# Stochastic flows in the Brownian web and net 

Emmanuel Schertzer Rongfeng Sun Jan M. Swart

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#### Abstract

It is known that certain one-dimensional nearest-neighbor random walks in i.i.d. random space-time environments have diffusive scaling limits. Here, in the continuum limit, the random environment is represented by a 'stochastic flows of kernels', which is a collection of random kernels that can be loosely interpreted as the transition probabilities of a Markov process in a random environment. The theory of stochastic flows of kernels was first developed by Le Jan and Raimond, who showed that each such flow is characterized by its $n$-point motions. Our work focuses on a class of stochastic flows of kernels with Brownian $n$-point motions which, after their inventors, will be called Howitt-Warren flows.

Our main result gives a graphical construction of general Howitt-Warren flows, where the underlying random environment takes on the form of a suitably marked Brownian web. This extends earlier work of Howitt and Warren who showed that a special case, the so-called 'erosion flow', can be constructed from two coupled 'sticky Brownian webs'. Our construction for general Howitt-Warren flows is based on a Poisson marking procedure developed by Newman, Ravishankar and Schertzer for the Brownian web. Alternatively, we show that a special subclass of the Howitt-Warren flows can be constructed as random flows of mass in a Brownian net, introduced by Sun and Swart.

Using these constructions, we prove some new results for the Howitt-Warren flows. In particular, we show that the kernels spread with a finite speed and have a locally finite support at deterministic times if and only if the flow is embeddable in a Brownian net. We show that the kernels are always purely atomic at deterministic times, but, with the exception of the erosion flows, exhibit random times when the kernels are purely nonatomic. We moreover prove ergodic statements for a class of measure-valued processes induced by the Howitt-Warren flows.

Our work also yields some new results in the theory of the Brownian web and net. In particular, we prove several new results about coupled sticky Brownian webs and about a natural coupling of a Brownian web with a Brownian net. We also introduce a 'finite graph representation' which gives a precise description of how paths in the Brownian net move between deterministic times.


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## 1 Introduction

### 1.1 Overview

In LR04a, Le Jan and Raimond introduced the notion of a stochastic flow of kernels, which is a collection of random probability kernels that can be loosely viewed as the transition kernels of a Markov process in a random space-time environment, where restrictions of the environment to disjoint time intervals are independent and the environment is stationary in time. For suitable versions of such a stochastic flow of kernels (when they exist), this loose interpretation is exact, see Definition 2.1 below and the remark following it. Given the environment, one can sample $n$ independent copies of the Markov process and then average over the environment. This defines the $n$-point motion for the flow, which satisfies a natural consistency condition: namely, the marginal distribution of any $k$ components of an $n$-point motion is necessarily a $k$-point motion. A fundamental result of Le Jan and Raimond [R04a shows that conversely, any family of Feller processes that is consistent in this way gives rise to an (essentially) unique stochastic flow of kernels.

As an example, in LR04b, the authors used Dirichlet forms to construct a consistent family of reversible $n$-point motions on the circle, which are $\alpha$-stable Lévy processes with some form of sticky interaction characterized by a real parameter $\theta$. In particular, for $\alpha=2$, these are sticky Brownian motions. Subsequently, Howitt and Warren HW09a used a martingale problem approach to construct a much larger class of consistent Feller processes on $\mathbb{R}$, which are Brownian motions with some form of sticky interaction characterized by a finite measure $\nu$ on $[0,1]$. In particular, if $\nu$ is a multiple of the Lebesgue measure, these are the sticky Brownian motions of Le Jan and Raimond. From now on, and throughout this paper, we specialize to the case of Browian underlying motions. By the general result of Le Jan and Raimond mentioned above, the sticky Brownian motions of Le Jan and Raimond, resp. Howitt and Warren, are the $n$-point motions of an (essentially) unique stochastic flow of kernels on $\mathbb{R}$, which we call a $L e$ Jan-Raimond flow, resp. Howitt-Warren flow (the former being a special case of the latter). It has been shown in [LL04, HW09a that these objects can be obtained as diffusive scaling limits of one-dimensional random walks in i.i.d. random space-time environments.

The main goal of the present paper is to give a graphical construction of Howitt-Warren flows that follows as closely as possible the discrete construction of random walks in an i.i.d. random environment. In particular, we want to make explicit what represents the random environment in the continuum setting. The original construction of Howitt-Warren flows using $n$-point motions does not tell us much about this. In HW09b, it was shown that the HowittWarren flow with $\nu=\delta_{0}+\delta_{1}$, known as the erosion flow, can be constructed using two coupled Brownian webs, where one Brownian web serves as the random space-time environment, while the conditional law of the second Brownian web determines the stochastic flow of kernels.

We will extend this construction to general Howitt-Warren flows, where in the general case, the random environment consists of a Brownian web together with a marked Poisson point process which is concentrated on the so-called points of type $(1,2)$ of the Brownian web. A central tool in this construction is a Poisson marking procedure invented by Newman, Ravishankar and Schertzer in [NRS10. Of course, we also make extensive use of the theory of the Brownian web developed in TW98, FINR04. For a special subclass of the HowittWarren flows, we will show that alternatively the random space-time environment can be represented as a Brownian net, plus a countable collection of i.i.d. marks attached to its socalled separation points. Here, we use the theory of the Brownian net, which was developed
in [S08 and [SSS09.
Using our graphical construction, we prove a number of new properties for the HowittWarren flows. In particular, we give necessary and sufficient conditions in terms of the measure $\nu$ for the random kernels to spread with finite speed, for their support to consist of isolated points at deterministic times, and for the existence of random times when the kernels are nonatomic (Theorems 2.5, 2.7 and 2.8 below). We moreover use our construction to prove the existence of versions of Howitt-Warren flows with nice regularity properties (Proposition 3.8 below), in particular, versions which can be interpreted as bona fide transition kernels in a random space-time environment. Lastly, we study the invariant laws for measure-valued processes associated with the Howitt-Warren flows (Theorem 2.11).

Our graphical construction of the Howitt-Warren flows is to a large extent motivated by its discrete space-time counterpart, i.e., random walks in i.i.d. random space-time environments on $\mathbb{Z}$. Many of our proofs will also be based on discrete approximation. Therefore, in the rest of the introduction, we will introduce a class of random walks in i.i.d. random spacetime environments and some related objects of interest, and sketch heuristically how the Brownian web and the Brownian net will arise in the representation of the random space-time environment for the Howitt-Warren flows. An outline of the rest of the paper will be given at the end of the introduction.

### 1.2 Discrete Howitt-Warren flows

Let $\mathbb{Z}_{\text {even }}^{2}:=\{(x, t): x, t \in \mathbb{Z}, x+t$ is even $\}$ be the even sublattice of $\mathbb{Z}^{2}$. We interpret the first coordinate $x$ as space and the second coordinate $t$ as time, which is plotted vertically in figures. Let $\omega:=\left(\omega_{z}\right)_{z \in \mathbb{Z}_{\text {even }}^{2}}$ be i.i.d. [0, 1]-valued random variables with common distribution $\mu$. We view $\omega$ as a random space-time environment for a random walk, such that conditional on the environment $\omega$, if the random walk is at time $t$ at the position $x$, then in the next unit time step the walk jumps to $x+1$ with probability $\omega_{(x, t)}$ and to $x-1$ with the remaining probability $1-\omega_{(x, t)}$ (see Figure 11).

To formalize this, let $\mathbb{P}$ denote the law of the environment $\omega$ and for each $(x, s) \in \mathbb{Z}_{\text {even }}^{2}$, let $\mathbf{Q}_{(x, s)}^{\omega}$ denote the conditional law, given the random environment $\omega$, of the random walk in random environment $X=(X(t))_{t \geq s}$ we have just described, started at time $s$ at position $X(s)=x$. Since parts of the random environment belonging to different times are independent, it is not hard to see that under the averaged (or 'annealed') law $\int \mathbb{P}(\mathrm{d} \omega) \mathbf{Q}_{(x, s)}^{\omega}$, the process $X$ is still a Markov chain, which in each time step jumps to the right with probability $\int \mu(\mathrm{d} q) q$ and to the left with the remaining probability $\int \mu(\mathrm{d} q)(1-q)$. Note that this is quite different from the usual random walk in random environment (RWRE) where the randomness is fixed for all time, and the averaged motion no longer has the Markov property.

We will be interested in three objects associated with the random walks in the i.i.d. random space-time environment $\omega$, namely: random transition kernels, $n$-point motions, and a measure-valued process. The law of each of these objects is uniquely characterized by $\mu$ and, conversely, uniquely determines $\mu$.

First of all, the random environment $\omega$ determines a family of random transition probability kernels,

$$
\begin{equation*}
K_{s, t}^{\omega}(x, y):=\mathbf{Q}_{(x, s)}^{\omega}\left[X_{t}=y\right] \quad\left(s \leq t,(x, s),(y, t) \in \mathbb{Z}_{\text {even }}^{2}\right) \tag{1.1}
\end{equation*}
$$

which satisfy
(i) $\sum_{y:(y, t) \in \mathbb{Z}_{\text {even }}^{2}} K_{s, t}^{\omega}(x, y) K_{t, u}^{\omega}(y, z)=K_{s, u}^{\omega}(x, z) \quad\left(s \leq t \leq u,(x, s),(z, u) \in \mathbb{Z}_{\text {even }}^{2}\right)$.


Figure 1: Random walk on $\mathbb{Z}_{\text {even }}^{2}$ in a random environment $\omega$.
(ii) For each $t_{0}<\cdots<t_{n}$, the random variables $\left(K_{t_{i-1}, t_{i}}^{\omega}\right)_{i=1, \ldots, n}$ are independent.
(iii) $K_{s, t}^{\omega}$ and $K_{s+u, t+u}^{\omega}$ are equal in law for each $u \in \mathbb{Z}_{\text {even }}:=\{2 x: x \in \mathbb{Z}\}$.

We call the collection of random probability kernels $\left(K_{s, t}^{\omega}\right)_{s \leq t}$ the discrete Howitt-Warren flow with characteristic measure $\mu$. Such a collection is a discrete time analogue of a stochastic flow of kernels as introduced by Le Jan and Raimond in [LR04a] (see Definition 2.1 below).

Next, given the environment $\omega$, we can sample a collection of independent random walks

$$
\begin{equation*}
(\vec{X}(t))_{t \geq 0}=\left(X_{1}(t), \ldots, X_{n}(t)\right)_{t \geq 0} \tag{1.2}
\end{equation*}
$$

in the random environment $\omega$, started at time zero from deterministic sites $x_{1}, \ldots, x_{n} \in \mathbb{Z}_{\text {even }}$, respectively. It is easy to see that under the averaged law

$$
\begin{equation*}
\int \mathbb{P}(\mathrm{d} \omega) \bigotimes_{i=1}^{n} \mathbf{Q}_{\left(x_{i}, 0\right)}^{\omega} \tag{1.3}
\end{equation*}
$$

the process $\vec{X}=(\vec{X}(t))_{t \geq 0}$ is still a Markov chain, which we call the discrete n-point motion. Its transition probabilities are given by

$$
\begin{equation*}
\mathbb{P}_{s, t}^{(n)}(\vec{x}, \vec{y})=\int \mathbb{P}(\mathrm{d} \omega) \prod_{i=1}^{n} K_{s, t}^{\omega}\left(x_{i}, y_{i}\right) \quad\left(s \leq t,\left(x_{i}, s\right),\left(y_{i}, t\right) \in \mathbb{Z}_{\mathrm{even}}^{2}, i=1, \ldots, n\right) \tag{1.4}
\end{equation*}
$$

Note that these discrete $n$-point motions are consistent in the sense that any $k$ coordinates of $\vec{X}$ are distributed as a discrete $k$-point motion. Each coordinate $X_{i}$ is distributed as a nearestneighbor random walk thats makes jumps to the right with probability $\int \mu(\mathrm{d} q) q$. Because of the spatial independence of the random environment, the coordinates evolve independently when they are at different positions. To see that there is some nontrivial interaction when they are at the same position, note that if $k+l$ coordinates are at position $x$ at time $t$, then
the probability that in the next time step the first $k$ coordinates jump to $x+1$ while the last $l$ coordinates jump to $x-1$ equals $\int \mu(\mathrm{d} q) q^{k}(1-q)^{l}$, which in general does not factor into $\left(\int \mu(\mathrm{d} q) q\right)^{k}\left(\int \mu(\mathrm{~d} q)(1-q)\right)^{l}$. Note that the law of $\omega_{(0,0)}$ is uniquely determined by its moments, which are in turn determined by the transition probabilities of the discrete $n$-point motions (for each $n$ ).

Finally, based on the family of kernels $\left(K_{s, t}^{\omega}\right)_{s \leq t}$, we can define a measure-valued process

$$
\begin{equation*}
\rho_{t}(x)=\sum_{y \in \mathbb{Z}_{\text {even }}} \rho_{0}(y) K_{0, t}^{\omega}(y, x) \quad\left(t \geq 0,(x, t) \in \mathbb{Z}_{\text {even }}^{2}\right) \tag{1.5}
\end{equation*}
$$

where $\rho_{0}$ is any locally finite initial measure on $\mathbb{Z}_{\text {even }}$. Note that conditional on $\omega$, the process $\rho=\left(\rho_{t}\right)_{t \geq 0}$ evolves deterministically, with

$$
\begin{equation*}
\rho_{t+1}(x):=\omega_{(x-1, t)} \rho_{t}(x-1)+\left(1-\omega_{(x+1, t)}\right) \rho_{t}(x+1) \quad\left((x, t+1) \in \mathbb{Z}_{\text {even }}^{2}, t \geq 0\right) \tag{1.6}
\end{equation*}
$$

Under the law $\mathbb{P}$, the process $\rho$ is a Markov chain, taking values alternatively in the spaces of finite measures on $\mathbb{Z}_{\text {even }}$ and $\mathbb{Z}_{\text {odd }}:=\{2 x+1: x \in \mathbb{Z}\}$. Note that (1.6) says that in the time step from $t$ to $t+1$, an $\omega_{(x, t)}$-fraction of the mass at $x$ is sent to $x+1$ and the rest is sent to $x-1$. Obviously, this dynamics preserves the total mass. In particular, if $\rho_{0}$ is a probability measure, then $\rho_{t}$ is a probability measure for all $t \geq 0$. We call $\rho$ the discrete Howitt-Warren process.

We will be interested in the diffusive scaling limits of all these objects, which will be (continuum) Howitt-Warren flows and their associated $n$-point motions and measure-valued processes, respectively. Note that the discrete Howitt-Warren flow $\left(K_{s, t}^{\omega}\right)_{s \leq t}$ determines the random environment $\omega$ a.s. uniquely. The law of $\left(K_{s, t}^{\omega}\right)_{s \leq t}$ is uniquely determined by either the law of its $n$-point motions or the law of its associated measure-valued process.

### 1.3 Scaling limits of discrete Howitt-Warren flows

We now recall from HW09a the conditions under which the $n$-point motions of a sequence of discrete Howitt-Warren flows converge to the $n$-point motions of a (continuum) stochastic flow of kernels, which we call a Howitt-Warren flow. We will then use discrete approximation to sketch heuristically how such a Howitt-Warren flow can be constructed from a Brownian web or net.

Let $\left(\varepsilon_{k}\right)_{k \in \mathbb{N}}$ be positive constants tending to zero, and let $\left(\mu_{k}\right)_{k \in \mathbb{N}}$ be probability laws on $[0,1]$ satisfying ${ }^{1}$
(i) $\varepsilon_{k}^{-1} \int(2 q-1) \mu_{k}(\mathrm{~d} q) \underset{k \rightarrow \infty}{\longrightarrow} \beta$,

$$
\begin{equation*}
\varepsilon_{k}^{-1} q(1-q) \mu_{k}(\mathrm{~d} q) \underset{k \rightarrow \infty}{\Longrightarrow} \nu(\mathrm{~d} q) \tag{ii}
\end{equation*}
$$

for some $\beta \in \mathbb{R}$ and finite measure $\nu$ on $[0,1]$, where $\Rightarrow$ denotes weak convergence. Howitt and Warren HW09a ${ }^{2}$ proved that under condition (1.7), if we scale space by $\varepsilon_{k}$ and time by

[^0]$\varepsilon_{k}^{2}$, then the discrete $n$-point motions with characteristic measure $\mu_{k}$ converge to a collection of Brownian motions with drift $\beta$ and some form of sticky interaction characterized by the measure $\nu$. These Brownian motions form a consistent family of Feller processes, hence by the general result of Le Jan and Raimond mentioned in Section 1.1, they are the $n$-point motions of some stochastic flow of kernels, which we call the Howitt-Warren flow with drift $\beta$ and characteristic measure $\nu$. The definition of Howitt-Warren flows and their $n$-point motions will be given more precisely in Section 2.

Now let us use discrete approximation to explain heuristically how to construct a HowittWarren flow based on a Brownian web or net. The construction based on the Brownian net is conceptually easier, so we consider this case first.

Let $\beta \in \mathbb{R}$ and let $\nu$ be a finite measure on $[0,1]$. Assuming, as we must in this case, that $\int \frac{\nu(\mathrm{d} q)}{q(1-q)}<\infty$, we may define a sequence of probability measures $\mu_{k}$ on $[0,1]$ by

$$
\begin{align*}
& \mu_{k}:=b \varepsilon_{k} \bar{\nu}+\frac{1}{2}\left(1-(b+c) \varepsilon_{k}\right) \delta_{0}+\frac{1}{2}\left(1-(b-c) \varepsilon_{k}\right) \delta_{1} \\
& \text { where } \quad b:=\int \frac{\nu(\mathrm{d} q)}{q(1-q)}, \quad c:=\beta-\int(2 q-1) \frac{\nu(\mathrm{d} q)}{q(1-q)}, \quad \bar{\nu}(\mathrm{d} q):=\frac{\nu(\mathrm{d} q)}{b q(1-q)} . \tag{1.8}
\end{align*}
$$

Then $\mu_{k}$ is a probability measure on $[0,1]$ for $k$ sufficiently large (such that $1-(b+|c|) \varepsilon_{k} \geq 0$ ), and the $\mu_{k}$ satisfy 1.7 . Thus, when space is rescaled by $\varepsilon_{k}$ and time by $\varepsilon_{k}^{2}$, the discrete Howitt-Warren flow with characteristic measure $\mu_{k}$ approximates a Howitt-Warren flow with drift $\beta$ and characteristic measure $\nu$.

Let $\omega^{\langle k\rangle}:=\left(\omega_{z}^{\langle k\rangle}\right)_{z \in \mathbb{Z}_{\text {even }}^{2}}$ be i.i.d. with common law $\mu_{k}$, which serves as the random environment for a discrete Howitt-Warren flow with characteristic measure $\mu_{k}$. We observe that for large $k$, most of the $\omega_{z}^{\langle k\rangle}$ are either zero or one. In view of this, it is convenient to alternatively encode $\omega^{\langle k\rangle}$ as follows. For each $z=(x, t) \in \mathbb{Z}_{\text {even }}^{2}$, if $\omega_{z}^{\langle k\rangle} \in(0,1)$, then we call $z$ a separation point, set $\bar{\omega}_{z}^{\langle k\rangle}=\omega_{z}^{\langle k\rangle}$, and we draw two arrows from $z$, leading respectively to $(x \pm 1, t+1)$. When $\omega_{z}^{\langle k\rangle}=0$, resp. 1 , we draw a single arrow from $z$ to $(x-1, t+1)$, resp. $(x+1, t+1)$. Note that the collection of arrows generates a branching-coalescing structure $N^{\langle k\rangle}$ on $\mathbb{Z}_{\text {even }}^{2}$ (see Figure 2) and conditional on $N^{\langle k\rangle}$, the $\bar{\omega}_{z}^{\langle k\rangle}$ at separation points $z$ of $N^{\langle k\rangle}$ are independent with common law $\bar{\nu}$. Therefore the random environment $\omega^{\langle k\rangle}$ can be represented by the pair $\left(N^{\langle k\rangle}, \bar{\omega}^{\langle k\rangle}\right)$, where a walk in such an environment must navigate along $N^{\langle k\rangle}$, and when it encounters a separation point $z$, it jumps either left or right with probability $1-\bar{\omega}^{\langle k\rangle}$, resp. $\bar{\omega}^{\langle k\rangle}$.

It turns out that the pair $\left(N^{\langle k\rangle}, \bar{\omega}^{\langle k\rangle}\right)$ has a meaningful diffusive scaling limit. In particular, if space is scaled by $\varepsilon_{k}$ and time by $\varepsilon_{k}^{2}$, then $N^{\langle k\rangle}$ converges to a limiting branching-coalescing structure $\mathcal{N}$ called the Brownian net, the theory of which was developed in [SS08, SSS09]. In particular, the separation points of $N^{\langle k\rangle}$ have a continuum analogue, the so-called separation points of $\mathcal{N}$, where incoming trajectories can continue along two groups of outgoing trajectories. These separation points are dense in space and time, but countable. Conditional on $\mathcal{N}$, we can then assign i.i.d. random variables $\bar{\omega}_{z}$ with common law $\bar{\nu}$ to the separation points of $\mathcal{N}$. The pair $(\mathcal{N}, \bar{\omega})$ provides a representation for the random space-time environment underlying the Howitt-Warren flow with drift $\beta$ and characteristic measure $\nu$. A random motion in such a random environment must navigate along $\mathcal{N}$, and whenever it comes to a separation point $z$, with probability $1-\bar{\omega}_{z}$ resp. $\bar{\omega}_{z}$, it continues along the left resp. right of the two groups of outgoing trajectories in $\mathcal{N}$ at $z$. We will recall the formal definition of the Brownian net and give a rigorous construction of random walks navigating in $\mathcal{N}$ in Section 4 .


Figure 2: Representation of the random environment $\left(\omega_{z}^{\langle k\rangle}\right)_{z \in \mathbb{Z}_{\text {even }}^{2}}$ in terms of a marked discrete net $\left(N^{\langle k\rangle}, \bar{\omega}^{\langle k\rangle}\right)$.


Figure 3: Representation of the random environment $\left(\omega_{z}^{\langle k\rangle}\right)_{z \in \mathbb{Z}_{\text {even }}^{2}}$ in terms of a marked discrete web $\left(W_{0}^{\langle k\rangle}, \omega^{\langle k\rangle}\right)$.

We now consider Howitt-Warren flows whose characteristic measure is a general finite measure $\nu$. Let $\left(\mu_{k}\right)_{k \in \mathbb{N}}$ satisfy 1.7 and let $\omega^{\langle k\rangle}:=\left(\omega_{z}^{\langle k\rangle}\right)_{z \in \mathbb{Z}_{\text {even }}^{2}}$ be an i.i.d. random spacetime environment with common law $\mu_{k}$. Contrary to the previous situation, it will now in general not be true that the most of the $\omega_{z}^{\langle k\rangle}$, s are either zero or one. Nevertheless, it is still true that for large $k$, most of the $\omega_{z}^{\langle k\rangle}$ s are either close to zero or to one. To take advantage of this fact, conditional on $\omega^{\langle k\rangle}$, we sample independent $\{-1,+1\}$-valued random variables $\left(\alpha_{z}^{\langle k\rangle}\right)_{z \in \mathbb{Z}_{\text {even }}^{2}}$ such that $\alpha_{z}^{\langle k\rangle}=+1$ with probability $\omega_{z}^{\langle k\rangle}$. For each $z=(x, t) \in \mathbb{Z}_{\text {even }}^{2}$, we draw an arrow from $(x, t)$ to $\left(x+\alpha_{z}^{\langle k\rangle}, t+1\right)$. These arrows define a coalescing structure $W_{0}^{\langle k\rangle}$ on $\mathbb{Z}_{\text {even }}^{2}$ (see Figure 3). Think of these arrows as assigning to each point $z$ a preferred direction, which, in most cases, will be +1 if $\omega_{z}^{\langle k\rangle}$ is close to one and -1 if $\omega_{z}^{\langle k\rangle}$ is close to zero.

Now let us describe the joint law of $\left(\omega^{\langle k\rangle}, \alpha^{\langle k\rangle}\right)$ differently. First of all, if we forget about $\omega^{\langle k\rangle}$, then the $\left(\alpha_{z}^{\langle k\rangle}\right)_{z \in \mathbb{Z}_{\text {even }}^{2}}$ are just i.i.d. $\{-1,+1\}$-valued random variables which take the value +1 with probability $\int q \mu_{k}(\mathrm{~d} q)$. Second, conditional on $\alpha^{\langle k\rangle}$, the random variables $\left(\omega_{z}^{\langle k\rangle}\right)_{z \in \mathbb{Z}_{\text {even }}^{2}}$ are independent with distribution

$$
\begin{equation*}
\mu_{k}^{1}:=\frac{(1-q) \mu_{k}(\mathrm{~d} q)}{\int(1-q) \mu_{k}(\mathrm{~d} q)}, \quad \text { resp. } \quad \mu_{k}^{\mathrm{r}}:=\frac{q \mu_{k}(\mathrm{~d} q)}{\int q \mu_{k}(\mathrm{~d} q)} \tag{1.9}
\end{equation*}
$$

depending on whether $\alpha_{z}^{\langle k\rangle}=-1$ resp. +1 . Therefore, we can alternatively construct our random space-time environment $\omega^{\langle k\rangle}$ in such a way, that first we construct an i.i.d. collection $\alpha^{\langle k\rangle}$ as above, and then conditional on $\alpha^{\langle k\rangle}$, independently for each $z \in \mathbb{Z}_{\text {even }}^{2}$, we choose $\omega_{z}^{\langle k\rangle}$ with law $\mu_{k}^{1}$ if $\alpha_{z}^{\langle k\rangle}=-1$ and law $\mu_{k}^{\mathrm{r}}$ if $\alpha_{z}^{\langle k\rangle}=+1$.

Let $W_{0}^{\langle k\rangle}$ denote the coalescing structure on $\mathbb{Z}_{\text {even }}^{2}$ generated by the arrows associated with $\left(\alpha_{z}^{\langle k\rangle}\right)_{z \in \mathbb{Z}_{\text {even }}^{2}}$ (see Figure 3). Then $\left(W_{0}^{\langle k\rangle}, \omega^{\langle k\rangle}\right)$ gives an alternative representation of the random environment $\omega^{\langle k\rangle}$. A random walk in such an environment navigates in such a way that whenever it comes to a point $z \in \mathbb{Z}_{\text {even }}^{2}$, the walk jumps to the right with probability $\omega_{z}^{\langle k\rangle}$ and to the left with the remaining probability. The important thing to observe is that if $k$ is large, then $\omega_{z}^{\langle k\rangle}$ is with large probability close to zero if $\alpha^{\langle k\rangle}=-1$ and close to one if $\alpha^{\langle k\rangle}=+1$. In view of this, the random walk in the random environment $\left(W_{0}^{\langle k\rangle}, \omega^{\langle k\rangle}\right)$ will most of its time walk along paths in $W_{0}^{\langle k\rangle}$.

It turns out that ( $W_{0}^{\langle k\rangle}, \omega^{\langle k\rangle}$ ) has a meaningful diffusive scaling limit. In particular, if space is scaled by $\varepsilon_{k}$ and time by $\varepsilon_{k}^{2}$, then the coalescing structure $W_{0}^{\langle k\rangle}$ converges to a limit $\mathcal{W}_{0}$ called the Brownian web (with drift $\beta$ ), which loosely speaking is a collection of coalescing Brownian motions starting from every point in space and time. These provide the default paths a motion in the limiting random environment must follow. The i.i.d. random variables $\omega_{z}^{\langle k\rangle}$ turn out to converge to a marked Poisson point process which is concentrated on socalled points of type $(1,2)$ in $\mathcal{W}_{0}$, which are points where there is one incoming path and two outgoing paths. These points are divided into points of type $(1,2)_{1}$ and $(1,2)_{\mathrm{r}}$, depending on whether the incoming path continues on the left or right. A random motion in such an environment follows paths in $\mathcal{W}_{0}$ by default, but whenever it comes to a marked point $z$ of type ( 1,2 ), it continues along the left resp. right outgoing path with probability $1-\omega_{z}$ resp. $\omega_{z}$. We will give the rigorous construction in Section 3. The procedure of marking a Poisson set of points of type $(1,2)$ that we need here was first developed by Newman, Ravishankar and Schertzer in [NRS10], who used it (among other things) to give an alternative construction of the Brownian net.

### 1.4 Outline and discussion

The rest of the paper is organized as follows. Sections $2 \sqrt{4}$ provide an extended introduction where we rigorously state our results. In Section 2, we recall the notion of a stochastic flow of kernels, first introduced in LR04a, and Howitt and Warren's HW09a sticky Brownian motions, to give a rigorous definition of Howitt-Warren flows. We then state out main results for these Howitt-Warren flows, including properties for the kernels and results for the associated measure-valued processes. In Section 33, in particular in Theorem 3.7, we present our construction of Howitt-Warren flows based on a 'reference' Brownian web with a Poisson marking, which is the main result of this paper. Along the way, we will recall the necessary background on the Brownian web. In Section4, we show that a special subclass of the HowittWarren flows can be constructed alternatively as flows of mass in the Brownian net, making the heuristics in Section 1.3 rigorous. Along the way, we will recall the necessary background on the Brownian net and establish some new results on a coupling between a Brownian web and a Brownian net. Sections 510 are devoted to proofs. In particular, we refer to Section 5 for an outline of the proofs. The paper concludes with a number of appendices.

Our work leaves several open problems. One question, for example, is how to characterize the measure-valued processes associated with a Howitt-Warren flow (see (2.1) below) by means of a well-posed martingale problem. Other questions (martingale problem formulation, path properties) refer to the duals (in the sense of linear systems duality) of these measure-valued processes, introduced in (11.1) below, which we have not investigated in much detail.

Moving away from the Brownian case, we note that it is an open problem whether our methods can be generalized to other stochastic flows of kernels than those introduced by Howitt and Warren. In particular, this applies to the stochastic flows of kernels with $\alpha$-stable Lévy $n$-point motions introduced in LR04b for $\alpha<2$. A first step on this road would be the construction of an $\alpha$-stable Lévy web which should generalize the presently known Brownian web. Some first steps in this direction have recently been taken in [EMS09.

## 2 Results for Howitt-Warren flows

In this section, we recall the notion of a stochastic flow of kernels, define the Howitt-Warren flows, and state our results on these Howitt-Warren flows, which include almost sure path properties and ergodic theorems for the associated measure-valued processes. The proofs of these results are based on our graphical construction of the Howitt-Warren flows, which we postpone to Sections 34 due to the extensive background we need to recall.

### 2.1 Stochastic flows of kernels

In LR04a, Le Jan and Raimond developed a theory of stochastic flows of kernels, which may admit versions that can be interpreted as the random transition probability kernels of a Markov process in a stationary random space-time environment. The notion of a stochastic flow of kernels generalizes the usual notion of a stochastic flow, which is a family of random mappings $\left(\phi_{s, t}^{\omega}\right)_{s \leq t}$ from a space $E$ to itself. In the special case that all kernels are deltameasures, a stochastic flow of kernels reduces to a stochastic flow in the usual sense of the word.

Since stochastic flows of kernels play a central role in our work, we take some time to recall their defintion. For any Polish space $E$, we let $\mathcal{B}(E)$ denote the Borel $\sigma$-field on $E$
and write $\mathcal{M}(E)$ and $\mathcal{M}_{1}(E)$ for the spaces of finite measures and probability measures on $E$, respectively, equipped with the topology of weak convergence and the associated Borel $\sigma$-field. By definition, a probability kernel on $E$ is a function $K: E \times \mathcal{B}(E) \rightarrow \mathbb{R}$ such that the map $x \mapsto K(x, \cdot)$ from $E$ to $\mathcal{M}_{1}(E)$ is measurable. By a random probability kernel, defined on some probability space $(\Omega, \mathcal{F}, \mathbb{P})$, we will mean a function $K: \Omega \times E \times \mathcal{B}(E) \rightarrow \mathbb{R}$ such that the map $(\omega, x) \mapsto K^{\omega}(x, \cdot)$ from $\Omega \times E$ to $\mathcal{M}_{1}(E)$ is measurable. We say that two random probability kernels $K, K^{\prime}$ are equal in finite dimensional distributions if for each $x_{1}, \ldots, x_{n} \in E$, the $n$-tuple of random probability measures $\left(K\left(x_{1}, \cdot\right), \ldots, K\left(x_{n}, \cdot\right)\right)$ is equally distributed with $\left(K^{\prime}\left(x_{1}, \cdot\right), \ldots, K^{\prime}\left(x_{n}, \cdot\right)\right)$. We say that two or more random probability kernels are independent if their finite-dimensional distributions are independent.

Definition 2.1 (Stochastic flow of kernels) $A$ stochastic flow of kernels on $E$ is a collection $\left(K_{s, t}\right)_{s \leq t}$ of random probability kernels on $E$ such that ${ }^{3}$
(i) For all $s \leq t \leq u$ and $x \in E$, a.s. $K_{s, s}(x, A)=\delta_{x}(A)$ and $\int_{E} K_{s, t}(x, \mathrm{~d} y) K_{t, u}(y, A)=$ $K_{s, u}(x, A)$ for all $A \in \mathcal{B}(E)$.
(ii) For each $t_{0}<\cdots<t_{n}$, the random probability kernels $\left(K_{t_{i-1}, t_{i}}\right)_{i=1, \ldots, n}$ are independent.
(iii) $K_{s, t}$ and $K_{s+u, t+u}$ are equal in finite-dimensional distributions for each real $s \leq t$ and $u$.

The finite-dimensional distributions of a stochastic flow of kernels are the laws of n-tuples of random probability measures of the form $\left(K_{s_{1}, t_{1}}\left(x_{1}, \cdot\right), \ldots, K_{s_{n}, t_{n}}\left(x_{n}, \cdot\right)\right)$, where $x_{i} \in E$ and $s_{i} \leq t_{i}, i=1, \ldots, n$.

Remark. If the random set of probability 1 on which Definition 2.1 (i) holds can be chosen uniformly for all $s \leq t \leq u$ and $x \in E$, then we can interpret $\left(K_{s, t}\right)_{s \leq t}$ as bona fide transition kernels of a random motion in random environment. For the stochastic flows of kernels we are interested in, we will prove the existence of a version of $K$ which satisfies this property (see Proposition 2.3 below). To the best of our knowledge, it is not known whether such a version always exists for general stochastic flows of kernels, even if we restrict ourselves to those defined by a consistent family of Feller processes.

If $\left(K_{s, t}\right)_{s \leq t}$ is a stochastic flow of kernels and $\rho_{0}$ is a finite measure on $E$, then setting

$$
\begin{equation*}
\rho_{t}(\mathrm{~d} y):=\int \rho_{0}(\mathrm{~d} x) K_{0, t}(x, \mathrm{~d} y) \quad(t \geq 0) \tag{2.1}
\end{equation*}
$$

defines an $\mathcal{M}(E)$-valued Markov process $\left(\rho_{t}\right)_{t \geq 0}$. Moreover, setting

$$
\begin{equation*}
P_{t-s}^{(n)}(\vec{x}, \mathrm{~d} \vec{y}):=\mathbb{E}\left[K_{s, t}\left(x_{1}, \mathrm{~d} y_{1}\right) \cdots K_{s, t}\left(x_{n}, \mathrm{~d} y_{n}\right)\right] \quad\left(\vec{x} \in E^{n}, s \leq t\right) \tag{2.2}
\end{equation*}
$$

defines a Markov transition function on $E^{n}$. We call the Markov process with these transition probabilities the n-point motion associated with the stochastic flow of kernels $\left(K_{s, t}\right)_{s \leq t}$. We observe that the $n$-point motions of a stochastic flow of kernels satisfy a natural consistency condition: namely, the marginal distribution of any $k$ components of an $n$-point motion is

[^1]necessarily a $k$-point motion for the flow. A fundamental result of Le Jan and Raimond [LR04a, Thm 2.1] states that conversely, any consistent family of Feller processes on a locally compact space $E$ gives rise to a stochastic flow of kernels on $E$ which is unique in finitedimensional distributions, $4^{4}$

### 2.2 Howitt-Warren flows

As will be proved in Proposition A.5 below, under the condition (1.7), if space and time are rescaled respectively by $\varepsilon_{k}$ and $\varepsilon_{k}^{2}$, then the $n$-point motions associated with the discrete Howitt-Warren flow introduced in Section 1.2 with characteristic measure $\mu_{k}$ converge to a collection of Brownian motions with drift $\beta$ and some form of sticky interaction characterized by the measure $\nu$. These Brownian motions solve a well-posed martingale problem, which we formulate now.

Let $\beta \in \mathbb{R}, \nu$ a finite measure on $[0,1]$, and define constants $\left(\beta_{+}(m)\right)_{m \geq 1}$ by

$$
\begin{align*}
& \beta_{+}(1):=\beta \quad \text { and } \\
& \beta_{+}(m):=\beta+2 \int \nu(\mathrm{~d} q) \sum_{k=0}^{m-2}(1-q)^{k} \quad(m \geq 2) \tag{2.3}
\end{align*}
$$

We note that in terms of these constants, 1.7 is equivalent to

$$
\begin{equation*}
\varepsilon_{k}^{-1} \int\left(1-2(1-q)^{m}\right) \mu_{k}(\mathrm{~d} q) \underset{k \rightarrow \infty}{\longrightarrow} \beta_{+}(m) \quad(m \geq 1) \tag{2.4}
\end{equation*}
$$

For $\emptyset \neq \Delta \subset\{1, \ldots, n\}$, we define

$$
\begin{equation*}
f_{\Delta}(\vec{x}):=\max _{i \in \Delta} x_{i} \quad \text { and } \quad g_{\Delta}(\vec{x}):=\left|\left\{i \in \Delta: x_{i}=f_{\Delta}(\vec{x})\right\}\right| \quad\left(\vec{x} \in \mathbb{R}^{n}\right) \tag{2.5}
\end{equation*}
$$

where $|\cdot|$ denotes the cardinality of a set.
The martingale problem we are about to formulate was invented by Howitt and Warren HW09a. We have reformulated their definition in terms of the functions $f_{\Delta}$ in (2.5), which form a basis of the vector space of test functions used in [HW09a, Def 2.1] (see Appendix A for a proof). This greatly simplifies the statement of the martingale problem and also facilitates our proof of the convergence of the $n$-point motions of discrete Howitt-Warren flows.

Definition 2.2 (Howitt-Warren martingale problem) We say that an $\mathbb{R}^{n}$-valued process $\vec{X}=(\vec{X}(t))_{t \geq 0}$ solves the Howitt-Warren martingale problem with drift $\beta$ and characteristic measure $\nu$ if $\vec{X}$ is a continuous, square-integrable semimartingale, the covariance process between $X_{i}$ and $X_{j}$ is given by

$$
\begin{equation*}
\left\langle X_{i}, X_{j}\right\rangle(t)=\int_{0}^{t} 1_{\left\{X_{i}(s)=X_{j}(s)\right\}} \mathrm{d} s \quad(t \geq 0, i, j=1, \ldots, n) \tag{2.6}
\end{equation*}
$$

and, for each nonempty $\Delta \subset\{1, \ldots, n\}$,

$$
\begin{equation*}
f_{\Delta}(\vec{X}(t))-\int_{0}^{t} \beta_{+}\left(g_{\Delta}(\vec{X}(s))\right) \mathrm{d} s \tag{2.7}
\end{equation*}
$$

is a martingale with respect to the filtration generated by $\vec{X}$.

[^2]Remark. We could have stated a similar martingale problem where instead of the functions $f_{\Delta}$ from (2.5) we use the functions $\tilde{f}_{\Delta}(x):=\min _{i \in \Delta} x_{i}$ and we replace the $\beta_{+}(m)$ defined in (2.3) by

$$
\begin{equation*}
\beta_{-}(1):=\beta \quad \text { and } \quad \beta_{-}(m):=\beta-2 \int \nu(\mathrm{~d} q) \sum_{k=0}^{m-2} q^{k} \quad(m \geq 2) . \tag{2.8}
\end{equation*}
$$

It is not hard to prove that both martingale problems are equivalent.
Remark. In the special case $n=2$, condition (2.7) is equivalent to the condition that

$$
\begin{equation*}
X_{1}(t)-\beta t, \quad X_{2}(t)-\beta t, \quad\left|X_{1}(t)-X_{2}(t)\right|-4 \nu([0,1]) \int_{0}^{t} 1_{\left\{X_{1}(s)=X_{2}(s)\right\}} \mathrm{d} s \tag{2.9}
\end{equation*}
$$

are martingales. In HW09a, such $\left(X_{1}, X_{2}\right)$ are called $\theta$-coupled Brownian motions, with $\theta=2 \nu([0,1])$. In this case, $X_{1}(t)-X_{2}(t)$ is a Brownian motion with stickiness at the origin.

Howitt and Warren HW09a, Prop. 8.1] proved that their martingale problem is wellposed and its solutions form a consistent family of Feller processes. Therefore, by the already mentioned result of Le Jan and Raimond [LR04a, Thm 2.1], there exists a stochastic flow of kernels $\left(K_{s, t}\right)_{s \leq t}$ on $\mathbb{R}$, unique in finite-dimensional distributions, such that the $n$-point motions of $\left(K_{s, t}\right)_{s \leq t}$ (in the sense of $(\overline{2.2})$ ) are given by the unique solutions of the HowittWarren martingale problem. We call this stochastic flow of kernels the Howitt-Warren flow with drift $\beta$ and characteristic measure $\nu$. It can be shown that Howitt-Warren flows are the diffusive scaling limits, in the sense of weak convergence of finite dimensional distributions, of the discrete Howitt-Warren flows with characteristic measures $\mu_{k}$ satisfying (1.7). (Indeed, this is a direct consequence of Proposition A. 5 below on the convergence of $n$-point motions.)

We will show that it is possible to construct versions of Howitt-Warren flows which are bona fide transition probability kernels of a random motion in a random space-time environment, and the kernels have 'regular' parameter dependence.

Proposition 2.3 (Regular parameter dependence) For each $\beta \in \mathbb{R}$ and finite measure $\nu$ on $[0,1]$, there exists a version of the Howitt-Warren flow $\left(K_{s, t}\right)_{s \leq t}$ with drift $\beta$ and characteristic measure $\nu$ such that in addition to the properties (i)-(iii) from Definition 2.1:
(i)' A.s., $\int_{E} K_{s, t}(x, \mathrm{~d} y) K_{t, u}(y, A)=K_{s, u}(x, A)$ for all $s \leq t \leq u, x \in E$ and $A \in \mathcal{B}(E)$.
(iv) A.s., the map $t \mapsto K_{s, t}(x, \cdot)$ from $[s, \infty)$ to $\mathcal{M}_{1}(\mathbb{R})$ is continuous for each $(s, x) \in \mathbb{R}^{2}$.

When the characteristic measure $\nu=0$, solutions to the Howitt-Warren martingale problem are coalescing Brownian motions. In this case, the associated stochastic flow of kernels is a stochastic flow (in the usual sense), which is known as the Arratia flow. In the special case that $\beta=0$ and $\nu$ is Lebesgue measure, the Howitt-Warren flow and its $n$-point motions are reversible. This stochastic flow of kernels has been constructed before (on the unit circle instead of $\mathbb{R}$ ) by Le Jan and Raimond in [R04b] using Dirichlet forms. We will call any stochastic flow of kernels with $\nu(\mathrm{d} x)=c \mathrm{~d} x$ for some $c>0$ a Le Jan-Raimond flow. In HW09b, Howitt and Warren constructed a stochastic flow of kernels with $\beta=0$ and $\nu=\frac{1}{2}\left(\delta_{0}+\delta_{1}\right)$, which they called the erosion flow. In this paper, we will call this flow the symmetric erosion flow and more generally, we will say that a Howitt-Warren flow is an erosion flow if $\nu=c_{0} \delta_{0}+c_{1} \delta_{1}$ with $c_{0}+c_{1}>0$. The paper HW09b] gives an explicit construction of the symmetric erosion


Figure 4: Four examples of Howitt-Warren flows. All examples have drift $\beta=0$ and stickiness parameter $\theta=2$. From left to right and from top to bottom: 1 . the equal splitting flow $\nu=\delta_{1 / 2}, 2$. the 'parabolic' flow $\nu(\mathrm{d} q)=6 q(1-q) \mathrm{d} q$, 3. the Le Jan-Raimond flow $\nu(\mathrm{d} q)=\mathrm{d} q$, 4 . the symmetric erosion flow $\nu=\frac{1}{2}\left(\delta_{0}+\delta_{1}\right)$. The first two flows have left and right speeds $\beta_{-}, \beta_{+}= \pm 4$ and $\beta_{-}, \beta_{+}= \pm 6$, respectively, while the last two flows have $\beta_{-}, \beta_{+}= \pm \infty$. The spatial dimension of each picture is 1.4 , the temporal dimension is 0.2 , and the initial state is Lebesgue measure.
flow based on coupled Brownian webs. Their construction can actually be extended to any erosion flow and can be seen as a precursor and special case of our construction of general Howitt-Warren flows in this paper.

### 2.3 Path properties

In this subsection, we state a number of results on the almost sure path properties of the measure-valued Markov process $\left(\rho_{t}\right)_{t \geq 0}$ defined in terms of a Howitt-Warren flow by (2.1). Throughout this subsection, we will assume that $\rho_{0}$ is a finite measure, and $\rho_{t}$ is defined using a version of the Howitt-Warren flow $\left(K_{s, t}\right)_{s \leq t}$, which satisfies property (iv) in Proposition 2.3, but not necessarily property (i)'. Then it is not hard to see that for any $\rho_{0} \in \mathcal{M}(\mathbb{R})$, the Markov process $\left(\rho_{t}\right)_{t \geq 0}$ defined in (2.1) has continuous sample paths in $\mathcal{M}(\mathbb{R})$. We call this process the Howitt-Warren process with drift $\beta$ and characteristic measure $\nu$.

