) dr \\ &\leq e^{-(\alpha-\alpha_0)t} U_{\alpha_0} g(x) \xrightarrow[t \rightarrow \infty]{} 0 \quad \text{for } \alpha > \alpha_0, \end{aligned}$$

by taking the Laplace transform and using integration by parts we obtain, on the other hand, that for $x \in E_0$

$$\mathbb{E}_x \{M(t)\} = 0, \quad t \geq 0 \iff \alpha U_\alpha f(x) - f(x) - U_\alpha g(x) = 0, \quad \alpha_0 < \alpha \in \mathbb{Q}_+.$$

Note that since f is finely continuous and using also Remark 6.2, we deduce that

$$u_\alpha := \alpha U_\alpha f - f - U_\alpha g \quad \text{is finely continuous on } E_0, \quad \alpha > \alpha_0.$$

From (6.18), the fact that u_α is finely continuous, and Remark 6.22 we deduce

$$u_\alpha = 0 \text{ } \mu\text{-a.e. on } E_0, \quad \text{hence } E_\alpha := \{x \in E_0 : u_\alpha(x) = 0\} \text{ is } \mu\text{-polar.}$$

Consequently, by the same Remark 6.22 we deduce that there exists a μ -inessential set E'_α such that $E'_\alpha \supset E_\alpha$, $\alpha > \alpha_0$. Finally, we consider

$$E' := \bigcap_{\alpha_0 < \alpha \in \mathbb{Q}} E'_\alpha.$$

We can now conclude since by Remark 6.34 we get that E' is μ -inessential, and for every $x \in E'$ we have that $\mathbb{E}_x \{M(t)\} = 0, t \geq 0$, hence M is an (\mathcal{F}_t) -martingale under \mathbb{P}_x . \square

6.5 Adding jumps by perturbation with kernels

Let X be a right process on a Lusin measurable space (E, \mathcal{B}) which for simplicity is considered to be conservative, i.e. it has a.s. infinite lifetime. Let $(P_t)_{t \geq 0}$ denote its transition function, with resolvent $\mathcal{U} = (U_\alpha)_{\alpha > 0}$. In the sequel $(\mathcal{L}, \mathcal{D}_b(\mathcal{L}))$ will be the generator (on $b\mathcal{B}$) of X in the following sense:

$$\mathcal{D}_b(\mathcal{L}) := U_\alpha(b\mathcal{B}) \quad \text{and} \quad \mathcal{L}f := \alpha f - g, \quad f = U_\alpha g, \quad g \in b\mathcal{B}. \quad (6.19)$$

The following connection between $(\mathcal{L}, \mathcal{D}_b(\mathcal{L}))$ and the martingale problem associated to X is known. We present its proof for completeness.

Proposition 6.57. *The following assertions are equivalent for $f, g \in b\mathcal{B}$.*

(i) $f \in \mathcal{D}_b(\mathcal{L})$ and $\mathcal{L}f = g$.

(ii) *The process $(M_f(t) := f(X(t)) - f(X(0)) - \int_0^t g(X(s)) ds)_{t \geq 0}$ is a càdlàg $(\mathcal{F}(t))$ -martingale under \mathbb{P}_x for every $x \in E$.*

Proof. Assume (i). Since it is easy to see that

$$P_t f = f + \int_0^t P_s g \, ds, \quad t \geq 0,$$

it follows that $\mathbb{E}_x \left\{ f(X(t)) - f(X(0)) - \int_0^t g(X(s)) \, ds \right\} = 0$ for every $t \geq 0$. Assertion (ii) now follows by Lemma 6.54.

Assume now that (ii) holds. Then we clearly have

$$\mathbb{E}_x \{M_f(t)\} = 0, \quad t \geq 0,$$

hence

$$\int_0^\infty e^{-\alpha t} P_t f \, dt = \frac{1}{\alpha} f + \int_0^\infty e^{-\alpha t} \int_0^t P_s g \, ds \, dt, \quad \alpha > 0,$$

which means that

$$f = U_\alpha (\alpha f - g) \in \mathcal{D}_b(\mathcal{L}), \quad \text{and} \quad \mathcal{L}f = \alpha f - (\alpha - g) = g.$$

□

Remark 6.58. Assertion (ii) in Proposition 6.57 is precisely the property of the function $f \in b\mathcal{B}$ to belong to the domain of the extended generator considered in [48]; see also Remark 3.3 from [11]. Consequently, Proposition 6.57 shows that $\mathcal{D}_b(\mathcal{L})$ coincides with the domain from [48] of the extended generator acting on bounded functions.