See Figures 4 and 5 for some simulations of Howitt-Warren processes for various choices of the characteristic measure $\nu$. There are a number of parameters that are important for the


Figure 5: Example of an asymmetric flow: the one-sided erosion flow with $\beta=0$ and $\nu=\delta_{1}$. This flow has left and right speeds $\beta_{-}=-\infty$ and $\beta_{+}=2$, respectively. The spatial dimension of this picture is 3.2 and the temporal dimension is 0.4 . The initial state is Lebesgue measure up to the point 1.9 and zero from there onwards.
behavior of these processes. First of all, following [HW09a], we define

$$
\begin{equation*}
\theta(k, l)=\int \nu(\mathrm{d} q) q^{k-l}(1-q)^{l-1} \quad(k, l \geq 1) \tag{2.10}
\end{equation*}
$$

In a certain excursion theoretic sense, $\theta(k, l)$ describes the rate at which a group of $k+$ $l$ coordinates of the $n$-point motion that are at the same position splits into two groups consisting of $k$ and $l$ specified coordinates, respectively. In particular, following again notation in HW09a, we set

$$
\begin{equation*}
\theta:=2 \theta(1,1)=2 \int_{[0,1]} \nu(\mathrm{d} q), \tag{2.11}
\end{equation*}
$$

and we call $\theta$ the stickiness parameter of the Howitt-Warren flow. Note that when $\theta$ is increased, particles separate with a higher rate, hence the flow is less sticky. The next proposition shows that with the exception of the Arratia flow, by a simple transformation of space-time, we can always scale our flow such that $\beta=0$ and $\theta=2$. Below, for any $A \subset \mathbb{R}$ and $a \in \mathbb{R}$ we write $a A:=\{a x: x \in A\}$ and $A+a:=\{x+a: x \in A\}$.

Proposition 2.4 (Scaling and removal of the drift) Let $\left(K_{s, t}\right)_{s \leq t}$ be a Howitt-Warren flow with drift $\beta$ and characteristic measure $\nu$. Then:
(a) For each $a>0$, the stochastic flow of kernels $\left(K_{s, t}^{\prime}\right)_{s \leq t}$ defined by $K_{a^{2} s, a^{2} t}^{\prime}(a x, a A):=$ $K_{s, t}(x, A)$ is a Howitt-Warren flow with drift $a^{-1} \beta$ and characteristic measure $a^{-1} \nu$.
(b) For each $a \in \mathbb{R}$, the stochastic flow of kernels $\left(K_{s, t}^{\prime}\right)_{s \leq t}$ defined by $K_{s, t}^{\prime}(x+a s, A+a t):=$ $K_{s, t}(x, A)$ is a Howitt-Warren flow with drift $\beta+a$ and characteristic measure $\nu$.

There are two more parameters that are important for the behavior of a Howitt-Warren
flow. We define

$$
\begin{align*}
& \beta_{-}:=\beta-2 \int \nu(\mathrm{~d} q)(1-q)^{-1},  \tag{2.12}\\
& \beta_{+}:=\beta+2 \int \nu(\mathrm{~d} q) q^{-1}
\end{align*}
$$

Note that $\beta_{+}=\lim _{m \rightarrow \infty} \beta_{+}(m)$, where $\left(\beta_{+}(m)\right)_{m \geq 1}$ are the constants defined in 2.3). We call $\beta_{-}$and $\beta_{+}$the left speed and right speed of a Howitt-Warren flow, respectively. The next theorem shows that these names are justified. Below, $\operatorname{supp}(\mu)$ denotes the support of a measure $\mu$, i.e., the smallest closed set that contains all mass.

Theorem 2.5 (Left and right speeds) Let $\left(\rho_{t}\right)_{t \geq 0}$ be a Howitt-Warren process with drift $\beta$ and characteristic measure $\nu$, and let $\beta_{-}, \beta_{+}$be defined as in 2.12). Set $r_{t}:=\sup \left(\operatorname{supp}\left(\rho_{t}\right)\right)$ $(t \geq 0)$. Then:
(a) If $\beta_{+}<\infty$ and $r_{0}<\infty$, then $\left(r_{t}\right)_{t \geq 0}$ is a Brownian motion with drift $\beta_{+}$. If $\beta_{+}<\infty$ and $r_{0}=\infty$, then $r_{t}=\infty$ for all $t \geq 0$.
(b) If $\beta_{+}=\infty$, then $r_{t}=\infty$ for all $t>0$.

Analogue statements hold for $l_{t}:=\inf \left(\operatorname{supp}\left(\rho_{t}\right)\right)$, with $\beta_{+}$replaced by $\beta_{-}$.
It turns out that the support of a Howitt-Warren process is itself a Markov process. Let $\operatorname{Closed}(\mathbb{R})$ be the space of closed subsets of $\mathbb{R}$. We equip $\operatorname{Closed}(\mathbb{R})$ with a topology such that $A_{n} \rightarrow A$ if and only if $\bar{A}_{n} \xrightarrow{\text { Haus }} \bar{A}$, where $\bar{A}$ denotes the closure of a set $A$ in $[-\infty, \infty]$ and $\xrightarrow{\text { Haus }}$ means convergence of compact subsets of $[-\infty, \infty]$ in the Hausdorff topology. The branchingcoalescing point set is a $\operatorname{Closed}(\mathbb{R})$-valued Markov process that has been introduced in [SS08, Thm 1.11]. Its definition involves the Brownian net; see formula (4.6) below. The following proposition, which we cite from [SS08, Thm 1.11 and Prop. 1.15] and [SSS09, Prop. 3.14], lists some of its elementary properties.

Proposition 2.6 (Properties of the branching-coalescing point set) Let $\xi=\left(\xi_{t}\right)_{t \geq 0}$ be the branching-coalescing point set defined in (4.6), started in any initial state $\xi_{0} \in \operatorname{Closed}(\mathbb{R})$. Then:
(a) The process $\xi$ is a Closed $(\mathbb{R})$-valued Markov process with continuous sample paths.
(b) If $\sup \left(\xi_{0}\right)<\infty$, then $\left(\sup \left(\xi_{t}\right)\right)_{t \geq 0}$ is a Brownian motion with drift +1 . Likewise, if $-\infty<\inf \left(\xi_{0}\right)$, then $\left(\inf \left(\xi_{t}\right)\right)_{t \geq 0}$ is a Brownian motion with drift -1 .
(c) The law of a Poisson point set with intensity 2 is a reversible invariant law for $\xi$ and the limit law of $\xi_{t}$ as $t \rightarrow \infty$ for any initial state $\xi_{0} \neq \emptyset$.
(d) For each deterministic time $t>0$, a.s. $\xi_{t}$ is a locally finite subset of $\mathbb{R}$.
(e) Almost surely, there exists a dense set $\mathcal{T} \subset(0, \infty)$ such that for each $t \in \mathcal{T}$, the set $\xi_{t}$ contains no isolated points.

Our next result shows how Howitt-Warren processes and the branching-coalescing point set are related. Note that this result covers all possible values of $\beta_{-}, \beta_{+}$, except the case $\beta_{-}=\beta_{+}$ which corresponds to the Arratia flow. In (2.13) below, we continue to use the notation $a A+b:=\{a x+b: x \in A\}$.

Theorem 2.7 (Support process) Let $\left(\rho_{t}\right)_{t \geq 0}$ be a Howitt-Warren process with drift $\beta$ and characteristic measure $\nu$ and let $\beta_{-}, \beta_{+}$be defined as in (2.12). Then:
(a) If $-\infty<\beta_{-}<\beta_{+}<\infty$, then for all $t>0$,

$$
\begin{equation*}
\operatorname{supp}\left(\rho_{t}\right)=\frac{1}{2}\left(\beta_{+}-\beta_{-}\right) \xi_{t}+\frac{1}{2}\left(\beta_{-}+\beta_{+}\right) t, \tag{2.13}
\end{equation*}
$$

where $\left(\xi_{t}\right)_{t \geq 0}$ is a branching-coalescing point set.
(b) If $\beta_{-}=-\infty$ and $\beta_{+}<\infty$, then $\operatorname{supp}\left(\rho_{t}\right)=\left(-\infty, r_{t}\right] \cap \mathbb{R}$ for all $t>0$, where $r_{t}:=$ $\sup \left(\operatorname{supp}\left(\rho_{t}\right)\right)$. An analogue statement holds when $\beta_{-}>-\infty$ and $\beta_{+}=\infty$.
(c) If $\beta_{-}=-\infty$ and $\beta_{+}=\infty$, then $\operatorname{supp}\left(\rho_{t}\right)=\mathbb{R}$ for all $t>0$.

Proposition 2.6 (d) and Theorem 2.7 (a) imply that if the left and right speeds of a Howitt-Warren process are finite, then at deterministic times the process is purely atomic. The next theorem generalizes this statement to any Howitt-Warren process, but shows that if the characteristic measure puts mass on the open interval $(0,1)$, then there are random times when the statement fails to hold.

Theorem 2.8 (Atomicness) Let $\left(\rho_{t}\right)_{t \geq 0}$ be a Howitt-Warren process with drift $\beta$ and characteristic measure $\nu$. Then:
(a) For each $t>0$, the measure $\rho_{t}$ is a.s. purely atomic.
(b) If $\int_{(0,1)} \nu(\mathrm{d} q)>0$, then a.s. there exists a dense set of random times $t>0$ when $\rho_{t}$ is purely non-atomic.
(c) If $\int_{(0,1)} \nu(\mathrm{d} q)=0$, then a.s. $\rho_{t}$ is purely atomic at all $t>0$.

In the special case that $\nu$ is (a multiple of) Lebesgue measure, a weaker version of part (a) has been proved in [LR04b, Prop. 9 (c)]. Part (b) is similar to Proposition 2.6 (e) and in fact, by Theorem 2.7 (a), implies the latter. Note that parts (b) and (c) of the theorem reveal an interesting dichotomy between erosion flows (where $\nu$ is nonzero and concentrated on $\{0,1\}$ ) and all other Howitt-Warren flows (except the Arratia flow, for which atomicness is trivial). For erosion flows, we have an exact description of the set of space-time points where $\left(\rho_{t}\right)_{t \geq 0}$ has an atom in terms of an underlying Brownian web, see Theorem 9.6 below.

### 2.4 Infinite starting measures and discrete approximation

The ergodic behavior of the branching-coalescing point set is well-understood (see Proposition 2.6 (c)). As a consequence, by Theorem 2.7 (a), it is known that if we start a Howitt-Warren process with left and right speeds $\beta_{-}=-1, \beta_{+}=1$ in any nonzero initial state, then its support will converge in law to a Poisson point process with intensity 2. This does not mean, however, that the Howitt-Warren process itself converges in law. Indeed, since its 1-point motion is Brownian motion, it is easy to see that any Howitt-Warren process started in a finite initial measure satisfies $\lim _{t \rightarrow \infty} \mathbb{E}\left[\rho_{t}(K)\right]=0$ for any compact $K \subset \mathbb{R}$. To find nontrivial invariant laws, we must start the process in infinite initial measures.

To this aim, let $\mathcal{M}_{\text {loc }}(\mathbb{R})$ denote the space of locally finite measures on $\mathbb{R}$, endowed with the vague topology. Let $\left(K_{s, t}\right)_{s \leq t}$ be a version of the Howitt-Warren flow with $-\infty<\beta_{-}$and $\beta_{+}<\infty$, which satisfies Proposition 2.3 (iv). We will prove that for any $\rho_{0} \in \mathcal{M}_{\text {loc }}(\mathbb{R})$,

$$
\begin{equation*}
\rho_{t}:=\int \rho_{0}(\mathrm{~d} x) K_{0, t}(x, \cdot) \quad(t \geq 0) \tag{2.14}
\end{equation*}
$$

defines an $\mathcal{M}_{\text {loc }}(\mathbb{R})$-valued Markov process. If $\beta_{+}-\beta_{-}=\infty$, then mass can spread infinitely fast, hence we cannot define the Howitt-Warren process $\left(\rho_{t}\right)_{t \geq 0}$ for arbitrary $\rho_{0} \in \mathcal{M}_{\text {loc }}(\mathbb{R})$. In this case, we will use the class

$$
\begin{equation*}
\mathcal{M}_{\mathrm{g}}(\mathbb{R}):=\left\{\rho \in \mathcal{M}_{\mathrm{loc}}(\mathbb{R}): \int_{\mathbb{R}} e^{-c x^{2}} \rho(\mathrm{~d} x)<\infty \text { for all } c>0\right\} \tag{2.15}
\end{equation*}
$$

endowed with the topology that $\mu_{n} \rightarrow \mu$ if and only if $e^{-c x^{2}} \mu_{n}(\mathrm{~d} x)$ converges weakly to $e^{-c x^{2}} \mu(\mathrm{~d} x)$ for all $c>0$, which can be seen to be equivalent to $\mu_{n} \rightarrow \mu$ in the vague topology plus $\int e^{-c x^{2}} \mu_{n}(\mathrm{~d} x) \rightarrow \int e^{-c x^{2}} \mu(\mathrm{~d} x)$ for all $c>0$. Note that $\mathcal{M}_{\mathrm{loc}}(\mathbb{R})$ and $\mathcal{M}_{\mathrm{g}}(\mathbb{R})$ are Polish spaces.

Observe that by Definition 2.1 (i), the Howitt-Warren process $\left(\rho_{t}\right)_{t \geq 0}$ defined in 2.14) satisfies

$$
\begin{equation*}
\rho_{t}=\int \rho_{s}(\mathrm{~d} x) K_{s, t}(x, \cdot) \quad \text { a.s. } \tag{2.16}
\end{equation*}
$$

for each deterministic $s<t$. We will also use (2.16) to define Howitt-Warren processes starting at any deterministic time $s \in \mathbb{R}$.

Theorem 2.9 (Infinite starting mass and continuous dependence) Let $\beta \in \mathbb{R}$, let $\nu$ be a finite measure on $[0,1]$, and let $\left(K_{s, t}\right)_{s \leq t}$ be a version of the Howitt-Warren flow with drift $\beta$ and characteristic measure $\nu$ satisfying property (iv) from Proposition 2.3. Then:
(a) For any $\rho_{0} \in \mathcal{M}_{\mathrm{g}}(\mathbb{R})$, formula 2.14 defines an $\mathcal{M}_{\mathrm{g}}(\mathbb{R})$-valued Markov process with continuous sample paths, satisfying

$$
\begin{equation*}
\mathbb{E}\left[\rho_{t}(K)\right]<\infty \quad(t \geq 0, K \subset \mathbb{R} \text { compact }) \tag{2.17}
\end{equation*}
$$

Moreover, if $\left(\rho_{t}^{\langle n\rangle}\right)_{t \geq s_{n}}$ are processes started at times $s_{n}$ with deterministic initial data $\rho_{s_{n}}^{\langle n\rangle}$, and $s_{n} \rightarrow 0$, then for any $t>0$ and $t_{n} \rightarrow t$,

$$
\begin{equation*}
\rho_{s_{n}}^{\langle n\rangle} \underset{n \rightarrow \infty}{\Longrightarrow} \rho_{0} \quad \text { implies } \quad \rho_{t_{n}}^{\langle n\rangle} \underset{n \rightarrow \infty}{\Longrightarrow} \rho_{t} \quad \text { a.s. } \tag{2.18}
\end{equation*}
$$

where $\Rightarrow$ denotes convergence in $\mathcal{M}_{\mathrm{g}}(\mathbb{R})$.
(b) Assume moreover that $\beta_{+}-\beta_{-}<\infty$. Then, for any $\rho_{0} \in \mathcal{M}_{\mathrm{loc}}(\mathbb{R})$, formula 2.14) defines an $\mathcal{M}_{\text {loc }}(\mathbb{R})$-valued Markov process with continuous sample paths. Moreover, formula 2.18) holds with convergence in $\mathcal{M}_{\mathrm{g}}(\mathbb{R})$ replaced by vague convergence in $\mathcal{M}_{\mathrm{loc}}(\mathbb{R})$.

Remark. The convergence in (2.18) implies the continuous dependence of the law of $\rho_{t}$ on the starting time and the initial law, which is known as the Feller property. Note that when $\rho_{0}$ is a finite measure, the continuity in $t$ of $\rho_{t}$ in the space $\mathcal{M}(\mathbb{R})$ already follows from Proposition 2.3 (iv). However, for our purposes, we will only consider the spaces $\mathcal{M}_{\mathrm{g}}$ and $\mathcal{M}_{\text {loc }}$.

Remark. When $\beta_{+}-\beta_{-}=\infty,\left(\rho_{t}\right)_{t>0}$ may not be well-defined if $\rho_{0} \notin \mathcal{M}_{\mathrm{g}}(\mathbb{R})$. Indeed, by Theorem [2.7, if $\beta_{+}-\beta_{-}=\infty$, then for any fixed $t>0$, we can find $x_{n} \in \mathbb{Z}$ with $\left|x_{n}\right| \rightarrow \infty$ such that $\mathbb{P}\left(K_{0, t}\left(x_{n},[0,1]\right)<\varepsilon_{n}\right)<2^{-n}$ for some $\varepsilon_{n}>0$. Therefore $\rho_{0}:=\varepsilon_{n}^{-1} \delta_{x_{n}}$ has $\rho_{0} \in \mathcal{M}_{\mathrm{loc}}(\mathbb{R})$, and almost surely, $\rho_{t}([0,1])=\infty$.
Remark. Theorems 2.5, 2.7, and 2.8 carry over without change to the case of infinite starting measures. To see this, note that it is easy to check from (2.14) that

$$
\begin{equation*}
\rho_{0} \ll \tilde{\rho}_{0} \quad \text { implies } \quad \rho_{t} \ll \tilde{\rho}_{t} \quad(t \geq 0) \tag{2.19}
\end{equation*}
$$

where $\ll$ denotes absolute continuity. Since for each $\rho_{0} \in \mathcal{M}_{\text {loc }}(\mathbb{R})$, we can find a finite measure $\rho_{0}^{\prime}$ that is equivalent to $\rho$, statements about the support of $\rho_{t}$ and atomicness immediately generalize to the case of locally finite starting measures.

We also collect here a discrete approximation result for Howitt-Warren processes.
Theorem 2.10 (Convergence of discrete Howitt-Warren processes) Let $\varepsilon_{k}$ be positive constants converging to zero, and let $\mu_{k}$ be probability measures on $[0,1]$ satisfying (1.7) for some real $\beta$ and finite measure $\nu$ on $[0,1]$. Let $\left(\rho_{t}^{\langle k\rangle}\right)_{t \geq 0}$ be a discrete Howitt-Warren process with characteristic measure $\mu_{k}$ defined as in (1.5), where $K_{0, t}^{\omega}(x, \cdot)$ therein is defined for all $t>0$ by letting the random walk $\left(X_{t}\right)_{t \geq 0}$ in 1.1 be linearly interpolated between integer times. Let $\bar{\rho}_{t}^{\langle k\rangle}(\mathrm{d} x):=\rho_{\varepsilon_{k}^{-2} t}^{\langle k\rangle}\left(\varepsilon_{k}^{-1} \mathrm{~d} x\right)$. If $\rho_{0}^{\langle k\rangle}$ is deterministic and $\bar{\rho}_{0}^{\langle k\rangle} \Rightarrow \rho_{0}$ in $\mathcal{M}_{g}(\mathbb{R})$, then for any $T>0$,

$$
\begin{equation*}
\left(\bar{\rho}_{t}^{(k\rangle}\right)_{0 \leq t \leq T} \underset{k \rightarrow \infty}{\Longrightarrow}\left(\rho_{t}\right)_{0 \leq t \leq T}, \tag{2.20}
\end{equation*}
$$

where $\rho_{t}$ is a Howitt-Warren process with drift $\beta$, characteristic measure $\nu$, and initial condition $\rho_{0}$, and $\Rightarrow$ denotes weak convergence in law of random variables taking values in $\mathcal{C}\left([0, T], \mathcal{M}_{g}(\mathbb{R})\right)$, the space of continuous functions from $[0, T]$ to $\mathcal{M}_{g}(\mathbb{R})$ equipped with the uniform topology.

### 2.5 Ergodic properties

We are now ready to discuss the ergodic behavior of Howitt-Warren processes. Note that for a given Howitt-Warren flow $\left(K_{s, t}\right)_{s \leq t}$, the right-hand side of (2.14) is a.s. a linear function of the starting measure $\rho_{0}$. In view of this, Howitt-Warren processes belong to the class of so-called linear systems. The theory of linear systems on $\mathbb{Z}^{d}$ has been developed by Liggett and Spitzer, see e.g. [LS81] and [Lig05, Chap. IX]. We will adapt this theory to the continuum setting here. First we define the necessary notion.

We let $\mathcal{I}$ denote the set of invariant laws of a given Howitt-Warren processes, i.e., $\mathcal{I}$ is the set of probability laws $\Lambda$ on $\mathcal{M}_{\text {loc }}(\mathbb{R})\left(\operatorname{resp} . \mathcal{M}_{\mathbf{g}}(\mathbb{R})\right.$ if $\left.\beta_{+}-\beta_{-}=\infty\right)$ such that $\mathbb{P}\left[\rho_{0} \in \cdot\right]=\Lambda$ implies $\mathbb{P}\left[\rho_{t} \in \cdot\right]=\Lambda$ for all $t \geq 0$. We let $\mathcal{T}$ denote the set of homogeneous (i.e., translation invariant) laws on $\mathcal{M}_{\text {loc }}(\mathbb{R})\left(\right.$ resp. $\left.\mathcal{M}_{\mathbf{g}}(\mathbb{R})\right)$, i.e., laws $\Lambda$ such that $\mathbb{P}[\rho \in \cdot]=\Lambda$ implies $\mathbb{P}\left[T_{a} \rho \in \cdot\right]=\Lambda$ for all $a \in \mathbb{R}$, where $T_{a} \rho(A):=\rho(A+a)$ denotes the spatial shift map. Note that $\mathcal{I}$ and $\mathcal{T}$ are both convex sets. We write $\mathcal{I}_{\mathrm{e}}, \mathcal{T}_{\mathrm{e}}$, and $(\mathcal{I} \cap \mathcal{T})_{\mathrm{e}}$ to denote respectively the set of extremal elements in $\mathcal{I}, \mathcal{T}$, and $\mathcal{I} \cap \mathcal{T}$. Below, $\mathcal{C}_{\mathrm{c}}(\mathbb{R})$ denotes the space of continuous real function on $\mathbb{R}$ with compact support.

Theorem 2.11 (Homogeneous invariant laws for Howitt-Warren processes) Let $\beta \in$ $\mathbb{R}$, let $\nu$ be a finite measure on $[0,1]$ with $\nu \neq 0$. Then for the corresponding Howitt-Warren process $\left(\rho_{t}\right)_{t \geq 0}$, we have:
(a) $(\mathcal{I} \cap \mathcal{T})_{\mathrm{e}}=\left\{\Lambda_{c}: c \geq 0\right\}$, where $\Lambda_{c}(\mathrm{~d}(c \rho))=\Lambda_{1}(\mathrm{~d} \rho)$ for all $c \geq 0$, and

$$
\begin{align*}
\int \Lambda_{1}(\mathrm{~d} \rho) \int \rho(\mathrm{d} x) \phi(x) & =\int \phi(x) \mathrm{d} x  \tag{2.21}\\
\int \Lambda_{1}(\mathrm{~d} \rho) \int \rho(\mathrm{d} x) \phi(x) \int \rho(\mathrm{d} y) \psi(y) & =\int \phi(x) \mathrm{d} x \int \psi(y) \mathrm{d} y+\frac{\int \phi(x) \psi(x) \mathrm{d} x}{2 \nu([0,1])} \tag{2.22}
\end{align*}
$$

for any $\phi, \psi \in \mathcal{C}_{\mathbf{c}}(\mathbb{R})$.
(b) If $\mathbb{P}\left[\rho_{0} \in \cdot\right] \in \mathcal{T}_{\mathrm{e}}$ and $\mathbb{E}\left[\rho_{0}([0,1])\right]=c \geq 0$, then $\mathbb{P}\left[\rho_{t} \in \cdot\right]$ converges weakly to $\Lambda_{c}$. Furthermore, if $\mathbb{E}\left[\rho_{0}([0,1])^{2}\right]<\infty$, then for any $\phi, \psi \in \mathcal{C}_{\mathbf{c}}(\mathbb{R})$,

$$
\begin{equation*}
\lim _{t \rightarrow \infty} \mathbb{E}\left[\int \rho_{t}(\mathrm{~d} x) \phi(x) \int \rho_{t}(\mathrm{~d} y) \psi(y)\right]=\int \Lambda_{c}(\mathrm{~d} \rho) \int \rho(\mathrm{d} x) \phi(x) \int \rho(\mathrm{d} y) \psi(y) \tag{2.23}
\end{equation*}
$$

(c) If $\mathbb{P}\left[\rho_{0} \in \cdot\right] \in \mathcal{T}_{\mathrm{e}}$ and $\mathbb{E}\left[\rho_{0}([0,1])\right]=\infty$, then the laws $\mathbb{P}\left[\rho_{t} \in \cdot\right]$ have no weak cluster point as $t \rightarrow \infty$ which is supported on $\mathcal{M}_{\mathrm{loc}}(\mathbb{R})$.
(d) If $\Lambda \in \mathcal{I} \cap \mathcal{T}$, then there exists a probability measure $\gamma$ on $[0, \infty)$ such that $\Lambda=$ $\int_{0}^{\infty} \gamma(\mathrm{d} c) \Lambda_{c}$.

Remark. When $\nu$ is Lebesgue measure, it is known that (see [LR04b, Prop. 9 (b)]) $\Lambda_{\mathrm{c}}$ is the law of $c \rho^{*}$, where $\rho^{*}=\sum_{(x, u) \in \mathcal{P}} u \delta_{x}$ for a Poisson point process $\mathcal{P}$ on $\mathbb{R} \times[0, \infty)$ with intensity measure $\mathrm{d} x \times u^{-1} e^{-u} \mathrm{~d} u$.

Theorem 2.11 shows that each Howitt-Warren process has a unique (modulo a constant multiple) homogeneous invariant law, which by (2.22) has zero off-diagonal correlations. Moreover, any ergodic law at time 0 with finite density converges under the dynamics to the unique homogeneous invariant law with the same density.

In line with Theorems 2.7 and 2.8 we have the following support properties for $\Lambda_{c}$.
Theorem 2.12 (Support of stationary process) Let $c>0$ and let $\rho$ be an $\mathcal{M}_{\mathrm{loc}}(\mathbb{R})$-valued random variable with law $\Lambda_{c}$, the extremal homogeneous invariant law defined in Theorem 2.11 for the Howitt-Warren process with drift $\beta$ and characteristic measure $\nu \neq 0$. Then:
(a) If $\beta_{+}-\beta_{-}<\infty$, then $\operatorname{supp}(\rho)$ is a Poisson point process with intensity $\beta_{+}-\beta_{-}$.
(b) If $\beta_{+}-\beta_{-}=\infty$, then $\rho$ is a.s. atomic with $\operatorname{supp}(\rho)=\mathbb{R}$.

## 3 Construction of Howitt-Warren flows in the Brownian web

In this section, we make the heuristics in Section 1.3 rigorous and give a graphical construction of the Howitt-Warren flows using a procedure of Poisson marking of the Brownian web invented by Newman, Ravishankar and Schertzer [NRS10. The random environment for the HowittWarren flow will turn out to be a Brownian web, which we call the reference web, plus a marked Poisson point process on the reference web. Given such an environment, we will then construct a second coupled Brownian web, which we call the sample web, which is constructed by modifying the reference web by switching the orientation of marked points of type (1,2). The kernels of the Howitt-Warren flow are then constructed from the quenched law of the sample web, conditional on the reference web and the associated marked Poisson point process.

This construction generalizes the construction of the erosion flow based on coupled Brownian webs given in HW09b. For erosion flows, the random environment consists only of a reference web (without marked points) and the construction of the sample web can be done by specifying the joint law of the reference web and the sample web by means of a martingale problem. This is the approach taken in HW09b]. In the general case, when the random environment also contains marked points, this approach does not work. Therefore, in our approach, even for erosion flows, we will give a graphical construction of the sample web by marking and switching paths in the reference web.

Discrete approximation will be an important tool in many of our proofs and is helpful for understanding the continuum models. Therefore, in Section 3.1, we will first formulate the notion of a quenched law of sample webs conditional on the random environment for discrete Howitt-Warren flows. In Section 3.2, we then recall the necessary background on the Brownian web and Poisson marking for the Brownian web. In Section 3.3, we show how coupled Brownian webs can be constructed by Poisson marking and switching paths in a reference web. In Section 3.4, we state our main result, Theorem 3.7, which is the construction of Howitt-Warren flows using the Poisson marking of a reference Brownian web, and we also state some regularity properties for the Howitt-Warren flows. Lastly, in Section 3.5, we state a convergence result on the quenched law of discrete webs, which will be used to identify the flows we construct in Theorem 3.7 as being, indeed, the Howitt-Warren flows defined in Section 2.2 earlier through their $n$-point motions. The statements of this section are proved in Sections 6 and 7.

### 3.1 A quenched law on the space of discrete webs

As in Section 1.2, let $\omega:=\left(\omega_{z}\right)_{z \in \mathbb{Z}_{\text {even }}^{2}}$ be i.i.d. [0, 1]-valued random variables with common distribution $\mu$. Instead of using $\omega$ as a random environment for a single random walk started from one fixed time and position, as we did in Section 1.2, we will now use $\omega$ as a random environment for a collection of coalescing random walks starting from each point in $\mathbb{Z}_{\text {even }}^{2}$. To this aim, conditional on $\omega$, let $\alpha=\left(\alpha_{z}\right)_{z \in \mathbb{Z}_{\text {even }}^{2}}$ be independent $\{-1,+1\}$-valued random variables such that $\alpha_{z}=+1$ with probability $\omega_{z}$ and $\alpha_{z}=-1$ with probability $1-\omega_{z}$. For each $(x, s) \in \mathbb{Z}_{\text {even }}^{2}$, we let $p_{(x, s)}^{\alpha}:\{s, s+1, \ldots\} \rightarrow \mathbb{Z}$ be the function $p_{(x, s)}^{\alpha}=p$ defined by

$$
\begin{equation*}
p(s):=x \quad \text { and } \quad p(t+1):=p(t)+\alpha_{(p(t), t)} \quad(t \geq s) \tag{3.1}
\end{equation*}
$$

Then $p_{(x, s)}^{\alpha}$ is the path of a random walk in the random environment $\omega$, started at time $s$ at position $x$. It is easy to see that paths $p_{z}^{\alpha}, p_{z^{\prime}}^{\alpha}$ starting at different points $z, z^{\prime}$ coalesce when they meet. We call the collection of paths

$$
\begin{equation*}
\mathcal{U}^{\alpha}:=\left\{p_{z}^{\alpha}: z \in \mathbb{Z}_{\text {even }}^{2}\right\} \tag{3.2}
\end{equation*}
$$

the discrete web associated with $\alpha$ (see Figure 6). Let $\mathbb{P}$ denote the law of $\omega$ and let

$$
\begin{equation*}
\mathbf{Q}^{\omega}:=\mathbb{P}\left[\mathcal{U}^{\alpha} \in \cdot \mid \omega\right] \tag{3.3}
\end{equation*}
$$

denote the conditional law of $\alpha$ given $\omega$. Then under the averaged law $\int \mathbb{P}(\mathrm{d} \omega) \mathbf{Q}^{\omega}$, paths in $\mathcal{U}^{\alpha}$ are coalescing random walks that in each time step jump to the right with probability $\int \mu(\mathrm{d} q) q$ and to the left with the remaining probability $\int \mu(\mathrm{d} q)(1-q)$.

We will be more interested in the quenched law $\mathbf{Q}^{\omega}$ defined in (3.3). One has

$$
\begin{equation*}
\mathbf{Q}_{z}^{\omega}=\mathbf{Q}^{\omega}\left[p_{z}^{\alpha} \in \cdot\right] \tag{3.4}
\end{equation*}
$$



Figure 6: A discrete web and its dual.
where $\mathbf{Q}_{z}^{\omega}$ is the conditional law of the random walk in random environment started from $z \in \mathbb{Z}_{\text {even }}^{2}$ defined in Section 1.2. In particular, by 1.1),

$$
\begin{equation*}
K_{s, t}^{\omega}(x, y)=\mathbf{Q}^{\omega}\left[p_{(x, s)}^{\alpha}(t)=y\right]=\mathbb{P}\left[p_{(x, s)}^{\alpha}(t)=y \mid \omega\right] \quad\left((x, s),(y, t) \in \mathbb{Z}_{\text {even }}^{2}\right), \tag{3.5}
\end{equation*}
$$

where $\left(K_{s, t}^{\omega}\right)_{s \leq t}$ is the discrete Howitt-Warren flow with characteristic measure $\mu$. In view of this, the random law $\mathbf{Q}^{\omega}$ contains all information that we are interested in. We call $\mathbf{Q}^{\omega}$ the discrete quenched law with characteristic measure $\mu$. In the next sections, we will construct a continuous analogue of this quenched law and use it to define Howitt-Warren flows.

### 3.2 The Brownian web

As pointed out in the previous subsection, under the averaged law $\int \mathbb{P}(\mathrm{d} \omega) \mathbf{Q}^{\omega}$, the discrete web $\mathcal{U}^{\alpha}$ is a collection of coalescing random walks, started from every point in $\mathbb{Z}_{\text {even }}^{2}$. It turns out that such discrete webs have a well-defined diffusive scaling limit, which is basically a collection of coalescing Brownian motions, starting from each point in space and time, and which is called a Brownian web. The Brownian web arose from the work of Arratia Arr79, Arr81 and has since been studied by Tóth and Werner [TW98]. More recently, Fontes, Isopi, Newman and Ravishankar [FINR04] have introduced a by now standard framework in which the Brownian web is regarded as a random compact set of paths in a suitable Polish space.

It turns out that associated to each Brownian web, there is a dual Brownian web, which is a collection of coalescing Brownian motions running backwards in time. To understand this on a heuristic level, let $\left(\alpha_{z}\right)_{z \in \mathbb{Z}_{\text {even }}^{2}}$ be an i.i.d. collection of $\{-1,+1\}$-valued random variables. If for each $z=(x, t) \in \mathbb{Z}_{\text {even }}^{2}$, we draw an arrow from $(x, t)$ to $\left(x+\alpha_{z}, t+1\right)$, then paths along these arrows form a discrete web as introduced in the previous section. Now, if for each $z=(x, t) \in \mathbb{Z}_{\text {even }}^{2}$, we draw in addition a dual arrow from $(x, t+1)$ to $\left(x-\alpha_{z}, t\right)$, then paths along these dual arrows form a dual discrete web of coalescing random walks running backwards in time, which do not cross paths in the forward web (see Figure 6). The dual Brownian web arises as the diffusive scaling limit of such a dual discrete web.

We now introduce these objects formally. Let $R_{\mathrm{c}}^{2}$ be the compactification of $\mathbb{R}^{2}$ obtained by equipping the set $R_{\mathrm{c}}^{2}:=\mathbb{R}^{2} \cup\{( \pm \infty, t): t \in \mathbb{R}\} \cup\{(*, \pm \infty)\}$ with a topology such that


Figure 7: The compactification $R_{\mathrm{c}}^{2}$ of $\mathbb{R}^{2}$.
$\left(x_{n}, t_{n}\right) \rightarrow( \pm \infty, t)$ if $x_{n} \rightarrow \pm \infty$ and $t_{n} \rightarrow t \in \mathbb{R}$, and $\left(x_{n}, t_{n}\right) \rightarrow(*, \pm \infty)$ if $t_{n} \rightarrow \pm \infty$ (regardless of the behavior of $x_{n}$ ). An explicit way to construct such a compactification is as follows. Let $\Theta: \mathbb{R}^{2} \rightarrow \mathbb{R}^{2}$ be defined by

$$
\begin{equation*}
\Theta(x, t)=\left(\Theta_{1}(x, t), \Theta_{2}(t)\right):=\left(\frac{\tanh (x)}{1+|t|}, \tanh (t)\right) \tag{3.6}
\end{equation*}
$$

and let $\Theta\left(\mathbb{R}^{2}\right)$ denote the image of $\mathbb{R}^{2}$ under $\Theta$. Then the closure of $\Theta\left(\mathbb{R}^{2}\right)$ in $\mathbb{R}^{2}$ is in a natural way isomorphic to $R_{\mathrm{c}}^{2}$ (see Figure 7 ).

By definition, a path $\pi$ in $R_{\mathrm{c}}^{2}$ with starting time $\sigma_{\pi}$ is a function $\pi:\left[\sigma_{\pi}, \infty\right] \rightarrow[-\infty, \infty] \cup\{*\}$ such that $t \mapsto(\pi(t), t)$ is a continuous map from $\left[\sigma_{\pi}, \infty\right]$ to $R_{\mathrm{c}}^{2}$. We will often view paths as subsets of $R_{\mathrm{c}}^{2}$, i.e., we identify a path $\pi$ with its graph $\left\{(\pi(t), t): t \in\left[\sigma_{\pi}, \infty\right]\right\}$. We let $\Pi$ denote the space of all paths in $R_{\mathrm{c}}^{2}$ with all possible starting times in $[-\infty, \infty]$, equipped with the metric

$$
\begin{equation*}
d\left(\pi_{1}, \pi_{2}\right):=\left|\Theta_{2}\left(\sigma_{\pi_{1}}\right)-\Theta_{2}\left(\sigma_{\pi_{2}}\right)\right| \vee \sup _{t \geq \sigma_{\pi_{1}} \wedge \sigma_{\pi_{2}}}\left|\Theta_{1}\left(\pi_{1}\left(t \vee \sigma_{\pi_{1}}\right), t\right)-\Theta_{1}\left(\pi_{2}\left(t \vee \sigma_{\pi_{2}}\right), t\right)\right| \tag{3.7}
\end{equation*}
$$

and we let $\mathcal{K}(\Pi)$ denote the space of all compact subsets $K \subset \Pi$, equipped with the Hausdorff metric

$$
\begin{equation*}
d_{\mathrm{H}}\left(K_{1}, K_{2}\right)=\sup _{x_{1} \in K_{1}} \inf _{x_{2} \in K_{2}} d\left(x_{1}, x_{2}\right) \vee \sup _{x_{2} \in K_{2}} \inf _{x_{1} \in K_{1}} d\left(x_{1}, x_{2}\right) . \tag{3.8}
\end{equation*}
$$

Both $\Pi$ and $\mathcal{K}(\Pi)$ are complete separable metric spaces. The set $\hat{\Pi}$ of all dual paths $\hat{\pi}$ : $\left[-\infty, \hat{\sigma}_{\hat{\pi}}\right] \rightarrow[-\infty, \infty] \cup\{*\}$ with starting time $\hat{\sigma}_{\hat{\pi}} \in[-\infty, \infty]$ is defined analoguously to $\Pi$.

We adopt the convention that if $f: R_{\mathrm{c}}^{2} \rightarrow R_{\mathrm{c}}^{2}$ and $A \subset R_{\mathrm{c}}^{2}$, then $f(A):=\{f(z): z \in A\}$ denotes the image of $A$ under $f$. Likewise, if $\mathcal{A}$ is a set of subsets of $R_{\mathrm{c}}^{2}$ (e.g. a set of paths), then $f(\mathcal{A}):=\{f(A): A \in \mathcal{A}\}$. This also applies to notation such as $-A:=\{-z: z \in A\}$. If $\mathcal{A} \subset \Pi$ is a set of paths and $A \subset R_{\mathrm{c}}^{2}$, then we let $\mathcal{A}(A):=\left\{\pi \in \mathcal{A}:\left(\pi\left(\sigma_{\pi}\right), \sigma_{\pi}\right) \in A\right\}$ denote the subspace of all paths in $\mathcal{A}$ with starting points in $A$, and for $z \in R_{\mathrm{c}}^{2}$ we write $\mathcal{A}(z):=\mathcal{A}(\{z\})$.

The next proposition, which follows from [FINR04, Theorem 2.1], [FINR06, Theorem 3.7], and [SS08, Theorem 1.9], gives a characterization of the Brownian web $\mathcal{W}$ and its dual $\hat{\mathcal{W}}$. Below, we say that a path $\pi \in \Pi$ crosses a dual path $\hat{\pi} \in \hat{\Pi}$ from left to right if there exist
$\sigma_{\pi} \leq s<t \leq \hat{\sigma}_{\hat{\pi}}$ such that $\pi(s)<\hat{\pi}(s)$ and $\hat{\pi}(t)<\hat{\pi}(t)$. Crossing from right to left is defined analogously.

Proposition 3.1 (Characterization of the Brownian web and its dual) For each $\beta \in$ $\mathbb{R}$, there exists a $\mathcal{K}(\Pi) \times \mathcal{K}(\hat{\Pi})$-valued random variable $(\mathcal{W}, \hat{\mathcal{W}})$, called the double Brownian web with drift $\beta$, whose distribution is uniquely determined by the following properties:
(a) For each deterministic $z \in \mathbb{R}^{2}$, almost surely there is a unique path $\pi_{z} \in \mathcal{W}(z)$ and a unique dual path $\hat{\pi}_{z} \in \hat{\mathcal{W}}(z)$.
(b) For any deterministic countable dense subset $\mathcal{D} \subset \mathbb{R}^{2}$, almost surely, $\mathcal{W}$ is the closure in $\Pi$ of $\left\{\pi_{z}: z \in \mathcal{D}\right\}$ and $\hat{\mathcal{W}}$ is the closure in $\hat{\Pi}$ of $\left\{\hat{\pi}_{z}: z \in \mathcal{D}\right\}$.
(c) For any finite deterministic set of points $z_{1}, \ldots, z_{k} \in \mathbb{R}^{2}$, the paths $\left(\pi_{z_{1}}, \ldots, \pi_{z_{k}}\right)$ are distributed as a collection of coalescing Brownian motions, each with drift $\beta$.
(d) For any deterministic $z \in \mathbb{R}^{2}$, the dual path $\hat{\pi}_{z}$ is the a.s. unique path in $\hat{\Pi}(z)$ that does not cross any path in $\mathcal{W}$.

If $(\mathcal{W}, \hat{\mathcal{W}})$ is a double Brownian web as defined in Proposition 3.1, then we call $\mathcal{W}$ a Brownian web and $\hat{\mathcal{W}}$ the associated dual Brownian web. Note that $\hat{\mathcal{W}}$ is a.s. uniquely determined by $\mathcal{W}$. Although this is not obvious from the definition, the dual Brownian web is indeed a Brownian web rotated by 180 degrees. Indeed, $(\mathcal{W}, \hat{\mathcal{W}})$ is equally distributed with $(-\hat{\mathcal{W}},-\mathcal{W})$.

Definition 3.2 (Incoming and outgoing paths) We say that a path $\pi \in \Pi$ is an incoming path at a point $z=(x, t) \in \mathbb{R}^{2}$ if $\sigma_{\pi}<t$ and $\pi(t)=x$. We say that $\pi$ is an outgoing path at $z$ if $\sigma_{\pi}=t$ and $\pi(t)=x$. We say that two incoming paths $\pi_{1}, \pi_{2}$ at $z$ are strongly equivalent, denoted as $\pi_{1}={ }_{\mathrm{in}}^{z} \pi_{2}$, if $\pi_{1}=\pi_{2}$ on $[t-\varepsilon, t]$ for some $\varepsilon>0$. For $z \in \mathbb{R}^{2}$, let $m_{\mathrm{in}}(z)$ denote the number of equivalence classes of incoming paths in $\mathcal{W}$ at $z$ and let $m_{\text {out }}(z)$ denote the cardinality of $\mathcal{W}(z)$. Then $\left(m_{\mathrm{in}}(z), m_{\text {out }}(z)\right)$ is called the type of the point $z$ in $\mathcal{W}$. The type ( $\left.\hat{m}_{\text {in }}(z), \hat{m}_{\text {out }}(z)\right)$ of a point $z$ in the dual Brownian web $\hat{\mathcal{W}}$ is defined analogously.

We cite the following result from TW98, Proposition 2.4] or [FINR06, Theorems 3.113.14]. See Figure 8 for an illustration.

Proposition 3.3 (Special points of the Brownian web) Almost surely, all points $z \in \mathbb{R}^{2}$ are of one of the following types in $\mathcal{W} / \hat{\mathcal{W}}$ : $(0,1) /(0,1),(0,2) /(1,1),(0,3) /(2,1),(1,1) /(0,2)$, $(1,2) /(1,2)$, and $(2,1) /(0,3)$. For each deterministic $t \in \mathbb{R}$, almost surely, each point in $\mathbb{R} \times\{t\}$ is of type $(0,1) /(0,1),(0,2) /(1,1)$, or $(1,1) /(0,2)$. Deterministic points $z \in \mathbb{R}^{2}$ are a.s. of type $(0,1) /(0,1)$.

For us, points of type $(1,2) /(1,2)$ are of special importance. Note that these are the only points at which there are incoming paths both in $\mathcal{W}$ and in $\hat{\mathcal{W}}$. Points of type $(1,2)$ in $\mathcal{W}$ are further distinguished into points of type $(1,2)_{1}$ and $(1,2)_{\mathrm{r}}$, according to whether the left or the right outgoing path in $\mathcal{W}$ is the continuation of the (up to equivalence unique) incoming path.