Let K be a bounded kernel on (E, \mathcal{B}) , $(P'_t)_{t \geq 0}$ be the transition function of the process X killed with the multiplicative functional induced by $K1$ (cf. Ch. III from [33]), and let $(\mathcal{L} - K1, \mathcal{D}_b(\mathcal{L} - K1))$ be its generator. Applying Lemma 1.4 from [31] on $b\mathcal{B}$ it follows that

$$\mathcal{D}_b(\mathcal{L} - K1) = \mathcal{D}_b(\mathcal{L}) \quad \text{and} \quad (\mathcal{L} - K1)u = \mathcal{L}u - K1u \quad \text{for every } u \in \mathcal{D}_b(\mathcal{L}).$$

The next proposition is a version of Proposition 4.5 from [27]; see also Proposition 3.4 from [28].

Proposition 6.59. *Let K be a bounded kernel on (E, \mathcal{B}) . Then the following assertions hold.*

(i) *There exists a right process $X^K = (\Omega^K, \mathcal{F}^K, \mathcal{F}_t^K, X^K(t), \theta^K(t), \mathbb{P}_x^K)$ with state space E , having the generator*

$$(\mathcal{L} - K1 + K, \mathcal{D}_b(\mathcal{L} - K1 + K)) \quad \text{and} \quad \mathcal{D}_b(\mathcal{L} - K1 + K) = \mathcal{D}_b(\mathcal{L}).$$

Moreover, the fine topology induced by (the resolvent of) X^K coincides with that induced by X .

(ii) *The process X^K solves the martingale problem for $(\mathcal{L} - K1 + K, \mathcal{D}_b(\mathcal{L}))$ under \mathbb{P}_x^K . More precisely, for every $f \in \mathcal{D}_b(\mathcal{L})$ and $x \in E$ the process*

$$\left(f(X^K(t)) - f(X^K(0)) - \int_0^t (\mathcal{L}f - K1f + Kf)(X^K(s)) \, ds \right)_{t \geq 0} \quad (6.20)$$

is a càdlàg (\mathcal{F}_t^K) -martingale under \mathbb{P}_x^K .

Proof. (i) We apply [27, Proposition 4.5] for $c := K1$ and for the sub-Markovian kernel $\frac{1}{c}K$, where the measure $\frac{1}{c}K(x, dy)$ is zero provided that $c(x) = 0$. It follows that for any $f \in bp\mathcal{B}$ the equation

$$r_t(x) = P'_t f(x) + \int_0^t P'_{t-u}(Kr_u)(x) \, du, \quad t \geq 0, \quad x \in E,$$

has a unique solution $Q_t f \in bp\mathcal{B}$, the function $(t, x) \mapsto Q_t f(x)$ is measurable, the family $(Q_t)_{t \geq 0}$ is a semigroup of sub-Markovian kernels on (E, \mathcal{B}) and it is the transition function of a right process $X^K =$

$(\Omega^K, \mathcal{F}^K, \mathcal{F}_t^K, X^K(t), \mathbb{P}_x^K)$ with state space E . Let $\mathcal{U}^K = (U_\alpha^K)_{\alpha>0}$ on (E, \mathcal{B}) be the resolvent of kernels induced by $(Q_t)_{t \geq 0}$, $U_\alpha^K = \int_0^\infty e^{-\alpha t} Q_t dt$, $\alpha > 0$. Then

$$U_\beta^K = U'_\beta + U'_\beta K U_\beta^K \quad \text{for all } \beta > 0, \quad (6.21)$$

where $\mathcal{U}' = (U'_\alpha)_{\alpha>0}$ is the resolvent of kernels induced by $(P'_t)_{t \geq 0}$. We have $\mathcal{E}(\mathcal{U}_\beta^K) \subset \mathcal{E}(\mathcal{U}'_\beta)$ and $[b\mathcal{E}(\mathcal{U}_\beta^K)] = [b\mathcal{E}(\mathcal{U}'_\beta)] = [b\mathcal{E}(\mathcal{U}_\beta)]$. In particular, the fine topology induced by X^K and the fine topology of X coincide.

Let $(\mathcal{L}^K, \mathcal{D}_b(\mathcal{L}^K))$ be the generator of X^K . By [31, Lemma 1.4] applied for the Banach space $b\mathcal{B}$, and using also (6.21), it follows that $\mathcal{D}_b(\mathcal{L}^K) = \mathcal{D}_b(\mathcal{L})$ and $\mathcal{L}^K = \mathcal{L} - K1 + K$. Hence $(\mathcal{L} - K1 + K, \mathcal{D}_b(\mathcal{L} - K1 + K))$ is the generator of X^K and we have $\mathcal{D}_b(\mathcal{L} - K1 + K) = \mathcal{D}_b(\mathcal{L})$ as claimed.

The proof of assertion (ii) is standard and we omit it; see e.g. Section 6.4. The right continuity of the martingale follows from Theorem 6.13 since $f \in \mathcal{D}_b(\mathcal{L})$, hence f is finely continuous; by the general theory, the martingale is in fact càdlàg. \square

Example: adding jumps to a path-continuous Markov process. The perturbation with a kernel of the generator was used in [9] and [31] in order to modify the jumps in the time evolution of a Markov process. However, the basic reference here is the pioneering articles [49] and [55]; see also [37] for a very recent related work.

Proposition 6.60 (cf. e.g. [55], [31]). *Assume that (E, τ) is locally compact with countable basis and \mathcal{B} is the Borel σ -algebra. Let X be the right process on E fixed above which is in addition assumed to have τ -continuous paths. Further, let K be a bounded kernel on (E, \mathcal{B}) . Then the right process X^K provided by Proposition 6.59 is a Hunt process on (E, τ) . Moreover, for every $f \in pb\mathcal{B} \otimes \mathcal{B}$, $f(y, y) = 0$, $y \in E$, we have*

$$\mathbb{E}_x \left\{ \sum_{s \leq t} f(X^K(s-), X^K(s)) \right\} = \mathbb{E}_x \left\{ \int_0^t \int_E f(X^K(s), y) K(X^K(s), dy) ds \right\}, \quad t \geq 0. \quad (6.22)$$

Proof. Let $X' = (\Omega', \mathcal{F}', \mathcal{F}'_t, X'(t), \mathbb{P}'^x)$ be the right process with state space E and lifetime ζ' , obtained by killing X with the multiplicative functional induced by $K1$. We consider (using the terminology from [55] and [31]) the (simple) Markov process X^K with τ -càdlàg paths obtained by resurrecting X' with the kernel R defined as

$$R(\omega', dy) = \frac{1_{[\zeta' < \infty]}(\omega')}{K1(X'(\zeta'-)(\omega'))} K(X'(\zeta'-)(\omega'), dy) + 1_{[\zeta' = \infty]}(\omega') \delta_\Delta(dy), \quad \omega' \in \Omega'.$$

The proof now finishes by noticing that X^K corresponds to the process constructed in [49, Theorem 2.5], hence it is in fact a Hunt process. \square

Remark 6.61. *Note that $f \in D(\mathcal{L})$ used as test function for the martingale problem (6.20) automatically satisfies the restriction $\mathcal{L}f \in b\mathcal{B}$. This, of course, fits very well with generators L in differential form that have locally bounded coefficients, which is the case of Lemma 3.12 and Proposition 3.13 in Section 3.4, or the case of the linearized operator \mathbf{L}_t^u given by (5.9) in Section 5.2. However, in Section 2 we are interested to cover more general coefficients, so in what follows we aim to present an L^1 -approach that allows to solve the perturbed martingale problem (6.20) for discontinuous diffusion coefficients and drift coefficients that are merely in L^2_{loc} .*

The L^1 -approach. In this paragraph we assume that there exists a σ -finite measure μ on E such that

$$(P_t)_{t \geq 0} \text{ can be extended to a } C_0\text{-semigroup on } L^1(\mu), \quad (6.23)$$

where recall that $(P_t)_{t \geq 0}$ is the transition function of the right process X fixed in the beginning of this part.

Remark 6.62. *Notice that, by [21, Proposition 2.1] it follows that (6.23) holds provided that μ is an excessive measure w.r.t. X , i.e., $\mu \circ P_t \leq \mu$ for all $t > 0$. The converse also holds: If a sub-Markovian C_0 -resolvent of contractions on $L^1(\mu)$ is given, then it is associated to a right process, however, one has to enlarge E by a zero set; cf. Theorem 2.2 from [17]. One can show that under certain additional conditions the above right process lives on E ; cf. [18].*

We further denote by $(L_1, D(L_1))$ the generator induced by $(P_t)_{t \geq 0}$ on $L^1(\mu)$.