Proposition 3.3 shows that although for each deterministic $z \in \mathbb{R}^{2}$, a.s. $\mathcal{W}(z)$ contains a single path, there exist random points $z$ where $\mathcal{W}(z)$ contains up to three paths. Sometimes, it will be necessary to choose a unique element of $\mathcal{W}(z)$ for each $z \in \mathbb{R}^{2}$. To that aim, for


Figure 8: Special points of the Brownian web. On the left: $(0,1) /(0,1)$. Top row: $(1,1) /(0,2)$, $(2,1) /(0,3),(1,2)_{1} /(1,2)_{1}$. Bottom row: $(0,2) /(1,1),(0,3) /(2,1),(1,2)_{\mathrm{r}} /(1,2)_{\mathrm{r}}$.
each $z \in \mathbb{R}^{2}$, we let $\pi_{z}^{+}$denote the right-most element of $\mathcal{W}(z)$. We define $\pi_{z}^{\uparrow}$ in the same way, except that at points of type $(1,2)_{1}$, we let $\pi_{z}^{\uparrow}$ be the left-most element of $\mathcal{W}(z)$. Note that as a consequence of this choice, whenever there are incoming paths at $z$, the path $\pi_{z}^{\uparrow}$ is the continuation of any incoming path at $z$.

The next proposition, which follows from [NRS10, Prop. 3.1], shows that it is possible to define something like the intersection local time of $\mathcal{W}$ and $\mathcal{W}$. Below, $|I|$ denotes the Lebesgue measure of a set $I \subset \mathbb{R}$.

Proposition 3.4 (Intersection local time) Let $(\mathcal{W}, \hat{\mathcal{W}})$ be the double Brownian web. Then a.s. there exists a unique measure $\ell$, concentrated on the set of points of type $(1,2)$ in $\mathcal{W}$, such that for each $\pi \in \mathcal{W}$ and $\hat{\pi} \in \hat{\mathcal{W}}$,

$$
\begin{align*}
& \ell\left(\left\{z=(x, t) \in \mathbb{R}^{2}: \sigma_{\pi}<t<\hat{\sigma}_{\hat{\pi}}, \pi(t)=x=\hat{\pi}(t)\right\}\right) \\
& \quad=\lim _{\varepsilon \downarrow 0} \varepsilon^{-1}\left|\left\{t \in \mathbb{R}: \sigma_{\pi}<t<\hat{\sigma}_{\hat{\pi}},|\pi(t)-\hat{\pi}(t)| \leq \varepsilon\right\}\right| \tag{3.9}
\end{align*}
$$

where the limit on the right-hand side exists and is finite. The measure $\ell$ is a.s. nonatomic and $\sigma$-finite. We let $\ell_{1}$ and $\ell_{\mathrm{r}}$ denote the restrictions of $\ell$ to the sets of points of type $(1,2)_{1}$ and $(1,2)_{\mathrm{r}}$, respectively.

We remark that $\ell(O)=\infty$ for every open nonempty subset $O \subset \mathbb{R}^{2}$, but $\ell$ is $\sigma$-finite. To see the latter, for any path $\pi \in \Pi$, let $\pi^{\circ}:=\left\{(\pi(t), t): t \in\left(\sigma_{\pi}, \infty\right)\right\}$ denote its interior, and define the interior $\hat{\pi}^{\circ}$ of a dual path $\hat{\pi} \in \hat{\Pi}$ analogously. Let $\mathcal{D} \subset \mathbb{R}^{2}$ be a deterministic countable dense set and for $z \in \mathcal{D}$, let $\pi_{z}$, resp. $\hat{\pi}_{z}$, denote the a.s. unique path in $\mathcal{W}$, resp. $\hat{\mathcal{W}}$, starting from $z$. Then by Proposition $3.4, \ell\left(\pi_{z}^{\circ} \cap \hat{\pi}_{\hat{z}}^{\circ}\right)<\infty$ for each $z, \hat{z} \in \mathcal{D}$, while by [SS08, Lemma 3.4 (b)], $\ell$ is concentrated on $\bigcup_{z, \hat{z} \in \mathcal{D}}\left(\pi_{z}^{\circ} \cap \hat{\pi}_{\tilde{z}}^{\circ}\right)$.

### 3.3 Sticky Brownian webs

We collect here some facts about a natural way to couple two Brownian webs. Such coupled Brownian webs will then be used in the next subsection to give a graphical construction of

Howitt-Warren flows. We first start with a 'reference' Brownian web $\mathcal{W}$, which is then used to construct a second, 'modified' or 'sample' Brownian web $\mathcal{W}$ ' by 'switching' a suitable Poisson subset of points of type $(1,2)_{1}$ of $\mathcal{W}$ into points of type $(1,2)_{\mathrm{r}}$, and vice versa, using a marking procedure developed in (NRS10.

To formulate this rigorously, let $z=(x, t)$ be a point of type $(1,2)_{1}$ in $\mathcal{W}$, and let

$$
\begin{equation*}
\mathcal{W}_{\text {in }}(z):=\left\{\pi \in \mathcal{W}: \sigma_{\pi}<t, \pi(t)=x\right\} \tag{3.10}
\end{equation*}
$$

denote the set of incoming paths in $\mathcal{W}$ at $z$. For any $\pi \in \mathcal{W}_{\text {in }}(z)$, let $\pi^{t}:=\left\{(\pi(s), s): \sigma_{\pi} \leq\right.$ $s \leq t\}$ denote the piece of $\pi$ leading up to $z$, and let $\mathcal{W}(z)=\{l, r\}$ be the outgoing paths in $\mathcal{W}$ at $z$, where $l<r$ on $(t, t+\varepsilon)$ for some $\varepsilon>0$. Since $z$ is of type $(1,2)_{1}$, identifying a path with its graph, we have $\pi=\pi^{t} \cup l$ for each $\pi \in \mathcal{W}_{\text {in }}(z)$. We define

$$
\begin{equation*}
\operatorname{switch}_{z}(\mathcal{W}):=\left(\mathcal{W} \backslash \mathcal{W}_{\text {in }}(z)\right) \cup\left\{\pi^{t} \cup r: \pi \in \mathcal{W}_{\text {in }}(z)\right\} \tag{3.11}
\end{equation*}
$$

Then $\operatorname{switch}_{z}(\mathcal{W})$ differs from $\mathcal{W}$ only in that $z$ is now of type $(1,2)_{\mathrm{r}}$ instead of $(1,2)_{1}$. In a similar way, if $z$ is of type $(1,2)_{\mathrm{r}}$ in $\mathcal{W}$, then we let $\operatorname{switch}_{z}(\mathcal{W})$ denote the web obtained from $\mathcal{W}$ by switching $z$ into a point of type $(1,2)_{1}$. If $z_{1}, \ldots, z_{n}$ are points of type $(1,2)$ in $\mathcal{W}$, then we let $\operatorname{switch}_{\left\{z_{1}, \ldots, z_{n}\right\}}(\mathcal{W}):=\operatorname{switch}_{z_{1}} \circ \cdots \circ \operatorname{switch}_{z_{n}}(\mathcal{W})$ denote the web obtained from $\mathcal{W}$ by switching the orientation of the points $z_{1}, \ldots, z_{n}$. Note that it does not matter in which order we perform the switching. Recall that a point $z$ is of type $(1,2)_{1}\left(\operatorname{resp} .(1,2)_{\mathrm{r}}\right)$ in the dual Brownian web $\hat{\mathcal{W}}$ if and only if it is of type $(1,2)_{1}\left(\operatorname{resp} .(1,2)_{\mathrm{r}}\right)$ in $\mathcal{W}$. We define switching in $\hat{\mathcal{W}}$ analogously to switching in $\mathcal{W}$.

The next theorem, which is similar to [NRS10, Prop. 6.1], shows how by switching the orientation of a countable Poisson set of points of type (1,2), we can obtain a well-defined modified Brownian web. Recall the definition of the intersection local time measure $\ell$ from Proposition 3.4 and note that since $\ell$ is $\sigma$-finite, the set $S$ below is a.s. a countable subset of the set of all points of type $(1,2)$.

Theorem 3.5 (Modified Brownian web) Let $\mathcal{W}$ be a Brownian web wih drift $\beta$, let $\ell$ be the intersection local time measure between $\mathcal{W}$ and its dual and let $\ell_{1}, \ell_{\mathrm{r}}$ denote the restrictions of $\ell$ to the sets of points of type $(1,2)_{1}$ and $(1,2)_{\mathrm{r}}$ in $\mathcal{W}$, respectively. Let $c_{1}, c_{r} \geq 0$ be constants and conditional on $\mathcal{W}$, let $S$ be a Poisson point set with intensity $c_{1} \ell_{1}+c_{\mathrm{r}} \ell_{\mathrm{r}}$. Then, for any sequence of finite sets $\Delta_{n} \uparrow S$, the limit

$$
\begin{equation*}
\left(\mathcal{W}^{\prime}, \hat{\mathcal{W}}^{\prime}\right):=\lim _{\Delta_{n} \uparrow S}\left(\operatorname{switch}_{\Delta_{n}}(\mathcal{W}), \operatorname{switch}_{\Delta_{n}}(\hat{\mathcal{W}})\right) \tag{3.12}
\end{equation*}
$$

exists in $\mathcal{K}(\Pi) \times \mathcal{K}(\hat{\Pi})$ a.s. and does not depend on the choice of the sequence $\Delta_{n} \uparrow S$. Moreover, $\mathcal{W}^{\prime}$ is a Brownian web with drift $\beta^{\prime}=\beta+c_{1}-c_{\mathrm{r}}$ and $\hat{\mathcal{W}}^{\prime}$ is its dual.

If $\left(\mathcal{W}, \mathcal{W}^{\prime}\right)$ are coupled as in Theorem 3.5, then we say that $\left(\mathcal{W}, \mathcal{W}^{\prime}\right)$ is a pair of sticky Brownian webs with drifts $\beta, \beta^{\prime}$ and coupling parameter $\kappa:=\min \left\{c_{1}, c_{\mathrm{r}}\right\}$. In the special case that $\kappa=0$ and $\beta \leq \beta^{\prime}$, we call $\left(\mathcal{W}, \mathcal{W}^{\prime}\right)$ a left-right Brownian web with drifts $\beta, \beta^{\prime}$. Leftright Brownian webs have been introduced with the help of a 'left-right stochastic differential equation' (instead of the marking construction above) in [SS08]. Pairs of sticky Brownian webs with general coupling parameters $\kappa \geq 0$ have been introduced by means of a martingale problem in HW09b, Section 7]. We will prove in Lemma 6.18 below that the constructions of left-right Brownian webs given above and in [SS08] are equivalent. We will not make use of the martingale formulation of sticky Brownian webs developed in HW09b.

For any point $z$ of type $(1,2)$ in some Brownian web $\mathcal{W}$, we call

$$
\operatorname{sign}_{\mathcal{W}}(z):= \begin{cases}-1 & \text { if } z \text { is of type }(1,2)_{1} \text { in } \mathcal{W}  \tag{3.13}\\ +1 & \text { if } z \text { is of type }(1,2)_{\mathrm{r}} \text { in } \mathcal{W}\end{cases}
$$

the $\operatorname{sign}$ of $z$ in $\mathcal{W}$. If $\left(\mathcal{W}, \mathcal{W}^{\prime}\right)$ is a pair of sticky Brownian webs, then it is known SSS09, Thm. 1.7] that the set of points of type $(1,2)$ in $\mathcal{W}$ in general does not coincide with the set of points of type $(1,2)$ in $\mathcal{W}^{\prime}$. The next proposition says that nevertheless, in the sense of intersection local time measure, almost all points of type $(1,2)$ in $\mathcal{W}$ are also of type $(1,2)$ in $\mathcal{W}^{\prime}$, and these point have the orientation one expects.

Proposition 3.6 (Change of reference web) In the setup of Theorem 3.5, let $\ell^{\prime}$ be the intersection local time measure between $\mathcal{W}^{\prime}$ and its dual and let $\ell_{1}^{\prime}, \ell_{\mathrm{r}}^{\prime}$ denote the restrictions of $\ell^{\prime}$ to the sets of points of type $(1,2)_{1}$ and $(1,2)_{\mathrm{r}}$ in $\mathcal{W}^{\prime}$, respectively. Then:
(i) Almost surely, $\ell_{1}^{\prime}=\ell_{1}$ and $\ell_{\mathrm{r}}^{\prime}=\ell_{\mathrm{r}}$.
(ii) $S=\left\{z \in \mathbb{R}^{2}: z\right.$ is of type $(1,2)$ in both $\mathcal{W}$ and $\mathcal{W}^{\prime}$, and $\left.\operatorname{sign}_{\mathcal{W}}(z) \neq \operatorname{sign}_{\mathcal{W}^{\prime}}(z)\right\}$ a.s.
(iii) Conditional on $\mathcal{W}^{\prime}$, the set $S$ is a Poisson point set with intensity $c_{\mathrm{r}} \ell_{1}^{\prime}+c_{1} \ell_{\mathrm{r}}^{\prime}$ and $\mathcal{W}=$ $\lim _{\Delta_{n} \uparrow S} \operatorname{switch}_{\Delta_{n}}\left(\mathcal{W}^{\prime}\right)$.

Let $\left(\mathcal{W}_{0}, \mathcal{W}\right)$ be a pair of sticky Brownian webs with drifts $\beta_{0}, \beta$ and coupling parameter $\kappa$, and let $\pi_{z}^{\uparrow}$ and $\pi_{z}^{+}$denote the special paths in $\mathcal{W}(z)$ defined below Proposition 3.3. Let

$$
\begin{equation*}
K_{s, t}^{\uparrow}(x, A):=\mathbb{P}\left[\pi_{(x, s)}^{\uparrow}(t) \in A \mid \mathcal{W}_{0}\right] \quad(s \leq t, x \in \mathbb{R}, A \in \mathcal{B}(\mathbb{R})) \tag{3.14}
\end{equation*}
$$

and let $K_{s, t}^{+}(x, A)$ be defined similarly, with $\pi_{z}^{\uparrow}$ replaced by $\pi_{z}^{+}$. Then, as we will see in Theorem 3.7 below, $\left(K_{s, t}^{\uparrow}\right)_{s \leq t}$ and $\left(K_{s, t}^{+}\right)_{s \leq t}$ are versions of the Howitt-Warren flow with drift $\beta$ and characteristic measure $\nu=c_{1} \delta_{1}+c_{\mathrm{r}} \delta_{0}$, where

$$
\begin{equation*}
c_{1}:=\kappa+\max \left\{0, \beta^{\prime}-\beta\right\} \quad \text { and } \quad c_{\mathrm{r}}:=\kappa+\max \left\{0, \beta-\beta^{\prime}\right\} \tag{3.15}
\end{equation*}
$$

In the special case that $c_{1}=c_{\mathrm{r}}$, this was proved in HW09b, formula (5)]. In the next subsection, we set out to give a similar construction for any Howitt-Warren flow.

### 3.4 Marking construction of Howitt-Warren flows

We now give a construction of a general Howitt-Warren flow based on two coupled webs, which is the central result of this paper. More precisely, the random environment of the HowittWarren flow will be represented by a reference Brownian web $\mathcal{W}_{0}$ plus a set $\mathcal{M}$ of marked points of type $(1,2)$. Conditional on $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, we will modify $\mathcal{W}_{0}$ in a similar way as in Theorem 3.5 to construct a sample Brownian web $\mathcal{W}$, whose law conditional on $\left(\mathcal{W}_{0}, \mathcal{M}\right)$ then defines the Howitt-Warren flow via the continuous analogue of (3.5). For erosion flows, the set $\mathcal{M}$ of marked points is empty, hence our representation reduces to $(3.14)$.

Let $\mathcal{W}_{0}$ be a Brownian web with drift $\beta_{0}$ and let $\nu_{1}$ and $\nu_{\mathrm{r}}$ be finite measures on $[0,1]$. Let $\ell, \ell_{1}$ and $\ell_{\mathrm{r}}$ be defined for $\mathcal{W}_{0}$ as in Proposition 3.4, and conditional on $\mathcal{W}_{0}$, let $\mathcal{M}$ be a Poisson point set on $\mathbb{R}^{2} \times[0,1]$ with intensity

$$
\begin{equation*}
\ell_{1}(\mathrm{~d} z) \otimes 21_{\{0<q\}} q^{-1} \nu_{1}(\mathrm{~d} q)+\ell_{\mathrm{r}}(\mathrm{~d} z) \otimes 21_{\{q<1\}}(1-q)^{-1} \nu_{\mathrm{r}}(\mathrm{~d} q) \tag{3.16}
\end{equation*}
$$

Elements of $\mathcal{M}$ are pairs $(z, q)$ where $z$ is a point of type $(1,2)$ in $\mathcal{W}_{0}$ and $q \in[0,1]$. Since $\ell$ is nonatomic, for each point $z$ of type $(1,2)$ there is at most one $q$ such that $(z, q) \in \mathcal{M}$, and we may write $\mathcal{M}=\left\{\left(z, \omega_{z}\right): z \in M\right\}$. We call points $z \in M$ marked points and we call $\omega_{z}$ the mark of $z$.

Conditional on the reference web $\mathcal{W}_{0}$ and the set of marked points $\mathcal{M}$, we construct independent $\{-1,+1\}$-valued random variables $\left(\alpha_{z}\right)_{z \in M}$ with $\mathbb{P}\left[\alpha_{z}=+1 \mid\left(\mathcal{\mathcal { W } _ { 0 }}, \mathcal{M}\right)\right]=\omega_{z}$, and we set

$$
\begin{equation*}
A:=\left\{z \in M: \alpha_{z} \neq \operatorname{sign}_{\mathcal{W}_{0}}(z)\right\} \tag{3.17}
\end{equation*}
$$

In addition, conditional on $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, we let $B$ be a Poisson point set with intensity $2 \nu_{1}(\{0\}) \ell_{1}+$ $2 \nu_{\mathrm{r}}(\{1\}) \ell_{\mathrm{r}}$, independent of $A$. We observe that conditional on $\mathcal{W}_{0}$, but integrating out the randomness of $\mathcal{M}$, the set $A \cup B$ is a Poisson point set with intensity

$$
\begin{align*}
& 2\left(\int_{(0,1]} q q^{-1} \nu_{1}(\mathrm{~d} q)+\nu_{1}(\{0\})\right) \ell_{1}+2\left(\int_{[0,1)}(1-q)(1-q)^{-1} \nu_{\mathrm{r}}(\mathrm{~d} q)+\nu_{\mathrm{r}}(\{1\})\right) \ell_{\mathrm{r}}  \tag{3.18}\\
& \quad \quad=2 \nu_{1}([0,1]) \ell_{1}+2 \nu_{\mathrm{r}}([0,1]) \ell_{\mathrm{r}}
\end{align*}
$$

Therefore, by Theorem 3.5 the limit

$$
\begin{equation*}
\mathcal{W}:=\lim _{\Delta_{n} \uparrow A \cup B} \operatorname{switch}_{\Delta_{n}}\left(\mathcal{W}_{0}\right) \tag{3.19}
\end{equation*}
$$

exists in $\mathcal{K}(\Pi)$ a.s. and is a Brownian web with drift

$$
\begin{equation*}
\beta:=\beta_{0}+2 \nu_{1}([0,1])-2 \nu_{\mathrm{r}}([0,1]) \tag{3.20}
\end{equation*}
$$

Theorem 3.7 (Construction of Howitt-Warren flows) Let $\beta \in \mathbb{R}$ and let $\nu$ be a finite measure on $[0,1]$. For any finite measures $\nu_{1}$ and $\nu_{\mathrm{r}}$ on $[0,1]$ satisfying

$$
\begin{equation*}
\nu(\mathrm{d} q):=(1-q) \nu_{1}(\mathrm{~d} q)+q \nu_{\mathrm{r}}(\mathrm{~d} q), \tag{3.21}
\end{equation*}
$$

let $\beta_{0}$ be determined from $\beta, \nu_{1}$ and $\nu_{\mathrm{r}}$ as in (3.20). Let $\mathcal{W}_{0}$ be a reference Brownian web with drift $\beta_{0}$ and define a set of marked points $\mathcal{M}$ and sample Brownian web $\mathcal{W}$ as in (3.16) and (3.19). Let $\pi_{z}^{\uparrow}$ and $\pi_{z}^{+}$denote the special paths in $\mathcal{W}(z)$ defined below Proposition 3.3. Set

$$
\begin{equation*}
K_{s, t}^{\uparrow}(x, A):=\mathbb{P}\left[\pi_{(x, s)}^{\uparrow}(t) \in A \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right] \quad(s \leq t, x \in \mathbb{R}, A \in \mathcal{B}(\mathbb{R})) \tag{3.22}
\end{equation*}
$$

and define $K_{s, t}^{+}(x, A)$ similarly with $\pi_{(x, s)}^{\uparrow}$ replaced by $\pi_{(x, s)}^{+}$. Then $\left(K_{s, t}^{\uparrow}\right)_{s \leq t}$ and $\left(K_{s, t}^{+}\right)_{s \leq t}$ are versions of the Howitt-Warren flow with drift $\beta$ and characteristic measure $\nu$. In the special case that $\nu_{1}=\nu_{\mathrm{r}}$, the triple $\left(\mathcal{W}_{0}, \mathcal{M}, \mathcal{W}\right)$ is equally distributed with $\left(\mathcal{W}, \mathcal{M}, \mathcal{W}_{0}\right)$.

Remark. In the special case that $\nu_{1}=\nu_{\mathrm{r}}$, the construction of the reference Brownian web and set of marked points arises as the scaling limit of the discrete construction outlined in Section 1.3, i.e., if $\omega^{\langle k\rangle}=\left(\omega_{z}^{\langle k\rangle}\right)_{z \in \mathbb{Z}_{\text {aven }}^{2}}$ is a collection of independent $[0,1]$-valued random variables with laws $\mu_{k}$ satisfying (1.7) and conditional on $\omega^{\langle k\rangle}$ we construct an independent collection $\alpha^{\langle k\rangle}=\left(\alpha_{z}^{\langle k\rangle}\right)_{z \in \mathbb{Z}_{\text {even }}^{2}}$ of $\{-1,+1\}$-valued random variables with $\mathbb{P}\left[\alpha_{z}^{\langle k\rangle}=+1 \mid \omega^{\langle k\rangle}\right]=$ $\omega_{z}^{\langle k\rangle}$, then the discrete reference web corresponding to $\alpha^{\langle k\rangle}$ converges after diffusive rescaling to $\mathcal{W}_{0}$ and $\left\{\left(z, \omega_{z}\right): 0<\omega_{z}<1\right\}$ converges to the set of marked points $\mathcal{M}$. The more general construction in Theorem 3.7 where possibly $\nu_{1} \neq \nu_{\mathrm{r}}$ and $\beta_{0}$ is possibly different from
$\beta$ arises as the diffusive scaling limit of discrete contructions where we first choose a reference collection of random variables $\alpha^{\langle k\rangle}$ with $\varepsilon_{k}^{-1} \mathbb{E}\left[\alpha_{z}^{\langle k\rangle}\right] \rightarrow \beta_{0}$ and then conditional on $\alpha^{\langle k\rangle}$, we choose independent $\left(\omega_{z}^{\langle k\rangle}\right)_{z \in \mathbb{Z}_{\text {even }}^{2}}$ with $\mathbb{P}\left[\omega_{z} \in \mathrm{~d} q \mid \alpha^{\langle k\rangle}=-1\right]=\mu_{k}^{1}(\mathrm{~d} q)$ and $\mathbb{P}\left[\omega_{z} \in \mathrm{~d} q \mid \alpha^{\langle k\rangle}=\right.$ $+1]=\mu_{k}^{\mathrm{r}}(\mathrm{d} q)$, where generalizing (1.9), $\mu_{k}^{\mathrm{l}}$ and $\mu_{k}^{\mathrm{r}}$ are any laws such that

$$
\begin{equation*}
\mathbb{P}\left[\alpha^{\langle k\rangle}=-1\right] \mu_{k}^{1}(\mathrm{~d} q)+\mathbb{P}\left[\alpha^{\langle k\rangle}=+1\right] \mu_{k}^{\mathrm{r}}(\mathrm{~d} q)=\mu_{k}(\mathrm{~d} q) . \tag{3.23}
\end{equation*}
$$

This more general construction will sometimes be handy. For example, for erosion flows where $\nu=c_{0} \delta_{0}+c_{1} \delta_{1}$ for some $c_{0}, c_{1} \geq 0$, it is most natural to choose $\nu_{1}=c_{0} \delta_{0}$ and $\nu_{\mathrm{r}}=c_{1} \delta_{1}$. Also, for Howitt-Warren flows where one or both of the speeds $\beta_{-}, \beta_{+}$defined in (2.12) are finite, it is sometimes handy to choose either $\nu_{1}=0$ or $\nu_{\mathrm{r}}=0$.

If $\left(\mathcal{W}_{0}, \mathcal{M}, \mathcal{W}\right)$ are a reference Brownian web, the set of marked points, and the sample Brownian web as defined above Theorem 3.7 , then we call the random probability measure $\mathbb{Q}$ on $\mathcal{K}(\Pi)$ defined by

$$
\begin{equation*}
\mathbb{Q}:=\mathbb{P}\left[\mathcal{W} \in \cdot \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right] \tag{3.24}
\end{equation*}
$$

the Howitt-Warren quenched law with drift $\beta$ and characteristic measure $\nu$. In Section 3.5 below, we will show that in some precisely defined way, these Howitt-Warren quenched laws are the diffusive scaling limits of the discrete quenched laws defined in Section 3.1.

Since at deterministic points in the Brownian web $\mathcal{W}$ there is a.s. only one outgoing path, the stochastic flows of kernels $\left(K_{s, t}^{\dagger}\right)_{s \leq t}$ and $\left(K_{s, t}^{+}\right)_{s \leq t}$ from Theorem 3.7 obviously have the same finite-dimensional distributions. They are, however, not the same. Each version has its own pleasant properties.

Proposition 3.8 (Regular parameter dependence) Let $\left(K_{s, t}^{\uparrow}\right)_{s \leq t}$ and $\left(K_{s, t}^{+}\right)_{s \leq t}$ be defined as in Theorem 3.7. Then, of the following properties, $\left(K_{s, t}^{+}\right)_{s \leq t}$ satisfies (a)-(c) and $\left(K_{s, t}^{\uparrow}\right)_{s \leq t}$ satisfies (a), (b) and (d).
(a) Setting $\mathbb{R}_{\leq}^{2}:=\left\{(s, t) \in \mathbb{R}^{2}: s \leq t\right\}$, the map $(s, t, x, \omega) \mapsto K_{s, t}(x, \cdot)(\omega)$ is a measurable map from $\mathbb{R}_{\leq}^{2} \times \mathbb{R} \times \Omega$ to $\mathcal{M}_{1}(\mathbb{R})$.
(b) A.s., the map $t \mapsto K_{s, t}(x, \cdot)$ from $[s, \infty)$ to $\mathcal{M}_{1}(\mathbb{R})$ is continuous for all $s \in \mathbb{R}$ and $x \in \mathbb{R}$.
(c) A.s., $x \mapsto K_{s, t}(x, A)$ is a càdlàg function from $\mathbb{R}$ to $\mathbb{R}$ for each $s<t$ and $A \in \mathcal{B}(\mathbb{R})$.
(d) A.s., $\int_{\mathbb{R}} K_{s, t}(x, \mathrm{~d} y) K_{t, u}(y, A)=K_{s, u}(x, A)$ for all $s \leq t \leq u, x \in \mathbb{R}$, and $A \in \mathcal{B}(\mathbb{R})$.

Proposition 3.8 (b) and (d) show that $\left(K_{s, t}^{\uparrow}\right)_{s \leq t}$ yields a version of a Howitt-Warren flow with the properties listed in Proposition 2.3. In particular, Proposition 3.8 (d) makes it a family of bona fide transition probability kernels of a Markov process in a random space-time environment.

### 3.5 Discrete approximation

Recall the definitions of the discrete quenched laws $\mathbf{Q}$ in (3.3) and the Howitt-Warren quenched laws $\mathbb{Q}$ in (3.24). In this subsection, we formulate a convergence result which says that if $\mu_{k}$ is a sequence of probability laws on $[0,1]$ satisfying (1.7), then the associated discrete quenched
laws $\mathbf{Q}_{\langle k\rangle}$, diffusively rescaled, converge to the Howitt-Warren quenched law $\mathbb{Q}$ with drift $\beta$ and characteristic measure $\nu$. This abstract result then implies other convergence results such as the convergence of Howitt-Warren flows, Howitt-Warren processes, and $n$-point motions. Since the $n$-point motions of discrete Howitt-Warren flows will be shown in Proposition A. 5 to converge to solutions of the Howitt-Warren martingale problems, this will also verify that the flows we constructed in Theorem 3.7 are indeed versions of the Howitt-Warren flow.

To formulate our convergence statement properly, we need to identify a discrete quenched law $\mathbf{Q}$ with a random probability law on the space $\mathcal{K}(\Pi)$ of compact subsets of the space $\Pi$ of paths defined in Section 3.2. Recall the definition of the paths $p_{z}^{\alpha}$ in (3.1) and the discrete webs $\mathcal{U}^{\alpha}$ in 3.2 . We wish to view $\mathcal{U}^{\alpha}$ as a random variable with values in $\mathcal{K}(\Pi)$. To this aim, we modify our definition of $\mathcal{U}^{\alpha}$ as follows. First, for each $z=(x, s) \in \mathbb{Z}_{\text {even }}^{2}$ we make $p_{z}^{\alpha}(t)$ into a path in $\Pi$ by linear interpolation between integer times and by setting $p_{z}^{\alpha}(\infty):=*$. Next, we add to $\mathcal{U}^{\alpha}$ all trivial paths $\pi$, with starting times $\sigma_{\pi} \in \mathbb{Z} \cup\{-\infty, \infty\}$, such that $\pi$ is identically $-\infty$ or $+\infty$ on $\left[\sigma_{\pi}, \infty\right) \cap \mathbb{R}$. With this modified definition, it can be checked that $\mathcal{U}^{\alpha}$ is indeed a random compact subset of $\Pi$, as desired.

For $\varepsilon>0$, we let $S_{\varepsilon}: R_{\mathrm{c}}^{2} \rightarrow R_{\mathrm{c}}^{2}$ denote the scaling map

$$
\begin{equation*}
S_{\varepsilon}(x, t):=\left(\varepsilon x, \varepsilon^{2} t\right) \quad\left((x, t) \in R_{\mathrm{c}}^{2}\right) \tag{3.25}
\end{equation*}
$$

As usual, we identify paths with their graphs; then $S_{\varepsilon}(\pi)$ is the path obtained by diffusively rescaling a path $\pi$ with $\varepsilon$, and $S_{\varepsilon}\left(\mathcal{U}^{\alpha}\right)$ is the random collection of paths obtained by diffusively rescaling paths in $\mathcal{U}^{\alpha}$. If $\mathbf{Q}^{\omega}$ is a discrete quenched law as defined in $(3.3)$ and $\varepsilon>0$, then we write

$$
\begin{equation*}
S_{\varepsilon}\left(\mathbf{Q}^{\omega}\right):=\mathbf{Q}^{\omega}\left[S_{\varepsilon}\left(\mathcal{U}^{\alpha}\right) \in \cdot\right] \tag{3.26}
\end{equation*}
$$

i.e., $S_{\varepsilon}\left(\mathbf{Q}^{\omega}\right)$ is the image under the scaling map $S_{\varepsilon}$ of the quenched law of $\mathcal{U}^{\alpha}$. Note that $S_{\varepsilon}\left(\mathbf{Q}^{\omega}\right)$, so defined, is a random probability law on the space of compact subsets of $\Pi$.

Theorem 3.9 (Convergence of quenched laws) Let $\varepsilon_{k}$ be positive constants, converging to zero and $\mu_{k}$ be probability measures on $[0,1]$ satisfying (1.7) for some real $\beta$ and finite measure $\nu$ on $[0,1]$. Let $\omega^{\langle k\rangle}=\left(\omega_{z}^{\langle k\rangle}\right)_{z \in \mathbb{Z}_{\text {even }}^{2}}$ be i.i.d. $[0,1]$-valued random variables with distibution $\mu_{k}$, let $\mathbf{Q}_{\langle k\rangle}:=\mathbf{Q}^{\omega^{\langle k\rangle}}$ be the discrete quenched law defined in (3.3), and let $\mathbb{Q}$ be the Howitt-Warren quenched law with drift $\beta$ and characteristic measure $\nu$ defined in (3.24). Then

$$
\begin{equation*}
\mathbb{P}\left[S_{\varepsilon_{k}}\left(\mathbf{Q}_{\langle k\rangle}\right) \in \cdot\right] \underset{k \rightarrow \infty}{\Longrightarrow} \mathbb{P}[\mathbb{Q} \in \cdot] \tag{3.27}
\end{equation*}
$$

where $\Rightarrow$ denotes weak convergence of probability laws on $\mathcal{M}_{1}(\mathcal{K}(\Pi))$.

## 4 Construction of Howitt-Warren flows in the Brownian net

In this section, we show that when a Howitt-Warren flow with drift $\beta$ and characteristic measure $\nu$ has finite left and right speeds, or equivalently, $b:=\int q^{-1}(1-q)^{-1} \nu(\mathrm{~d} q)<\infty$, then we can alternatively construct the flow as a random flow of mass in the Brownian net. Analogous to Theorem 3.7, the random environment will now be represented as a Brownian net $\mathcal{N}$ plus a set of i.i.d. marks $\bar{\omega}:=\left(\bar{\omega}_{z}\right)_{z \in S}$ attached to the separation points $S$ of $\mathcal{N}$, each with law $\bar{\nu}(\mathrm{d} q):=b^{-1} q^{-1}(1-q)^{-1} \nu(\mathrm{~d} q)$. Conditional on $(\mathcal{N}, \bar{\omega})$, we can construct the sample web $\mathcal{W}$ by choosing trajectories in $\mathcal{N}$ that turn in the 'right' way at separation points. The Howitt-Warren flow is then defined from the law of $\mathcal{W}$ conditional on $(\mathcal{N}, \bar{\omega})$ as in (3.22).

In Sections 4.1 4.2, we recall the necessary background on the Brownian net and properties of its separation points. In Section 4.3, we first state some coupling results between the Brownian web and Brownian net, which will help shed more light on the marking constructions of sticky Brownian webs in Theorem 3.5. In Section 4.4 we then give our main result on the alternative construction of Howitt-Warren flows with finite left and right speeds using the Brownian net. Lastly in Section 4.5, we formulate what we call Brownian half-nets, and state some properties for the support of the Howitt-Warren quenched law defined in (3.24), which will imply Theorems 2.5 and 2.7 . We note that, apart from being used to construct HowittWarren flows with finite left and right speeds, the theory of the Brownian net will also play an important role for Howitt-Warren flows with infinite left or right speed, such as in the proof of results in Section 4.5, as well as in the proof of Theorems 3.5 and 3.9 .

### 4.1 The Brownian net

The Brownian net arises as the diffusive scaling limit of branching-coalescing random walks in the limit of small branching probability. It was first introduced by Sun and Swart in [S08] and independently by Newman, Ravishankar and Schertzer in [NRS10]. A further study of its properties was carried out in [SSS09]. We now recall the definition of the Brownian net given in SS08.

Recall that in Section 3.3, we defined a left-right Brownian web to be a pair of sticky Brownian webs $\left(\mathcal{W}^{l}, \mathcal{W}^{\mathrm{r}}\right)$ with drifts $\beta_{-} \leq \beta_{+}$and coupling parameter $\kappa=0$. At present, we will need the original definition of a left-right Brownian web given in [SS08]. In Lemma 6.18 below, we will prove that both definitions are equivalent.

Following [SS08], we call $\left(l_{1}, \ldots, l_{m} ; r_{1}, \ldots, r_{n}\right)$ a collection of left-right coalescing Brownian motions with drifts $\beta_{-} \leq \beta_{+}$, if $\left(l_{1}, \ldots, l_{m}\right)$ and $\left(r_{1}, \ldots, r_{n}\right)$ are distributed as collections of coalescing Brownian motions with drift $\beta_{-}$and $\beta_{+}$, respectively, if paths in $\left(l_{1}, \ldots, l_{m} ; r_{1}, \ldots, r_{n}\right)$ evolve independently when they are apart, and the interaction between $l_{i}$ and $r_{j}$ when they meet is described by the two-dimensional stochastic differential equation

$$
\begin{align*}
\mathrm{d} L_{t} & =1_{\left\{L_{t} \neq R_{t}\right\}} \mathrm{d} B_{t}^{\mathrm{l}}+1_{\left\{L_{t}=R_{t}\right\}} \mathrm{d} B_{t}^{\mathrm{s}}+\beta_{-} \mathrm{d} t \\
\mathrm{~d} R_{t} & =1_{\left\{L_{t} \neq R_{t}\right\}} \mathrm{d} B_{t}^{\mathrm{r}}+1_{\left\{L_{t}=R_{t}\right\}} \mathrm{d} B_{t}^{\mathrm{s}}+\beta_{+} \mathrm{d} t \tag{4.1}
\end{align*}
$$

where $B_{t}^{\mathrm{l}}, B_{t}^{\mathrm{r}}, B_{t}^{\mathrm{s}}$ are independent standard Brownian motions, and $(L, R)$ are subject to the constraint that

$$
\begin{equation*}
L_{t} \leq R_{t} \text { for all } t \geq \inf \left\{s: L_{s}=R_{s}\right\} \tag{4.2}
\end{equation*}
$$

It can be shown that subject to the condition 4.2 , solutions to the $\mathrm{SDE}(4.1)$ are unique in distribution [SS08, Proposition 2.1].

Let $\mathcal{W}^{l}, \mathcal{W}^{\mathrm{r}}$ be two Brownian webs with drifts $\beta_{-} \leq \beta_{+}$, and for determinstic $z \in \mathbb{R}^{2}$, let $l_{z}$ resp. $r_{z}$ denote the a.s. unique path in $\mathcal{W}^{l}$ resp. $\mathcal{W}^{r}$ starting from $z$. Following [SS08], we say that $\left(\mathcal{W}^{l}, \mathcal{W}^{\mathrm{r}}\right)$ is a left-right Brownian web if for any finite deterministic set of points $z_{1}, \ldots, z_{m}, z_{1}^{\prime}, \ldots, z_{n}^{\prime} \in \mathbb{R}^{2}$, the collection $\left(l_{z_{1}}, \ldots, l_{z_{m}} ; r_{z_{1}^{\prime}}, \ldots, r_{z_{n}^{\prime}}\right)$ is distributed as left-right coalescing Brownian motions. Elements of $\mathcal{W}^{1}$ (resp. $\mathcal{W}^{\mathrm{r}}$ ) are called left-most (resp. right-most) paths.

It was shown in SS08 that each left-right Brownian web a.s. determines an associated Brownian net and vice versa. There, three different ways were given to construct a Brownian net from its associated left-right Brownian web, which are known as the hopping construction and the constructions using wedges and meshes, which we recall now.


Figure 9: A mesh $M(r, l)$ with bottom point $z$ and a wedge $W(\hat{r}, \hat{l})$ with bottom point $z$.
Hopping: We call $t \in \mathbb{R}$ an intersection time of two paths $\pi, \pi^{\prime} \in \Pi$ if $\sigma_{\pi}, \sigma_{\pi^{\prime}}<t<\infty$ and $\pi(t)=\pi^{\prime}(t)$. If $t$ is an intersection time of $\pi$ and $\pi^{\prime}$, then we can define a new path $\pi^{\prime \prime}$ by concatenating the piece of $\pi$ before $t$ with the piece of $\pi^{\prime}$ after $t$, i.e., by setting $\pi^{\prime \prime}:=\left\{(\pi(s), s): s \in\left[\sigma_{\pi}, t\right]\right\} \cup\left\{\left(\pi^{\prime}(s), s\right): s \in[t, \infty]\right\}$. For any collection of paths $\mathcal{A} \subset \Pi$, we let $\mathcal{H}_{\text {int }}(\mathcal{A})$ denote the smallest set of paths containing $\mathcal{A}$ that is closed under such 'hopping' from one path onto another at intersection times, i.e., $\mathcal{H}_{\text {int }}(\mathcal{A})$ is the set of all paths $\pi \in \Pi$ of the form

$$
\begin{equation*}
\pi=\bigcup_{k=1}^{m}\left\{\left(\pi_{k}(s), s\right): s \in\left[t_{k-1}, t_{k}\right]\right\}, \tag{4.3}
\end{equation*}
$$

where $\pi_{1}, \ldots, \pi_{m} \in \mathcal{A}, \sigma_{\pi_{1}}=t_{0}<\cdots<t_{m}=\infty$, and $t_{k}$ is an intersection time of $\pi_{k}$ and $\pi_{k+1}$ for each $k=1, \ldots, m-1$.
Wedges: Let $\hat{\mathcal{W}}^{1}, \hat{\mathcal{W}}^{\mathrm{r}}$ be the dual Brownian webs associated with a left-right Brownian web $\left(\mathcal{W}^{\mathrm{l}}, \mathcal{W}^{\mathrm{r}}\right)$. Any pair $\hat{l} \in \hat{\mathcal{W}}^{\mathrm{l}}, \hat{r} \in \hat{\mathcal{W}}^{\mathrm{r}}$ with $\hat{r}\left(\hat{\sigma}_{\hat{l}} \wedge \hat{\sigma}_{\hat{r}}\right)<\hat{l}\left(\hat{\sigma}_{\hat{l}} \wedge \hat{\sigma}_{\hat{r}}\right)$ defines an open set (see Figure 9)

$$
\begin{equation*}
W(\hat{r}, \hat{l}):=\left\{(x, u) \in \mathbb{R}^{2}: \hat{\tau}_{\hat{r}, \hat{l}}<u<\hat{\sigma}_{\hat{l}} \wedge \hat{\sigma}_{\hat{r}}, \hat{r}(u)<x<\hat{l}(u)\right\}, \tag{4.4}
\end{equation*}
$$

where $\tau_{\hat{r}, \hat{l}}:=\sup \left\{t<\hat{\sigma}_{\hat{l}} \wedge \hat{\sigma}_{\hat{r}}: \hat{r}(t)=\hat{l}(t)\right\}$ is the first (backward) hitting time of $\hat{r}$ and $\hat{l}$, which might be $-\infty$. Such an open set is called a wedge of $\left(\hat{\mathcal{W}}^{1}, \hat{\mathcal{W}}^{\mathrm{r}}\right)$. If $\tau_{\hat{r}, \hat{l}}>-\infty$, then we call $\tau_{\hat{r}, \hat{l}}$ the bottom time, and $\left(\hat{l}\left(\tau_{\hat{r}, \hat{l}}\right), \tau_{\hat{r}, \hat{l}}\right)$ the bottom point of the wedge $W(\hat{r}, \hat{l})$.
Meshes: By definition, a mesh of $\left(\mathcal{W}^{\mathrm{l}}, \mathcal{W}^{\mathrm{r}}\right)$ (see Figure 9) is an open set of the form

$$
\begin{equation*}
M=M(r, l)=\left\{(x, t) \in \mathbb{R}^{2}: \sigma_{l}<t<\tau_{l, r}, r(t)<x<l(t)\right\}, \tag{4.5}
\end{equation*}
$$

where $l \in \mathcal{W}^{\mathrm{l}}, r \in \mathcal{W}^{\mathrm{r}}$ are paths such that $\sigma_{l}=\sigma_{r}, l\left(\sigma_{l}\right)=r\left(\sigma_{r}\right)$ and $r(s)<l(s)$ on $\left(\sigma_{l}, \sigma_{l}+\varepsilon\right)$ for some $\varepsilon>0$, and $\tau_{l, r}:=\inf \left\{t>\sigma_{l}: l(t)=r(t)\right\}$. We call $\left(l\left(\sigma_{l}\right), \sigma_{l}\right)$ the bottom point, $\sigma_{l}$ the bottom time, $\left(l\left(\tau_{l, r}\right), \tau_{l, r}\right)$ the top point, $\tau_{l, r}$ the top time, $r$ the left boundary, and $l$ the right boundary of $M$.

Given an open set $A \subset \mathbb{R}^{2}$ and a path $\pi \in \Pi$, we say $\pi$ enters $A$ if there exist $\sigma_{\pi}<s<t$ such that $\pi(s) \notin A$ and $\pi(t) \in A$. We say $\pi$ enters $A$ from outside if there exists $\sigma_{\pi}<s<t$ such that $\pi(s) \notin \bar{A}$ and $\pi(t) \in A$. We now recall the following characterization of the Brownian net from [SS08, Theorems 1.3, 1.7 and 1.10]. Below, $\overline{\mathcal{A}}$ denotes the closure of a set of paths $\mathcal{A} \subset \Pi$ in the topology on $\Pi$.

Theorem 4.1 (Brownian net associated with a left-right Brownian web) Let $\left(\mathcal{W}^{l}, \mathcal{W}^{\mathrm{r}}\right)$ be a left-right Brownian web with drifts $\beta_{-} \leq \beta_{+}$and let $\hat{\mathcal{W}}^{\mathrm{l}}, \hat{\mathcal{W}}^{\mathrm{r}}$ be the dual Brownian webs associated with $\mathcal{W}^{1}, \mathcal{W}^{\mathrm{r}}$. Then there exists a random compact set of paths $\mathcal{N} \in \mathcal{K}(\Pi)$, called the Brownian net, that is a.s. uniquely determined by any of the following equivalent conditions:
(i) $\mathcal{N}=\overline{\mathcal{H}_{\text {int }}\left(\mathcal{W}^{\mathrm{l}} \cup \mathcal{W}^{\mathrm{r}}\right)}$ a.s.
(ii) $\mathcal{N}=\left\{\pi \in \Pi: \pi\right.$ does not enter any wedge of $\left(\hat{\mathcal{W}}^{1}, \hat{\mathcal{W}}^{\mathrm{r}}\right)$ from outside $\}$ a.s.
(iii) $\mathcal{N}=\left\{\pi \in \Pi: \pi\right.$ does not enter any mesh of $\left.\left(\mathcal{W}^{l}, \mathcal{W}^{\mathrm{r}}\right)\right\}$ a.s.