Let K be a bounded kernel on (E, \mathcal{B}) as above, and consider the following assumption on K :

(\mathbf{H}_K^μ) K induces a bounded operator on $L^1(\mu)$, namely

$$\exists c \in (0, \infty) : \int_E |Kf(x)| \mu(dx) \leq c \int_E |f(x)| \mu(dx), \quad f \in L^1(\mu).$$

Remark 6.63. Let $v \in bp\mathcal{B}$. Then the multiplication by v is trivially a bounded operator on $L^p(\mu)$. Moreover, if K satisfies (\mathbf{H}_K^μ) , then vK trivially satisfies (\mathbf{H}_{vK}^μ) for any $v \in pb\mathcal{B}$, in particular for $v = K1$.

Proposition 6.64. Let X be the right process fixed above with its transition function satisfying (6.23), and K be a bounded kernel on (E, \mathcal{B}) which satisfies (\mathbf{H}_K^μ) , and set $c := K1$. Let X^K be the right process provided by Proposition 6.59. Then the following assertions hold.

(i) The transition function of X^K extends to a C_0 -semigroup on $L^1(\mu)$, whose generator is

$$L_1 - c + K, \quad D(L_1 - c + K) = D(L_1).$$

(ii) Let $f \in D(L_1) \cap b\mathcal{B}$ be finely continuous. Then there exists a set $\tilde{E} \in \mathcal{B}$ such that $E \setminus \tilde{E}$ is μ -inessential and for every $x \in \tilde{E}$

$$\mathbb{E}_x^K \left\{ \int_0^t |L_1 f - cf + Kf|(X^K(r)) dr \right\} < \infty, \quad t \geq 0,$$

and the process

$$\left(f(X^K(t)) - f(X^K(0)) - \int_0^t (L_1 f - cf + Kf)(X^K(s)) ds \right)_{t \geq 0} \quad (6.24)$$

is a càdlàg (\mathcal{F}_t^K) -martingale under \mathbb{P}_x^K .

Proof. (i). First of all, we can apply [43, Chapter III, 1.3] two times in a row, first for the perturbation $L_1 - c$ and then for the perturbation $L_1 - c + K$, to ensure that $(L_1 - c + K, D(L_1))$ generates a C_0 -semigroup on $L^1(\mu)$, which we denote by $(S_t)_{t \geq 0}$. Furthermore, by [43, Chapter III, 1.7] we have

$$S_t = P_t^c + \int_0^t P_{t-s}^c cK S_s ds \quad \text{as operators on } L^1(\mu), \quad t \geq 0.$$

It is then straightforward to deduce that

$$S_t f = Q_t f \quad \mu\text{-a.e.}, \quad f \in L^1(\mu) \cap b\mathcal{B}, \quad t \geq 0,$$

where recall that $(Q_t)_{t \geq 0}$ is the transition function of X^K . Thus, assertion (i) is proved.

(ii). Let $f \in L^1(\mu) \cap b\mathcal{B}$ be finely continuous and consider

$$M(t) := f(X^K(t)) - f(X^K(0)) - \int_0^t (L_1 - c + K)u(X^K(r))dr, \quad t \geq 0,$$

which is well defined \mathbb{P}_μ -a.e. and from $L^1(\mathbb{P}_\mu^K)$. Since $(Q_t)_{t \geq 0}$ is a C_0 -semigroup on $L^1(\mu)$ with generator $L_1 - c + K$, we get that $\mathbb{E}_\nu^K \{M(t)\} = 0$, $t \geq 0$ for every $\nu \in \mathcal{P}(E)$ which satisfy $\nu \leq \delta\mu$ for some constant $\delta \in (0, \infty)$. Now, assertion (ii) is entailed by Proposition 6.56. \square

The Lyapunov approach. In this paragraph we let X be the right process fixed above with transition function $(P_t)_{t \geq 0}$ and resolvent $\mathcal{U} := (U_\alpha)_{\alpha > 0}$, and assume

$$\exists V : E \rightarrow [0, \infty) \text{ } \mathcal{B}\text{-measurable and } \omega > 0 \text{ such that } P_t V \leq e^{\omega t} V, t \geq 0. \quad (6.25)$$

That is, using the terminology from Definition 6.1, V is \mathcal{U}_ω -supermedian.