The set $\mathcal{N}$ is closed under hopping, i.e., $\mathcal{N}=\mathcal{H}_{\mathrm{int}}(\mathcal{N})$ a.s. Moreover, if $\mathcal{D} \subset \mathbb{R}^{2}$ is a deterministic countable dense set, then a.s., for each $z \in \mathcal{D}$, the set $\mathcal{N}(z)$ contains a minimal element $l_{z}$ and a maximal element $r_{z}$, and one has $\mathcal{W}^{1}=\overline{\left\{l_{z}: z \in \mathcal{D}\right\}}$ and $\mathcal{W}^{r}=\overline{\left\{r_{z}: z \in \mathcal{D}\right\}}$ a.s.

If $\left(\mathcal{W}^{\mathrm{l}}, \mathcal{W}^{\mathrm{r}}\right)$ and $\mathcal{N}$ are coupled as in Theorem 4.1, then we call $\mathcal{N}$ the Brownian net associated with $\left(\mathcal{W}^{\mathrm{l}}, \mathcal{W}^{\mathrm{r}}\right)$. We also call $\beta_{-}, \beta_{+}$the left and right speed of $\mathcal{N}$. The Brownian net with left and right speeds $\beta_{-}=-1$ and $\beta_{+}=+1$ is called the standard Brownian net. We note that if $\beta_{-}=\beta_{+}$, then $\mathcal{W}^{\mathrm{l}}=\mathcal{W}^{\mathrm{r}}=\mathcal{N}$, i.e., the Brownian net reduces to a Brownian web. It is known [SS08, formula (1.22)] that if $\left(\mathcal{W}^{l}, \mathcal{W}^{r}\right)$ is a left-right Brownian web and $\hat{\mathcal{W}}^{1}, \hat{\mathcal{W}}^{\text {r }}$ are the dual Brownian webs associated with $\mathcal{W}^{\mathrm{l}}, \mathcal{W}^{\mathrm{r}}$, then $\left(-\hat{\mathcal{W}}^{\mathrm{l}},-\hat{\mathcal{W}}^{\mathrm{r}}\right)$ is equally distributed with $\left(\mathcal{W}^{\mathrm{l}}, \mathcal{W}^{\mathrm{r}}\right)$. Therefore, $\left(\hat{\mathcal{W}}^{\mathrm{l}}, \hat{\mathcal{W}}^{\mathrm{r}}\right)$ defines an a.s. unique dual Brownian net $\hat{\mathcal{N}}$ in the same way as $\left(\mathcal{W}^{\mathrm{l}}, \mathcal{W}^{\mathrm{r}}\right)$ defines $\mathcal{N}$.

If $A$ is any closed subset of $\mathbb{R}$ and $\mathcal{N}$ is a standard Brownian net, then setting

$$
\begin{equation*}
\xi_{t}:=\{\pi(t): \pi \in \mathcal{N}(A \times\{0\})\} \quad(t \geq 0) \tag{4.6}
\end{equation*}
$$

defines a Markov process taking values in the space of closed subsets of $\mathbb{R}$, called the branchingcoalescing point set. We refer to Proposition 2.6 for some of its basic properties.

### 4.2 Separation points

Loosely speaking, the separation points of a Brownian net are the limits of separation points of the approximating branching-coalescing random walks, i.e., they are points where paths in the Brownian net have a choice whether to 'turn left' or 'turn right'. These points play an important role in our proofs. In this subsection, we recall some basic facts about them.

Recall from Definition 3.2 the definition of strong equivalence of incoming paths. Following [SSS09], we adopt the following definition.

Definition 4.2 (Equivalence of incoming and outgoing paths) We call two incoming paths $\pi_{1}, \pi_{2} \in \Pi$ at a point $z=(x, t) \in \mathbb{R}^{2}$ equivalent paths entering $z$, denoted by $\pi_{1} \sim_{\text {in }}^{z} \pi_{2}$, if $\pi_{1}\left(t-\varepsilon_{n}\right)=\pi_{2}\left(t-\varepsilon_{n}\right)$ for a sequence $\varepsilon_{n} \downarrow 0$. We call two outgoing paths $\pi_{1}, \pi_{2}$ at a point $z$ equivalent paths leaving $z$, denoted by $\pi_{1} \sim_{\text {out }}^{z} \pi_{2}$, if $\pi_{1}\left(t+\varepsilon_{n}\right)=\pi_{2}\left(t+\varepsilon_{n}\right)$ for a sequence $\varepsilon_{n} \downarrow 0$.

In spite of the suggestive notation, these are not equivalence relations on the spaces of all paths in $\Pi$ entering resp. leaving a point. However, it is known that:
(i) If $\mathcal{W}$ is a Brownian web and $\pi_{1}, \pi_{2} \in \mathcal{W}$ satisfy $\pi_{1}(t)=\pi_{2}(t)$ for some $\sigma_{\pi_{1}}, \sigma_{\pi_{2}}<t$, then $\pi_{1}=\pi_{2}$ on $[t, \infty]$.
(ii) If $\left(\mathcal{W}^{l}, \mathcal{W}^{\mathrm{r}}\right)$ is a left-right Brownian web and $l \in \mathcal{W}^{\mathrm{l}}, r \in \mathcal{W}^{\mathrm{r}}$ satisfy $l(t)=r(t)$ for some $\sigma_{l}, \sigma_{r}<t$, then $l \leq r$ on $[t, \infty]$.


Figure 10: Structure of separation points and a path $\pi$ in the Brownian net turning left at a separation point $z$.

Using (4.7), it is easy to see that if $\left(\mathcal{W}^{l}, \mathcal{W}^{l}\right)$ is a left-right Brownian web, then a.s. for all $z \in \mathbb{R}^{2}$, the relations $\sim_{\text {in }}^{z}$ and $\sim_{\text {out }}^{z}$ are equivalence relations on the set of paths in $\mathcal{W}^{l} \cup \mathcal{W}^{\mathrm{r}}$ entering resp. leaving $z$, and the equivalence classes of paths in $\mathcal{W}^{l} \cup \mathcal{W}^{r}$ entering resp. leaving $z$ are naturally ordered from left to right.

In previous work [SSS09, Theorem 1.7], we have given a complete classification of points $z \in \mathbb{R}^{2}$ according to the structure of the equivalence classes in $\mathcal{W}^{\mathrm{l}} \cup \mathcal{W}^{\mathrm{r}}$ entering resp. leaving $z$, in the spirit of the classification of special points of the Brownian web in Proposition 3.3. It turns out there are 20 types of special points in a left-right Brownian web. Here, we will only be interested in separation points.

By definition, we say that a point $z=(x, t) \in \mathbb{R}^{2}$ is a separation point of two paths $\pi_{1}, \pi_{2} \in \Pi$ if $\sigma_{\pi_{1}}, \sigma_{\pi_{2}}<t, \pi_{1}(t)=x=\pi_{2}(t)$, and $\pi_{1} \not \chi_{\text {out }}^{z} \pi_{2}$. We say that $z$ is a separation point of some collection of paths $\mathcal{A}$ if there exist $\pi_{1}, \pi_{2} \in \mathcal{A}$ such that $z$ is a separation point of $\pi_{1}$ and $\pi_{2}$. Recall the definition of the dual Brownian net below Theorem 4.1. We cite the following proposition from [SSS09, Prop. 2.6 and Thm. 1.12(a)]. See Figure 10.

Proposition 4.3 (Separation points) Let $\mathcal{N}$ be a Brownian net with left and right speeds $\beta_{-}<\beta_{+}$and let $\left(\mathcal{W}^{\mathrm{l}}, \mathcal{W}^{\mathrm{r}}\right)$ be its associated left-right Brownian web. Then:
(a) A.s., $S:=\left\{z \in \mathbb{R}^{2}: z\right.$ is a separation point of $\left.\mathcal{N}\right\}=\left\{z \in \mathbb{R}^{2}: z\right.$ is a separation point of $\hat{\mathcal{N}}\}$, and $S$ is countable.
(b) A.s., $S=\left\{z \in \mathbb{R}^{2}: z\right.$ is of type $(1,2)_{1}$ in $\mathcal{W}^{1}$ and of type $(1,2)_{\mathrm{r}}$ in $\left.\mathcal{W}^{\mathrm{r}}\right\}$.
(c) For given $z \in S$, let $l_{z}$ and $r_{z}$ denote the, up to strong equivalence unique, incoming paths in $\mathcal{W}^{\mathrm{l}}$ resp. $\mathcal{W}^{\mathrm{r}}$ at $z$ and let $l_{z}^{\prime}$ and $r_{z}^{\prime}$ be the elements of $\mathcal{W}^{\mathrm{l}}(z)$ resp. $\mathcal{W}^{\mathrm{r}}(z)$ that are not continuations of $l_{z}$ resp. $r_{z}$. Then, a.s. for each $z \in S$, one has $l_{z} \sim_{\mathrm{in}}^{z} r_{z}, l_{z} \sim_{\text {out }}^{z} r_{z}^{\prime}$ and $l_{z}^{\prime} \sim_{\text {out }}^{z} r_{z}$.
(d) With the same notation as in (c), a.s. for each $z=(x, t) \in S$ and for each incoming path $\pi \in \mathcal{N}$ at $z$, there exists some $\varepsilon>0$ such that $l_{z} \leq \pi \leq r_{z}$ on $[t-\varepsilon, \infty)$. Moreover, each path $\pi \in \mathcal{N}$ leaving $z$ satisfies either $l_{z} \leq \pi \leq r_{z}^{\prime}$ on $[t, \infty)$ or $l_{z}^{\prime} \leq \pi \leq r_{z}$ on $[t, \infty)$.

Note that part (c) of this proposition says that at each separation point $z$ there is one pair of equivalent (in the sense of $\sim_{\text {in }}^{z}$ ) incoming paths $\left\{l_{z}, r_{z}\right\}$ and there are two pairs of equivalent (in the sense of $\sim_{\text {out }}^{z}$ ) outgoing paths: $\left\{l_{z}, r_{z}^{\prime}\right\}$ and $\left\{l_{z}^{\prime}, r_{z}\right\}$. By part (d), whenever a path $\pi \in \mathcal{N}$ enters $z$, it must do so squeezed between $\left\{l_{z}, r_{z}\right\}$ and it must leave $z$ squeezed either between the pair $\left\{l_{z}, r_{z}^{\prime}\right\}$ or between the pair $\left\{l_{z}^{\prime}, r_{z}\right\}$ (see Figure 10). By part (b), these are the only points in $\mathbb{R}^{2}$ where paths in $\mathcal{N}$ can separate from each other, and by part (a), there are only countable many of these points.

### 4.3 Switching and hopping inside a Brownian net

In this subsection, we show how it is possible to construct a Brownian web inside a Brownian net by turning separation points into points of type $(1,2)$ with i.i.d. orientations. In the next subsection, this will be used to state the main result of this section, which is an analogue of Theorem 3.7 and gives an alternative construction of Howitt-Warren flows with finite left and right speeds based on a reference Brownian net.

Recall from Proposition 4.3 (d) (see also Figure 10) that if $\mathcal{N}$ is a Brownian net and $\pi \in \mathcal{N}$ is some path entering a separation point $z=(x, t)$ of $\mathcal{N}$, then $\pi$ must leave $z$ squeezed between one of the two outgoing pairs $\left\{l_{z}, r_{z}^{\prime}\right\}$ or $\left\{l_{z}^{\prime}, r_{z}\right\}$. We write

$$
\operatorname{sign}_{\pi}(z):= \begin{cases}-1 & \text { if } l_{z} \leq \pi \leq r_{z}^{\prime} \text { on }[t, \infty)  \tag{4.8}\\ +1 & \text { if } l_{z}^{\prime} \leq \pi \leq r_{z} \text { on }[t, \infty)\end{cases}
$$

Recall the definition of the dual Brownian net below Theorem 4.1 and recall from Proposition 4.3 (a) that the set $S$ of separation points of $\mathcal{N}$ coincides with the set of separation points of $\hat{\mathcal{N}}$. For $\hat{\pi} \in \hat{\mathcal{N}}$, we define $\operatorname{sign}_{\hat{\pi}}(z)$ to be the sign of $-z$ in $-\hat{\mathcal{N}}$. The next theorem, which will be proved by discrete approximation, shows that it is possible to define a Brownian web 'inside' a Brownian net.

Theorem 4.4 (Brownian web inside a Brownian net) Let $\mathcal{N}$ be a Brownian net with left and right speeds $\beta_{-} \leq \beta_{+}$, let $\hat{\mathcal{N}}$ be its associated dual Brownian net, and let $r \in[0,1]$. Let $S$ be the set of separation points of $\mathcal{N}$ and conditional on $\mathcal{N}$, let $\alpha=\left(\alpha_{z}\right)_{z \in S}$ be a collection of i.i.d. $\{-1,+1\}$-valued random variables such that $\mathbb{P}\left[\alpha_{z}=+1 \mid \mathcal{N}\right]=r$. Then

$$
\begin{align*}
& \mathcal{W}:=\left\{\pi \in \mathcal{N}: \operatorname{sign}_{\pi}(z)=\alpha_{z} \forall z \in S \text { s.t. } \pi \text { enters } z\right\}, \\
& \hat{\mathcal{W}}:=\left\{\hat{\pi} \in \hat{\mathcal{N}}: \operatorname{sign}_{\hat{\pi}}(z)=\alpha_{z} \forall z \in S \text { s.t. } \hat{\pi} \text { enters } z\right\} \tag{4.9}
\end{align*}
$$

defines a Brownian web $\mathcal{W}$ with drift $\beta:=(1-r) \beta_{-}+r \beta_{+}$and its associated dual Brownian web $\hat{\mathcal{W}}$. In particular, if $r=0$ resp. $r=1$, then $\mathcal{W}$ is the left (resp. right) Brownian web associated with $\mathcal{N}$. In general, if $\ell$ denotes the intersection local time measure between $\mathcal{W}$ and its dual and $\ell_{1}, \ell_{\mathrm{r}}$ are the restrictions of $\ell$ to the sets of points of type $(1,2)_{1}$ resp. $(1,2)_{\mathrm{r}}$, then, conditional on $\mathcal{W}$, the sets $S_{1}:=\left\{z \in S: \alpha_{z}=-1\right\}$ and $S_{\mathrm{r}}:=\left\{z \in S: \alpha_{z}=+1\right\}$ are independent Poisson point sets with intensities $\left(\beta_{+}-\beta\right) \ell_{1}$ and $\left(\beta-\beta_{-}\right) \ell_{\mathrm{r}}$, respectively.

For any Brownian web $\mathcal{W}$, we may without loss of generality assume that $\mathcal{W}$ is constructed 'inside' some Brownian net $\mathcal{N}$ as in Theorem 4.4. This will be very helpful in understanding marking constructions based on $\mathcal{W}$ such as the 'switching' construction of sticky Brownian webs in Theorem 3.5 or the marking construction of a Brownian net. To reap the full profit of Theorem 4.4 , we need one more result, which we formulate next.

Let $\mathcal{W}$ be a Brownian web with drift $\beta$. For each point $z$ of type $(1,2)$, let $\operatorname{switch}_{z}(\mathcal{W})$ denote the web obtained from $\mathcal{W}$ by switching the orientation of $z$ as in (3.11), and let

$$
\begin{equation*}
\operatorname{hop}_{z}(\mathcal{W}):=\mathcal{W} \cup \operatorname{switch}_{z}(\mathcal{W}) \tag{4.10}
\end{equation*}
$$

be the compact set of paths obtained from $\mathcal{W}$ by allowing hopping at $z$, i.e., by allowing incoming paths at $z$ to continue along any of the outgoing paths. More generally, if $\Delta$ is a finite set of points of type $(1,2)$ in $\mathcal{W}$, then we set

$$
\begin{equation*}
\operatorname{hop}_{\Delta}(\mathcal{W}):=\bigcup_{\Delta^{\prime} \subset \Delta} \operatorname{switch}_{\Delta^{\prime}}(\mathcal{W}), \tag{4.11}
\end{equation*}
$$

where the union ranges over all subsets $\Delta^{\prime} \subset \Delta$, with $\operatorname{switch}_{\emptyset}(\mathcal{W}):=\mathcal{W}$.
Proposition 4.5 (Switching and hopping inside a Brownian net) Let $\mathcal{N}$ be a Brownian net with left and right speeds $\beta_{-} \leq \beta_{+}$and set of separation points $S$. Conditional on $\mathcal{N}$, let $\alpha=\left(\alpha_{z}\right)_{z \in S}$ be a collection of i.i.d. $\{-1,+1\}$-valued random variables such that $\mathbb{P}\left[\alpha_{z}=\right.$ $+1 \mid \mathcal{N}]=r$ and let $\mathcal{W}$ be a Brownian web with drift $\beta:=(1-r) \beta_{-}+r \beta_{+}$defined inside $\mathcal{N}$ as in (4.9). Then, a.s. for each subset $S^{\prime} \subset S$ and for each sequence of finite sets $\Delta_{n} \uparrow S^{\prime}$, the limits

$$
\begin{align*}
\mathcal{W}^{\prime} & :=\lim _{\Delta_{n} \uparrow S^{\prime}} \operatorname{switch}_{\Delta_{n}}(\mathcal{W}),  \tag{4.12}\\
\mathcal{N}^{\prime} & =\lim _{\Delta_{n} \uparrow S^{\prime}} \operatorname{lop}_{\Delta_{n}}(\mathcal{W})
\end{align*}
$$

exist in $\mathcal{K}(\Pi)$ and are given by

$$
\begin{align*}
& \mathcal{N}^{\prime}=\left\{\pi \in \mathcal{N}: \operatorname{sign}_{\pi}(z)=\alpha_{z} \forall z \in S \backslash S^{\prime} \text { s.t. } \pi \text { enters } z\right\}, \\
& \mathcal{W}^{\prime}=\mathcal{N}^{\prime} \cap\left\{\pi \in \mathcal{N}: \operatorname{sign}_{\pi}(z)=-\alpha_{z} \forall z \in S^{\prime} \text { s.t. } \pi \text { enters } z\right\} . \tag{4.13}
\end{align*}
$$

### 4.4 Construction of Howitt-Warren flows inside a Brownian net

By combining Theorem 4.4 and Proposition 4.5, one can give short proofs of some of the results we have seen before, such as the marking construction of sticky Brownian webs (Theorem 3.5), Proposition 3.6 on changing the reference web, and the equivalence of the definitions of a leftright Brownian web given in Sections 3.3 and 4.1. From Theorem 4.4 and Proposition 4.5 , one moreover easily deduces the following result, which is similar to the marking construction of the Brownian net given in [NRS10, Sec. 3.3.1 and Thm. 5.5]. For the proofs of all these results, we refer to Section 6.6.

Theorem 4.6 (Marking construction of the Brownian net) Let $\mathcal{W}$ be a Brownian web with drift $\beta$ and let $c_{1}, c_{\mathrm{r}} \geq 0$. Let $\ell$ denote the intersection local time measure between $\mathcal{W}$ and its dual, and let $\ell_{1}$ and $\ell_{\mathrm{r}}$ denote the restrictions of $\ell$ to the sets of points of type $(1,2)_{1}$ and $(1,2)_{\mathrm{r}}$, respectively. Conditional on $\mathcal{W}$, let $S_{1}$ and $S_{\mathrm{r}}$ be independent Poisson point sets
with intensities $c_{1} \ell_{1}$ and $c_{\mathrm{r}} \ell_{\mathrm{r}}$, respectively. Then, for any sequence of finite sets $\Delta_{n}^{1} \uparrow S_{1}$ and $\Delta_{n}^{\mathrm{r}} \uparrow S_{\mathrm{r}}$, the limits

$$
\begin{align*}
\text { (i) } \mathcal{N}:=\lim _{n \rightarrow \infty} \operatorname{hop}_{\Delta_{n}^{1} \cup \Delta_{n}^{\mathrm{r}}}(\mathcal{W}), \\
\text { (ii) } \mathcal{W}^{\mathrm{l}}:=\lim _{n \rightarrow \infty}{\operatorname{switch} \Delta_{n}^{\mathrm{r}}}^{(\mathcal{W})},  \tag{4.14}\\
\text { (iii) } \mathcal{W}^{\mathrm{r}}:=\lim _{n \rightarrow \infty} \operatorname{switch}_{\Delta_{n}^{1}}(\mathcal{W})
\end{align*}
$$

exist in $\mathcal{K}(\Pi)$ a.s. and do not depend on the choice of the sequences $\Delta_{n}^{1} \uparrow S_{1}$ and $\Delta_{n}^{\mathrm{r}} \uparrow S_{\mathrm{r}}$. Moreover, $\mathcal{N}$ is a Brownian net with left and right speeds $\beta_{-}:=\beta-c_{r}$ and $\beta_{+}:=\beta+c_{1}$, $\left(\mathcal{W}^{1}, \mathcal{W}^{\mathrm{r}}\right)$ is its associated left-right Brownian web, and $S:=S_{1} \cup S_{\mathrm{r}}$ is its set of separation points. If $c_{1}+c_{\mathrm{r}}>0$, then conditional on $\mathcal{N}$, the random variables $\left(\operatorname{sign}_{\mathcal{W}}(z)\right)_{z \in S}$ are i.i.d. with $\mathbb{P}\left[\operatorname{sign}_{\mathcal{W}}(z)=+1 \mid \mathcal{N}\right]=c_{\mathrm{r}} /\left(c_{1}+c_{\mathrm{r}}\right)$.

Our final result of this section shows how the construction of Brownian webs inside a Brownian net given in Theorem 4.4 can be used to construct Howitt-Warren flows with finite left and right speeds, providing an alternative to Theorem 3.7. Recall that a Howitt-Warren quenched law with drift $\beta$ and characteristic measure $\nu$ is a random probability measure $\mathbb{Q}$ on $\mathcal{K}(\Pi)$ with law as defined in (3.24).

Theorem 4.7 (Construction of Howitt-Warren flows with finite speeds) Let $\beta \in \mathbb{R}$ and let $\nu$ be a finite measure on $[0,1]$ such that the speeds $\beta_{-}, \beta_{+}$defined in (2.12) are finite. Let $\mathcal{N}$ be a Brownian net with left and right speeds $\beta_{-}, \beta_{+}$and let $S$ be its set of separation points. Conditional on $\mathcal{N}$, let $\omega:=\left(\omega_{z}\right)_{z \in S}$ be a collection of i.i.d. $[0,1]$-valued random variables with law $\bar{\nu}(\mathrm{d} q):=b^{-1} q^{-1}(1-q)^{-1} \nu(\mathrm{~d} q)$, where $b:=\int q^{-1}(1-q)^{-1} \nu(\mathrm{~d} q)$, and conditional on $(\mathcal{N}, \omega)$, let $\left(\alpha_{z}\right)_{z \in S}$ be a collection of independent $\{-1,+1\}$-valued random variables such that $\mathbb{P}\left[\alpha_{z}=1 \mid(\mathcal{N}, \omega)\right]=\omega_{z}$. Set

$$
\begin{equation*}
\mathcal{W}:=\left\{\pi \in \mathcal{N}: \operatorname{sign}_{\pi}(z)=\alpha_{z} \forall z \in S \text { s.t. } \pi \text { enters } z\right\} . \tag{4.15}
\end{equation*}
$$

Then setting

$$
\begin{equation*}
\mathbb{Q}:=\mathbb{P}[\mathcal{W} \in \cdot \mid(\mathcal{N}, \omega)] \tag{4.16}
\end{equation*}
$$

yields a Howitt-Warren quenched law with drift $\beta$ and characteristic measure $\nu$. In particular, setting

$$
\begin{equation*}
K_{s, t}^{\uparrow}(x, A):=\mathbb{P}\left[\pi_{(x, s)}^{\uparrow}(t) \in A \mid(\mathcal{N}, \omega)\right] \quad(s \leq t, x \in \mathbb{R}, A \in \mathcal{B}(\mathbb{R})) \tag{4.17}
\end{equation*}
$$

and defining $K_{s, t}^{+}(x, A)$ similarly with $\pi_{(x, s)}^{\uparrow}$ replaced by $\pi_{(x, s)}^{+}$yields versions of the HowittWarren flow with drift $\beta$ and characteristic measure $\nu$ with properties as described in Proposition 3.8.

### 4.5 Support of the quenched law

In this subsection, we formulate a theorem on the support of Howitt-Warren quenched laws, which will imply Theorems 2.5 and 2.7. Before we can do this, we need to introduce Brownian half-nets, which are basically Brownian nets with either infinite left speed and finite right speed, or vice versa. Recall that a path $\pi \in \Pi$ crosses a dual path $\hat{\pi} \in \hat{\Pi}$ from left to right if there exist $\sigma_{\pi} \leq s<t \leq \hat{\sigma}_{\hat{\pi}}$ such that $\pi(s)<\hat{\pi}(s)$ and $\pi(t)>\hat{\pi}(t)$. Crossing from right to left and crossing of forward paths are defined analogously. We will prove the following analogue of Theorem 4.1.

Theorem 4.8 (Brownian half-net associated with a Brownian web) Let $\mathcal{W}$ be a Brownian web with drift $\beta$ and let $\hat{\mathcal{W}}$ be its dual. Then there exists a random closed set of paths $\mathcal{H}_{-} \subset \Pi$ that is a.s. uniquely determined by any of the following equivalent conditions:
(i) $\mathcal{H}_{-}=\{\pi \in \Pi: \pi$ does not cross any path of $\mathcal{W}$ from left to right $\}$ a.s.
(ii) $\mathcal{H}_{-}=\{\pi \in \Pi: \pi$ does not cross any path of $\hat{\mathcal{W}}$ from left to right $\}$ a.s.

Moreover, if $\mathcal{D} \subset \mathbb{R}^{2}$ is a deterministic countable dense set, then a.s., for each $z \in \mathcal{D}$, the set $\mathcal{H}_{-}(z)$ contains a maximal element $\pi_{z}$, and one has $\mathcal{W}=\overline{\left\{\pi_{z}: z \in \mathcal{D}\right\}}$. Analogue statements hold with $\mathcal{H}_{-}$replaced by $\mathcal{H}_{+}$, 'from left to right' replaced by 'from right to left' and 'maximal element' replaced by 'minimal element'.

If $\mathcal{H}_{-}\left(\right.$resp. $\left.\mathcal{H}_{+}\right)$and $\mathcal{W}$ are coupled as in Theorem 4.8, then we call $\mathcal{H}_{-}\left(\right.$resp. $\left.\mathcal{H}_{+}\right)$a Brownian half-net with left and right speeds $-\infty, \beta$ (resp. $\beta,+\infty$ ), and we call $\mathcal{W}$ the right (resp. left) Brownian web associated with $\mathcal{H}_{-}$(resp. $\mathcal{H}_{+}$).

Let $\mathbb{Q}$ be a Howitt-Warren quenched law with drift $\beta$ and characteristic measure $\nu$ as defined as in (3.24), or alternatively, in the case of finite speeds, as in 4.16). Then $\mathbb{Q}$ is a random probability law on the space of webs. In particular, if $\mathcal{W}$ is a $\mathcal{K}(\Pi)$-valued random variable with (random) law $\mathbb{Q}$, then for each $z \in \mathbb{R}^{2}$ we can define special paths $\pi_{z}^{\uparrow}$ and $\pi_{z}^{+}$ in $\mathcal{W}(z)$ as below Proposition 3.3. In analogy with the conditional law $\mathbf{Q}_{(x, s)}^{\omega}$ of the random walk in random environment defined in Section 1.2 , in the continuum setting, we define

$$
\begin{equation*}
\mathbb{Q}_{z}^{+}:=\mathbb{Q}\left[\pi_{z}^{+} \in \cdot\right] \quad\left(z \in \mathbb{R}^{2}\right) \tag{4.18}
\end{equation*}
$$

and we define $\mathbb{Q}_{z}^{\uparrow}$ similarly, with $\pi_{z}^{+}$replaced by $\pi_{z}^{\uparrow}$. In particular, if $\mathbb{Q}$ is defined as in 3.24) or as in 4.16), this says that $\mathbb{Q}_{z}^{+}:=\mathbb{P}\left[\pi_{z}^{+} \in \cdot \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right]$ resp. $\mathbb{Q}_{z}^{+}:=\mathbb{P}\left[\pi_{z}^{+} \in \cdot \mid(\mathcal{N}, \omega)\right]$. We note that since typical points in $\mathbb{R}^{2}$ are of type $(0,1)$ in $\mathcal{W}$, for deterministic $z \in \mathbb{R}^{2}$, the random variables $\mathbb{Q}_{z}^{+}$and $\mathbb{Q}_{z}^{\uparrow}$ are equal a.s. It follows that for any deterministic finite measure $\mu$ on $\mathbb{R}^{2}$, one has $\int \mu(\mathrm{d} z) \mathbb{Q}_{z}^{+}=\int \mu(\mathrm{d} z) \mathbb{Q}_{z}^{\dagger}$.

Theorem 4.9 (Support property) Let $\mathbb{Q}$ be a Howitt-Warren quenched law with drift $\beta$ and characteristic measure $\nu$, and let $\beta_{-}, \beta_{+}$be the left and right speeds defined in (2.12). Then there exists a random, closed subset $\mathcal{N} \subset \Pi$ such that for any deterministic finite measure $\mu$ on $\mathbb{R}^{2}$,

$$
\begin{equation*}
\operatorname{supp}\left(\int \mu(\mathrm{d} z) \mathbb{Q}_{z}^{+}\right)=\overline{\mathcal{N}(\operatorname{supp}(\mu))} \quad \text { a.s. } \tag{4.19}
\end{equation*}
$$

If $-\infty<\beta_{-} \leq \beta_{+}<+\infty$, then $\mathcal{N}$ is Brownian net with left and right speeds $\beta_{-}, \beta_{+}$. If either $-\infty=\beta_{-}<\beta_{+}<+\infty$ or $-\infty<\beta_{-}<\beta_{+}=+\infty$, then $\mathcal{N}$ is a Brownian half-net with left and right speeds $\beta_{-}, \beta_{+}$. If $-\infty=\beta_{-}<\beta_{+}=+\infty$, then $\mathcal{N}=\Pi$.

Note that above, $\operatorname{supp}(\mu)$ is a closed subset of $\mathbb{R}^{2}$, but not necessarily of $R_{\mathrm{c}}^{2}$, which is why in general we need to take the closure of $\mathcal{N}(\operatorname{supp}(\mu))$ in the space of paths $\Pi$. If $(*,-\infty) \notin$ $\overline{\operatorname{supp}(\mu)}$ or if $\operatorname{supp}(\mu)=\mathbb{R}^{2}$, then it is moreover true that $\overline{\mathcal{N}(\operatorname{supp}(\mu))}=\mathcal{N}(\overline{\operatorname{supp}(\mu)})$, where $\overline{\operatorname{supp}(\mu)}$ denotes the closure of $\operatorname{supp}(\mu)$ in $R_{\mathrm{c}}^{2}$; see Lemma 8.5 below.

We note, without proof, that in the setup of Theorem 4.9, it can be shown that $\mathcal{N}=$ $\cup \operatorname{supp}(\mathbb{Q})$, where $\cup \operatorname{supp}(\mathbb{Q}):=\{\pi: \pi \in \mathcal{A}$ for some $\mathcal{A} \in \operatorname{supp}(\mathbb{Q})\}$ denotes the union of all elements of $\operatorname{supp}(\mathbb{Q}) \subset \mathcal{K}(\Pi)$. We state as an open problem to characterize $\operatorname{supp}(\mathbb{Q})$ itself $($ rather than just $\cup \operatorname{supp}(\mathbb{Q}))$.

## 5 Outline of the proofs

Our results are proved in Sections $6 \sqrt{11}$ below. In Section 6 we collect some well-known and less well-known facts about the Brownian web and net and prove some new results that we will need further on. In particular, in Section 6.2 we prove a 'finite graph representation' that gives a precise description of how paths in the Brownian net move between deterministic times. Sections 6.3 6.5 then culminate in Theorem 6.15, the central result of the section, which is about discrete approximation of a Brownian web embedded in a Brownian net and implies Theorem 4.4. In Section 6.6, this is then used, together with the finite graph representation, to prove Theorem 3.5 and Proposition 3.6 on the construction of sticky Brownian webs and related results such as Proposition 4.5 and Theorem 4.6.

In Section 7 we prove our main results: Theorem 3.9 on the convergence of the quenched laws on the space of webs, and Theorems 3.7 and 4.7 on the construction of Howitt-Warren flows using a marked reference Brownian web or net. Here we also harvest some immediate consequences of our construction, such as the existence of regular versions of Howitt-Warren flows (Proposition 2.3 and 3.8) and scaling (Proposition 2.4).

In Section 8 we prove our results on the support of Howitt-Warren flows. In Section 8.1, we prove a number of preparatory results about generalized Brownian nets with possibly infinite left and right speeds. In particular, we prove Theorem 4.8 on Brownian half-nets and prepare for the proof of Theorem 4.9 on the support of the quenched law on the space of webs. In Section 8.2, we prove Theorem 4.9 and use it to deduce Theorems 2.5 and 2.7 on the left and right speeds and the support of Howitt-Warren processes.

In Section 9, we address questions of atomicness. In particular, parts (a), (b) and (c) of Theorem 2.8 are proved in Sections 9.1, 9.2 and 9.3 , respectively.

In Section 10 we prove Theorems 2.9 and 2.10 on Howitt-Warren processes with infinite starting mass and the convergence of rescaled discrete Howitt-Warren processes, while Section 11 contains the proofs of Theorems 2.11 and 2.12 on homogeneous invariant laws.

The paper concludes with four appendices on the Howitt-Warren martingale problem and some other technical issues.

The table below gives a complete overview of where the proofs can be found of all results stated so far. Further results stated in the following sections will be proved on the spot. Below, cited means that the proof of the listed result is cited from other sources.

| Result | Proved in | Result | Proved in | Result | Proved in |
| :---: | :---: | :---: | :---: | :---: | :---: |
| Prop. 2.3 | Sect. 7.3 | Thm. 2.12 | Sect. 11.4 | Thm. 4.1 | cited |
| Prop. 2.4 | Sect. 7.3 | Prop. 3.1 | cited | Prop. 4.3 | cited |
| Thm. 2.5 | Sect. $\overline{8.2}$ | Prop. 3.3 | cited | Thm. 4.4 | Sect. 6.5 |
| Prop. 2.6 | cited | Prop. 3.4 | cited | Prop. 4.5 | Sect. 6.6 |
| Thm. 2.7 | Sect. 8.2 | Thm. 3.5 | Sect. 6.6 | Thm. 4.6 | Sect. $\overline{6.6}$ |
| Thm. 2.8 | Sect. 9 | Prop. 3.6 | Sect. $\overline{6.6}$ | Thm. 4.7 | Sect. 7.2 |
| Thm. 2.9 | Sect. 10.1 | Thm. 3.7 | Sect. 7.2 | Thm. 4.8 | Sect. 8.1 |
| Thm. 2.10 | Sect. 10.2 | Prop. 3.8 | Sect. 7.3 | Thm. $\overline{4.9}$ | Sect. $\overline{8.2}$ |
| Thm. 2.11 | Sect. 11.4 | Thm. 3.9 | Sect. 7.1 |  |  |

## 6 Coupling of the Brownian web and net

The main aim of this section is to prove Theorem 4.4 and Proposition 4.5, which will be our main tools for constructing modified Brownian webs and nets by switching or hopping inside a reference Brownian web or net. In particular, after proving these theorems, we will apply them to prove Theorems 3.5 and 4.6 on the switching construction of sticky Brownian webs and the marking construction of the Brownian net.

In order to prepare for the proofs of Theorem 4.4 and Proposition 4.5, we first need to take a closer look at the separation points of a Brownian net, introduced in Section 4.2, It has been proved in [SSS09] that for deterministic times $S<U$, there are only locally finitely many ' $S, U$-relevant' separation points that decide where paths in the Brownian net started at time $S$ end up at time $U$. After recalling some basic facts about these relevant separation points in Section 6.1, we use them in Section 6.2 to give a rather precise description, by means of a 'finite graph representation', of the way paths in the Brownian net move between time $S$ and $U$.

Since discrete approximation will play an important role in our proofs, Sections 6.36 .4 are devoted to discrete approximation of the Brownian web and net, and related objects such as intersection local times and separation points. In Section 6.5, we then use these results to prove a result about the convergence of a discrete web embedded in a discrete net to analogue Brownian objects. This result then immediately yields Theorem 4.4 on the construction of a Brownian web inside a Brownian net. In addition, it lays the basis for proofs of other convergence results such as Theorem 3.9 on the convergence of quenched laws. In Section 6.6, finally, we use the finite graph representation developed in Section 6.2 together with Theorem 4.4 to prove Proposition 4.5 and we combine Theorem 4.4 and Proposition 4.5 to prove Theorems 3.5 and 4.6 and some related results.

### 6.1 Relevant separation points

The set of separation points of a Brownian net $\mathcal{N}$ is dense in $\mathbb{R}^{2}$ and also along any path $\pi \in \mathcal{N}$. It turns out, however, that for given deterministic times $S<U$, the set of separation points that are relevant for deciding where paths in the Brownian net started at time $S$ end up at time $U$ is a locally finite subset of $\mathbb{R} \times[S, U]$.

Following [SSS09], we say that a separation point $z=(x, t)$ of a Brownian net $\mathcal{N}$ is $S, U$ relevant for some $-\infty \leq S<t<U \leq \infty$, if there exists $\pi \in \mathcal{N}$ such that $\sigma_{\pi}=S$ and $\pi(t)=x$, and there exist $l \in \mathcal{W}^{\mathrm{l}}(z)$ and $r \in \mathcal{W}^{\mathrm{r}}(z)$ such that $l<r$ on $(t, U)$. (Note that since we are assuming that $z$ is a separation point, $l$ and $r$ have to be the paths $l_{z}$ and $r_{z}$ from Proposition 4.3 (c). In particular, $l$ and $r$ are continuations of incoming paths at $z$.) The next proposition follows easily, by Brownian scaling, from [SSS09, Lemma 2.8 and Prop. 2.9]. Part (a) says that the definition of relevant separation points is symmetric with respect to duality; see also Figure 11.

Proposition 6.1 (Relevant separation points) Let $\mathcal{N}$ be a Brownian net with left and right speeds $\beta_{-} \leq \beta_{+}$. Then:
(a) A.s. for each $-\infty \leq S<U \leq \infty$, a separation point $z=(x, t)$ with $S<t<U$ is $S, U$-relevant in $\mathcal{N}$ if and only if $-z$ is $-U,-S$-relevant in the rotated dual Brownian net $-\hat{\mathcal{N}}$.


Figure 11: An $S, U$-relevant separation point.
(b) For each deterministic $-\infty \leq S<U \leq \infty$, if $R_{S, U}$ denotes the set of $S, U$-relevant separation points, then

$$
\begin{equation*}
\mathbb{E}\left[\left|R_{S, U} \cap A\right|\right]=2 b \int_{A} \Psi_{b}(t-S) \Psi_{b}(U-t) \mathrm{d} x \mathrm{~d} t \quad(A \in \mathcal{B}(\mathbb{R} \times(S, U))) \tag{6.1}
\end{equation*}
$$

where $b:=\left(\beta_{+}-\beta_{-}\right) / 2$,

$$
\begin{equation*}
\Psi_{b}(t):=\frac{e^{-b^{2} t}}{\sqrt{\pi t}}+2 b \Phi(b \sqrt{2 t}) \quad(0<t \leq \infty) \tag{6.2}
\end{equation*}
$$

and $\Phi(x):=\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{x} e^{-y^{2} / 2} \mathrm{~d} y$. In particular, if $-\infty<S, U<\infty$, then $R_{S, U}$ is a.s. a locally finite subset of $\mathbb{R} \times[S, U]$.

We will need yet another characterization of relevant separation points. To formulate this, we first need to recall the definition of crossing times from [SSS09, Def. 2.4].

Definition 6.2 (Crossing and crossing points) We say that a forward path $\pi \in \Pi$ crosses a dual path $\hat{\pi} \in \hat{\Pi}$ from left to right at time $t$ if there exist $\sigma_{\pi} \leq t_{-}<t<t_{+} \leq \hat{\sigma}_{\hat{\pi}}$ such that $\pi\left(t_{-}\right)<\hat{\pi}\left(t_{-}\right), \hat{\pi}\left(t_{+}\right)<\pi\left(t_{+}\right)$, and $t=\inf \left\{s \in\left(t_{-}, t_{+}\right): \hat{\pi}(s)<\pi(s)\right\}=\sup \left\{s \in\left(t_{-}, t_{+}\right):\right.$ $\pi(s)<\hat{\pi}(s)\}$. Crossing from right to left is defined analogously. We call $z=(x, t) \in \mathbb{R}^{2} a$ crossing point of $\pi \in \Pi$ and $\hat{\pi} \in \hat{\Pi}$ if $\pi(t)=x=\hat{\pi}(t)$ and $\pi$ crosses $\hat{\pi}$ either from left to right or from right to left at time $t$.

Lemma 6.3 (Relevant separation points and crossing points) Almost surely for each $-\infty \leq S<U \leq \infty$ and $z \in \mathbb{R} \times(S, U)$, the following statements are equivalent:
(i) $z$ is an $S, U$-relevant separation point.
(ii) $z$ is a crossing point of some $\pi \in \mathcal{N}$ and $\hat{\pi} \in \hat{\mathcal{N}}$ with $\sigma_{\pi}=S$ and $U=\hat{\sigma}_{\hat{\pi}}$.

Proof. If $z=(x, t)$ is an $S, U$-relevant separation point, then by Proposition 6.1 there exist $\pi^{\prime} \in \mathcal{N}$ starting at time $\sigma_{\pi^{\prime}}=S$ and $\hat{\pi}^{\prime} \in \hat{\mathcal{N}}$ starting at time $\hat{\sigma}_{\hat{\pi}^{\prime}}=U$ such that $\pi^{\prime}$ and $\hat{\pi}^{\prime}$ enter $z$. By [SSS09, Prop. 2.6], $z$ is a crossing point of some $r \in \mathcal{W}^{\mathrm{r}}$ and $\hat{l} \in \hat{\mathcal{W}}^{\mathrm{l}}$. Let $\pi$ be the concatenation of $\pi^{\prime}$ on $[S, t]$ and $r$ on $[t, U]$ and likewise, let $\hat{\pi} \in \mathcal{N}$ be the concatenation of $\hat{\pi}^{\prime}$ on $[t, U]$ and $\hat{l}$ on $[S, t]$. Since by Theorem 4.1, $\mathcal{N}$ is closed under hopping, we see that $\pi \in \mathcal{N}$ and $\hat{\pi} \in \hat{\mathcal{N}}$. By the structure of separation points (Proposition 4.3 (d)), $z$ is a crossing point of $\pi$ and $\hat{\pi}$, proving the implication (i) $\Rightarrow$ (ii).

Conversely, if $z$ is a crossing point of some $\pi \in \mathcal{N}$ and $\hat{\pi} \in \hat{\mathcal{N}}$ with $\sigma_{\pi}=S$ and $U=\hat{\sigma}_{\hat{\pi}}$, then by the classification of special points of the Brownian net [SSS09, Thm. 1.7] and their structure [SSS09, Thm. 1.12 (d)], $z$ must be a separation point of $\mathcal{N}$. By [SSS09, Lemma 2.7 (a)], the presence of the dual path $\hat{\pi}$ implies the existence of $l \in \mathcal{W}^{\mathrm{l}}(z), r \in \mathcal{W}^{\mathrm{r}}(z)$ such that $l<r$ on $(t, U)$, hence $z$ is $S, U$-relevant.

### 6.2 Finite graph representation

In this section, we give a rather precise description of how paths in a Brownian net move between deterministic times $S, U$. In particular, we will construct an oriented graph whose internal vertices are relevant separation points and whose directed edges are pairs consisting of a left-most and right-most path, such that each path in the Brownian net starting at time $S$ must between times $S$ and $U$ move through an oriented path in this graph, and conversely, for each oriented path in the graph there exist paths in the Brownian net following this path.

As a preparation, we need some results from [SSS09] on the special points of the Brownian net. Almost surely, there are 20 types of special points in the Brownian net, but we will only need those that occur at deterministic times, of which there are only three. Let $\mathcal{N}$ be a Brownian net with associated left-right Brownian web $\left(\mathcal{W}^{l}, \mathcal{W}^{r}\right)$. Recall the notion of strong equivalence of paths from Definition 3.2 and the relations $\sim_{\text {in }}^{z}$ and $\sim_{\text {out }}^{z}$ from Definition 4.2. As remarked there, these are equivalence relations on the set of paths in $\mathcal{W}^{1} \cup \mathcal{W}^{\mathrm{r}}$ entering resp. leaving a point $z$, and the corresponding equivalence classes are naturally ordered from left to right. In general, such an equivalence class may be of three types. If it contains only paths in $\mathcal{W}^{1}$ then we say it is of type l, if it contains only paths in $\mathcal{W}^{r}$ then we say it is of type r , and if it contains both paths in $\mathcal{W}^{\mathrm{l}}$ and $\mathcal{W}^{\mathrm{r}}$ then we say it is of type p , standing for pair. To denote the type of a point $z \in \mathbb{R}^{2}$ in a Brownian net $\mathcal{N}$, we first list the incoming equivalence classes in $\mathcal{W}^{1} \cup \mathcal{W}^{\mathrm{r}}$ from left to right and then, separated by a comma, the outgoing equivalence classes.

In our case, there are only three types of points of interest, namely the types (o, p), (p, p) and $(\mathrm{o}, \mathrm{pp})$, where a $o$ means that there are no incoming paths in $\mathcal{N}$ at $z$. We note that by property (4.7) (i), an outgoing equivalence class of type p at a point $z$ contains exactly one path in $\mathcal{W}^{1}(z)$ and one path in $\mathcal{W}^{\mathrm{r}}(z)$. By the same property, at points of type $(\mathrm{p}, \mathrm{p})$, all incoming paths in $\mathcal{W}^{1}$ are strongly equivalent and likewise all incoming paths in $\mathcal{W}^{\mathrm{r}}$ are strongly equivalent. We cite the following result from [SSS09, Thms. 1.7 and 1.12] and [SS08, Prop. 1.8]. Recall the definition of the dual Brownian net below Theorem 4.1.

Proposition 6.4 (Special points at deterministic times) Let $\mathcal{N}$ be a Brownian net, let $\hat{\mathcal{N}}$ be its dual, and let $\left(\mathcal{W}^{1}, \mathcal{W}^{\mathrm{r}}\right)$ and $\left(\hat{\mathcal{W}}^{\mathrm{l}}, \hat{\mathcal{W}}^{\mathrm{r}}\right)$ be the left-right Brownian web and the dual left-right Brownian web associated with $\mathcal{N}$ and $\hat{\mathcal{N}}$. Then:
(a) For each deterministic $t \in \mathbb{R}$, almost surely, each point in $\mathbb{R} \times\{t\}$ is either of type $(\mathrm{o}, \mathrm{p}) /(\mathrm{o}, \mathrm{p}),(\mathrm{p}, \mathrm{p}) /(\mathrm{o}, \mathrm{pp})$ or $(\mathrm{o}, \mathrm{pp}) /(\mathrm{p}, \mathrm{p})$ in $\mathcal{N} / \hat{\mathcal{N}}$, and all of these type occur.


Figure 12: Finite graph representation.
(b) Almost surely, for each point $z=(x, t)$ of type $(\mathrm{o}, \mathrm{p}),(\mathrm{p}, \mathrm{p})$ or $(\mathrm{o}, \mathrm{pp})$ in $\mathcal{N}$ and $\pi \in \mathcal{N}(z)$, there exist $l \in \mathcal{W}^{\mathrm{l}}(z)$ and $r \in \mathcal{W}^{\mathrm{r}}(z)$ such that $l \sim_{\text {out }}^{z} r$ and $l \leq \pi \leq r$ on $[t, \infty)$.
(c) Almost surely, for each point $z=(x, t)$ of type $(\mathrm{p}, \mathrm{p})$ in $\mathcal{N}$, for each $l \in \mathcal{W}^{\mathrm{l}}, r \in \mathcal{W}^{\mathrm{r}}$ and $\pi \in \mathcal{N}$ entering $z$, there exists an $\varepsilon>0$ such that $l \leq \pi \leq r$ on $[t-\varepsilon, \infty)$.

Let $-\infty<S<U<\infty$ be deterministic times, let $R_{S, U}$ be the set of $S, U$-relevant separation points of $\mathcal{N}$ and set

$$
\begin{align*}
& R_{S}:=\mathbb{R} \times\{S\}, \\
& R_{U}:=\left\{(x, U): x \in \mathbb{R}, \exists \pi \in \mathcal{N} \text { with } \sigma_{\pi}=S \text { s.t. } \pi(U)=x\right\} . \tag{6.3}
\end{align*}
$$

We make the set $R:=R_{S} \cup R_{S, U} \cup R_{U}$ into an oriented graph by writing $z \rightarrow_{l, r} z^{\prime}$ if $z, z^{\prime} \in R$, $z \neq z^{\prime}, l \in \mathcal{W}^{\mathrm{l}}(z), r \in \mathcal{W}^{\mathrm{r}}(z), l \sim_{\text {out }}^{z} r$, and $l \sim_{\text {in }}^{z^{\prime}} r$.

Proposition 6.5 (Finite graph representation) Let $\mathcal{N}$ be a Brownian net with associated left-right Brownian web $\left(\mathcal{W}^{1}, \mathcal{W}^{r}\right)$ and let $-\infty<S<U<\infty$ be deterministic times. Let $R:=R_{S} \cup R_{S, U} \cup R_{U}$ and the relation $\rightarrow_{l, r}$ be defined as above. Then, a.s. (see Figure 12):
(a) For each $z \in R_{S}$ that is not of type ( $\mathrm{o}, \mathrm{pp}$ ), there exist unique $l \in \mathcal{W}^{\mathrm{l}}(z), r \in \mathcal{W}^{\mathrm{r}}(z)$ and $z^{\prime} \in R$ such that $z \rightarrow_{l, r} z^{\prime}$.
(b) For each $z=(x, t)$ such that either $z \in R_{S, U}$ or $z \in R_{S}$ is of type ( $\mathrm{o}, \mathrm{pp}$ ), there exist unique $l, l^{\prime} \in \mathcal{W}^{\mathrm{l}}(z), r, r^{\prime} \in \mathcal{W}^{\mathrm{r}}(z)$ and $z^{\prime}, z^{\prime \prime} \in R$ such that $l \leq r^{\prime}<l^{\prime} \leq r$ on $(t, t+\varepsilon)$ for some $\varepsilon>0, z \rightarrow_{l, r^{\prime}} z^{\prime}$ and $z \rightarrow_{l^{\prime}, r} z^{\prime \prime}$. For $z \in R_{S, U}$ one has $z^{\prime} \neq z^{\prime \prime}$. For $z \in R_{S}$ of type ( $\mathrm{o}, \mathrm{pp}$ ), one has $z^{\prime} \neq z^{\prime \prime}$ if and only if there exists a dual path $\hat{\pi} \in \hat{\mathcal{N}}$ with $\hat{\sigma}_{\hat{\pi}}=U$ such that $\hat{\pi}$ enters $z$.
(c) For each $\pi \in \mathcal{N}$ with $\sigma_{\pi}=S$, there exist $z_{i}=\left(x_{i}, t_{i}\right) \in R(i=0, \ldots, n)$ and $l_{i} \in \mathcal{W}^{1}\left(z_{i}\right)$, $r_{i} \in \mathcal{W}^{1}\left(z_{i}\right)(i=0, \ldots, n-1)$ such that $z_{0} \in R_{S}, z_{n} \in R_{U}, z_{i} \rightarrow l_{i}, r_{i} z_{i+1}$ and $l_{i} \leq \pi \leq r_{i}$ on $\left[t_{i}, t_{i+1}\right](i=0, \ldots, n-1)$.
(d) If $z_{i}=\left(x_{i}, t_{i}\right) \in R(i=0, \ldots, n)$ and $l_{i} \in \mathcal{W}^{\mathrm{l}}\left(z_{i}\right), r_{i} \in \mathcal{W}^{\mathrm{l}}\left(z_{i}\right)(i=0, \ldots, n-1)$ satisfy $z_{0} \in R_{S}, z_{n} \in R_{U}$, and $z_{i} \rightarrow_{l_{i}, r_{i}} z_{i+1}(i=0, \ldots, n-1)$, then there exists a $\pi \in \mathcal{N}$ with $\sigma_{\pi}=S$ such that $l_{i} \leq \pi \leq r_{i}$ on $\left[t_{i}, t_{i+1}\right]$.
Proof. By Proposition 6.4 (a), each $z=(x, t) \in R_{S}$ that is not of type ( $\mathrm{o}, \mathrm{pp}$ ) must be of type $(\mathrm{o}, \mathrm{p})$ or $(\mathrm{p}, \mathrm{p})$, hence there exists a unique pair $(l, r)$ consisting of one left-most path $l \in \mathcal{W}^{\mathrm{l}}(z)$ and one right-most path $r \in \mathcal{W}^{\mathrm{r}}(z)$, such that $l \sim_{\text {out }}^{z} r$. Likewise, by Proposition 4.3, for each $z \in R_{S, U}$ there exist exactly two pairs $\left(l, r^{\prime}\right)$ and $\left(l^{\prime}, r\right)$ such that $l, l^{\prime} \in \mathcal{W}^{\mathrm{l}}(z), r, r^{\prime} \in \mathcal{W}^{\mathrm{r}}(z)$, $l \sim_{\text {out }}^{z} r^{\prime}$ and $l^{\prime} \sim_{\text {out }}^{z} r$, and the same is true for $z \in R_{S}$ that is of type (o, pp), by the properties of such points. Therefore, in order to prove parts (a) and (b), assume that $z \in R_{S} \cup R_{S, U}$ and that $l \in \mathcal{W}^{\mathrm{l}}(z)$ and $r \in \mathcal{W}^{\mathrm{r}}(z)$ satisfy $l \sim_{\text {out }}^{z} r$. We claim that there exists a unique $z^{\prime} \in R_{S, U} \cup R_{U}$ such that $z \rightarrow l_{l, r} z^{\prime}$.

To see this, let $\tau:=\sup \{u \in(t, U): l(u)=r(u)\}$ be the last time $l$ and $r$ separate before time $U$. If $\tau=U$, then there exists some $z^{\prime} \in \mathbb{R} \times\{U\}$ such that $l$ and $r$ enter $z^{\prime}$ and hence, by Proposition 6.4 (a), $l \sim_{\text {in }}^{z^{\prime}} r$. On the other hand, if $\tau<U$, then we claim that $z^{\prime}=\left(x^{\prime}, t^{\prime}\right):=(l(\tau), \tau)$ is an $S, U$-relevant separation point. To prove this, we must show that there exists some $\pi \in \mathcal{N}$ with $\sigma_{\pi}=S$ and $\pi\left(t^{\prime}\right)=x^{\prime}$, the other parts of the definition being obviously satisfied. If $z \in \mathbb{R}_{S}$ we may take $\pi=l$. If $z \in R_{S, U}$, then there exists some $\pi \in \mathcal{N}$ with $\sigma_{\pi}=S$ such that $\pi(t)=x$. By Theorem 4.1, the Brownian net is closed under hopping, therefore we may concatenate $\pi$ with $l$ to find a path in $\mathcal{N}$ starting at time $S$ and entering $z^{\prime}$. This proves that $z^{\prime}$ is an $S, U$-relevant separation point. By Proposition 4.3 (b), the left-most and right-most paths entering a separation point are up to strong equivalence unique, and $l \sim_{\mathrm{in}}^{z^{\prime}} r$. This proves the existence of a $z^{\prime} \in R_{S, U} \cup R_{U}$ such that $z \rightarrow_{l, r} z^{\prime}$. The uniqueness of $z^{\prime}$ follows from the fact that only the last separation point of $l$ and $r$ before time $U$ can be $S, U$-relevant. This completes the proof of part (a).

To complete the proof of part (b), it suffices to show that $z^{\prime} \neq z^{\prime \prime}$ if and only if there exists a dual path $\hat{\pi} \in \hat{\mathcal{N}}$ with $\hat{\sigma}_{\hat{\pi}}=U$ such that $\hat{\pi}$ enters $z$. To see this, we note that if $z^{\prime}=z^{\prime \prime}$, then the paths $l$ and $r$ starting at $z=(x, t)$ meet before time $U$. Conversely, if $\tau_{l, r}:=\inf \{u>t: l(u)=r(u)\}<U$, then $l$ and $r$ cannot enter a $S, U$-relevant separation point before time $\tau_{l, r}$, while after time $\tau_{l, r}$, by the arguments above, $l$ and $r$ must lead to the same point in $R$. The statement now follows from the fact that by [SSS09, Lemma 2.7], there exists a dual path $\hat{\pi} \in \hat{\mathcal{N}}$ with $\hat{\sigma}_{\hat{\pi}}=U$ entering $z$ if and only if there exist $l \in \mathcal{W}^{\mathrm{l}}(z)$ and $r \in \mathcal{W}^{\mathrm{r}}(z)$ such that $l<r$ on $(t, U)$.

To prove part (c), set $z_{0}:=(\pi(S), S)$. By Proposition 6.4 (b) there exist unique $l_{0} \in \mathcal{W}^{\mathrm{l}}\left(z_{0}\right)$ and $r_{0} \in \mathcal{W}^{1}\left(z_{0}\right)$ such that $l_{0} \sim_{\text {out }}^{z_{0}} r_{0}$ and $l_{0} \leq \pi \leq r_{0}$, hence by what we have just proved one has $z_{0} \rightarrow l_{0}, r_{0} z_{1}$ for some unique $z_{1}=\left(x_{1}, t_{1}\right) \in R$. If $t_{1}=U$ we are done. Otherwise, by what we have just proved, there exist $l, l^{\prime} \in \mathcal{W}^{l}\left(z_{1}\right), r, r^{\prime} \in \mathcal{W}^{\mathrm{l}}\left(z_{1}\right)$, and $z^{\prime}, z^{\prime \prime} \in R$ such that $z_{0} \rightarrow_{l, r^{\prime}} z^{\prime}$ and $z_{0} \rightarrow_{l^{\prime}, r} z^{\prime \prime}$. By Proposition 4.3 , the path $\pi$ must either turn left or right at $z_{1}$, so setting either $\left(l_{1}, r_{1}\right)=\left(l, r^{\prime}\right)$ or $\left(l_{1}, r_{1}\right)=\left(l^{\prime}, r\right)$ and $z_{2}=\left(x_{2}, t_{2}\right):=z^{\prime}$ or $z^{\prime \prime}$ we have that $l_{1} \leq \pi \leq r_{1}$ on $\left[t_{1}, t_{2}\right]$. Continuing this process, which terminates after a finite number of steps by Proposition 6.1, we find a sequence of points $z_{0}, \ldots, z_{n}$ and paths $l_{0}, r_{0}, \ldots, l_{n-1}, r_{n-1}$ with the desired properties.

Finally, part (d) follows from [SSS09, Thm. 1.12 (d)] which implies that the concatenation of the paths $l_{0}, \ldots, l_{n-1}$ defines a path $\pi \in \mathcal{N}$ with all the desired properties.

We will sometimes need the following extension of Proposition 6.5.
Corollary 6.6 (Steering paths between deterministic times) Let $\mathcal{N}$ be a Brownian net with associated left-right Brownian web $\left(\mathcal{W}^{\mathrm{l}}, \mathcal{W}^{\mathrm{r}}\right)$ and let $-\infty<T_{1}<\cdots<T_{m}<\infty$ be
deterministic times. Set

$$
\begin{align*}
R_{T_{1}} & :=\mathbb{R} \times\left\{T_{1}\right\}, \\
R_{T_{m}} & :=\left\{\left(x, T_{m}\right): x \in \mathbb{R}, \exists \pi \in \mathcal{N} \text { with } \sigma_{\pi}=T_{1} \text { s.t. } \pi\left(T_{m}\right)=x\right\},  \tag{6.4}\\
R_{T_{k}, T_{k+1}} & :=\left\{z \in \mathbb{R}^{2}: z \text { is a } T_{k}, T_{k+1}\right. \text {-relevant separation point, } \\
& \left.\exists \pi \in \mathcal{N} \text { with } \sigma_{\pi}=T_{1} \text { s.t. } \pi \text { enters } z\right\} .
\end{align*}
$$

Then all conclusions of Proposition 6.5 remain valid with $R_{S}$ replaced by $R_{T_{1}}, R_{S, U}$ replaced by $\bigcup_{k=1}^{m-1} R_{T_{k}, T_{k+1}}$, and $R_{U}$ replaced by $R_{T_{m}}$, except that in part (b), it may happen that $z^{\prime}=z^{\prime \prime}$ for some $z \in \bigcup_{k=1}^{m-1} R_{T_{k}, T_{k+1}}$ or $z^{\prime} \neq z^{\prime \prime}$ for some $z \in R_{S}$ even though there is no path $\hat{\pi} \in \hat{\mathcal{N}}$ starting at time $T_{m}$ entering $z$.

Moreover, a.s. for each $\pi, \pi^{\prime} \in \mathcal{N}$ satisfying $\sigma_{\pi}, \sigma_{\pi^{\prime}} \leq T_{1}, \sigma_{\pi} \wedge \sigma_{\pi^{\prime}}<T_{1}, \pi\left(T_{1}\right) \leq \pi^{\prime}\left(T_{1}\right)$ and $\operatorname{sign}_{\pi}(z) \leq \operatorname{sign}_{\pi^{\prime}}(z)$ for all $z \in \bigcup_{k=1}^{m-1} R_{T_{k}, T_{k+1}}$ such that both $\pi$ and $\pi^{\prime}$ enter $z$, one has $\pi\left(T_{k}\right) \leq \pi^{\prime}\left(T_{k}\right)$ for $k=1, \ldots, m$.

Proof. This generalization of Proposition 6.5 follows by 'pasting together' the finite graph representations for the consecutive time intervals $\left[T_{k}, T_{k+1}\right]$, where we use that by Proposition 6.4 (a), if $\pi \in \mathcal{N}$ satisfies $\sigma_{\pi}=T_{1}$, then the points $\left(\pi\left(T_{k}\right), T_{k}\right)(k=2, \ldots, m)$ must be of type ( $\mathrm{p}, \mathrm{p}$ ) in $\mathcal{N}$.

To prove the statement about the paths $\pi, \pi^{\prime}$, by symmetry, we may assume without loss of generality that $\sigma_{\pi}<T_{1}$. In this case, by Proposition 6.4 (a), the point $\left(\pi\left(T_{1}\right), T_{1}\right)$ must be of type $(\mathrm{p}, \mathrm{p})$ in $\mathcal{N}$. Let $z_{i}=\left(x_{i}, t_{i}\right)(i=1, \ldots, n)$ and $l_{i}, r_{i}(i=1, \ldots, n-1)$ be defined for $\pi$ as in Proposition 6.5 (c), and define $\pi_{-}, \pi_{+}:\left[T_{1}, T_{m}\right] \rightarrow \mathbb{R}$ by

$$
\begin{equation*}
\pi_{-}:=l_{i} \quad \text { and } \quad \pi_{+}:=r_{i} \text { on }\left[t_{i}, t_{i+1}\right] \quad(i=1, \ldots, n-1) \tag{6.5}
\end{equation*}
$$

Then $\pi_{-}\left(T_{k}\right)=\pi\left(T_{k}\right)=\pi_{+}\left(T_{k}\right)$ for $k=1, \ldots, m$ and $\pi_{-} \leq \pi^{\prime}$ on $\left[T_{1}, T_{m}\right]$.

### 6.3 Discrete approximation of the Brownian web

In this and the next section, we recall known results about convergence of dicrete webs and nets to Brownian webs and nets and prove some related, new results about convergence of intersection local times and relevant separation points. In Section 6.5, we then use these results to prove a new convergence result about Brownian webs embedded in Brownian nets, which will form the basis for the proof of Theorem 3.9 on the convergence of quenched laws, which will be proved in Section 7.1 .

Before we turn our attention to the details of the Brownian web, we first explain two simple, general principles that we will be using several time in what follows.

## Lemma 6.7 (Weak convergence of coupled random variables)

(a) Let $E$ be a Polish space, let $\left(F_{i}\right)_{i \in I}$ be a finite or countable collection of Polish spaces and for each $i \in I$, let $f_{i}: E \rightarrow F_{i}$ be a measurable function. Let $X, X_{k}, Y_{k, i}$ be random variables $(k \geq 1, i \in I)$ such that $X, X_{k}$ take values in $E$ and $Y_{k, i}$ takes values in $F_{i}$. Then

$$
\begin{align*}
& \mathbb{P}\left[\left(X_{k}, Y_{k, i}\right) \in \cdot\right] \underset{k \rightarrow \infty}{\Longrightarrow} \mathbb{P}\left[\left(X, f_{i}(X)\right) \in \cdot\right] \quad \forall i \in I \\
& \quad \text { implies } \mathbb{P}\left[\left(X_{k},\left(Y_{k, i}\right)_{i \in I}\right) \in \cdot\right] \underset{k \rightarrow \infty}{\Longrightarrow} \mathbb{P}\left[\left(X,\left(f_{i}(X)\right)_{i \in I}\right) \in \cdot\right] \tag{6.6}
\end{align*}
$$

where $\Rightarrow$ denotes weak convergence of probability laws on the Polish spaces $E \times F_{i}$ and $E \times \prod_{i \in I} F_{i}$, respectively .
(b) Let $E, F, G$ be Polish spaces, let $f: E \rightarrow F$ and $g: F \rightarrow G$ be measurable functions, let $X, X_{k}, Y, Y_{k}$ and $Z_{k}$ be random variables taking values in $E, F$ and $G$, respectively ( $k \geq 1$ ). Then

$$
\begin{align*}
& \mathbb{P}\left[\left(X_{k}, Y_{k}\right) \in \cdot\right] \underset{k \rightarrow \infty}{\Longrightarrow} \mathbb{P}[(X, f(X)) \in \cdot] \text { and } \mathbb{P}\left[\left(Y_{k}, Z_{k}\right) \in \cdot\right] \underset{k \rightarrow \infty}{\Longrightarrow} \mathbb{P}[(Y, g(Y)) \in \cdot] \\
& \quad \text { implies } \mathbb{P}\left[\left(X_{k}, Y_{k}, Z_{k}\right) \in \cdot\right] \underset{k \rightarrow \infty}{\Longrightarrow} \mathbb{P}[(X, f(X), g(f(X))) \in \cdot] . \tag{6.7}
\end{align*}
$$

where $\Rightarrow$ denotes weak convergence of probability laws on the Polish spaces $E \times F, F \times G$, and $E \times F \times G$, respectively.

Proof. For part (a), we observe that the assumed weak convergence of ( $X_{k}, Y_{k, i}$ ) for each $i \in I$ implies tightness of the laws of $\left(X_{k},\left(Y_{k, i}\right)_{i \in I}\right)$. Let $\left(X,\left(Y_{i}\right)_{i \in I}, \ldots\right)$ be any weak subsequential limit. Then $\left(X, Y_{i}\right)$ is equally distributed with $\left(X, f_{i}(X)\right)$, hence $Y_{i}=f_{i}(X)$ a.s. for each $i \in I$. Similarly, in the set-up of part (b), the weak convergence of $\left(X_{k}, Y_{k}\right)$ and $\left(Y_{k}, Z_{k}\right)$ implies tightness of the laws of $\left(X_{k}, Y_{k}, Z_{k}\right)$, while for each weak subsequential limit $(X, Y, Z)$ one has $Y=f(X)$ and $Z=g(Y)$ a.s.

We note that by Skorohod's representation theorem (see e.g. [Bil99, Theorem 6.7]) the left-hand side of (6.6) implies that for each $i \in I$, we can find a coupling of the $X_{k}, Y_{k, i}$ and $X$ such that $\left(X_{k}, Y_{k, i}\right) \rightarrow\left(X, f_{i}(X)\right)$ a.s. By the right-hand side of 6.6), we can find a coupling that works for all $i \in I$ simultaneously. We will apply this principle many times, e.g. when $X$ is a Browian web, $f(X)$ is its associated dual Brownian web, $g(X)$ is the set of paths starting at a given time etc. We will not always be explicit in our choice of the measurable maps $f, g$ but it is clear from the context that they can be constructed.

Recall from Section 3.1 that each i.i.d. collection $\alpha=\left(\alpha_{z}\right)_{z \in \mathbb{Z}_{\text {evenn }}^{2}}$ of $\{-1,+1\}$-valued random variables defines a discrete web $\mathcal{U}^{\alpha}=\left\{p_{z}^{\alpha}: z \in \mathbb{Z}_{\text {even }}^{2}\right\}$ as in (3.2). As in Section 3.5, by linear interpolation and by adding trivial paths that are constantly $-\infty$ or $+\infty$, we view $\mathcal{U}^{\alpha}$ as a random compact subset of the space of paths $\Pi$ introduced in Section 3.2.

Let $\mathbb{Z}_{\text {odd }}^{2}:=\{(x, t): x, t \in \mathbb{Z}, x+t$ is odd $\}$ be the odd sublattice of $\mathbb{Z}^{2}$. For each $(x, s) \in$ $\mathbb{Z}_{\text {odd }}^{2}$, we let $\hat{p}_{(x, s)}^{\alpha}=\hat{p}$, defined by (compare 3.1p)

$$
\begin{equation*}
\hat{p}(s):=x \quad \text { and } \quad \hat{p}(t-1):=\hat{p}(t)-\alpha_{(\hat{p}(t), t-1)} \quad(t \leq s) \tag{6.8}
\end{equation*}
$$

denote the dual path started at $(x, s)$ and we let $\hat{\mathcal{U}}^{\alpha}=\left\{\hat{p}_{z}^{\alpha}: z \in \mathbb{Z}_{\text {even }}^{2}\right\}$ denote the dual discrete web associated with $\mathcal{U}^{\alpha}$. We view $\hat{\mathcal{U}}^{\alpha}$ as a random compact subset of the space of dual paths $\hat{\Pi}$. In line with earlier notation, for any $A \subset \mathbb{Z}_{\text {even }}^{2}$ (resp. $A \subset \mathbb{Z}_{\text {odd }}^{2}$ ), we let $\mathcal{U}^{\alpha}(A)$ (resp. $\left.\hat{\mathcal{U}}^{\alpha}(A)\right)$ denote the set of paths in $\mathcal{U}^{\alpha}$ (resp. $\hat{\mathcal{U}}^{\alpha}$ ) starting from $A$. We define diffusive scaling maps $S_{\varepsilon}$ as in (3.25) and use $S_{\varepsilon}\left(\mathcal{A}_{1}, \ldots, \mathcal{A}_{n}\right)$ as a shorthand for $\left(S_{\varepsilon}\left(\mathcal{A}_{1}\right), \ldots, S_{\varepsilon}\left(\mathcal{A}_{n}\right)\right)$.

The following result follows easily from [FINR04, Theorem 6.1] on the convergence of discrete webs to the Brownian web and Proposition 3.1 on the characterization of the dual Brownian web.

Theorem 6.8 (Convergence to the double Brownian web) Let $\varepsilon_{k}$ be positive constants, tending to zero. For each $k$, let $\alpha^{\langle k\rangle}=\left(\alpha_{z}^{\langle k\rangle}\right)_{z \in \mathbb{Z}_{\text {even }}^{2}}$ be an i.i.d. collection of $\{-1,+1\}$-valued random variables, let $\mathcal{U}_{\langle k\rangle}:=\mathcal{U}^{\alpha^{\langle k\rangle}}$ and $\hat{\mathcal{U}}_{\langle k\rangle}:=\hat{\mathcal{U}}^{\alpha^{(k)}}$ be the discrete web and dual discrete web associated with $\alpha^{\langle k\rangle}$, and assume that $\lim _{k \rightarrow \infty} \varepsilon_{k}^{-1} \mathbb{E}\left[\alpha_{z}^{\langle k\rangle}\right]=\beta$ for some $\beta \in \mathbb{R}$. Then

$$
\begin{equation*}
\mathbb{P}\left[S_{\varepsilon_{k}}\left(\mathcal{U}_{\langle k\rangle}, \hat{\mathcal{U}}_{\langle k\rangle}\right) \in \cdot\right] \underset{k \rightarrow \infty}{\Longrightarrow} \mathbb{P}[(\mathcal{W}, \hat{\mathcal{W}}) \in \cdot], \tag{6.9}
\end{equation*}
$$

where $\Rightarrow$ denotes weak convergence of probability laws on $\mathcal{K}(\Pi) \times \mathcal{K}(\hat{\Pi}), \mathcal{W}$ is a Brownian web with drift $\beta$ and $\hat{\mathcal{W}}$ is its dual.

For notational convenience, let us write

$$
\begin{equation*}
\Sigma_{T}:=\left\{(x, t) \in R_{\mathrm{c}}^{2}: t=T\right\} \tag{6.10}
\end{equation*}
$$

so that, e.g., $\mathcal{N}\left(\Sigma_{T}\right)=\left\{\pi \in \mathcal{N}: \sigma_{\pi}=T\right\}$. We use similar notation for sets of discrete paths. The following strengthening of Theorem 6.8 is sometimes handy.

Lemma 6.9 (Convergence of paths started at given times) In the setup of Theorem 6.8. let $T_{k} \in \mathbb{Z} \cup\{-\infty, \infty\}$ be times such that $\varepsilon_{k}^{2} T_{k} \rightarrow T$ for some $T \in[-\infty,+\infty]$. Then

$$
\begin{equation*}
\mathbb{P}\left[S_{\varepsilon_{k}}\left(\mathcal{U}_{\langle k\rangle}, \mathcal{U}_{\langle k\rangle}\left(\Sigma_{T_{k}}\right)\right) \in \cdot\right] \underset{k \rightarrow \infty}{\Longrightarrow} \mathbb{P}\left[\left(\mathcal{W}, \mathcal{W}\left(\Sigma_{T}\right)\right) \in \cdot\right] . \tag{6.11}
\end{equation*}
$$

Proof. It follows from the tightness of the $S_{\varepsilon_{k}}\left(\mathcal{U}_{\langle k\rangle}\right)$ and Lemma B. 4 that also the laws of the $\mathcal{K}(\Pi)^{2}$-valued random variables $S_{\varepsilon_{k}}\left(\mathcal{U}_{\langle k\rangle}, \mathcal{U}_{\langle k\rangle}\left(K_{k}\right)\right)$ are tight. By going to a subsequence if necessary and invoking Skorohod's representation theorem, we may assume that they converge to an a.s. limit $(\mathcal{W}, \mathcal{A})$. It is easy to see that $\mathcal{A} \subset \mathcal{W}\left(\Sigma_{T}\right)$. If $T= \pm \infty$, then $\mathcal{W}\left(\Sigma_{T}\right)$ contains only trivial paths and it is easy to check that also $\mathcal{A} \supset \mathcal{W}\left(\Sigma_{T}\right)$. To get this inclusion for $-\infty<T<\infty$, let $\mathcal{D}$ be a deterministic countable dense subset of $\mathbb{R} \times\{T\}$. Since $\mathcal{A}(z)$ is nonempty for each $z \in \mathcal{D}$ and since $\mathcal{W}(z)$ contains a single path for each $z \in \mathcal{D}$, we conclude that $\mathcal{A} \supset \mathcal{W}(\mathcal{D})$. Since $\mathcal{A}$ is compact and $\mathcal{W}\left(\Sigma_{T}\right)$ is the closure of $\mathcal{W}(\mathcal{D})$, it follows that $\mathcal{A}=\mathcal{W}\left(\Sigma_{T}\right)$.

We next formulate a result which says that the intersection local time measure $\ell$ between a forward and dual Brownian web as defined in Proposition 3.4 is the limit of the intersection local time measures between approximating forward and dual discrete webs. Since $\ell$ is locally infinite, such a statement on its own cannot make sense. Rather, we will show that the restriction of $\ell$ to the intersection of finitely many forward and dual paths is a.s. the weak limit of the analogue discrete object.

For any $K \subset \mathcal{K}\left(R_{\mathrm{c}}^{2}\right)$, we let

$$
\begin{equation*}
\operatorname{Img}(K):=\left\{z \in R_{\mathrm{c}}^{2}: \exists A \in K \text { s.t. } z \in A\right\} \tag{6.12}
\end{equation*}
$$

denote the union of all sets in $K$. We call $\operatorname{Img}(K)$ the image set (or trace) of $K$. In particular, if $\mathcal{A}$ is a set of paths (which, as usual, we identify with their graphs), then $\operatorname{Img}(\mathcal{A})=\{(\pi(t), t)$ : $\left.t \geq \sigma_{\pi}, \pi \in \mathcal{A}\right\}$. Similarly, if $\mathcal{A}$ is a set of discrete paths, then

$$
\begin{equation*}
\operatorname{Img}(\mathcal{A}):=\left\{(x, t) \in \mathbb{Z}_{\text {even }}^{2}: t \geq \sigma_{\pi}, \pi \in \mathcal{A}\right\} \tag{6.13}
\end{equation*}
$$

and we use similar notation for a set $\hat{\mathcal{A}}$ of discrete dual paths, where in this case $\operatorname{Img}(\hat{\mathcal{A}})$ is a subset of $\mathbb{Z}_{\text {odd }}^{2}$.

Proposition 6.10 (Convergence of intersection local time) Let $\varepsilon_{k}$ be positive constants, tending to zero. Let $\alpha^{\langle k\rangle}$ be collections of i.i.d. $\{-1,+1\}$-valued random variables satisfying $\lim _{k \rightarrow \infty} \varepsilon_{k}^{-1} \mathbb{E}\left[\alpha_{z}^{\langle k\rangle}\right]=\beta$ for some $\beta \in \mathbb{R}$, and let $\mathcal{U}_{\langle k\rangle}, \hat{\mathcal{U}}_{\langle k\rangle}$ be the discrete web and its dual associated with $\alpha^{\langle k\rangle}$. Let $\mathcal{W}$ be a Brownian web with drift $\beta$, $\hat{\mathcal{W}}$ be its dual, $\ell$ be the intersection
local time measure between $\mathcal{W}$ and $\hat{\mathcal{W}}$, and $\ell_{\mathrm{r}}$ be the restriction of $\ell$ to the set of points of type $(1,2)_{\mathrm{r}}$. Let

$$
\begin{align*}
\Delta_{k} & =\left\{z_{1}^{k}, \ldots, z_{m}^{k}\right\} \subset \mathbb{Z}_{\text {even }}^{2}, & \hat{\Delta}_{k} & =\left\{\hat{z}_{1}^{k}, \ldots, \hat{z}_{n}^{k}\right\} \subset \mathbb{Z}_{\text {odd }}^{2},  \tag{6.14}\\
\Delta & =\left\{z_{1}, \ldots, z_{m}\right\} \subset \mathbb{R}^{2}, & & \hat{\Delta}
\end{align*}=\left\{\hat{z}_{1}, \ldots, \hat{z}_{n}\right\} \subset \mathbb{R}^{2}, ~ l
$$

be finite sets such that $S_{\varepsilon_{k}}\left(z_{i}^{k}\right) \rightarrow z_{i}$ and $S_{\varepsilon_{k}}\left(\hat{z}_{j}^{k}\right) \rightarrow \hat{z}_{j}$ as $k \rightarrow \infty$ for each $i=1, \ldots, m$ and $j=1, \ldots, n$. Set

$$
\ell_{\mathrm{r}}^{\langle k\rangle}:=\varepsilon_{k} \sum_{z \in I_{k} \cap Z_{\mathrm{r}}^{(k)}} \delta_{S_{\varepsilon_{k}}(z)},
$$

where

$$
\begin{align*}
Z_{\mathrm{r}}^{\langle k\rangle} & :=\left\{z \in \mathbb{Z}_{\text {even }}^{2}: \alpha_{(x, t)}^{\langle k\rangle}=+1\right\},  \tag{6.15}\\
I_{k} & :=\left\{(x, t) \in \mathbb{Z}_{\text {even }}^{2}:(x, t) \in \operatorname{Img}\left(\mathcal{U}_{\langle k\rangle}\left(\Delta_{k}\right)\right),(x, t+1) \in \operatorname{Img}\left(\hat{\mathcal{U}}_{\langle k\rangle}\left(\hat{\Delta}_{k}\right)\right)\right\} .
\end{align*}
$$

Let $\ell_{\mathrm{r}}(\Delta, \hat{\Delta})$ denote the restriction of $\ell_{\mathrm{r}}$ to the set $\operatorname{Img}(\mathcal{W}(\Delta)) \cap \operatorname{Img}(\hat{\mathcal{W}}(\hat{\Delta}))$. Then

$$
\begin{equation*}
\mathbb{P}\left[\left(S_{\varepsilon_{k}}\left(\mathcal{U}_{\langle k\rangle}\right), \ell_{\mathrm{r}}^{\langle k\rangle}\right) \in \cdot\right] \underset{k \rightarrow \infty}{\Longrightarrow} \mathbb{P}\left[\left(\mathcal{W}, \ell_{\mathrm{r}}(\Delta, \hat{\Delta})\right) \in \cdot\right] \tag{6.16}
\end{equation*}
$$

where $\Rightarrow$ denotes weak convergence of probability laws on $\mathcal{K}(\Pi) \times \mathcal{M}\left(\mathbb{R}^{2}\right)$, and $\mathcal{M}\left(\mathbb{R}^{2}\right)$ is the space of finite measures on $\mathbb{R}^{2}$ equipped with the topology of weak convergence.

Proof. Since

$$
\mathbb{P}\left[S_{\varepsilon_{k}}\left(\mathcal{U}_{\langle k\rangle}, \hat{\mathcal{U}}_{\langle k\rangle}, \mathcal{U}_{\langle k\rangle}\left(\Delta_{k}\right), \hat{\mathcal{U}}_{\langle k\rangle}\left(\hat{\Delta}_{k}\right)\right) \in \cdot\right] \underset{k \rightarrow \infty}{\Longrightarrow} \mathbb{P}[(\mathcal{W}, \hat{\mathcal{W}}, \mathcal{W}(\Delta), \hat{\mathcal{W}}(\hat{\Delta})) \in \cdot]
$$

and a.s. $(\mathcal{W}(\Delta), \hat{\mathcal{W}}(\hat{\Delta}))$ determines $\ell_{\mathrm{r}}(\Delta, \hat{\Delta})$, by Lemma 6.7 (b), proving 6.16) reduces to proving

$$
\begin{equation*}
\left.\mathbb{P}\left[\left(S_{\varepsilon_{k}}\left(\mathcal{U}_{\langle k\rangle}\left(\Delta_{k}\right), \hat{\mathcal{U}}_{\langle k\rangle}\left(\hat{\Delta}_{k}\right)\right), \ell_{\mathrm{r}}^{\langle k\rangle}\right) \in \cdot\right] \underset{k \rightarrow \infty}{\Longrightarrow} \mathbb{P}\left[(\mathcal{W}(\Delta), \hat{\mathcal{W}}(\hat{\Delta})), \ell_{\mathrm{r}}(\Delta, \hat{\Delta})\right) \in \cdot\right] . \tag{6.17}
\end{equation*}
$$

We will make a further reduction.
For $k \in \mathbb{N}, 1 \leq i \leq m$ and $1 \leq j \leq n$, let $\mathcal{U}_{\langle k\rangle}\left(z_{i}^{k}\right)=\left\{p_{i}^{k}\right\}$ and $\hat{\mathcal{U}}_{\langle k\rangle}\left(\hat{z}_{j}^{k}\right)=\left\{\hat{p}_{j}^{k}\right\}$. Let $t_{i}^{k}$ and $\hat{t}_{j}^{k}$ denote respectively the starting time of $p_{i}^{k}$ and $\hat{p}_{j}^{k}$. Similarly, let $\mathcal{W}\left(z_{i}\right)=\left\{\pi_{i}\right\}$ and $\hat{\mathcal{W}}\left(\hat{z}_{j}\right)=\left\{\hat{\pi}_{j}\right\}$, with starting time $t_{i}$ for $\pi_{i}$ and $\hat{t}_{j}$ for $\hat{\pi}_{j}$. Let $\left(\tau_{u v}^{k}\right)_{1 \leq u<v \leq m}$ be the time of coalescence between $p_{u}^{k}$ and $p_{v}^{k}$, and let $\left(\hat{\tau}_{u v}^{k}\right)_{1 \leq u<v \leq n}$ be the time of coalescence between $\hat{p}_{u}^{k}$ and $\hat{p}_{v}^{k}$. Define $\left(\tau_{u v}\right)_{1 \leq u<v \leq m}$ and $\left(\hat{\tau}_{u v}\right)_{1 \leq u<v \leq n}$ similarly for $\mathcal{W}(\Delta)$ and $\hat{\mathcal{W}}(\hat{\Delta})$. For $1 \leq u \leq m$ and $1 \leq v \leq n$, let

$$
\ell_{\mathrm{r}, u v}^{\langle k\rangle}:=\varepsilon_{k} \sum_{z:=(x, t) \in \mathbb{Z}_{\text {even }}^{2}} \delta_{S_{\varepsilon_{k}}(z)} 1_{\left\{p_{u}^{k}(t)=\hat{p}_{v}^{k}(t+1)=\hat{p}_{v}^{k}(t)+1=x\right\}}
$$

be the rescaled intersection local time measure of $p_{u}^{k}$ and $\hat{p}_{v}^{k}$ on points with $\alpha_{(x, t)}^{\langle k\rangle}=1$, and similarly let $\ell_{\mathrm{r}, u v}$ be the intersection local time measure of $\pi_{u}$ and $\hat{\pi}_{v}$ on points of type $(1,2)_{\mathrm{r}}$. We note that $\ell_{\mathrm{r}}(\Delta, \hat{\Delta})$ can be uniquely constructed from $\left(\tau_{u v}\right)_{1 \leq u<v \leq m},\left(\hat{\tau}_{u v}\right)_{1 \leq u<v \leq n}$, and $\left(\ell_{r, u v}\right)_{1 \leq u \leq m, 1 \leq v \leq n}$. For example, we can go through the indices $(u v)_{1 \leq u \leq m, 1 \leq v}$ in numeric order, and at each step, we add the proper restriction of $\ell_{\mathrm{r}, u v}$ to the construction of $\ell_{\mathrm{r}}(\Delta, \hat{\Delta})$ so as to exclude overlaps among $\left(\ell_{\mathrm{r}, u v}\right)_{1 \leq u \leq m, 1 \leq v \leq n}$ due to the coalescence of paths. By the same
procedure, $\ell_{\mathrm{r}}^{\langle k\rangle}$ can be constructed from $\left(\tau_{u v}^{k}\right)_{1 \leq u<v \leq m},\left(\hat{\tau}_{u v}^{k}\right)_{1 \leq u<v \leq n}$, and $\left(\ell_{\mathrm{r}, u v}^{\langle k\rangle}\right)_{1 \leq u \leq m, 1 \leq n \leq v}$. To prove (6.17), it then suffices to prove

$$
\begin{align*}
& \mathbb{P}\left[\left(S_{\varepsilon_{k}}\left(\left(p_{i}^{k}\right)_{1 \leq i \leq m},\left(\hat{p}_{j}^{k}\right)_{1 \leq j \leq n}\right),\left(\varepsilon_{k}^{2} \tau_{u v}^{k}\right)_{1 \leq u<v \leq m},\left(\varepsilon_{k}^{2} \hat{\tau}_{u v}^{k}\right)_{1 \leq u<v \leq n},\left(\ell_{\mathrm{r}, u v}^{(k\rangle}\right)_{1 \leq u \leq m, 1 \leq n \leq v}\right) \in \cdot\right] \\
\Longrightarrow & \mathbb{P}\left[\left(\left(\pi_{i}\right)_{1 \leq i \leq m},\left(\hat{\pi}_{j}\right)_{1 \leq j \leq n},\left(\tau_{u v}\right)_{1 \leq u<v \leq m},\left(\hat{\tau}_{u v}\right)_{1 \leq u<v \leq n},\left(\ell_{\mathrm{r}, u v}\right)_{1 \leq u \leq m, 1 \leq v \leq n}\right) \in \cdot\right] . \tag{6.18}
\end{align*}
$$

It has been shown in the proof of [STW00, Thm. 8] that

$$
\begin{align*}
& \mathbb{P}\left[S_{\varepsilon_{k}}\left(\left(p_{i}^{k}\right)_{1 \leq i \leq m},\left(\hat{p}_{j}^{k}\right)_{1 \leq j \leq n},\left(\tau_{u v}^{k}\right)_{1 \leq u<v \leq m},\left(\hat{\tau}_{u v}^{k}\right)_{1 \leq u<v \leq n}\right) \in \cdot\right] \\
\Longrightarrow & \mathbb{P}\left[\left(\left(\pi_{i}\right)_{1 \leq i \leq m},\left(\hat{\pi}_{j}\right)_{1 \leq j \leq n},\left(\tau_{u v}\right)_{1 \leq u<v \leq m},\left(\hat{\tau}_{u v}\right)_{1 \leq u<v \leq n}\right) \in \cdot\right], \tag{6.19}
\end{align*}
$$

where $\left(p_{i}^{k}\right)_{1 \leq i \leq m}$ and $\left(\hat{p}_{j}^{k}\right)_{1 \leq j \leq n}$ were constructed as a deterministic transformation (via Skorohod reflection and colaescence) of a collection of independent random walks ( $\left.W_{i}^{k}\right)_{1 \leq i \leq m}$ and $\left(\hat{W}_{j}^{k}\right)_{1 \leq j \leq n}$. The same transformation was used to construct $\left(\pi_{i}\right)_{1 \leq i \leq m}$ and $\left(\hat{\pi}_{j}\right)_{1 \leq j \leq n}$ from a collection of independent Brownian motions $\left(B_{i}\right)_{1 \leq i \leq m}$ and $\left(\hat{B}_{j}\right)_{1 \leq j \leq n}$. Furthermore, this transformation together with the times of coalescence $\left(\tau_{u v}\right)_{1 \leq u<v<m}$ and $\left(\hat{\tau}_{u v}\right)_{1 \leq u<v \leq n}$ are a.s. continuous in $\left(B_{i}\right)_{1 \leq i \leq m}$ and $\left(\hat{B}_{j}\right)_{1 \leq j \leq n}$. The convergence in 6.19) then follows from Donsker's invariance principle. Since $\ell_{\mathrm{r}, u v}$ is uniquely determined by $\pi_{u}$ and $\hat{\pi}_{v}$, by Lemma 6.7, to prove (6.18), it then suffices to show that for each $1 \leq u \leq m$ and $1 \leq v \leq m$,

$$
\begin{equation*}
\mathbb{P}\left[\left(S_{\varepsilon_{k}}\left(p_{u}^{k}, \hat{p}_{v}^{k}\right), \ell_{\mathrm{r}, u v}^{\langle k\rangle}\right) \in \cdot\right] \underset{k \rightarrow \infty}{\Longrightarrow} \mathbb{P}\left[\left(\pi_{u}, \hat{\pi}_{v}, \ell_{\mathrm{r}, u v}\right) \in \cdot\right] \tag{6.20}
\end{equation*}
$$

Without loss of generality, we may assume $u=v=1$ in 6.20). We may also assume that $z_{1}=\left(x_{1}, t_{1}\right)$ and $\hat{z}_{1}=\left(\hat{x}_{1}, \hat{t}_{1}\right)$ satisfy $t_{1}<\hat{t}_{1}$, so that $\ell_{\mathrm{r}, 11}$ is not a.s. the zero measure, in which case (6.20) is trivial. We recall from [TW00, Thm. 3] (see also [FINR06, (3.6) and Thm. 3.7]) that, conditional on $\hat{\pi}_{1}, \pi_{1}$ is distributed as an independent Brownian motion $B_{1}$ with drift $\beta$, starting from $z_{1}$, and Skorohod reflected away from $\hat{\pi}_{1}$. More precisely, $\pi_{1}$ admits the representation

$$
\pi_{1}(t)= \begin{cases}B_{1}(t)+L_{\mathrm{r}}(t) & \text { if } \hat{\pi}_{1}\left(t_{1}\right)<x_{1},  \tag{6.21}\\ B_{1}(t)-L_{1}(t) & \text { if } x_{1}<\hat{\pi}_{1}\left(t_{1}\right),\end{cases}
$$

where

$$
\begin{array}{lll}
L_{\mathrm{r}}(t)=\sup _{t_{1} \leq s \leq t} \max \left\{0, \hat{\pi}_{1}(s)-B_{1}(s)\right\} & \text { for } t \in\left[t_{1}, \hat{t}_{1}\right], & L_{\mathrm{r}}(t)=L_{\mathrm{r}}\left(\hat{t}_{1}\right) \text { for } t \geq \hat{t}_{1}, \\
L_{\mathrm{l}}(t)=\sup _{t_{1} \leq s \leq t} \max \left\{0, B_{1}(s)-\hat{\pi}_{1}(s)\right\} & \text { for } t \in\left[t_{1}, \hat{t}_{1}\right], & L_{\mathrm{l}}(t)=L_{\mathrm{l}}\left(\hat{t}_{1}\right) \text { for } t \geq \hat{t}_{1} \tag{6.22}
\end{array}
$$

It was shown in NRS10, Prop. 3.1] and its proof that, with the construction of $\pi_{1}$ as in (6.21), almost surely $\ell_{\mathrm{r}, 11}(\mathbb{R} \times \cdot)=\mathrm{d} L_{\mathrm{r}}(\cdot)$, or equivalently,

$$
\begin{equation*}
\ell_{\mathrm{r}, 11}\left(\mathbb{R} \times\left[t_{1}, t\right]\right)=L_{\mathrm{r}}(t) \quad \text { for all } t \in\left[t_{1}, \hat{t}_{1}\right] . \tag{6.23}
\end{equation*}
$$

Since $\ell_{\mathrm{r}, 11}$ is concentrated on the graph of $\pi_{1}$, it follows that

$$
\begin{equation*}
\ell_{\mathrm{r}, 11}=\mathrm{d} L_{r} \circ \pi_{1}^{-1}, \tag{6.24}
\end{equation*}
$$

[^3]i.e., $\ell_{\mathrm{r}, 11}$ is the image of the measure $\mathrm{d} L_{r}$ under the map $\pi_{1}$.