Furthermore, we consider

$$\mathcal{B}_{\lesssim V} := \{f \in \mathcal{B} : \text{there exists } \delta \in (0, \infty) \text{ such that } |f| \leq \delta V\},$$

and set

$$\mathcal{D}_V(L) := U_\alpha(\mathcal{B}_{\lesssim V}), \alpha > \omega, \quad \text{and} \quad \mathcal{L}f := \alpha f - g, f = U_\alpha g, g \in \mathcal{B}_{\lesssim V}. \quad (6.26)$$

Remark 6.65. Note that $\mathcal{D}_V(\mathcal{L})$ is independent of $\alpha > \omega$.

Let us endow $\mathcal{B}_{\lesssim V}$ with the following weighted supremum norm

$$|f|_V := \sup_{x \in E} \frac{|f(x)|}{1 + V(x)}, \quad f \in \mathcal{B}_{\lesssim V}.$$

It is easy to check that $(\mathcal{B}_{\lesssim V}, |\cdot|_V)$ is a Banach space.

Further, we consider the following condition on K :

(H_K^V) K is a bounded kernel on (E, \mathcal{B}) and

$$\exists \delta \in (0, \infty) : KV \leq \delta V.$$

Proposition 6.66. Let X be the right process fixed above with its transition function satisfying (6.23), K be a bounded kernel on (E, \mathcal{B}) which satisfies **(H_K^V)**, and set $c := K1$. Let X^K be the right process provided by Proposition 6.59, with transition function $(Q_t)_{t \geq 0}$ and resolvent $\mathcal{U}^K := (U_\alpha^K)_{\alpha > 0}$. Then the following assertions hold.

(i) We have $\alpha U_{\alpha + \delta + \omega}^K V \leq V$, $\alpha > 0$, and if we set

$$\mathcal{D}_V(\mathcal{L} - c + K) := U_\alpha^K(\mathcal{B}_{\lesssim V}), \alpha > \delta + \omega, \quad (\mathcal{L} - c + K)f := \alpha f + g, \quad f = U_\alpha^K g, g \in \mathcal{B}_{\lesssim V},$$

then

$$\mathcal{D}_V(\mathcal{L} - c + K) = \mathcal{D}_V(\mathcal{L}).$$

(ii) For every $f \in \mathcal{D}_V(\mathcal{L})$ and every $x \in E$, the process

$$\left(f(X^K(t)) - f(X^K(0)) - \int_0^t (\mathcal{L}f - cf + Kf)(X^K(s)) ds \right)_{t \geq 0}$$

is a càdlàg (\mathcal{F}_t^K) -martingale under \mathbb{P}_x^K .

Proof. (i). Note that by (6.25) we deduce

$$\alpha U_{\alpha + \omega}^c V \leq \alpha U_{\alpha + \omega} V \leq V, \quad \alpha > 0$$

Also, iteration (6.21) we have

$$U_\alpha^K = U_\alpha^c \left(\sum_{n \geq 0} (KU_\alpha^c)^n \right) \quad \text{for all } \alpha > 0.$$

Moreover,

$$KU_{\alpha + \omega}^c V \leq \frac{\delta}{\alpha} V, \quad \alpha > 0.$$

By setting $V_n := V \wedge n$, $n \geq 1$, we have

$$\begin{aligned} U_{\alpha+\delta+\omega}^K V &= \sup_n U_{\alpha+\delta+\omega}^K V_n \leq \sup_n U_{\alpha+\delta+\omega}^c \left(\sum_{l \geq 0} (KU_{\alpha+\delta+\omega}^c)^l V_n \right) \\ &\leq U_{\alpha+\delta+\omega}^c \left(\sum_{l \geq 0} (KU_{\alpha+\delta+\omega}^c)^l V \right) \leq U_{\alpha+\delta+\omega}^c \left(V \sum_{l \geq 0} \left(\frac{\delta}{\alpha + \delta} \right)^l \right) \\ &= \frac{1}{\alpha} V \quad \text{for all } \alpha > 0. \end{aligned}$$

In order to prove the rest of assertion (i) we simply apply [31, Lemma 1.4] for the resolvents $(U_{\alpha+\delta+\omega})_{\alpha>0}$, $(U_{\alpha+\delta+\omega}^K)_{\alpha>0}$, and the Banach space $(\mathcal{B}_{\leq V}, |\cdot|_V)$.

(ii). This is now standard, see e.g. Section 6.4, noticing that the integrability of the martingale is ensured by (i). \square

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