There is a similar representation for $p_{1}^{k}, \hat{p}_{1}^{k}$, and $\ell_{\mathrm{r}, 11}^{\langle k}$. Indeed, if $W_{1}^{k}$ is an independent simple random walk on $\mathbb{Z}$ with drift $\mathbb{E}\left[\alpha_{z}^{\langle k\rangle}\right]$ and starting from $z_{1}^{k}$, then conditional on $\hat{p}_{1}^{k}$, we can construct $p_{1}^{k}$ as (see e.g. [STW00, Sec. 2.2.2] or the proof of [SSS09, Lemma 2.1])

$$
p_{1}^{k}(t)= \begin{cases}W_{1}^{k}(t)+L_{\mathrm{r}}^{k}(t) & \text { if } \hat{p}_{1}^{k}\left(t_{1}^{k}\right)<x_{1}^{k},  \tag{6.25}\\ W_{1}^{k}(t)-L_{1}^{k}(t) & \text { if } x_{1}^{k}<\hat{p}_{1}^{k}\left(t_{1}^{k}\right),\end{cases}
$$

where

$$
\begin{array}{ll}
L_{\mathrm{r}}^{k}(t)=\sup _{t_{1}^{k} \leq s \leq t} \max \left\{0,1+\hat{p}_{1}^{k}(s)-W_{1}^{k}(s)\right\} \text { for } t \in\left[t_{1}^{k}, \hat{t}_{1}^{k}\right], & L_{\mathrm{r}}^{k}(t)=L_{\mathrm{r}}^{k}\left(\hat{t}_{1}^{k}\right) \text { for } t \geq \hat{t}_{1}^{k},  \tag{6.26}\\
L_{1}^{k}(t)=\sup _{t_{1}^{k} \leq s \leq t} \max \left\{0,1+W_{1}^{k}(s)-\hat{p}_{1}^{k}(s)\right\} \text { for } t \in\left[t_{1}^{k}, \hat{t}_{1}^{k}\right], & L_{1}^{k}(t)=L_{1}^{k}\left(\hat{t}_{1}^{k}\right) \text { for } t \geq \hat{t}_{1}^{k}
\end{array}
$$

The constant 1 arises because $\hat{p}_{1}^{k}$ is a walk on $\mathbb{Z}_{\text {odd }}^{2}$, and $W_{1}^{k}$ is a walk on $\mathbb{Z}_{\text {even }}^{2}$. Let

$$
\bar{L}_{\mathrm{r}}^{k}(t)=\sum_{i=t_{1}^{k}}^{t} 1_{\left\{p_{1}^{k}(i)=\hat{p}_{1}^{k}(i+1)=\hat{p}^{k}(i)+1\right\}} \quad \text { for } t_{1}^{k} \leq t<\hat{t}_{1}^{k}
$$

Then, in analogy with (6.24),

$$
\begin{equation*}
\ell_{\mathrm{r}, 11}^{\langle k\rangle}=\mathrm{d}\left(S_{\varepsilon_{k}} \bar{L}_{\mathrm{r}}^{k}\right) \circ\left(p_{1}^{k}\right)^{-1} \tag{6.27}
\end{equation*}
$$

where $\mathrm{d}\left(S_{\varepsilon_{k}} \bar{L}_{\mathrm{r}}^{k}\right)(t)=\varepsilon_{k} \mathrm{~d} \overline{\mathrm{~L}}_{\mathrm{r}}^{k}\left(\varepsilon_{k}^{-2} t\right)$. To relate $L_{\mathrm{r}}^{k}$ and $\bar{L}_{\mathrm{r}}^{k}$, note that conditional on $\hat{p}_{1}^{k}$,

$$
\begin{equation*}
L_{\mathrm{r}}^{k}(t+1)-\left(1-\mathbb{E}\left[\alpha_{z}^{\langle k\rangle}\right]\right) \bar{L}_{\mathrm{r}}^{k}(t)=\sum_{i=t_{1}^{k}}^{t}\left(L_{\mathrm{r}}^{k}(i+1)-L_{\mathrm{r}}^{k}(i)-\left(1-\mathbb{E}\left[\alpha_{z}^{\langle k\rangle}\right]\right) 1_{\left\{\bar{L}_{\mathrm{r}}^{k}(i)-\bar{L}_{\mathrm{r}}^{k}(i-1)=1\right\}}\right) \tag{6.28}
\end{equation*}
$$

is a martingale, because $L_{\mathrm{r}}^{k}(i+1)-L_{\mathrm{r}}^{k}(i) \neq 0$ only when $p_{1}^{k}(i)=\hat{p}_{1}^{k}(i+1)=\hat{p}^{k}(i)+1$, and conditional on the later event, $L_{\mathrm{r}}^{k}(i+1)-L_{\mathrm{r}}^{k}(i)=0$ with probability $\frac{1+\mathbb{E}\left[\alpha_{z}^{\langle k\rangle}\right]}{2}$ and $L_{\mathrm{r}}^{k}(i+1)-L_{\mathrm{r}}^{k}(i)=2$ with probability $\frac{1-\mathbb{E}\left[\alpha_{z}^{\langle k\rangle}\right]}{2}$. By Doob's maximal inequality, conditional on $\hat{p}_{1}^{k}$,

$$
\begin{align*}
\mathbb{E}\left[\sup _{t_{1}^{k} \leq t<\hat{t}_{1}^{k}}\left|L_{\mathrm{r}}^{k}(t+1)-\left(1-\mathbb{E}\left[\alpha_{z}^{\langle k\rangle}\right]\right) \bar{L}_{\mathrm{r}}^{k}(t)\right|^{2}\right] & \leq 4 \mathbb{E}\left[\left|L_{\mathrm{r}}^{k}\left(\hat{t}_{1}^{k}\right)-\left(1-\mathbb{E}\left[\alpha_{z}^{\langle k\rangle}\right]\right) \bar{L}_{\mathrm{r}}^{k}\left(\hat{t}_{1}^{k}\right)\right|^{2}\right] \\
& \leq 16 \mathbb{E}\left[\bar{L}_{\mathrm{r}}^{k}\left(\hat{t}_{1}^{k}\right)\right] \tag{6.29}
\end{align*}
$$

We are now ready to prove 6.20. By 6.24) and 6.27), we can replace $\ell_{\mathrm{r}, 11}$ by $L_{r}$, and $\ell_{\mathrm{r}, 11}^{\langle k\rangle}$ by $S_{\varepsilon_{k}} \bar{L}_{\mathrm{r}}^{k}$. First let us extend the definition of all processes in discrete time to continuous time by linear interpolation. Note that (6.25) and (6.26) remain valid. By Donsker's invariance principle, the pair of independent processes $S_{\varepsilon_{k}}\left(\hat{p}_{1}^{k}, W_{1}^{k}\right)$ converge weakly to ( $\hat{\pi}_{1}, B_{1}$ ). By Skorohod's representation, we may assume from now on this convergence is almost sure by suitable coupling. If $x_{1}<\hat{\pi}_{1}\left(t_{1}\right)$, then trivially

$$
\begin{equation*}
S_{\varepsilon_{k}}\left(p_{1}^{k}, \hat{p}_{1}^{k}, L_{\mathrm{r}}^{k}\right) \underset{k \rightarrow \infty}{\longrightarrow}\left(\pi_{1}, \hat{\pi}_{1}, L_{r}\right) \tag{6.30}
\end{equation*}
$$

since $L_{\mathrm{r}}=0$, and so is $L_{\mathrm{r}}^{k}$ for all $k$ large. If $\hat{\pi}_{1}\left(t_{1}\right)<x_{1}$, then $S_{\varepsilon_{k}}\left(L_{\mathrm{r}}^{k}\right) \rightarrow L_{\mathrm{r}}$ uniformly on compacts, because the Skorohod reflection map which defines both $L_{\mathrm{r}}^{k}$ and $L_{\mathrm{r}}$ in 6.26 and respectively (6.22) is continuous in its arguments. Therefore,

$$
\begin{equation*}
\mathbb{P}\left[S_{\varepsilon_{k}}\left(p_{1}^{k}, \hat{p}_{1}^{k}, L_{\mathrm{r}}^{k}\right) \in \cdot\right] \underset{k \rightarrow \infty}{\Longrightarrow} \mathbb{P}\left[\left(\pi_{1}, \hat{\pi}_{1}, L_{\mathrm{r}}\right) \in \cdot\right] \tag{6.31}
\end{equation*}
$$

On the other hand, by (6.29),

$$
\begin{equation*}
\mathbb{E}\left[\sup _{t_{1}^{k} \leq t<t_{1}^{\hat{k}}} \varepsilon_{k}^{2}\left|L_{\mathrm{r}}^{k}(t+1)-\left(1-\mathbb{E}\left[\alpha_{z}^{\langle k\rangle}\right]\right) \bar{L}_{\mathrm{r}}^{k}(t)\right|^{2}\right] \leq 16 \varepsilon_{k}^{2} \mathbb{E}\left[\bar{L}_{\mathrm{r}}^{k}\left(\hat{t}_{1}^{k}\right)\right] \underset{k \rightarrow \infty}{\longrightarrow} 0 \tag{6.32}
\end{equation*}
$$

because the above inequality implies by triangle inequality that

$$
\left(1-\mathbb{E}\left[\alpha_{z}^{\langle k\rangle}\right]\right) \mathbb{E}\left[\left(\varepsilon_{k} \bar{L}_{\mathrm{r}}^{k}\left(\hat{t}_{1}^{k}\right)\right)^{2}\right]^{\frac{1}{2}} \leq \mathbb{E}\left[\left(\varepsilon_{k} L_{\mathrm{r}}^{k}\left(\hat{t}_{1}^{k}\right)\right)^{2}\right]^{\frac{1}{2}}+4 \sqrt{\varepsilon_{k} \mathbb{E}}\left[\varepsilon_{k} \bar{L}_{\mathrm{r}}^{k}\left(\hat{t}_{1}^{k}\right)\right]^{\frac{1}{2}}
$$

and since $\mathbb{E}\left[\left(\varepsilon_{k} L_{\mathrm{r}}^{k}\left(\hat{t}_{1}^{k}\right)\right)^{2}\right]$ is uniformly bounded in $k$ as easily seen from the definition of $L_{\mathrm{r}}^{k}$, so are $\mathbb{E}\left[\left(\varepsilon_{k} \bar{L}_{\mathrm{r}}^{k}\left(\hat{t}_{1}^{k}\right)\right)^{2}\right]$ and $\mathbb{E}\left[\varepsilon_{k} \bar{L}_{\mathrm{r}}^{k}\left(\hat{t}_{1}^{k}\right)\right]$. Since $\mathbb{E}\left[\alpha_{z}^{\langle k\rangle}\right] \rightarrow 0$, we conclude from 6.32 that

$$
\begin{equation*}
\mathbb{P}\left[S_{\varepsilon_{k}}\left(L_{r}^{k}, \bar{L}_{r}^{k}\right) \in \cdot\right] \underset{k \rightarrow \infty}{\Longrightarrow} \mathbb{P}\left[\left(L_{\mathrm{r}}, L_{\mathrm{r}}\right) \in \cdot\right] . \tag{6.33}
\end{equation*}
$$

By Lemma 6.7, (6.31) and (6.33) imply that

$$
\begin{equation*}
\mathbb{P}\left[S_{\varepsilon_{k}}\left(p_{1}^{k}, \hat{p}_{1}^{k}, \bar{L}_{\mathrm{r}}^{k}\right) \in \cdot\right] \underset{k \rightarrow \infty}{\Longrightarrow} \mathbb{P}\left[\left(\pi_{1}, \hat{\pi}_{1}, L_{\mathrm{r}}\right) \in \cdot\right] \tag{6.34}
\end{equation*}
$$

which in turn implies (6.20) and concludes our proof.

### 6.4 Discrete approximation of the Brownian net

It has been shown in [S08] that the Brownian net arises as the limit of systems of branchingcoalescing random walks, in the limit of small branching probability and after diffusive rescaling. In this section, we review this result and add some additional results on the approximation of (relevant) separation points by discrete separation points.

Let $\beta_{-} \leq \beta_{+}$be real constants. Let $\varepsilon_{k}$ be positive constants, converging to zero, and for each $k$, let

$$
\begin{equation*}
\left(\alpha^{\langle k\rangle 1}, \alpha^{\langle k\rangle \mathrm{r}}\right)=\left(\alpha_{z}^{\langle k\rangle 1}, \alpha_{z}^{\langle k\rangle \mathrm{r}}\right)_{z \in \mathbb{Z}_{\text {even }}^{2}} \tag{6.35}
\end{equation*}
$$

be an i.i.d. collection of $\{-1,+1\}^{2}$-valued random variables such that $\alpha_{z}^{\langle k\rangle 1} \leq \alpha_{z}^{\langle k\rangle \mathrm{r}}$ and

$$
\begin{equation*}
\varepsilon_{k}^{-1} \mathbb{E}\left[\alpha_{z}^{\langle k\rangle}\right] \underset{k \rightarrow \infty}{\longrightarrow} \beta_{-} \quad \text { and } \quad \varepsilon_{k}^{-1} \mathbb{E}\left[\alpha_{z}^{\langle k\rangle \mathrm{r}}\right] \underset{k \rightarrow \infty}{\longrightarrow} \beta_{+} . \tag{6.36}
\end{equation*}
$$

We let $\mathcal{U}_{\langle k\rangle}^{1}$ and $\mathcal{U}_{\langle k\rangle}^{\mathrm{r}}$ denote the discrete webs associated with $\alpha^{\langle k\rangle 1}$ and $\alpha^{\langle k\rangle \mathrm{r}}$, respectively. Then $\left(\mathcal{U}_{\langle k\rangle}^{1}, \mathcal{U}_{\langle k\rangle}^{\mathrm{r}}\right)$ is a discrete analogue of a left-right Brownian web as introduced in Section 3.3. We call the collection of discrete paths

$$
\begin{equation*}
\mathcal{V}_{\langle k\rangle}:=\left\{p: p(t+1)-p(t) \in\left\{\alpha_{(p(t), t)}^{\langle k\rangle}, \alpha_{(p(t), t)}^{\langle k\rangle \mathrm{r}}\right\} \forall t \geq \sigma_{p}\right\} \tag{6.37}
\end{equation*}
$$

the discrete net associated with $\left(\mathcal{U}_{\langle k\rangle}^{1}, \mathcal{U}_{\langle k\rangle}^{\mathrm{r}}\right)$. We observe that except for a rotation by 180 degrees and a shift from $\mathbb{Z}_{\text {even }}^{2}$ to $\mathbb{Z}_{\text {odd }}^{2}$, the discrete dual left-right web $\left(\hat{\mathcal{U}}_{\langle k\rangle}^{1}, \hat{\mathcal{U}}_{\langle k\rangle}^{\mathrm{r}}\right)$ is equally distributed with $\left(\mathcal{U}_{\langle k\rangle}^{1}, \mathcal{U}_{\langle k\rangle}^{\mathrm{r}}\right)$. In view of this, we define a dual discrete net $\hat{\mathcal{V}}_{\langle k\rangle}$ analogously to $\mathcal{V}_{\langle k\rangle}$. As in Section 3.5, we view the sets of discrete paths $\mathcal{U}_{\langle k\rangle}^{1}, \mathcal{U}_{\langle k\rangle}^{\mathrm{r}}, \mathcal{V}_{\langle k\rangle}$ as random compact subsets of the space of continuous paths $\Pi$.

We cite the following result from [SS08, Thm. 5.4].

## Theorem 6.11 (Convergence to the Brownian net)

Let $\varepsilon_{k}$ and $\mathcal{V}_{\langle k\rangle}, \mathcal{U}_{\langle k\rangle}^{1}, \mathcal{U}_{\langle k\rangle}^{\mathrm{r}}, \hat{\mathcal{V}}_{\langle k\rangle}, \hat{\mathcal{U}}_{\langle k\rangle}^{1}, \hat{\mathcal{U}}_{\langle k\rangle}^{\mathrm{r}}$ be as above and let $\mathcal{N}, \mathcal{W}^{1}, \mathcal{W}^{\mathrm{r}}, \hat{\mathcal{N}}, \hat{\mathcal{W}}^{1}, \hat{\mathcal{W}}^{\mathrm{r}}$ be a Brownian net with left and right speeds $\beta_{-} \leq \beta_{+}$, its associated left-right Brownian web, and their duals. Then

$$
\begin{equation*}
\mathbb{P}\left[S_{\varepsilon_{k}}\left(\mathcal{V}_{\langle k\rangle}, \mathcal{U}_{\langle k\rangle}^{\mathrm{l}}, \mathcal{U}_{\langle k\rangle}^{\mathrm{r}}, \hat{\mathcal{V}}_{\langle k\rangle}, \hat{\mathcal{U}}_{\langle k\rangle}^{\mathrm{l}}, \hat{\mathcal{U}}_{\langle k\rangle}^{\mathrm{r}}\right) \in \cdot\right] \underset{k \rightarrow \infty}{\Longrightarrow} \mathbb{P}\left[\left(\mathcal{N}, \mathcal{W}^{\mathrm{l}}, \mathcal{W}^{\mathrm{r}}, \hat{\mathcal{N}}, \hat{\mathcal{W}}^{1}, \hat{\mathcal{W}}^{\mathrm{r}}\right) \in \cdot\right] \tag{6.38}
\end{equation*}
$$

where $\Rightarrow$ denotes weak convergence of probability laws on $\mathcal{K}(\Pi)^{3} \times \mathcal{K}(\hat{\Pi})^{3}$.
The following analogue of Lemma 6.9 is sometimes handy.
Lemma 6.12 (Convergence of paths started at given time) In the setup of Theorem 6.11. let $T_{k} \in \mathbb{Z} \cup\{-\infty,+\infty\}$ satisfy $\varepsilon_{k}^{2} T_{k} \rightarrow T$ for some $T \in[-\infty, \infty]$. Then

$$
\begin{equation*}
\mathbb{P}\left[S_{\varepsilon_{k}}\left(\mathcal{V}_{\langle k\rangle}, \mathcal{V}_{\langle k\rangle}\left(\Sigma_{T_{k}}\right)\right) \in \cdot\right] \underset{k \rightarrow \infty}{\Longrightarrow} \mathbb{P}\left[\left(\mathcal{N}, \mathcal{N}\left(\Sigma_{T}\right)\right) \in \cdot\right] \tag{6.39}
\end{equation*}
$$

Proof. By Lemma B. 4 in the appendix, the tightness of the $S_{\varepsilon_{k}}\left(\mathcal{V}_{\langle k\rangle}\right)$ implies that also the laws of the $\mathcal{K}(\Pi)^{2}$-valued random variables $S_{\varepsilon_{k}}\left(\left(\mathcal{V}_{\langle k\rangle}, \mathcal{V}_{\langle k\rangle}\left(\Sigma_{T_{k}}\right)\right)\right.$ are tight. By going to a subsequence if necessary and invoking Skorohod's representation theorem, we may assume that they converge to an a.s. limit $(\mathcal{N}, \mathcal{A})$. It is easy to see that $\mathcal{A} \subset \mathcal{N}\left(\Sigma_{T}\right)$. To get the other inclusion, we distinguish three cases. The case $T=+\infty$ is trivial. If $-\infty<T<+\infty$, let $\mathcal{D}^{1}, \mathcal{D}^{\mathrm{r}}$ be deterministic countable dense subsets of $\mathbb{R}^{2}$ such that $\mathcal{D}^{1}$ is also dense in $\mathbb{R} \times\{T\}$. Let $\mathcal{H}_{\text {cros }}\left(\mathcal{W}^{\mathrm{l}}\left(\mathcal{D}^{\mathrm{l}}\right), \mathcal{W}^{\mathrm{r}}\left(\mathcal{D}^{\mathrm{r}}\right)\right)$ be the set of paths that can be obtained by concatenating finitely many paths in $\mathcal{W}^{\mathrm{l}}\left(\mathcal{D}^{\mathrm{l}}\right)$ and $\mathcal{W}^{\mathrm{r}}\left(\mathcal{D}^{\mathrm{r}}\right)$ at crossing times between left and right paths. Arguing as in the proof of [SS08, Thm. 5.4], we obtain that

$$
\begin{equation*}
\mathcal{A} \supset \Pi\left(\Sigma_{T}\right) \cap \mathcal{H}_{\text {cros }}\left(\mathcal{W}^{\mathrm{l}}\left(\mathcal{D}^{\mathrm{l}}\right), \mathcal{W}^{\mathrm{r}}\left(\mathcal{D}^{\mathrm{r}}\right)\right), \tag{6.40}
\end{equation*}
$$

hence by [SS08, Lemma 8.1] we conclude that $\mathcal{A} \supset \mathcal{N}\left(\Sigma_{T}\right)$. Finally, if $T=-\infty$, then let $\left.\mathcal{V}_{\langle k\rangle}\left(\Sigma_{-\infty}\right)\right|_{T_{k}} ^{\infty}$ denote the set of all restrictions of paths in $\mathcal{V}_{\langle k\rangle}\left(\Sigma_{-\infty}\right)$ to the time interval $\left[T_{k}, \infty\right]$. Since $\left.\mathcal{V}_{\langle k\rangle}\left(\Sigma_{-\infty}\right)\right|_{T_{k}} ^{\infty} \subset \mathcal{V}_{\langle k\rangle}\left(\Sigma_{T_{k}}\right)$, it then suffices to prove the claim if $T_{k}=-\infty$ for all $k$. But this is just [SS08, Lemma 9.2].

We will need one more result that is very close in spirit to Lemma 6.12 and can in fact be seen as a strengthening of the latter. If $\mathcal{N}$ is a Brownian net with left and right speeds $\beta_{-} \leq \beta_{+}$, then, generalizing (4.6), for any closed $A \subset \mathbb{R}$, we may define a Markov process taking values in the closed subsets of the real line by

$$
\begin{equation*}
\xi_{t}^{A}:=\{\pi(t): \pi \in \mathcal{N}(A \times\{0\})\} \quad(t \geq 0) \tag{6.41}
\end{equation*}
$$

We call $\xi^{A}$ the branching-coalescing point set with left and right speeds $\beta_{-}, \beta_{+}$. By combining [SS08, Prop. 1.12], Brownian scaling, and the well-known density of the Arratia flow (see [FINR02, equation (1.6)]), it is easy to check that the density of $\xi_{t}^{\mathbb{R}}$ is given by

$$
\begin{equation*}
\mathbb{E}\left[\left|\xi^{\mathbb{R}} \cap[x, y]\right|\right]=(y-x) \Psi_{b}(t) \quad(t>0, x<y) \tag{6.42}
\end{equation*}
$$

where $b:=\left(\beta_{+}-\beta_{-}\right) / 2$ and $\Psi_{b}$ is the function in (6.2).

If $\mathcal{V}$ is a discrete net defined from an i.i.d. collection of random variables $\left(\alpha_{z}^{1}, \alpha_{z}^{\mathrm{r}}\right)_{z \in \mathbb{Z}^{2}}$ as in (6.37), then we can define a discrete branching-coalescing point set in analogy with (6.41). In particular, we let

$$
\begin{equation*}
\Psi_{b_{-}, b_{+}}(t):=\mathbb{P}\left[\exists \pi \in \mathcal{V} \text { s.t. } \sigma_{\pi}=0, \pi(t)=x\right] \quad\left(t \geq 0,(x, t) \in \mathbb{Z}_{\text {even }}^{2}\right) \tag{6.43}
\end{equation*}
$$

denote its density, which is a function of $t$ and the speeds $b_{-}:=\mathbb{E}\left[\alpha_{z}^{1}\right]$ and $b_{+}:=\mathbb{E}\left[\alpha_{z}^{\mathrm{r}}\right]$ of the discrete net $\mathcal{V}$. In what follows, we will need the following fact.

Lemma 6.13 (Convergence of the density) Let $\varepsilon_{k}$ be positive constants, converging to zero and assume that $-1 \leq b_{k,-} \leq b_{k,+} \leq 1$ satisfy

$$
\begin{equation*}
\varepsilon_{k}^{-1} b_{k,-} \underset{k \rightarrow \infty}{\longrightarrow} \beta_{-} \quad \text { and } \quad \varepsilon_{k}^{-1} b_{k,+} \underset{k \rightarrow \infty}{\longrightarrow} \beta_{+} \tag{6.44}
\end{equation*}
$$

for some $\beta_{-} \leq \beta_{+}$. Let $\Psi_{b}(t)$ be the function in (6.2) with $b:=\left(\beta_{+}-\beta_{-}\right) / 2$. Then

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \sup _{\delta \leq t \leq \infty}\left|\left(2 \varepsilon_{k}\right)^{-1} \Psi_{b_{k,-}, b_{k,+}}\left(\left\lfloor\varepsilon_{k}^{-2} t\right\rfloor\right)-\Psi_{b}(t)\right|=0 \quad(\delta>0) \tag{6.45}
\end{equation*}
$$

and

$$
\begin{equation*}
\lim _{\delta \downarrow 0} \limsup _{k \rightarrow \infty} \int_{0}^{\delta} \mathrm{d} t\left(2 \varepsilon_{k}\right)^{-1} \Psi_{b_{k,-}, b_{k,+}}\left(\left\lfloor\varepsilon_{k}^{-2} t\right\rfloor\right)=0 . \tag{6.46}
\end{equation*}
$$

Proof. Fix $\delta>0$. First we derive a formula for $\Psi_{b_{-}, b_{+}}(t)$, defined as in 6.43). For $(x, t) \in$ $\mathbb{Z}_{\text {even }}^{2}$, let $\hat{p}^{\text {r }}$ (resp. $\hat{p}^{1}$ ) be the path starting from $(x-1, t)$ (resp. $\left.(x+1, t)\right)$ in the dual discrete rightmost (resp. leftmost) web $\hat{\mathcal{U}}^{\mathrm{r}}$ (resp. $\hat{\mathcal{U}}^{\mathrm{r}}$ ) associated with the discrete dual net $\hat{\mathcal{V}}$. Then by the discrete analogue of the wedge characterization of the Brownian net in Theorem 4.1, the event in the RHS of (6.43) occurs if and only if $\hat{p}^{\mathrm{r}}$ and $\hat{p}^{1}$ do not intersect on the time interval $[0, t]$. Before $\hat{p}^{r}$ and $\hat{p}^{1}$ intersect, the two paths evolve independently, with $\frac{\hat{p}^{1}(t-\cdot)-\hat{p}^{r}(t-\cdot)}{2}$ distributed as a random walk $\left(D_{i}\right)_{i \geq 0}$ with $D_{0}=1$ and increment distribution $\mathbb{P}(\Delta D=1)=$ $\gamma_{+}:=\frac{\left(1-b_{-}\right)\left(1+b_{+}\right)}{4}, \mathbb{P}(\Delta D=-1)=\gamma_{-}:=\frac{\left(1+b_{-}\right)\left(1-b_{+}\right)}{4}$, and $\mathbb{P}(\Delta D=0)=\gamma_{0}:=1-\gamma_{-}-\gamma_{+}$. Therefore

$$
\begin{equation*}
\Psi_{b_{-}, b_{+}}(t)=\mathbb{P}_{1}^{D}\left(\tau_{0}>t\right) \tag{6.47}
\end{equation*}
$$

where $\mathbb{P}_{1}^{D}(\cdot)$ denotes probability w.r.t. $D$ with $D_{0}=1$, and $\tau_{0}:=\inf \left\{i \geq 0: D_{i}=0\right\}$.
Let $N_{t}$ be the number of non-zero increments of $D$ up to time $t$, and let $\bar{D}$ be a random walk on $\mathbb{Z}$ with $\bar{D}_{0}=1$ and increment distribution $\mathbb{P}(\Delta \bar{D}= \pm 1)=\bar{\gamma}_{ \pm}:=\frac{\gamma_{ \pm}}{\gamma_{+}+\gamma_{-}}$. Then

$$
\begin{equation*}
\Psi_{b_{-}, b_{+}}(t)=\mathbb{E}_{1}^{D}\left[\mathbb{P}_{1}^{D}\left(\tau_{0}>t \mid N_{t}\right)\right]=\mathbb{E}_{1}^{D}\left[\mathbb{P}_{1}^{\bar{D}}\left(\tau_{0}>N_{t}\right)\right] . \tag{6.48}
\end{equation*}
$$

Note that for any $n \in \mathbb{N}$, the law of $\left(\bar{D}_{i}-1\right)_{0 \leq i \leq n}$ is absolutely continuous w.r.t. the law of a simple symmetric random walk $\left(X_{i}\right)_{0 \leq i \leq n}$ with $X_{0}=0$, and the Radon-Nikodym derivative
is given by $\left(2 \bar{\gamma}_{+}\right)^{\frac{n+X_{n}}{2}}\left(2 \bar{\gamma}_{-}\right)^{\frac{n-X_{n}}{2}}$. Therefore

$$
\begin{align*}
& \mathbb{P}_{1}^{\bar{D}}\left(\tau_{0}>n\right)=1-\mathbb{P}_{1}^{\bar{D}}\left(\tau_{0} \leq n\right)=1-\mathbb{P}_{1}^{\bar{D}}\left(\bar{D}_{n} \leq 0\right)-\mathbb{P}_{1}^{\bar{D}}\left(\bar{D}_{n} \geq 1, \tau_{0} \leq n\right) \\
= & \sum_{m=0}^{\infty}\left(\mathbb{P}_{1}^{\bar{D}}\left(\bar{D}_{n}=m+1\right)-\mathbb{P}_{1}^{\bar{D}}\left(\bar{D}_{n}=m+1, \tau_{0} \leq n\right)\right) \\
= & \sum_{m=0}^{\infty} \mathbb{E}_{0}^{X}\left[\left(2 \bar{\gamma}_{+}\right)^{\frac{n+m}{2}}\left(2 \bar{\gamma}_{-}\right)^{\frac{n-m}{2}}\left(1_{\left\{X_{n}=m\right\}}-1_{\left\{X_{n}=m, \tau_{-1} \leq n\right\}}\right)\right] \\
= & \sum_{m=0}^{\infty} \mathbb{E}_{0}^{X}\left[\left(2 \bar{\gamma}_{+}\right)^{\frac{n+m}{2}}\left(2 \bar{\gamma}_{-}\right)^{\frac{n-m}{2}}\left(1_{\left\{X_{n}=m\right\}}-1_{\left\{X_{n}=-m-2\right\}}\right)\right] \\
= & \sum_{m=0}^{1}\left(2 \bar{\gamma}_{+}\right)^{\frac{n+m}{2}}\left(2 \bar{\gamma}_{-}\right)^{\frac{n-m}{2}} \mathbb{P}_{0}^{X}\left(X_{n}=m\right)  \tag{6.49}\\
& +\sum_{m=2}^{\infty}\left(2 \bar{\gamma}_{+}\right)^{\frac{n+m}{2}}\left(2 \bar{\gamma}_{-}\right)^{\frac{n-m}{2}}\left(1-\frac{\bar{\gamma}_{-}}{\bar{\gamma}_{+}}\right) \mathbb{P}_{0}^{X}\left(X_{n}=m\right), \tag{6.50}
\end{align*}
$$

where we applied the reflection principle to $X$ and used $\mathbb{P}_{0}^{X}\left(X_{n}=-m-2\right)=\mathbb{P}_{0}^{X}\left(X_{n}=m+2\right)$.
We now specialize to the calculation of $\Psi_{b_{k,-}, b_{k,+}}\left(t_{k}\right)$ for $t_{k}:=\left\lfloor\varepsilon_{k}^{-2} t\right\rfloor$, where $b_{k,-}, b_{k,+}, \varepsilon_{k}$ satisfy 6.44. Note that to prove 6.46, it suffices to restrict the integral to $t \in\left[\varepsilon_{k}^{3 / 2}, \delta^{-1}\right]$. Therefore we assume $t \in\left[\varepsilon_{k}^{3 / 2}, \delta^{-1}\right]$ from now on, which implies in particular that $t_{k} \rightarrow \infty$ uniformly in $t$ as $k \rightarrow \infty$. By (6.48), where we replace $D, \bar{D}, N$. by $D^{k}, \bar{D}^{k}, N^{k}$, we have

$$
\begin{equation*}
\Psi_{b_{k,-}, b_{k,+}}\left(t_{k}\right)=\sum_{\left|n-t_{k} / 2\right| \leq t_{k}^{3 / 4}} \mathbb{P}_{1}^{D^{k}}\left(N_{t_{k}}^{k}=n\right) \mathbb{P}_{1}^{\bar{D}^{k}}\left(\tau_{0}>n\right)+\sum_{\left|n-t_{k} / 2\right|>t_{k}^{3 / 4}} \mathbb{P}_{1}^{D^{k}}\left(N_{t_{k}}^{k}=n\right) \mathbb{P}_{1}^{\bar{D}^{k}}\left(\tau_{0}>n\right), \tag{6.51}
\end{equation*}
$$

where for $k$ large, the second term is bounded by

$$
\begin{equation*}
\mathbb{P}_{1}^{D^{k}}\left(\left|N_{t_{k}}^{k}-t_{k} / 2\right|>t_{k}^{3 / 4}\right) \leq \mathbb{P}_{1}^{D^{k}}\left(\left|N_{t_{k}}^{k}-\left(\gamma_{k,+}+\gamma_{k,-}\right) t_{k}\right|>t_{k}^{2 / 3}\right) \leq 2 e^{-\frac{t_{k}^{4 / 3}}{2 t_{k}}}=2 e^{-\frac{1}{2} t_{k}^{1 / 3}} \tag{6.52}
\end{equation*}
$$

where we applied Hoeffding's concentration of measure inequality Hoe63 to $N_{t_{k}}^{k}$, which is a sum of $t_{k}$ i.i.d. $\{0,1\}$-valued random variables with mean $\gamma_{k,+}+\gamma_{k,-}=\frac{1}{2}+\frac{\beta_{+} \beta_{-}}{2} \varepsilon_{k}^{2}(1+o(1))$. Since as $k \rightarrow \infty, \varepsilon_{k}^{-1} e^{-\frac{1}{2} t_{k}^{1 / 3}} \rightarrow 0$ uniformly in $t \in\left[\varepsilon_{k}^{3 / 2}, \delta^{-1}\right]$, we can safely neglect the second term in (6.51) when proving (6.45)-(6.46).

Note that in the first sum in 6.51, $n=\varepsilon_{k}^{-2} t\left(\frac{1}{2}+o(1)\right)$ uniformly in $n$ and $t \in\left[\varepsilon_{k}^{3 / 2}, \delta^{-1}\right]$ as $k \rightarrow \infty$. For $n_{k}:=\varepsilon_{k}^{-2} t\left(\frac{1}{2}+o(1)\right)$, we have a representation for $\mathbb{P}_{1}^{\bar{D}^{k}}\left(\tau_{0}>n_{k}\right)$ as in 6.50, where the first sum in 6.50 gives

$$
\begin{align*}
\sum_{m=0}^{1}\left(2 \bar{\gamma}_{k,+}\right)^{\frac{n_{k}+m}{2}}\left(2 \bar{\gamma}_{k,-}\right)^{\frac{n_{k}-m}{2}} \mathbb{P}_{0}^{X}\left(X_{n_{k}}=m\right) & =\left(\frac{4 \gamma_{k,+} \gamma_{k,-}}{\left(\gamma_{k,+}+\gamma_{k,-}\right)^{2}}\right)^{\frac{n_{k}}{2}} \frac{2}{\sqrt{2 \pi n_{k}}}(1+o(1)) \\
& =\frac{2 \varepsilon_{k} e^{-b^{2} t(1+o(1))}}{\sqrt{\pi t}}(1+o(1)) \tag{6.53}
\end{align*}
$$

where we used

$$
\begin{aligned}
& 4 \gamma_{k, \pm}=\left(1 \pm b_{k,+}\right)\left(1 \mp b_{k,-}\right)=1 \pm \varepsilon_{k}(2 b+o(1)), \\
& 16 \gamma_{k,+} \gamma_{k,-}=\left(1-b_{k,+}^{2}\right)\left(1-b_{k,-}^{2}\right)=1-\varepsilon_{k}^{2}\left(\beta_{-}^{2}+\beta_{+}^{2}+o(1)\right), \\
& 2 \gamma_{k,+}+2 \gamma_{k,-}=1-b_{k,+} b_{k,-}=1-\varepsilon_{k}^{2}\left(\beta_{+} \beta_{-}+o(1)\right),
\end{aligned}
$$

and we applied the local central limit theorem, a strong version of which we need later is

$$
\begin{equation*}
\mathbb{P}_{0}^{X}\left(X_{s}=x\right)=1_{\{s+x \text { is even }\}} \frac{2 e^{-\frac{x^{2}}{2 s}}}{\sqrt{2 \pi s}}(1+o(1)) \tag{6.54}
\end{equation*}
$$

uniformly for all $|x| \leq s^{\frac{3}{4}}$ as $s \rightarrow \infty$. This can be deduced from Sto67, Theorem 3].
Analogously, the second term in 6.50 gives a contribution to $\mathbb{P}_{1}^{\bar{D}^{k}}\left(\tau_{0}>n_{k}\right)$ of

$$
\begin{align*}
& \left(\frac{4 \gamma_{k,+} \gamma_{k,-}}{\left(\gamma_{k,+}+\gamma_{k,-}\right)^{2}}\right)^{\frac{n_{k}}{2}}\left(1-\frac{\gamma_{k,-}}{\gamma_{k,+}}\right) \sum_{m=2}^{\infty}\left(\frac{\gamma_{k,+}}{\gamma_{k,-}}\right)^{\frac{m}{2}} \mathbb{P}_{0}^{X}\left(X_{n_{k}}=m\right) \\
= & \varepsilon_{k}(4 b+o(1)) e^{-b^{2} t(1+o(1))} \sum_{m=2}^{\infty} e^{\varepsilon_{k} m(2 b+o(1))} \mathbb{P}_{0}^{X}\left(X_{n_{k}}=m\right) \tag{6.55}
\end{align*}
$$

where we note that the sum is bounded by

$$
\sum_{m \in \mathbb{Z}} e^{3 b \varepsilon_{k} m} \mathbb{P}_{0}^{X}\left(X_{n_{k}}=m\right)=\mathbb{E}_{0}^{X}\left[e^{3 b \varepsilon_{k} X_{n_{k}}}\right]=\left(\frac{e^{3 b \varepsilon_{k}}+e^{-3 b \varepsilon_{k}}}{2}\right)^{\varepsilon_{k}^{-2} t\left(\frac{1}{2}+o(1)\right)}=O(1)
$$

uniformly in $t \in\left[\varepsilon_{k}^{\frac{3}{2}}, \delta^{-1}\right]$ as $k \rightarrow \infty$. Combined with 6.51 6.53, this implies 6.46.
To prove (6.45), we now restrict to $t \in\left[\delta, \delta^{-1}\right]$ and estimate the sum in (6.55) more preicsely. By Hoeffding's inequality Hoe63,

$$
\begin{equation*}
\mathbb{P}_{0}^{X}\left(X_{s} \geq m\right) \leq e^{-\frac{m^{2}}{2 s}} . \tag{6.56}
\end{equation*}
$$

Substituting this bound into (6.55) then gives

$$
\sum_{m>n_{k}^{3 / 4}}^{\infty} e^{\varepsilon_{k} m(2 b+o(1))} \mathbb{P}_{0}^{X}\left(X_{n_{k}}=m\right) \leq \sum_{m>n_{k}^{3 / 4}}^{\infty} e^{\varepsilon_{k} m(2 b+o(1))-\frac{m^{2}}{2 n_{k}}} \leq \sum_{m>n_{k}^{3 / 4}}^{\infty} e^{\varepsilon_{k} m(2 b+o(1))-\delta \varepsilon_{k}^{2} m^{2}(1+o(1))}=o(1)
$$

uniformly in $t \in\left[\delta, \delta^{-1}\right]$ as $k \rightarrow \infty$. On the other hand, by 6.54),

$$
\begin{align*}
& \sum_{2 \leq m \leq n_{k}^{3 / 4}}^{\infty} e^{\varepsilon_{k} m(2 b+o(1))} \mathbb{P}_{0}^{X}\left(X_{n_{k}}=m\right)=(1+o(1)) \sum_{\substack{2 \leq m \leq n_{k}^{3 / 4} \\
2 \mid\left(m+n_{k}\right)}}^{\infty} \frac{2 e^{\varepsilon_{k} m(2 b+o(1))-\frac{m^{2}}{2 n_{k}}}}{\sqrt{2 \pi n_{k}}} \\
= & (1+o(1)) \sum_{\substack{2 \leq m \leq n_{k}^{3 / 4} \\
2 \mid\left(m+n_{k}\right)}}^{\infty} \frac{e^{(2 b \sqrt{t}+o(1)) \frac{\varepsilon_{k}}{\sqrt{t}} m-(1+o(1))\left(\frac{\varepsilon_{k}}{\sqrt{t}} m\right)^{2}}}{\sqrt{\pi}} \frac{2 \varepsilon_{k}}{\sqrt{t}} \\
= & \frac{1+o(1)}{\sqrt{\pi}} \int_{0}^{\infty} e^{2 b \sqrt{t x}-x^{2}} \mathrm{~d} x=(1+o(1)) e^{b^{2} t} \int_{-b \sqrt{2 t}}^{\infty} \frac{e^{-\frac{x^{2}}{2}}}{\sqrt{2 \pi}} \mathrm{~d} x \tag{6.57}
\end{align*}
$$

by Riemann sum approximation. Substituting the last two estimates into 6.55 and combining with (6.51)-(6.53) then gives (6.45).

Proposition 6.14 (Convergence of relevant separation points) Let $\mathcal{V}_{\langle k\rangle}$ be a sequence of discrete nets as defined in (6.35)-(6.37) and let $\mathcal{N}$ be a Brownian net with left and right speeds $\beta_{-} \leq \beta_{+}$. Let $-\infty \leq S<U \leq \infty$ and let $S_{k}, U_{k} \in \mathbb{Z} \cup\{-\infty,+\infty\}$ be such that $\varepsilon_{k}^{2} S_{k} \rightarrow S$ and $\varepsilon_{k}^{2} U_{k} \rightarrow U$. Let $R_{S, U}$ denote the set of $S, U$-relevant separation points of $\mathcal{N}$ and let $R_{S_{k}, U_{k}}^{\langle k\rangle}$ denote the set of $S_{k}, U_{k}$-relevant separation points of $\mathcal{V}_{\langle k\rangle}$. Then it is possible to couple the $\mathcal{V}_{\langle k\rangle}$ and $\mathcal{N}$ in such a way that

$$
\begin{equation*}
S_{\varepsilon_{k}}\left(\mathcal{V}_{\langle k\rangle}\right) \underset{k \rightarrow \infty}{\longrightarrow} \mathcal{N} \quad \text { a.s. } \tag{6.58}
\end{equation*}
$$

and moreover

$$
\begin{equation*}
\sum_{z \in R_{S_{k}, U_{k}}^{(k)}} \delta_{S_{\varepsilon_{k}}(z)}^{\Longrightarrow} \underset{z \rightarrow \infty}{\Longrightarrow} \sum_{z \in R_{S, U}} \delta_{z} \quad \text { a.s. } \tag{6.59}
\end{equation*}
$$

where $\Rightarrow$ denotes vague convergence of locally finite measures on $\mathbb{R}^{2}$.
Remark. The convergence in 6.59 is stronger than the statement that for each $z \in R_{S, U}$ there exist $z_{k} \in R_{S_{k} . U_{k}}^{\langle k\rangle}$ such that $S_{\varepsilon_{k}}\left(z_{k}\right) \rightarrow z$. Indeed, since the counting measure on the right-hand side of (6.59) has no double points, such an approximating sequence is eventually unique, a fact that wil be important in the proof of Theorem 6.15 below.
Proof of Proposition 6.14. By Theorem 6.11, Lemma 6.12 , Lemma 6.7 (a) and the remarks below it, we can couple our random variables such that

$$
\begin{equation*}
S_{\varepsilon_{k}}\left(\mathcal{V}_{\langle k\rangle}, \hat{\mathcal{V}}_{\langle k\rangle}, \mathcal{V}_{\langle k\rangle}\left(\Sigma_{S_{k}}\right), \hat{\mathcal{V}}_{\langle k\rangle}\left(\Sigma_{U_{k}}\right)\right) \underset{k \rightarrow \infty}{\longrightarrow}\left(\mathcal{N}, \hat{\mathcal{N}}, \mathcal{N}\left(\Sigma_{S}\right), \hat{\mathcal{N}}\left(\Sigma_{U}\right)\right) \quad \text { a.s. } \tag{6.60}
\end{equation*}
$$

We claim that with this coupling, for each $z \in R_{S, U}$ there exist $z_{k} \in R_{S_{k}, U_{k}}^{\langle k\rangle}$ with $S_{\varepsilon_{k}}\left(z_{k}\right) \rightarrow z$. To see this, note that by Lemma 6.3, each $z \in R_{S, U}$ is a crossing point of some $\pi \in \mathcal{N}$ and $\hat{\pi} \in \hat{\mathcal{N}}$ with $\sigma_{\pi}=S$ and $U=\hat{\sigma}_{\hat{\pi}}$. By (6.60), there exist $p_{k} \in \mathcal{V}_{\langle k\rangle}\left(\Sigma_{S_{k}}\right)$ and $\hat{p}_{k} \in \hat{\mathcal{V}}_{\langle k\rangle}\left(\Sigma_{U_{k}}\right)$ such that $p_{k} \rightarrow \pi$ and $\hat{p}_{k} \rightarrow \hat{\pi}$. It follows from the definition of crossing points that for $k$ sufficiently large, there must exist points $z_{k} \in \mathbb{Z}_{\text {even }}^{2}$ such that $S_{\varepsilon_{k}}\left(z_{k}\right) \rightarrow z$ and $p_{k}$ crosses $\hat{p}_{k}$ in $z_{k}$. In particular, this implies that the $z_{k}$ must be $S_{k}, U_{k}$-relevant in $\mathcal{V}_{\langle k\rangle}$.

We next claim that for each $-\infty<T_{-}<T_{+}<\infty$ and $-\infty<x_{-}<x_{+}<+\infty$,

$$
\begin{equation*}
\mathbb{E}\left[\left|S_{\varepsilon_{k}}\left(R_{S_{k}, U_{k}}^{\langle k}\right) \cap A\right|\right] \underset{k \rightarrow \infty}{\longrightarrow} \mathbb{E}\left[\left|R_{S, U} \cap A\right|\right], \quad \text { where } \quad A=\left(x_{-}, x_{+}\right) \times\left(T_{-}, T_{+}\right) . \tag{6.61}
\end{equation*}
$$

To see this, recall that the discrete nets $\mathcal{V}_{\langle k\rangle}$ are defined from i.i.d. collections of random variables $\left(\alpha_{z}^{\langle k\rangle 1}, \alpha_{z}^{\langle k\rangle \mathrm{r}}\right)_{z \in \mathbb{Z}_{\text {even }}^{2}}$. We observe that for all $z=(x, t) \in \mathbb{Z}_{\text {even }}^{2}$ with $S_{k} \leq t<U_{k}$,

$$
\begin{align*}
\mathbb{P} & {\left[z \text { is } S_{k}, U_{k} \text {-relevant in } \mathcal{V}_{\langle k\rangle}\right] } \\
& =\mathbb{P}\left[\alpha_{z}^{1}<\alpha_{z}^{\mathrm{r}} \mathbb{P}\left[\exists \pi \in \mathcal{V}_{\langle k\rangle} \text { s.t. } \sigma_{\pi}=S_{k}, \pi(t)=x\right] \mathbb{P}\left[\exists \hat{\pi} \in \hat{\mathcal{V}}_{\langle k\rangle} \text { s.t. } \hat{\sigma}_{\hat{\pi}}=U_{k}, \hat{\pi}(t+1)=x\right]\right. \\
& =\frac{1}{2}\left(b_{k,+}-b_{k,-}\right) \Psi_{b_{k,-}, b_{k,+}}\left(t-S_{k}\right) \Psi_{b_{k,-}, b_{k,+}}\left(U_{k}-(t+1)\right), \tag{6.62}
\end{align*}
$$

where

$$
\begin{equation*}
b_{k,-}:=\mathbb{E}\left[\alpha_{z}^{\langle k\rangle}\right] \quad \text { and } \quad b_{k,+}:=\mathbb{E}\left[\alpha_{z}^{\langle k\rangle \mathrm{r}}\right], \tag{6.63}
\end{equation*}
$$

and $\Psi_{b_{-}, b_{+}}(t)$ is the function in (6.43). We claim that (6.61) now follows from Proposition 6.1 (b), Lemma 6.13, and Riemann sum approximation. Without going through the
details, note that after diffusive rescaling, the per unit density of points of $S_{\varepsilon}\left(\mathbb{Z}_{\text {even }}^{2}\right)$ in the plane is $\frac{1}{2} \varepsilon_{k}^{-3}$, and therefore, by Lemma 6.13 formula (6.62) says that after diffusive rescaling, the per unit density of relevant separation points at time $S<t<U$ is approximately given by

$$
\begin{align*}
& \frac{1}{2} \varepsilon_{k}^{-3} \cdot \frac{1}{2}\left(\varepsilon_{k} \beta_{+}-\varepsilon_{k} \beta_{-}\right) \cdot\left(2 \varepsilon_{k}\right) \Psi_{b}(t-S) \cdot\left(2 \varepsilon_{k}\right) \Psi_{b}(U-t)  \tag{6.64}\\
& \quad=2 b \Psi_{b}(t-S) \Psi_{b}(u-t) \quad \text { where } \quad b:=\left(\beta_{+}-\beta_{-}\right) / 2
\end{align*}
$$

which agrees with 6.1).
To prove the existence of a coupling such that (6.59) holds, let

$$
\begin{equation*}
\nu:=\sum_{z \in R_{S, U} \cap A} \delta_{z} \quad \text { and } \quad \nu_{k}:=\sum_{z \in S_{\varepsilon_{k}}\left(R_{S_{k}, U_{k}}^{k j)}\right) \cap A} \delta_{z} \tag{6.65}
\end{equation*}
$$

be random counting measures with atoms at the positions of the sets in (6.61). By (6.61), the laws of the $\nu_{k}$ 's are tight, so by going to a subsequence if necessary and invoking Skorohod's representation theorem, we can find a coupling such that in addition to (6.60), also $\nu_{k} \Rightarrow \nu^{*}$, where $\Rightarrow$ denotes weak convergence and $\nu^{*}$ is some finite counting measure on the closure $\bar{A}$ of $A$. Since for each $z \in R_{S, U}$ there exist $z_{k} \in R_{S_{k}, U_{k}}^{\langle k\rangle}$ such that $S_{\varepsilon_{k}}\left(z_{k}\right) \rightarrow z$, we know that $\nu \leq \nu^{*}$. By 6.61, we see that moreover $\mathbb{E}[\nu(A)]=\mathbb{E}\left[\nu^{*}(\bar{A})\right]$, so we conclude that $\nu=\nu^{*}$.

By Lemma 6.7 (a) and the remarks below it, we can find a coupling such that the measures in (6.59) converge weakly on ( $\left.x_{n,-}, x_{n,+}\right) \times\left(T_{n,-}, T_{n,+}\right)$ for each $n$, where $x_{n,-}, T_{n,-} \downarrow-\infty$ and $x_{n,+}, T_{n,+} \uparrow+\infty$, proving the vague convergence in (6.59).

### 6.5 Discrete approximation of a coupled Brownian web and net

In this section, we prove a convergence result for discrete webs that are defined 'inside' a discrete net. As a result, we will obtain Theorem 4.4. Our convergence result also prepares for the proof of Theorem 3.9 which will be given in Section 7.1 .

For $k \geq 1$, let $\left(\alpha^{\langle k\rangle}, \alpha^{\langle k\rangle \mathrm{r}}\right)$ be a collection of $\{-1,+1\}^{2}$-valued random variables indexed by $\mathbb{Z}_{\text {even }}^{2}$ as in $\left.6.35-6.36\right)$ and let $\mathcal{V}_{\langle k\rangle}$ and $\mathcal{U}_{\langle k\rangle}^{1}, \mathcal{U}_{\langle k\rangle}^{\mathrm{r}}$ be the associated discrete net (as defined in (6.37) and discrete left-right web.

Let $r \in[0,1]$ and, conditional on $\left(\alpha^{\langle k\rangle}, \alpha^{\langle k\rangle}\right)$, let $\alpha^{\langle k\rangle}=\left(\alpha_{z}^{\langle k\rangle}\right)_{z \in \mathbb{Z}_{\text {even }}^{2}}$ be a collection of independent $\{-1,+1\}$-valued random variables such that $\alpha_{z}^{\langle k\rangle 1} \leq \alpha_{z}^{\langle k\rangle} \leq \alpha_{z}^{\langle k\rangle \mathrm{r}}$ a.s. and

$$
\begin{equation*}
\mathbb{P}\left[\alpha_{z}^{\langle k\rangle}=\alpha_{z}^{\langle k\rangle \mathrm{r}} \mid\left(\alpha^{\langle k\rangle 1}, \alpha^{\langle k\rangle \mathrm{r}}\right)\right]=r \quad\left(z \in \mathbb{Z}_{\mathrm{even}}^{2}\right) \tag{6.66}
\end{equation*}
$$

Then, obviously, under the unconditioned law the collection $\alpha^{\langle k\rangle}=\left(\alpha_{z}^{\langle k\rangle}\right)_{z \in \mathbb{Z}_{\text {even }}^{2}}$ is i.i.d. with

$$
\begin{equation*}
\varepsilon_{k}^{-1} \mathbb{E}\left[\alpha_{z}^{\langle k\rangle}\right] \underset{k \rightarrow \infty}{\longrightarrow} \beta:=(1-r) \beta_{-}+r \beta_{+} . \tag{6.67}
\end{equation*}
$$

We let $\mathcal{U}_{\langle k\rangle}$ denote the discrete web associated with $\alpha^{\langle k\rangle}$. The following theorem implies Theorem 4.4.

Theorem 6.15 (Converge to a coupled Brownian web and net) Let $\mathcal{U}_{\langle k\rangle}$ and $\mathcal{V}_{\langle k\rangle}$ be coupled discrete webs and nets as above. Then

$$
\begin{equation*}
\mathbb{P}\left[S_{\varepsilon_{k}}\left(\mathcal{U}_{\langle k\rangle}, \mathcal{V}_{\langle k\rangle}\right) \in \cdot\right] \underset{k \rightarrow \infty}{\Longrightarrow} \mathbb{P}[(\mathcal{W}, \mathcal{N}) \in \cdot] \tag{6.68}
\end{equation*}
$$

where $\mathcal{N}$ is a Brownian net with left and right speeds $\beta_{-} \leq \beta_{+}$and $\mathcal{W}$ is a Brownian web with drift $\beta$. Letting $S$ denote the set of separation points of $\mathcal{N}$, one has a.s.:
(i) $\mathcal{W} \subset \mathcal{N}$ and each point $z \in S$ is of type $(1,2)$ in $\mathcal{W}$.
(ii) Conditional on $\mathcal{N}$, the random variables $\left(\operatorname{sign}_{\mathcal{W}}(z)\right)_{z \in S}$ are i.i.d. with $\mathbb{P}\left[\operatorname{sign}_{\mathcal{W}}(z)=\right.$ $+1 \mid \mathcal{N}]=r$.
(iii) Conditional on $\mathcal{W}$, the sets $S_{1}:=\left\{z \in S: \operatorname{sign}_{\mathcal{W}}(z)=-1\right\}$ and $S_{\mathrm{r}}:=\{z \in S$ : $\left.\operatorname{sign}_{\mathcal{W}}(z)=+1\right\}$ are independent Poisson point sets with intensities $\left(\beta_{+}-\beta\right) \ell_{1}$ and $\left(\beta-\beta_{-}\right) \ell_{\mathrm{r}}$, respectively.

Moreover,

$$
\begin{equation*}
\mathcal{W}=\left\{\pi \in \mathcal{N}: \operatorname{sign}_{\pi}(z)=\operatorname{sign}_{\mathcal{W}}(z) \forall z \in S \text { s.t. } \pi \text { enters } z\right\} \tag{6.69}
\end{equation*}
$$

In the special case that $r=0$ (resp. $r=1$ ), the Brownian web $\mathcal{W}$ is the left (resp. right) Brownian web associated with $\mathcal{N}$.

Proof. By Theorems 6.8 and 6.11 , the random variables $\mathcal{U}_{\langle k\rangle}$ and $\mathcal{V}_{\langle k\rangle}$, diffusively rescaled with $\varepsilon_{k}$, converge weakly in law to a Brownian web with drift $\beta$ and Brownian net with left and right speeds $\beta_{-}, \beta_{+}$, respectively. It follows that the laws on the left-hand side of (6.68) are tight, so by going to a subsequence if necessary we can assume that they converge weakly in law to a limit $(\mathcal{W}, \mathcal{N})$. We will show that each such limit point has the properties (i)-(iii) and satisfies moreover (6.69). Since property (ii) and formula (6.69) determine the joint law of $(\mathcal{W}, \mathcal{N})$ uniquely, this then proves the convergence in (6.68). Note that if $r=0$ (resp. $r=1$ ), then it follows from Theorem 6.11 that $\mathcal{W}$ is the left (resp. right) Brownian web associated with $\mathcal{N}$.

Proof of property (i). The fact that $\mathcal{W} \subset \mathcal{N}$ is immediate from the fact that $\mathcal{U}_{\langle k\rangle} \subset \mathcal{V}_{\langle k\rangle}$ for each $k$. To prove that each point $z \in S$ is of type $(1,2)$ in $\mathcal{W}$, we first claim that if $l \in \mathcal{W}^{1}\left(z^{\prime}\right)$ and $r \in \mathcal{W}^{\mathrm{r}}\left(z^{\prime}\right)$ satisfy $l \sim_{\text {out }}^{z^{\prime}} r$ for some $z^{\prime}=\left(x^{\prime}, t^{\prime}\right)$, then there exists a path $\pi \in \mathcal{W}\left(z^{\prime}\right)$ such that $l \leq \pi \leq r$ on $\left[t^{\prime}, \infty\right)$. This follows from the fact that by [SS08, Prop. 3.6 (b)], there exist $t_{n} \downarrow t^{\prime}$ such that $l\left(t_{n}\right)<r\left(t_{n}\right)$, while by [SS08, Prop. 1.8], any path in $\mathcal{N}$ started at time $t_{n}$ at a position in $\left(l\left(t_{n}\right), r\left(t_{n}\right)\right)$ is contained between $l$ and $r$. Using the compactness of $\mathcal{W}$, we find a path $\pi \in \mathcal{W}\left(z^{\prime}\right)$ with the desired property. Applying our claim to the one incoming and two outgoing left-right pairs at a separation point $z$ of $\mathcal{N}$, we find that there must be at least one incoming path and at least two outgoing paths in $\mathcal{W}$ at each such point. By the classification of special points in $\mathcal{W}$ (Proposition 3.3), it follows that $z$ is of type $(1,2)$ in $\mathcal{W}$.

Intermezzo. Before we turn to the proofs of properties (ii) and (iii), we first prove some prepatarory results. Since we are assuming that $S_{\varepsilon_{k}}\left(\mathcal{U}_{\langle k\rangle}, \mathcal{V}_{\langle k\rangle}\right)$ converges weakly in law to $(\mathcal{W}, \mathcal{N})$, by Skorohod's representation theorem, we can find a coupling such that the convergence is a.s. Let $\mathcal{T} \subset \mathbb{R}$ be a deterministic countable dense set of times and for $T \in \mathcal{T}$, set $T_{[k]}:=\left\lfloor\varepsilon_{k}^{-2} T\right\rfloor$. By Proposition 6.14. Lemma 6.7 (a) and the remarks below it, we can improve our coupling such that

$$
\begin{equation*}
\sum_{z \in R_{S_{[k]}, U}^{(k)}[k]} \delta_{S_{\varepsilon_{k}}(z)} \Longrightarrow \sum_{k \rightarrow \infty}^{\Longrightarrow} \delta_{z \in R_{S, U}} \quad \forall S, U \in \mathcal{T}, S<U, \tag{6.70}
\end{equation*}
$$

where $\Rightarrow$ denotes vague convergence of locally finite measures on $\mathbb{R}^{2}$, and $R_{S_{[k]}, U_{[k]}}^{\langle k}$ and $R_{S, U}$ denote the sets of $S_{[k]}, U_{[k]}$-relevant and $S, U$-relevant separation points of $\mathcal{V}_{\langle k\rangle}$ and $\mathcal{N}$, respectively.

It follows from 6.70 that for each $z \in R_{S, U}$, there exist $z_{k} \in R_{S_{[k],}, U_{[k]}}^{\langle k\rangle}$ such that $S_{\varepsilon_{k}}\left(z_{k}\right) \rightarrow$ $z$. We claim that such an approximating sequence is eventually unique. To see this, assume
that $z_{k}^{\prime} \in R_{S_{[k]}, U_{[k]}}^{\langle k}$ satisfy $S_{\varepsilon_{k}}\left(z_{k}^{\prime}\right) \rightarrow z$. We can choose $\delta>0$ such that the ball of radius $\delta$ around $z$ does not contain any other $S, U$-relevant separation points of $\mathcal{N}$ except for $z$. Then (6.70) shows that for $k$ sufficiently large, there is exactly one $S_{[k]}, U_{[k]}$-relevant separation point in the ball of radius $\delta$ around $z$, hence $z_{k}=z_{k}^{\prime}$ for $k$ sufficiently large.

Now let $S, U \in \mathcal{T}, S<U, z \in R_{S, U}$, and let $z_{k}$ be the eventually unique sequence of points in $R_{S_{[k]}, U_{[k]}}^{\langle k\rangle}$ such that $S_{\varepsilon_{k}}\left(z_{k}\right) \rightarrow z$. We claim that

$$
\begin{equation*}
\operatorname{sign}_{\mathcal{W}}(z)=+1 \quad \text { if and only if } \quad \alpha_{z_{k}}^{\langle k\rangle}=+1 \quad \text { eventually } \tag{6.71}
\end{equation*}
$$

(Here, " $\alpha_{z_{k}}^{\langle k\rangle}=+1$ eventually" means that there exists a $K$ such that $\alpha_{z_{k}}^{\langle k\rangle}=+1$ for all $k \geq K$.) It suffices to prove that $\operatorname{sign}_{\mathcal{W}}(z)=+1$ implies that $\alpha_{z_{k}}^{\langle k\rangle}=+1$ eventually. By symmetry, this then also shows that $\operatorname{sign}_{\mathcal{W}}(z)=-1$ implies that $\alpha_{z_{k}}^{\langle k\rangle}=-1$ eventually, proving 6.71.

If $\operatorname{sign}_{\mathcal{W}}(z)=+1$, then there exist $\pi \in \mathcal{W}$ and $\hat{l} \in \hat{\mathcal{W}}^{1}$ such that $\pi$ crosses $\hat{l}$ in $z=(x, t)$. It follows that there exist $p_{k} \in \mathcal{U}_{\langle k\rangle}$ and $\hat{l}_{k} \in \hat{\mathcal{U}}_{\langle k\rangle}^{1}$ such that $S_{\varepsilon_{k}}\left(p_{k}\right) \rightarrow \pi$ and $S_{\varepsilon_{k}}\left(\hat{l}_{k}\right) \rightarrow \hat{l}$, and points $z_{k}^{\prime} \in \mathbb{Z}_{\text {even }}^{2}$ with $S_{\varepsilon_{k}}\left(z_{k}^{\prime}\right) \rightarrow z$ such that $p_{k}$ crosses $\hat{l}_{k}$ in $z_{k}^{\prime}$. Let $\sigma_{\pi}<S^{\prime}<t<U^{\prime}<\hat{\sigma}_{\hat{l}}$. Then the $z_{k}^{\prime}$ are $S_{[k]}^{\prime}, U_{[k]}^{\prime}$-relevant for $k$ large enough, so by the principle of eventual uniqueness applied to the times $S^{\prime}, U^{\prime}$ we see that $z_{k}^{\prime}=z_{k}$ eventually. Since $\alpha_{z_{k}^{\prime}}^{\langle k\rangle}=+1$ for each $k$, this proves 6.71).

Next, let $\tilde{z} \in \mathbb{R}^{2}$ and $\tilde{z}_{k} \in \mathbb{Z}_{\text {even }}^{2}$ be deterministic points such that $S_{\varepsilon_{k}}\left(\tilde{z}_{k}\right) \rightarrow \tilde{z}$, let $p_{k}$ be the unique element of $\mathcal{U}_{\langle k\rangle}\left(\tilde{z}_{k}\right)$ and let $\pi$ be the a.s. unique element of $\mathcal{W}(\tilde{z})$. Since $S_{\varepsilon_{k}}\left(\mathcal{U}_{\langle k\rangle}, p_{k}\right)$ converges weakly in law to ( $\mathcal{W}, \pi$ ), by Lemma 6.7(a) and the remarks below it, we can improve our coupling such that $S_{\varepsilon_{k}}\left(p_{k}\right) \rightarrow \pi$ a.s. Now as before let $S, U \in \mathcal{T}, S<U, z \in R_{S, U}$, and let $z_{k}$ be the eventually unique sequence of points in $R_{S_{[k]}, U_{[k]}}^{k\rangle}$ such that $S_{\varepsilon_{k}}\left(z_{k}\right) \rightarrow z$. We claim that

$$
\begin{equation*}
\pi \text { enters } z \text { if and only if } p_{k} \text { enters } z_{k} \text { eventually as } k \rightarrow \infty \tag{6.72}
\end{equation*}
$$

Indeed, if $\pi$ does not enter $z$, then $\pi$ dos not enter some open ball around $z$, so it is clear that for $k$ sufficiently large, $p_{k}$ does not enter $z_{k}$. On the other hand, if $\pi$ enters $z$, then since by property (i), $z$ is either of type $(1,2)_{1}$ or of type $(1,2)_{\mathrm{r}}$ in $\mathcal{W}$, there must be either some $\hat{r} \in \hat{\mathcal{W}}^{\text {r }}$ such that $\pi$ crosses $\hat{r}$ from right to left or some $\hat{l} \in \hat{\mathcal{W}}^{\mathrm{l}}$ such that $\pi$ crosses $\hat{k}$ from left to right. By symmetry, it suffices to consider only the first case. In this case, there must exist $\hat{r}_{k} \in \hat{\mathcal{U}}_{\langle k\rangle}^{\mathrm{r}}$ and $z_{k}^{\prime} \in \mathbb{Z}_{\text {even }}^{2}$ with $S_{\varepsilon_{k}}\left(z_{k}^{\prime}\right) \rightarrow z$ such that for $k$ sufficiently large, $p_{k}$ crosses $\hat{r}_{k}$ in $z_{k}^{\prime}$. By the same argument as in the proof of (6.71) we see that $z_{k}^{\prime}=z_{k}$ eventually, hence $p_{k}$ enters $z_{k}$ eventually.

Proof of Property (ii). Let $\Theta \subset \mathbb{R}$ be a deterministic finite set, say $\Theta=\left\{T_{1}, \ldots, T_{m}\right\}$ with $T_{1}<\cdots<T_{m}$, and set

$$
\begin{align*}
& R_{\Theta}:=\left\{z \in \mathbb{R}^{2}: z \text { is a } T_{i}, T_{i+1} \text {-relevant separation point for some } 1 \leq i \leq m-1\right\}, \\
& R_{\Theta}^{+}:=\left\{z \in \mathbb{R}_{\Theta}: \operatorname{sign}_{\mathcal{W}}(z)=+1\right\} . \tag{6.73}
\end{align*}
$$

We let

$$
\begin{equation*}
\nu_{\Theta}:=\sum_{z \in R_{\Theta}} \delta_{z} \quad \text { and } \quad \nu_{\Theta}^{+}:=\sum_{z \in R_{\Theta}^{+}} \delta_{z} \tag{6.74}
\end{equation*}
$$

be counting measures with atoms at each point of $R_{\Theta}$ and $R_{\Theta}^{+}$, respectively. We wish to show that $R_{\Theta}^{+}$is an $r$-thinning of $R_{\Theta}$. By formula (D.4) and Lemma C.1 of Appendix D, it suffices
to show that (in notation introduced there)

$$
\begin{equation*}
\mathbb{E}\left[(1-f)^{\nu_{\Theta}^{+}} \mid \mathcal{N}\right]=(1-r f)^{\nu_{\Theta}} \quad \text { a.s. } \tag{6.75}
\end{equation*}
$$

for each deterministic $f: \mathbb{R}^{2} \rightarrow[0,1]$ that is continuous and has compact support. Equivalently, we may show that

$$
\begin{equation*}
\mathbb{E}\left[(1-f)^{\nu_{\Theta}^{+}} g(\mathcal{N})\right]=\mathbb{E}\left[(1-r f)^{\nu_{\Theta}} g(\mathcal{N})\right] \tag{6.76}
\end{equation*}
$$

for each $f$ as before and bounded continuous $g: \mathcal{K}(\Pi) \rightarrow \mathbb{R}$. We now choose deterministic $T_{k, i} \in \mathbb{Z}$ with $\varepsilon_{k}^{2} T_{k, i} \rightarrow T_{i}$ for each $1 \leq i \leq m$, and we set

$$
\begin{align*}
& R^{\langle k\rangle}:=\left\{z \in \mathbb{Z}_{\text {even }}^{2}: z\right. \text { is a } T_{k, i}, T_{k, i+1} \text {-relevant separation point } \\
&\text { of } \left.\mathcal{V}_{\langle k\rangle} \text { for some } 1 \leq i \leq m-1\right\},  \tag{6.77}\\
& R^{+\langle k\rangle}:=\left\{z \in \mathbb{R}^{\langle k\rangle}: \alpha_{z}=+1\right\},
\end{align*}
$$

and

$$
\begin{equation*}
\nu^{\langle k\rangle}:=\sum_{z \in R^{\langle k\rangle}} \delta_{S_{\varepsilon_{k}}(z)} \quad \text { and } \quad \nu^{+\langle k\rangle}:=\sum_{z \in R^{+}\langle k\rangle} \delta_{S_{\varepsilon_{k}}(z)} . \tag{6.78}
\end{equation*}
$$

By construction, $\nu^{+\langle k\rangle}$ is an $r$-thinning of $\nu^{\langle k\rangle}$, so

$$
\begin{equation*}
\mathbb{E}\left[(1-f)^{\nu^{+\langle k\rangle}} g\left(\mathcal{V}_{\langle k\rangle}\right)\right]=\mathbb{E}\left[(1-r f)^{\nu^{\langle k\rangle}} g\left(\mathcal{V}_{\langle k\rangle}\right)\right] \tag{6.79}
\end{equation*}
$$

Now (6.76) will follow by taking the limit in (6.79), provided we show that there exists a coupling such that

$$
\begin{equation*}
\text { (i) } \nu^{\langle k\rangle} \underset{k \rightarrow \infty}{\Longrightarrow} \nu_{\Theta}, \quad \text { (ii) } \nu^{+\langle k\rangle} \underset{k \rightarrow \infty}{\Longrightarrow} \nu_{\Theta}^{+}, \tag{6.80}
\end{equation*}
$$

where $\Rightarrow$ denotes vague convergence on $\mathbb{R}^{2}$. The existence of a coupling such that (6.80) (i) holds follows from Proposition 6.14. By $\sqrt{6.71)}$, we can improve this coupling such that also 6.80) (ii) holds. This completes our proof that $R_{\Theta}^{+}$is an $r$-thinning of $R_{\Theta}$. Since $\Theta$ is arbitrary and since each separation point ( $x, t$ ) is $S, U$-relevant for some $S<t<U$, Property (ii) follows.

Proof of Property (iii). Let $\Delta, \hat{\Delta} \subset \mathbb{R}^{2}$ be deterministic, finite sets and let $\ell_{1}(\Delta, \hat{\Delta})$ and $\ell_{\mathrm{r}}(\Delta, \hat{\Delta})$ be the restrictions of $\ell_{1}$ and $\ell_{\mathrm{r}}$, respectively, to the set $I:=\operatorname{Img}(\mathcal{W}(\Delta)) \cap \operatorname{Img}(\hat{\mathcal{W}}(\hat{\Delta}))$. We note that if $\Delta_{n}, \hat{\Delta}_{n}$ are finite sets increasing to countable limits $\Delta_{\infty}$ and $\hat{\Delta}_{\infty}$ that are dense in $\mathbb{R}^{2}$, then $\operatorname{Img}\left(\mathcal{W}\left(\Delta_{n}\right)\right) \cap \operatorname{Img}\left(\hat{\mathcal{W}}\left(\hat{\Delta}_{n}\right)\right)$ increases to the set of all points of type $(1,2)$ in $\mathcal{W}$, so Property (iii) will follow provided we show that for any deterministic finite $\Delta, \hat{\Delta} \subset \mathbb{R}^{2}$, conditional on $\mathcal{W}$, the sets $S_{1} \cap I$ and $S_{\mathrm{r}} \cap I$ are independent Poisson point sets with intensities $\left(\beta_{+}-\beta\right) \ell_{1}(\Delta, \hat{\Delta})$ and $\left(\beta-\beta_{-}\right) \ell_{r}(\Delta, \hat{\Delta})$, respectively.

Equivalently, this says that the set $\left\{(z,-1): z \in S_{1} \cap I\right\} \cup\left\{(z,+1): z \in S_{\mathrm{r}} \cap I\right\}$ is a Poisson point set on $\mathbb{R}^{2} \times\{-1,+1\}$ with intensity $\int \ell_{1}(\mathrm{~d} z) \delta_{(z,-1)}+\int \ell_{\mathrm{r}}(\mathrm{d} z) \delta_{(z,+1)}$. Thus, by formula (D.2) and Lemma C.1 of Appendix D, it suffices to show that (in notation introduced there)

$$
\begin{equation*}
\mathbb{E}\left[(1-f)^{\nu_{1}}(1-g)^{\nu_{\mathrm{r}}} \mid \mathcal{W}\right]=e^{-\left(\beta_{+}-\beta\right) \int f \mathrm{~d} \ell_{1}(\Delta, \hat{\Delta})-\left(\beta-\beta_{-}\right) \int g \mathrm{~d} \ell_{\mathrm{r}}(\Delta, \hat{\Delta})} \tag{6.81}
\end{equation*}
$$

for any deterministic, continuous $f, g: \mathbb{R}^{2} \rightarrow[0,1]$, where

$$
\begin{equation*}
\nu_{1}:=\sum_{z \in S_{1} \cap I} \delta_{z} \quad \text { and } \quad \nu_{\mathrm{r}}:=\sum_{z \in S_{\mathrm{r}} \cap I} \delta_{z} . \tag{6.82}
\end{equation*}
$$

Equivalently, we may show that

$$
\begin{equation*}
\mathbb{E}\left[(1-f)^{\nu_{1}}(1-g)^{\nu_{\mathrm{r}}} h(\mathcal{W})\right]=\mathbb{E}\left[e^{-\left(\beta_{+}-\beta\right) \int f \mathrm{~d} \ell_{1}(\Delta, \hat{\Delta})-\left(\beta-\beta_{-}\right) \int g \mathrm{~d} \ell_{\mathrm{r}}(\Delta, \hat{\Delta})} h(\mathcal{W})\right] \tag{6.83}
\end{equation*}
$$

for each $f, g$ as before and bounded continuous $h: \mathcal{K}(\Pi) \rightarrow \mathbb{R}$. Let $\Delta_{k} \subset \mathbb{Z}_{\text {even }}^{2}$ and $\hat{\Delta}_{k} \subset \mathbb{Z}_{\text {odd }}^{2}$ approximate $\Delta$ and $\hat{\Delta}$ as in Proposition 6.10, let $Z_{\mathrm{r}}^{\langle k\rangle}, I_{k}$ and $\ell_{\mathrm{r}}^{\langle k\rangle}$ be as defined in 6.15) and let $\ell_{1}^{\langle k\rangle}$ be defined similarly, with $Z_{\mathrm{r}}^{\langle k\rangle}$ replaced by $Z_{1}^{\langle k\rangle}:=\left\{z \in \mathbb{Z}_{\text {even }}^{2}: \alpha_{z}^{\langle k\rangle}=-1\right\}$. Let $S_{\langle k\rangle}$ be the set of separation points of $\mathcal{V}_{\langle k\rangle}$, and set

$$
\begin{equation*}
\nu_{1}^{\langle k\rangle}:=\sum_{z \in S_{\langle k\rangle} \cap I_{k} \cap Z_{1}^{(k\rangle}} \delta_{S_{\varepsilon_{k}}(z)} \quad \text { and } \quad \nu_{\mathrm{r}}^{\langle k\rangle}:=\sum_{z \in S_{\langle k\rangle} \cap I_{k} \cap Z_{\mathrm{r}}^{(k\rangle}} \delta_{S_{\varepsilon_{k}}(z)} . \tag{6.84}
\end{equation*}
$$

We know that conditional on $\mathcal{U}_{\langle k\rangle}$, the sets $S_{\langle k\rangle} \cap I_{k} \cap Z_{1}^{\langle k\rangle}$ and $S_{\langle k\rangle} \cap I_{k} \cap Z_{\mathrm{r}}{ }^{\langle k\rangle}$ are independent thinnings of the sets $I_{k} \cap Z_{1}^{\langle k\rangle}$ and $I_{k} \cap Z_{\mathrm{r}}^{\langle k\rangle}$, with thinning probabilities $b_{k, 1}$ and $b_{k, \mathrm{r}}$, respectively, which satisfy

$$
\begin{align*}
& b_{k, 1}:=\mathbb{P}\left[\alpha_{z}^{\langle k\rangle \mathrm{r}}=+1 \mid \alpha_{z}^{\langle k\rangle}=-1\right]=\frac{\mathbb{E}\left[\alpha_{z}^{\langle k\rangle \mathrm{r}}-\alpha_{z}^{\langle k\rangle}\right]}{2 \mathbb{P}\left[\alpha_{z}^{\langle k\rangle}=-1\right]} \underset{k \rightarrow \infty}{\sim}\left(\beta_{+}-\beta\right) \varepsilon_{k}  \tag{6.85}\\
& b_{k, \mathrm{r}}:=\mathbb{P}\left[\alpha_{z}^{\langle k\rangle}=-1 \mid \alpha_{z}^{\langle k\rangle}=+1\right]=\frac{\mathbb{E}\left[\alpha_{z}^{\langle k\rangle}-\alpha_{z}^{\langle k\rangle}\right]}{2 \mathbb{P}\left[\alpha_{z}^{\langle k\rangle}=+1\right]} \underset{k \rightarrow \infty}{\sim}\left(\beta-\beta_{-}\right) \varepsilon_{k} .
\end{align*}
$$

Therefore, by formula (D.4) of Appendix D, we have that

$$
\begin{equation*}
\mathbb{E}\left[(1-f)^{\nu_{1}^{\langle k\rangle}}(1-g)^{\nu_{\mathrm{r}}^{\langle k\rangle}} h\left(S_{\varepsilon_{k}}\left(\mathcal{U}_{\langle k\rangle}\right)\right]=\mathbb{E}\left[\left(1-b_{k, 1} f\right)^{\varepsilon_{k}^{-1} \ell_{1}^{\langle k\rangle}}\left(1-b_{k, \mathrm{r}} g\right)^{\varepsilon_{k}^{-1} \ell_{\mathrm{r}}^{\langle k\rangle}} h\left(S_{\varepsilon_{k}}\left(\mathcal{U}_{\langle k\rangle}\right)\right] .\right.\right. \tag{6.86}
\end{equation*}
$$

Recall that we are assuming throughout that our random variables are coupled in such a way that $\mathcal{U}_{\langle k\rangle}$ and $\mathcal{V}_{\langle k\rangle}$, diffusively rescaled with $\varepsilon_{k}$, converge to a Brownian web $\mathcal{W}$ and Brownian net $\mathcal{N}$, respectively. By Proposition 6.10, Lemma 6.7 (a) and the remarks below it we can improve our coupling such that moreover

$$
\begin{equation*}
\ell_{1}^{\langle k\rangle} \underset{k \rightarrow \infty}{\Longrightarrow} \ell_{1}(\Delta, \hat{\Delta}) \quad \text { and } \quad \ell_{\mathrm{r}}^{\langle k\rangle} \underset{k \rightarrow \infty}{\Longrightarrow} \ell_{\mathrm{r}}(\Delta, \hat{\Delta}) \tag{6.87}
\end{equation*}
$$

Thus, (6.83) will follow by taking the limit $k \rightarrow \infty$ in (6.86), provided we show that our coupling can be further improved such that also

$$
\begin{equation*}
\nu_{1}^{\langle k\rangle} \underset{k \rightarrow \infty}{\Longrightarrow} \nu_{1} \quad \text { and } \quad \nu_{\mathrm{r}}^{\langle k\rangle} \underset{k \rightarrow \infty}{\Longrightarrow} \nu_{\mathrm{r}}, \tag{6.88}
\end{equation*}
$$

where $\Rightarrow$ denotes weak convergence of finite measures on $\mathbb{R}^{2}$. Since $z \in S_{1} \cap I$ if and only if $z$ is a separation point of $\mathcal{N}, \operatorname{sign}_{\mathcal{W}}(z)=-1$, and $z$ is entered by a path $\pi \in \mathcal{W}(\Delta)$ and a path $\hat{\pi} \in \hat{\mathcal{W}}(\hat{\Delta})$, formula (6.88) follows from (6.71) and (6.72).

Proof of formula (6.69) Let $\tilde{\mathcal{W}}$ be defined by

$$
\begin{equation*}
\tilde{\mathcal{W}}:=\left\{\pi \in \mathcal{N}: \operatorname{sign}_{\pi}(z)=\operatorname{sign}_{\mathcal{W}}(z) \forall z \in S \text { s.t. } \pi \text { enters } z\right\} \tag{6.89}
\end{equation*}
$$

Then $\mathcal{W} \subset \tilde{\mathcal{W}}$ by the fact that $\mathcal{W} \subset \mathcal{N}$. To prove the other inclusion, let $\mathcal{T}$ be some deterministic countable dense subset of $\mathbb{R}$. Fix $\tilde{\pi} \in \tilde{\mathcal{W}}$. Choose $\sigma_{\tilde{\pi}}<s_{n} \in \mathcal{T}$ with $s_{n} \downarrow \sigma_{\tilde{\pi}}$.

For each $n$, we may choose some $\pi_{n} \in \mathcal{W}$ with $\sigma_{\pi_{n}}=s_{n}$ and $\pi_{n}\left(s_{n}\right)=\tilde{\pi}\left(s_{n}\right)$. Since $\tilde{\pi} \in \mathcal{N}$ is an incoming path at $\left(\tilde{\pi}\left(s_{n}\right), s_{n}\right)$, by Proposition 6.4 (a), this point is of type ( $\mathrm{p}, \mathrm{p}$ ) in $\mathcal{N}$. Using this and the finite graph representation (Proposition 6.5), we see that $\pi_{n}(t)=\tilde{\pi}(t)$ for each $s_{n} \leq t \in \mathcal{T}$, hence $\pi_{n}=\tilde{\pi}$ on $\left[s_{n}, \infty\right)$. It follows that $\pi_{n} \rightarrow \tilde{\pi}$, and therefore, by the compactness of $\mathcal{W}$, that $\tilde{\pi} \in \mathcal{W}$.

Proof of Theorem 4.4. Let $\mathcal{U}_{\langle k\rangle}$ and $\mathcal{V}_{\langle k\rangle}$ be as in Theorem 6.15 and let $\hat{\mathcal{U}}_{\langle k\rangle}$ and $\hat{\mathcal{V}}_{\langle k\rangle}$ be their associated dual discrete web and net. Then, by Theorems 6.8 and 6.11

$$
\begin{equation*}
\mathbb{P}\left[S_{\varepsilon_{k}}\left(\mathcal{U}_{\langle k\rangle}, \mathcal{V}_{\langle k\rangle}, \hat{\mathcal{U}}_{\langle k\rangle}, \hat{\mathcal{V}}_{\langle k\rangle}\right) \in \cdot\right] \underset{k \rightarrow \infty}{\Longrightarrow} \mathbb{P}[(\mathcal{W}, \mathcal{N}, \hat{\mathcal{W}}, \hat{\mathcal{N}}) \in \cdot] \tag{6.90}
\end{equation*}
$$

where $(\mathcal{W}, \mathcal{N})$ are a coupled Brownian web and net as in Theorem 6.15 and $\hat{\mathcal{W}}, \hat{\mathcal{N}}$ are the duals of $\mathcal{W}, \mathcal{N}$. Since $\left(-\hat{\mathcal{U}}_{\langle k\rangle},-\hat{\mathcal{V}}_{\langle k\rangle}\right)$ is equally distributed with $\left(\mathcal{U}_{\langle k\rangle}, \mathcal{V}_{\langle k\rangle}\right)$, we see that $(-\hat{\mathcal{W}},-\hat{\mathcal{N}})$ is equally distributed with $(\mathcal{W}, \mathcal{N})$. Now all statements in Theorem 4.4 follow from Theorem 6.15.

### 6.6 Switching and hopping in the Brownian web and net

In this section, we apply Theorem 4.4 together with the finite graph representation developed in Section 6.2 to prove Proposition 4.5 on switching and hopping inside a Brownian net. We then apply Theorem 4.4 and Proposition 4.5 to give short proofs of the marking construction of sticky Brownian webs (Theorem 3.5) and the Brownian net (Theorem 4.6), Proposition 3.6 on changing the reference web, and the equivalence of the definitions of a left-right Brownian web given in Sections 3.3 and 4.1. We also formulate and prove a result on the construction of sticky Brownian webs inside a Brownian net, analogous to Theorem 4.4.
Proof of Proposition 4.5 For each set $\Delta \subset S$, set

$$
\begin{align*}
\mathcal{N}_{\Delta} & =\left\{\pi \in \mathcal{N}: \operatorname{sign}_{\pi}(z)=\alpha_{z} \forall z \in S \backslash \Delta \text { s.t. } \pi \text { enters } z\right\} \\
\mathcal{W}_{\Delta} & =\left\{\pi \in \mathcal{N}: \operatorname{sign}_{\pi}(z)=-\alpha_{z} \forall z \in \Delta \text { s.t. } \pi \text { enters } z\right\} \cap \mathcal{N}_{\Delta} . \tag{6.91}
\end{align*}
$$

Since $\mathcal{N}_{S^{\prime}}, \mathcal{N}_{\Delta_{n}}, \mathcal{W}_{S^{\prime}}, \mathcal{W}_{\Delta_{n}}$ are contained in the compact set $\mathcal{N}$, it suffices to prove the following statements:

1. $\mathcal{W}_{\Delta_{n}}=\operatorname{switch}_{\Delta_{n}}(\mathcal{W})$ and $\mathcal{N}_{\Delta_{n}}=\operatorname{hop}_{\Delta_{n}}(\mathcal{W})$.
2. $\mathcal{N}_{S^{\prime}}, \mathcal{N}_{\Delta_{n}}, \mathcal{W}_{S^{\prime}}, \mathcal{W}_{\Delta_{n}}$ are closed sets.
3. $\mathcal{N}_{\Delta_{n}} \rightarrow \mathcal{N}_{S^{\prime}}$ and $\mathcal{W}_{\Delta_{n}} \rightarrow \mathcal{W}_{S^{\prime}}$.
4. Since $\operatorname{hop}_{\Delta_{n}}(\mathcal{W})=\bigcup_{\Delta^{\prime} \subset \Delta_{n}} \operatorname{switch}_{\Delta_{n}}(\mathcal{W})$ and $\mathcal{N}_{\Delta_{n}}=\bigcup_{\Delta^{\prime} \subset \Delta_{n}} \mathcal{W}_{\Delta_{n}}$ it suffices to prove that $\mathcal{W}_{\Delta}=\operatorname{switch}_{\Delta}(\mathcal{W})$ for each finite $\Delta \subset S$. By induction, it suffices to prove that $\mathcal{W}_{\Delta \cup\{z\}}=$ $\operatorname{switch}_{z}\left(\mathcal{W}_{\Delta}\right)$ for each finite $\Delta \subset S$ and $z \in S$. Here, by induction, we have that $\mathcal{W}_{\Delta}$ is a subset of $\mathcal{N}$ such that for each $z \in S$ the set $\mathcal{W}_{\Delta}(z)$ contains exactly two paths, say $\pi_{1}, \pi_{2}$, of which exactly one, say $\pi_{1}$, is the continuation of a path in the set $\mathcal{W}_{\Delta, \text { in }}(z)$ of paths in $\mathcal{W}_{\Delta}$ entering $z$. By definition, writing $z=(x, t)$, one has

$$
\begin{equation*}
\operatorname{switch}_{z}\left(\mathcal{W}_{\Delta}\right)=\left(\mathcal{W}_{\Delta} \backslash \mathcal{W}_{\Delta, \text { in }}(z)\right) \cup\left\{\pi^{t} \cup \pi_{2}: \pi \in \mathcal{W}_{\Delta, \text { in }}(z)\right\} \tag{6.92}
\end{equation*}
$$

By the structure of separation points (Proposition 4.3 (c) and (d)) and the fact that the net is closed under hopping between paths at intersection times [S508, Prop. 1.4], it follows that
each path of the form $\pi^{\prime}:=\pi_{1}^{t} \cup \pi_{2}$ with $\pi \in \mathcal{W}_{\Delta, \text { in }}(z)$ is an element of $\mathcal{N}$ and satisfies $\operatorname{sign}_{\pi^{\prime}}(z)=-\operatorname{sign}_{\pi_{1}}(z)$, proving that $\operatorname{switch}_{z}\left(\mathcal{W}_{\Delta}\right) \subset \mathcal{W}_{\Delta \cup\{z\}}$. Conversely, each $\pi^{\prime} \in \mathcal{W}_{\Delta \cup\{z\}}$ that enters $z$ is of the form $\pi^{\prime}=\pi^{t} \cup \pi_{2}$ where $\pi:=\pi^{\prime t} \cup \pi_{1} \in \mathcal{W}_{\Delta, \text { in }}(z)$, showing that $\operatorname{switch}_{z}\left(\mathcal{W}_{\Delta}\right) \supset \mathcal{W}_{\Delta \cup\{z\}}$.
2. It suffices to prove that if $\pi \in \mathcal{N}$ enters some point $z=(x, t) \in S$ and $\pi_{n} \in \mathcal{N}$ satisfy $\pi_{n} \rightarrow \pi$ and $\operatorname{sign}_{\pi_{n}}(z)=\alpha$ whenever $\pi_{n}$ enters $z$, then $\operatorname{sign}_{\pi}(z)=\alpha$. By symmetry, it suffices to treat the case $\alpha=-1$. We start by noting that there exists an $N$ such that for each $n \geq N$, the path $\pi_{n}$ enters $z$. This follows from the fact that, by Proposition 4.3, there exist dual paths $\hat{l}_{z}^{\prime}$ and $\hat{r}_{z}^{\prime}$ forming a dual mesh $\hat{M}\left(\hat{l}_{z}^{\prime}, \hat{r}_{z}^{\prime}\right)$, and each path in $\mathcal{N}$ starting in $\hat{M}\left(\hat{l}_{z}^{\prime}, \hat{r}_{z}^{\prime}\right)$ must enter $z$ [SSS09, Lemma 3.3]. Since $\operatorname{sign}_{\pi_{n}}(z)=-1$ for all $n \geq N$, we have that $\pi_{n} \leq r_{z}^{\prime}$ on $[t, \infty)$ for all $n \geq N$, hence the same holds for $\pi$ and $\operatorname{sign}_{\pi}(z)=-1$.
3. Since $\mathcal{N}_{\Delta_{n}}, \mathcal{W}_{\Delta_{n}} \subset \mathcal{N}$ and $\mathcal{N}$ is compact, by Lemma B. 3 in the appendix, the sets $\left\{\mathcal{N}_{\Delta_{n}}\right\}$ and $\left\{\mathcal{W}_{\Delta_{n}}\right\}$ are precompact, so by going to a subsequence if necessary, we may assume that $\mathcal{N}_{\Delta_{n}} \rightarrow \mathcal{N}^{*}$ and $\mathcal{W}_{\Delta_{n}} \rightarrow \mathcal{W}^{*}$ for some $\mathcal{N}^{*}, \mathcal{W}^{*} \in \mathcal{K}(\Pi)$. We need to show that $\mathcal{N}^{*}=\mathcal{N}_{S^{\prime}}$ and $\mathcal{W}^{*}=\mathcal{W}_{S^{\prime}}$. We observe that $\mathcal{W}_{\Delta_{n}}, \mathcal{N}_{\Delta_{n}} \subset \mathcal{N}_{S^{\prime}}$ so $\mathcal{W}^{*}, \mathcal{N}^{*} \subset \mathcal{N}_{S^{\prime}}$. Set

$$
\begin{equation*}
\tilde{\mathcal{W}}_{\Delta}:=\left\{\pi \in \mathcal{N}: \operatorname{sign}_{\pi}(z)=-\alpha_{z} \forall z \in \Delta \text { s.t. } \pi \text { enters } z\right\} \cap \mathcal{N}_{S^{\prime}} . \tag{6.93}
\end{equation*}
$$

Then $\mathcal{W}_{\Delta_{n}} \subset \tilde{\mathcal{W}}_{\Delta_{m}}$ for all $n \geq m$, so letting $n \rightarrow \infty$ we see that $\mathcal{W}^{*} \subset \tilde{\mathcal{W}}_{\Delta_{m}}$ for each $m$, hence $\mathcal{W}^{*} \subset \bigcap_{m} \mathcal{W}_{\Delta_{m}}=\mathcal{W}_{S^{\prime}}$.

To prove the opposite inclusions, we must show that for each $\pi \in \mathcal{N}_{S^{\prime}}$ there exist $\pi_{n} \in \mathcal{N}_{\Delta_{n}}$ such that $\pi_{n} \rightarrow \pi$ and likewise, for each $\pi \in \mathcal{W}_{S^{\prime}}$ there exist $\pi_{n} \in \mathcal{W}_{\Delta_{n}}$ such that $\pi_{n} \rightarrow \pi$. Let $\mathcal{T}$ be some deterministic countable dense subset of $\mathbb{R}$ and let $\pi \in \mathcal{N}_{S^{\prime}}$. Let $T_{1}, \ldots, T_{m} \in \mathcal{T}$ be such that $\sigma_{\pi}<T_{1}<\cdots<T_{m}$ and let $R_{T_{k}, T_{k+1}}$ be as in (6.4). We observe that at each point in $\mathbb{R}^{2}$ there starts at least one path in $\mathcal{N}_{\Delta_{n}}=\bigcup_{\Delta^{\prime} \subset \Delta_{n}}$ switch $_{\Delta^{\prime}}(\mathcal{W})$. Therefore, for each $n$ we can find some $\pi_{n} \in \mathcal{N}_{\Delta_{n}}$ such that $\sigma_{\pi_{n}}=T_{1}$ and $\pi_{n}\left(T_{1}\right)=\pi\left(T_{1}\right)$. Provided $n$ is sufficiently large, we may moreover choose $\pi_{n}$ with the property that $\operatorname{sign}_{\pi_{n}}(z)=\operatorname{sign}_{\pi}(z)$ for each point $z \in \bigcup_{k=1}^{m-1} R_{T_{k}, T_{k+1}}$ such that both $\pi_{n}$ and $\pi$ enter $z$, hence by Corollary 6.6. we conclude that $\pi_{n}\left(T_{k}\right)=\pi\left(T_{k}\right)$ for $k=1, \ldots, m$.

Thus, we have shown that for each finite set $T \subset\left(\sigma_{\pi}, \infty\right) \cap \mathcal{T}$ there exists an $N$ such that for all $n \geq N$ there exists some $\pi_{n} \in \mathcal{N}_{\Delta_{n}}$ with $\pi_{n}=\pi$ on $T$. Choosing $T_{m} \uparrow\left(\sigma_{\pi}, \infty\right) \cap \mathcal{T}$, using the compactness of $\mathcal{N}$, going to a subsequence if necessary, we can find $\pi_{n_{m}} \in \mathcal{N}_{\Delta_{n_{m}}}$ such that $\pi_{n_{m}} \rightarrow \pi$ locally uniformly on $\left(\sigma_{\pi}, \infty\right)$. Cutting off a piece of $\pi_{n_{m}}$ if necessary to make the starting times converge, we have found $\mathcal{N}_{\Delta_{n_{m}}} \ni \pi_{n_{m}} \rightarrow \pi$, proving that $\mathcal{N}^{*} \supset \mathcal{N}_{S^{\prime}}$. The proof that $\mathcal{W}^{*} \supset \mathcal{W}_{S^{\prime}}$ is completely analogous.

Proof of Theorem 3.5. Let $\beta \in \mathbb{R}$ and $c_{1}, c_{\mathrm{r}} \geq 0$. In Theorem 4.4, set $\beta_{-}:=\beta-c_{\mathrm{r}}$, $\beta_{+}:=\beta+c_{1}$, let $r:=c_{\mathrm{r}} /\left(c_{1}+c_{\mathrm{r}}\right)$ if $c_{1}+c_{\mathrm{r}}>0$, and choose some arbitrary $r \in[0,1]$ otherwise. Then $\mathcal{W}$, defined in (4.9), is a Brownian web with drift $\beta$ and conditional on $\mathcal{W}$, the set $S$ is a Poisson point set with intensity $c_{1} \ell_{1}+c_{\mathrm{r}} \ell_{\mathrm{r}}$. In Theorem 3.5, we may without loss of generality assume that $\mathcal{W}$ and $S$ are constructed in this way. Then Proposition 4.5 tells us that the limit $\mathcal{W}^{\prime}=\lim _{\Delta_{n} \uparrow S} \operatorname{switch}_{\Delta_{n}}(\mathcal{W})$ exists, does not depend on the choice of the $\Delta_{n}$, and is given by

$$
\begin{equation*}
\mathcal{W}^{\prime}=\left\{\pi \in \mathcal{N}: \operatorname{sign}_{\hat{\pi}}(z)=-\alpha_{z} \forall z \in S \text { s.t. } \hat{\pi} \text { enters } z\right\} . \tag{6.94}
\end{equation*}
$$

By Theorem 4.4, the dual webs $\hat{\mathcal{W}}, \hat{\mathcal{W}}^{\prime}$ associated with $\mathcal{W}, \mathcal{W}^{\prime}$ are given by

$$
\begin{align*}
\hat{\mathcal{W}} & =\left\{\hat{\pi} \in \hat{\mathcal{N}}: \operatorname{sign}_{\hat{\pi}}(z)=\alpha_{z} \forall z \in S \text { s.t. } \hat{\pi} \text { enters } z\right\}, \\
\hat{\mathcal{W}}^{\prime} & =\left\{\hat{\pi} \in \hat{\mathcal{N}}: \operatorname{sign}_{\hat{\pi}}(z)=\alpha_{z}^{\prime} \forall z \in S \text { s.t. } \hat{\pi} \text { enters } z\right\}, \tag{6.95}
\end{align*}
$$

so Proposition 4.5 tells us that $\hat{\mathcal{W}}^{\prime}=\lim _{\Delta_{n} \uparrow S} \operatorname{switch}_{\Delta_{n}}(\hat{\mathcal{W}})$. Since conditional on $\mathcal{N}$, the $\left(-\alpha_{z}\right)_{z \in S}$ are i.i.d. with parameter $1-r$, by Theorem 4.4. the Brownian web $\mathcal{W}^{\prime}$ has drift $\beta^{\prime}=r \beta_{-}+(1-r) \beta_{+}=\beta+c_{1}-c_{\mathrm{r}}$.

Proof of Proposition 3.6 (ii) and (iii). We continue to assume that $\mathcal{W}$ and $\mathcal{W}^{\prime}$ are defined inside a Brownian net $\mathcal{N}$ as in the proof of Theorem 3.5. Set $S_{1}:=\left\{z \in S: \alpha_{z}=-1\right\}$ and $S_{\mathrm{r}}:=\left\{z \in S: \alpha_{z}=+1\right\}$. Then, by Theorem 4.4, conditional on $\mathcal{W}$, the set $S_{1}$ is a Poisson point set with intensity $c_{1} \ell_{1}$ and the set $S_{\mathrm{r}}$ is a Poisson point set with intensity $c_{\mathrm{r}} \ell_{\mathrm{r}}$, and likewise, conditional on $\mathcal{W}^{\prime}$, the set $S_{1}$ is Poisson with intensity $c_{\mathrm{r}} \ell_{1}^{\prime}$ and $S_{\mathrm{r}}$ is Poisson with intensity $c_{1} \ell_{\mathrm{r}}^{\prime}$. In particular, this implies that a.s., each point $z \in S_{1}$ is of type $(1,2)_{1}$ in $\mathcal{W}$ and of type $(1,2)_{\mathrm{r}}$ in $\mathcal{W}^{\prime}$. Conversely, if $z \in \mathbb{R}^{2}$ is of type $(1,2)_{1}$ in $\mathcal{W}$ and of type $(1,2)_{\mathrm{r}}$ in $\mathcal{W}^{\prime}$, then $z$ is a separation point of some paths $\pi \in \mathcal{W}$ and $\pi^{\prime} \in \mathcal{W}^{\prime}$ and therefore, by the definition of separation points of $\mathcal{N}$ given before Proposition 4.3, $z \in S$.

Proof of Theorem 4.6. As in the previous two proofs, without loss of generality, we assume that $\mathcal{W}$ is embedded in a Brownian net $\mathcal{N}$ as in (4.9) and that $S_{1}:=\left\{z \in S: \alpha_{z}=-1\right\}$ and $S_{\mathrm{r}}:=\left\{z \in S: \alpha_{z}=+1\right\}$. Then, by Proposition 4.5, (4.14) (i) holds and the limits in (4.14) (ii) and (iii) exist and are given by

$$
\begin{align*}
& \mathcal{W}^{\mathrm{l}}=\left\{\pi \in \mathcal{N}: \operatorname{sign}_{\hat{\pi}}(z)=-1 \forall z \in S \text { s.t. } \hat{\pi} \text { enters } z\right\},  \tag{6.96}\\
& \mathcal{W}^{\mathrm{r}}=\left\{\pi \in \mathcal{N}: \operatorname{sign}_{\hat{\pi}}(z)=+1 \forall z \in S \text { s.t. } \hat{\pi} \text { enters } z\right\} .
\end{align*}
$$

By Theorem 4.4, $\left(\mathcal{W}^{1}, \mathcal{W}^{\mathrm{r}}\right)$ is the left-right Brownian web associated with $\mathcal{N}$. Since $S_{1}$ and $S_{\mathrm{r}}$ are Poisson point sets with intensities $c_{1} \ell_{1}$ and $c_{\mathrm{r}} \ell_{\mathrm{r}}$, respectively, each $z \in S$ is of type $(1,2)$ in $\mathcal{W}$ and $\operatorname{sign}_{\mathcal{W}}(z)=\alpha_{z}$, so by construction, conditional on $\mathcal{N}$, the random variables $\left(\operatorname{sign}_{\mathcal{W}}(z)\right)_{z \in S}$ are i.i.d. with $\mathbb{P}\left[\operatorname{sign}_{\mathcal{W}}(z)=+1 \mid \mathcal{N}\right]=r=c_{\mathrm{r}} /\left(c_{1}+c_{\mathrm{r}}\right)$.

To prepare for the proof of Proposition 3.6 (i), we need a lemma.
Lemma 6.16 (Sticky Brownian webs inside a Brownian net) Let $\mathcal{N}^{*}$ be a Brownian net with left and right speeds $\beta_{-}^{*} \leq \beta_{+}^{*}$ and set of separation points $S^{*}$. Conditional on $\mathcal{N}^{*}$, let $\left(\alpha_{z}, \alpha_{z}^{\prime}\right)_{z \in S}$ be an i.i.d. collection of random variables with values in $\{-1,+1\}^{2}$. Set

$$
\begin{align*}
p_{--} & :=\mathbb{P}\left[\left(\alpha_{z}, \alpha_{z}^{\prime}\right)=(-1,-1) \mid \mathcal{N}\right],  \tag{6.97}\\
S_{--} & :=\left\{z \in S^{*}:\left(\alpha_{z}, \alpha_{z}^{\prime}\right)=(-1,-1)\right\},
\end{align*}
$$

and let $p_{-+}, p_{+-}, p_{++}$and $S_{-+}, S_{+-}, S_{++}$be defined analogously. Set

$$
\text { (i) } \mathcal{W}:=\left\{\pi \in \mathcal{N}^{*}: \operatorname{sign}_{\pi}(z)=\alpha_{z} \forall z \in S^{*} \text { s.t. } \pi \text { enters } z\right\} \text {, }
$$

$$
\begin{equation*}
\text { (ii) } \mathcal{W}^{\prime}:=\left\{\pi \in \mathcal{N}^{*}: \operatorname{sign}_{\pi}(z)=\alpha_{z}^{\prime} \forall z \in S^{*} \text { s.t. } \pi \text { enters } z\right\} \text {. } \tag{6.98}
\end{equation*}
$$

Then $\mathcal{W}$ is a Brownian web with drift $\beta:=\left(p_{--}+p_{-+}\right) \beta_{-}^{*}+\left(p_{+-}+p_{++}\right) \beta_{+}^{*}$ and $\mathcal{W}^{\prime}$ is a Brownian web with drift $\beta^{\prime}:=\left(p_{--}+p_{+-}\right) \beta_{-}^{*}+\left(p_{-+}+p_{++}\right) \beta_{+}^{*}$.

Let $\ell$ denote the intersection local time measure between $\mathcal{W}$ and its dual, let $\ell_{1}, \ell_{\mathrm{r}}$ denote the restrictions of $\ell$ to the sets of points of type $(1,2)_{1}$ and $(1,2)_{\mathrm{r}}$ in $\mathcal{W}$, respectively, and let $\ell^{\prime}, \ell_{1}^{\prime}, \ell_{\mathrm{r}}^{\prime}$ be the same objects defined for $\mathcal{W}^{\prime}$. Then, conditional on $\mathcal{W}$, the sets

$$
\begin{equation*}
S_{--}, \quad S_{-+}, \quad S_{+-}, \quad S_{++} \tag{6.99}
\end{equation*}
$$

are independent Poisson point sets with respective intensities

$$
\begin{equation*}
p_{--}\left(\beta_{+}^{*}-\beta_{-}^{*}\right) \ell_{1}, \quad p_{-+}\left(\beta_{+}^{*}-\beta_{-}^{*}\right) \ell_{1}, \quad p_{+-}\left(\beta_{+}^{*}-\beta_{-}^{*}\right) \ell_{\mathrm{r}}, \quad p_{++}\left(\beta_{+}^{*}-\beta_{-}^{*}\right) \ell_{\mathrm{r}}, \tag{6.100}
\end{equation*}
$$

while conditional on $\mathcal{W}^{\prime}$, the sets in (6.99) are independent Poisson point sets with respective intensities

$$
\begin{equation*}
p_{--}\left(\beta_{+}^{*}-\beta_{-}^{*}\right) \ell_{1}^{\prime}, \quad p_{-+}\left(\beta_{+}^{*}-\beta_{-}^{*}\right) \ell_{\mathrm{r}}^{\prime}, \quad p_{+-}\left(\beta_{+}^{*}-\beta_{-}^{*}\right) \ell_{1}^{\prime}, \quad p_{++}\left(\beta_{+}^{*}-\beta_{-}^{*}\right) \ell_{\mathrm{r}}^{\prime} . \tag{6.101}
\end{equation*}
$$

Moreover, one has

$$
\begin{equation*}
\text { (i) } \mathcal{W}^{\prime}=\lim _{\Delta_{n} \uparrow S_{-+} \cup S_{+-}} \operatorname{switch}_{\Delta_{n}}(\mathcal{W}) \tag{6.102}
\end{equation*}
$$

$$
\text { (ii) } \mathcal{W}=\lim _{\Delta_{n} \uparrow S_{-+} \cup S_{+-}} \operatorname{switch}_{\Delta_{n}}\left(\mathcal{W}^{\prime}\right)
$$

Proof. By Theorem 4.4, formulas (6.98) (i) and (ii) define Brownian webs with drifts as claimed. By Proposition 4.5, the limits in (6.102) exist and coincide with the objects defined in (6.98). By Theorem 4.4, conditional on $\mathcal{W}$, the sets $S_{1}:=S_{--} \cup S_{-+}$and $S_{\mathrm{r}}:=S_{+-} \cup S_{++}$are independent Poisson point sets with intensities $\left(p_{--}+p_{-+}\right)\left(\beta_{+}^{*}-\beta_{-}^{*}\right) \ell_{1}$ and $\left(p_{+-}+p_{++}\right)\left(\beta_{+}^{*}-\right.$ $\left.\beta_{-}^{*}\right) \ell_{\mathrm{r}}$, respectively. In particular, this implies that each $z \in S_{\mathrm{l}}$ (resp. $z \in S_{\mathrm{r}}$ ) is of type (1,2) (resp. $\left.(1,2)_{\mathrm{r}}\right)$ in $\mathcal{W}$.

We claim that the $\sigma$-fields generated by, on the one hand, $\mathcal{W}$ and $S^{*}$, and, on the other hand, $\mathcal{N}^{*}$ and the collection of random variables $\alpha=\left(\alpha_{z}\right)_{z \in S^{*}}$ are identical. To see this, we note that by Proposition 4.5, $\mathcal{N}=\lim _{\Delta_{n} \uparrow S^{*}} \operatorname{hop}_{\Delta_{n}}(\mathcal{W})$. Since moreover $\alpha_{z}=\operatorname{sign}_{\mathcal{W}}(z)$ for all $z \in S^{*}$, this shows that $\overline{\mathcal{N}^{*}}$ and $\alpha$ are a.s. uniquely determined by $\mathcal{W}$ and $S^{*}$. Conversely, since $\mathcal{W}$ is given by (6.98) (i) and $S^{*}$ is the set of separation points of $\mathcal{N}^{*}$, we see that $\mathcal{W}$ and $S^{*}$ are a.s. uniquely determined by $\mathcal{N}^{*}$ and $\alpha$.

Conditional on $\mathcal{N}^{*}$ and $\alpha$, the random variables $\left(\alpha_{z}^{\prime}\right)_{z \in S^{*}}$ are independent, where $\mathbb{P}\left[\alpha_{z}^{\prime}=\right.$ $\left.+1 \mid\left(\mathcal{N}^{*}, \alpha\right)\right]$ equals $p_{-+} /\left(p_{--}+p_{-+}\right)$if $\alpha_{z}=-1$ and $p_{++} /\left(p_{+-}+p_{++}\right)$if $\alpha_{z}=+1$. It follows that conditional on $\mathcal{W}$ and $S^{*}$, the set $S_{-+}$is obtained from $S_{--} \cup S_{-+}$by independent thinning with probability $p_{-+} /\left(p_{--}+p_{-+}\right)$and likewise, the set $S_{++}$is obtained from $S_{+-} \cup S_{++}$by independent thinning with probability $p_{++} /\left(p_{+-}+p_{++}\right)$. Since independent thinning splits a Poisson point set in two independent Poisson point sets, we conclude that conditional on $\mathcal{W}$, the sets in (6.99) are independent Poisson point sets with intensities given in (6.100). By symmetry, an analogue statement holds for $\mathcal{W}^{\prime}$, i.e., conditional on $\mathcal{W}^{\prime}$, the sets in 6.99) are independent Poisson point sets with intensities given in 6.101.

Proof of Proposition 3.6 (i). By symmetry, it suffices to show that $\ell_{1}=\ell_{1}^{\prime}$. Let $\beta \in \mathbb{R}$ and $c_{1}, c_{\mathrm{r}} \geq 0$. In Lemma 6.16. set $\beta_{-}^{*}:=\beta-c_{\mathrm{r}}, \beta_{+}^{*}:=\beta+c_{1}+1$, and let $p_{--}:=1 /\left(1+c_{1}+c_{\mathrm{r}}\right)$, $p_{-+}:=c_{1} /\left(1+c_{1}+c_{\mathrm{r}}\right), p_{+-}:=c_{\mathrm{r}} /\left(1+c_{1}+c_{\mathrm{r}}\right)$, and $p_{++}:=0$. Let $\mathcal{W}, \mathcal{W}^{\prime}$ be as in 66.98) and set $S:=S_{-+} \cup S_{+--}$. Then conditional on $\mathcal{W}$, the set $S$ is a Poisson point set with intensity $c_{1} \ell_{1}+c_{\mathrm{r}} \ell_{\mathrm{r}}$ and $\mathcal{W}^{\prime}=\lim _{\Delta_{n} \uparrow S} \operatorname{switch}_{\Delta_{n}}(\mathcal{W})$. Without loss of generality, we may assume that the sticky Brownian webs in Proposition 3.6 are constructed in this way.

It follows from (6.102) that the $\sigma$-fields generated by, on the one hand $\mathcal{W}$ and $S$, and, on the other hand, $\mathcal{W}^{\prime}$ and $S$ coincide. By 6.100 and 6.101, conditional on this $\sigma$-field, the set $S_{--}$is a Poisson point set with intensity $\ell_{1}$ and also a Poisson point set with intensity $\ell_{1}^{\prime}$, i.e., the conditional law $\mathbb{P}\left[S_{--} \in \cdot \mid \mathcal{W}, S\right]$ is the law of a Poisson point set with intensity $\ell_{1}$ and also the law of a Poisson point set with intensity $\ell_{1}^{\prime}$. This is possible only if $\ell_{1}=\ell_{1}^{\prime}$.

The following lemma sometimes comes in handy.

Lemma 6.17 (Commutativity of switching) Let $\mathcal{W}$ be a Brownian web wih drift $\beta$, let $\ell$ be the intersection local time measure between $\mathcal{W}$ and its dual and let $\ell_{1}, \ell_{\mathrm{r}}$ denote the restrictions of $\ell$ to the sets of points of type $(1,2)_{1}$ and $(1,2)_{\mathrm{r}}$ in $\mathcal{W}$, respectively. Let $c_{1}, c_{\mathrm{r}}, c_{1}^{\prime}, c_{\mathrm{r}}^{\prime} \geq 0$ be constants and conditional on $\mathcal{W}$, let $S, S^{\prime}$ be independent Poisson point sets with intensities $c_{1} \ell_{1}+c_{\mathrm{r}} \ell_{\mathrm{r}}$ and $c_{1}^{\prime} \ell_{1}+c_{\mathrm{r}}^{\prime} \ell_{\mathrm{r}}$, respectively. Then

$$
\begin{equation*}
\lim _{\Delta_{n} \uparrow S} \operatorname{switch}_{\Delta_{n}}\left(\lim _{\Delta_{m}^{\prime} \uparrow S^{\prime}} \operatorname{switch}_{\Delta_{m}^{\prime}}(\mathcal{W})\right)=\lim _{\Delta_{k}^{\prime \prime} \uparrow S \cup S^{\prime}} \operatorname{switch}_{\Delta_{k}^{\prime \prime}}(\mathcal{W}) . \tag{6.103}
\end{equation*}
$$

Proof. Choose $\beta_{-}^{*} \leq \beta_{+}^{*}$ and $p_{--}, \ldots, p_{++}$, summing up to one, such that $c_{1}=p_{--}\left(\beta_{+}^{*}-\beta_{-}^{*}\right)$, $c_{1}^{\prime}=p_{-+}\left(\beta_{+}^{*}-\beta_{-}^{*}\right), c_{\mathrm{r}}=p_{+-}\left(\beta_{+}^{*}-\beta_{-}^{*}\right), c_{1}^{\prime}=p_{++}\left(\beta_{+}^{*}-\beta_{-}^{*}\right)$, and $\beta=\left(p_{--}+p_{-+}\right) \beta_{-}^{*}+$ $\left(p_{+-}+p_{++}\right) \beta_{+}^{*}$. Then, without loss of generality, we may assume that $\mathcal{W}$ is constructed inside a Brownian net $\mathcal{N}^{*}$ as in Lemma 6.16 and that $S=S_{--} \cup S_{+-}$and $S^{\prime}=S_{-+} \cup S_{++}$. Now Proposition 4.5 tells us that both sides of (6.103) are well-defined and given by

$$
\begin{equation*}
\mathcal{W}^{\prime \prime}=\left\{\pi \in \mathcal{N}^{*}: \operatorname{sign}_{\pi}(z) \neq \alpha_{z} \forall z \in S^{*} \text { s.t. } \pi \text { enters } z\right\} . \tag{6.104}
\end{equation*}
$$

The following Lemma has been announced in Section 3.3.
Lemma 6.18 (Equivalent definitions of left-right Brownian web) A pair of Brownian webs $\left(\mathcal{W}^{\mathrm{l}}, \mathcal{W}^{\mathrm{r}}\right)$ is a left-right Brownian web with drifts $\beta_{-}, \beta_{+}$as defined in Section 4.1 if and only if $\left(\mathcal{W}^{1}, \mathcal{W}^{\mathrm{r}}\right)$ is a pair of sticky Brownian webs with drifts $\beta_{-}, \beta_{+}$and coupling parameter $\kappa=0$, as defined in Section 3.3.

Proof. Let $\mathcal{N}$ be a Brownian net with left and right speeds $\beta_{-}, \beta_{+}$and let $S$ be its set of separation points. Then, by Theorem 4.4, the left-right Brownian web $\left(\mathcal{W}^{1}, \mathcal{W}^{\mathrm{r}}\right)$ associated with $\mathcal{N}$ is given by

$$
\begin{align*}
\mathcal{W}^{\mathrm{l}} & =\left\{\pi \in \mathcal{N}: \operatorname{sign}_{\pi}(z)=-1 \forall z \in S \text { s.t. } \pi \text { enters } z\right\},  \tag{6.105}\\
\mathcal{W}^{\mathrm{r}} & =\left\{\pi \in \mathcal{N}: \operatorname{sign}_{\pi}(z)=+1 \forall z \in S \text { s.t. } \pi \text { enters } z\right\} .
\end{align*}
$$

Moreover, by the same theorem, if $\ell$ denotes the intersection local time measure between $\mathcal{W}^{1}$ and its dual, and let $\ell_{1}$ and $\ell_{\mathrm{r}}$ denote the restrictions of $\ell$ to the sets of points of type $(1,2)_{1}$ and $(1,2)_{\mathrm{r}}$ in $\mathcal{W}^{\mathrm{l}}$, respectively, then conditional on $\mathcal{W}^{1}$, the set $S$ is a Poisson point set with intensity $\left(\beta_{+}-\beta_{-}\right) \ell_{1}$. By Proposition 4.5, it follows that $\mathcal{W}^{\mathrm{r}}=\lim _{\Delta_{n} \uparrow S} \operatorname{switch}_{\Delta_{n}}\left(\mathcal{W}^{\mathrm{l}}\right)$, hence $\left(\mathcal{W}^{\mathrm{l}}, \mathcal{W}^{\mathrm{r}}\right)$ is a pair of sticky Brownian webs with drifts $\beta_{-}, \beta_{+}$and coupling parameter $\kappa=0$ as defined in Section 3.3.

Conversely, if $\mathcal{W}^{1}$ is a Brownian web with drift $\beta_{-}$and if conditional on $\mathcal{W}^{1}$, the set $S$ is a Poisson point set with intensity $\left(\beta_{+}-\beta_{-}\right) \ell_{1}$, and $\mathcal{W}^{\mathrm{r}}=\lim _{\Delta_{n} \uparrow S} \operatorname{switch}_{\Delta_{n}}\left(\mathcal{W}^{\mathrm{l}}\right)$, then by Theorem 4.4, we may assume without loss of generality that $\mathcal{W}^{1}$ is defined inside a Brownian net $\mathcal{N}$ such that $S$ is the set of separation points of $\mathcal{N}$. Now Proposition 4.5 tells us that $\mathcal{W}^{\mathrm{r}}$ has the representation in 6.105), hence by Theorem $4.4\left(\mathcal{W}^{\mathrm{l}}, \mathcal{W}^{\mathrm{r}}\right)$ is the left-right Brownian web associated with $\mathcal{N}$.

## 7 Construction and convergence of Howitt-Warren flows

In this section, we prove our main results. We start in Section 7.1 with the proof of Theorem 3.9 on the convergence of the quenched laws on the space of webs. In Section 7.2, we then use this to show that the $n$-point motions of the sample web constructed in Theorem 3.7 solve the Howitt-Warren martingale problem, thereby identifying the stochastic flow of kernels there as a Howitt-Warren flow. Here we also prove the construction of Howitt-Warren flows inside a Brownian net (Theorem 4.7) and a result on the exchangeability of the reference and sample Brownian webs from Theorem 3.7 if $\nu_{1}=\nu_{\mathrm{r}}$. In Section 7.3, finally, we harvest some immediate consequences of our construction, such as scaling (Proposition 2.4) and the existence of regular versions of Howitt-Warren flows (Proposition 2.3 and 3.8).

### 7.1 Convergence of quenched laws

In this section, we prove Theorem 3.9. The measures $S_{\varepsilon_{k}}\left(\mathbf{Q}_{\langle k\rangle}\right)$ and $\mathbb{Q}$ from Theorem 3.9 are random probability measures on the Polish space $\mathcal{K}(\Pi)$. Therefore, by Daw91, Thm. 3.2.9], the convergence in $(3.27)$ is equivalent to the convergence of the moment measures of $S_{\varepsilon_{k}}\left(\mathbf{Q}_{\langle k\rangle}\right)$ to the moment measures of $\mathbb{Q}$.

We start by describing these moment measures. Let $\left(\mathcal{W}_{0}, \mathcal{M}\right)$ be a marked reference web as in Section 3.4 and conditional on $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, let $\mathcal{W}_{1}, \mathcal{W}_{2}, \ldots$ be an i.i.d. sequence of sample webs constructed as in (3.19). Then the unconditional law

$$
\begin{equation*}
\mathbb{P}\left[\left(\mathcal{W}_{1}, \ldots, \mathcal{W}_{n}\right) \in \cdot\right] \tag{7.1}
\end{equation*}
$$

is the $n$-th moment measure of $\mathbb{Q}$. Similarly, for each $k$, conditional on an i.i.d. collection of $[0,1]$-valued random variables $\omega^{\langle k\rangle}=\left(\omega_{z}^{\langle k\rangle}\right)_{z \in \mathbb{Z}_{\text {even }}^{2}}$ with law $\mu_{k}$, let $\alpha^{\langle k\rangle}, \ldots, \alpha^{\langle k\rangle n}$ be independent collections $\alpha^{\langle k\rangle i}=\left(\alpha_{z}^{\langle k\rangle i}\right)_{z \in \mathbb{Z}_{\text {even }}^{2}}$ of $\{-1,+1\}$-valued random variables with $\mathbb{P}\left[\alpha_{z}^{\langle k\rangle i}=+1 \mid \omega^{\langle k\rangle}\right]=\omega_{z}^{\langle k\rangle}$, and let $\mathcal{U}_{\langle k\rangle}^{i}:=\mathcal{U}^{\alpha^{\langle k\rangle i}}$ be the discrete web associated with $\alpha^{\langle k\rangle i}$ as defined in (3.2). Then the averaged law

$$
\begin{equation*}
\mathbb{P}\left[S_{\varepsilon_{k}}\left(\mathcal{U}_{\langle k\rangle}^{1}, \ldots, \mathcal{U}_{\langle k\rangle}^{n}\right) \in \cdot\right] \tag{7.2}
\end{equation*}
$$

is the $n$-th moment measure of $S_{\varepsilon_{k}}\left(\mathbf{Q}_{\langle k\rangle}\right)$. We need to prove weak convergence of the laws in (7.2) to those in (7.1).

Our strategy will be to embed the Brownian webs $\mathcal{W}_{1}, \ldots, \mathcal{W}_{n}$ in a Brownian net $\mathcal{N}$, and similarly for the rescaled discrete webs. We will then prove weak convergence in law for the discrete net and webs to $\left(\mathcal{N}, \mathcal{W}_{1}, \ldots, \mathcal{W}_{n}\right)$ much in the same way as we have proved Theorem 6.15

We start by recalling how the sample Brownian webs $\mathcal{W}_{1}, \mathcal{W}_{2}, \ldots$ are constructed in terms of the marked reference Brownian web $\left(\mathcal{W}_{0}, \mathcal{M}\right)$. The basic ingredients of the construction are the drift $\beta_{0}$ of the reference web $\mathcal{W}_{0}$ and finite measures $\nu_{1}, \nu_{\mathrm{r}}$ on $[0,1]$. Given $\mathcal{W}_{0}$, the set of marked points $\mathcal{M}=\left\{\left(z, \omega_{z}\right): z \in M\right\}$ is then a Poisson point set with intensity as in (3.16). To construct $\mathcal{W}_{1}, \mathcal{W}_{2}, \ldots$, conditional on $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, independently for $i=1,2, \ldots$, we let $\left(\alpha_{z}^{i}\right)_{z \in M}$ be a collection of independent $\{-1,+1\}$-valued random variables with $\mathbb{P}\left[\alpha_{z}^{i}=\right.$ $\left.+1 \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right]=\omega_{z}$, we set $A_{i}:=\left\{z \in M: \alpha_{z}^{i} \neq \operatorname{sign}_{\mathcal{W}_{0}}(z)\right\}$, we let $B_{i}$ be a Poisson point set with intensity $2 \nu_{1}(\{0\}) \ell_{1}+2 \nu_{\mathrm{r}}(\{1\}) \ell_{\mathrm{r}}$, independent of $A_{i}$, and as in 3.19), we set

$$
\begin{equation*}
\mathcal{W}_{i}:=\lim _{\Delta_{n} \uparrow A_{i} \cup B_{i}} \operatorname{switch}_{\Delta_{n}}\left(\mathcal{W}_{0}\right) \quad(i=1,2, \ldots) \tag{7.3}
\end{equation*}
$$

Then the $\mathcal{W}_{1}, \mathcal{W}_{2}, \ldots$ are conditionally i.i.d. given $\left(\mathcal{W}_{0}, \mathcal{M}\right)$ and $\mathbb{P}\left[\mathcal{W}_{i} \in \cdot \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right]$ is the Howitt-Warren quenched law with drift $\beta$ and characteristic measure $\nu$ given by (3.20) and (3.21).

We wish to show that for each $n \geq 1$, the Brownian webs $\mathcal{W}_{1}, \ldots, \mathcal{W}_{n}$ from (7.3) are in a natural way embedded in a Brownian net. To that aim, for any set of paths $\mathcal{A} \subset \Pi$ and set of times $\mathcal{T} \subset \mathbb{R}$, we let $\mathcal{H}_{\mathcal{T}}(\mathcal{A})$ denote the set of paths that can be obtained from $\mathcal{A}$ by hopping finitely often at times in $\mathcal{T}$, i.e., $\mathcal{H}_{\mathcal{T}}(\mathcal{A})$ contains all paths of the form

$$
\begin{align*}
& \pi=\bigcup_{i=1}^{n}\left\{\left(\pi_{i}(t), t\right): t_{i-1} \leq t \leq t_{i}\right\} \quad \text { where } \quad \pi_{1}, \ldots, \pi_{n} \in \mathcal{A}, t_{1}, \ldots, t_{n-1} \in \mathcal{T} \\
& \quad t_{0}<\cdots<t_{n}=\infty, \sigma_{\pi_{1}}=t_{0}, \sigma_{\pi_{i+1}} \leq t_{i}, \pi_{i+1}\left(t_{i}\right)=\pi_{i}\left(t_{i}\right) \quad(1 \leq i \leq n-1) \tag{7.4}
\end{align*}
$$

where as usual we identify a path $\pi$ with its graph $\left\{(\pi(t), t): t \geq \sigma_{\pi}\right\}$. Moreover, we set

$$
\begin{equation*}
\Gamma_{n}:=\{-1,+1\}^{n} \backslash\{(-1, \ldots,-1),(+1, \ldots,+1)\} \tag{7.5}
\end{equation*}
$$

and for each $n \geq 2$, we define a probability measure $\Lambda_{n}$ on $\Gamma_{n}$ by

$$
\begin{equation*}
\Lambda_{n}(\vec{\gamma}):=\frac{1}{Z} \int q^{k-1}(1-q)^{n-k-1} \nu(\mathrm{~d} q), \quad \text { where } \quad k:=\left|\left\{i: \gamma^{i}=+1\right\}\right| \tag{7.6}
\end{equation*}
$$

where $Z$ is the normalization constant given by

$$
\begin{equation*}
Z:=\int \frac{1-q^{n}-(1-q)^{n}}{q(1-q)} \nu(\mathrm{d} q) \tag{7.7}
\end{equation*}
$$

with the convention that the integrand in (7.7) takes on the value $n$ at the points $q=0,1$.
Below, if $\mathcal{A}$ is a set of paths, then $\overline{\mathcal{A}}$ denotes the closure of $\mathcal{A}$ in the topology on the path space $\Pi$. We note that in 7.8 , if $\mathcal{D}$ is moreover dense in $\mathbb{R}^{2}$, then $\overline{\mathcal{N}_{n}(\mathcal{D})}=\mathcal{N}_{n}$. (This follows, for example, from [SS08, Thm. 1.3].)

Lemma 7.1 (Construction of moment measures) Conditional on the marked reference Brownian web $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, let $\mathcal{W}_{1}, \mathcal{W}_{2}, \ldots$ be an i.i.d. sequence of sample Brownian webs defined as in (7.3). Then, for each $n \geq 1$, there exists an a.s. unique Brownian net $\mathcal{N}_{n}$ with left and right speeds $\beta_{-}(n), \beta_{+}(n)$ given by (2.3) such that for any deterministic countable set $\mathcal{D} \subset \mathbb{R}^{2}$ and countable dense set of times $\mathcal{T} \subset \mathbb{R}$,

$$
\begin{equation*}
\overline{\mathcal{N}_{n}(\mathcal{D})}=\overline{\mathcal{H}_{\mathcal{T}}\left(\mathcal{W}_{1} \cup \cdots \cup \mathcal{W}_{n}\right)(\mathcal{D})} \quad \text { a.s. } \tag{7.8}
\end{equation*}
$$

Let $S_{n}$ be the set of separation points of $\mathcal{N}_{n}$. Then each $z \in S_{n}$ is of type $(1,2)$ in $\mathcal{W}_{1}, \ldots, \mathcal{W}_{n}$ and conditional on $\mathcal{N}_{n}$, the random variables $\left(\vec{\alpha}_{z}\right)_{z \in S_{n}}$ defined by

$$
\begin{equation*}
\vec{\alpha}_{z}:=\left(\operatorname{sign}_{\mathcal{W}_{1}}(z), \ldots, \operatorname{sign}_{\mathcal{W}_{n}}(z)\right) \tag{7.9}
\end{equation*}
$$

are i.i.d. with law $\Lambda_{n}$ defined in (7.6). Moreover, one has

$$
\begin{equation*}
\mathcal{W}_{i}=\left\{\pi \in \mathcal{N}_{n}: \operatorname{sign}_{\pi}(z)=\alpha_{z}^{i} \forall z \in S_{n} \text { s.t. } \pi \text { enters } z\right\} \quad(i=1, \ldots, n) \tag{7.10}
\end{equation*}
$$

Proof. Set $C:=\bigcup_{i=1}^{n}\left(A_{i} \cup B_{i}\right)$ and let $C_{\mathrm{l}}$ and $C_{\mathrm{r}}$ denote the restrictions of $C$ to the sets of points of type $(1,2)_{1}$ and $(1,2)_{\mathrm{r}}$ in $\mathcal{W}_{0}$, respectively. For each $\vec{\gamma} \in\{-1,+1\}^{n}$, set

$$
\begin{equation*}
C(\vec{\gamma}):=\left\{z \in C: \operatorname{sign}_{\mathcal{W}_{i}}(z)=\gamma_{i} \forall i=1, \ldots, n\right\} \tag{7.11}
\end{equation*}
$$

and define $C_{1}(\vec{\gamma})$ and $C_{\mathrm{r}}(\vec{\gamma})$ similarly, with $C$ replaced by $C_{1}$ resp. $C_{\mathrm{r}}$. By our definition of the sample Brownian webs, conditional on $\mathcal{W}_{0}$, the sets $\left\{C(\vec{\gamma}): \vec{\gamma} \in\{-1,+1\}^{n}\right\}$ are independent Poisson point sets with intensity $c_{1}(\vec{\gamma}) \ell_{1}+c_{\mathrm{r}}(\vec{\gamma}) \ell_{\mathrm{r}}$, where (compare (3.18))

$$
\left.\begin{array}{rl}
c_{1}(\vec{\gamma}) & =21_{\{0<k\}} \int_{(0,1]} q^{k}(1-q)^{n-k} q^{-1} \nu_{1}(\mathrm{~d} q)+21_{\{k=1\}} \nu_{1}(\{0\}) \\
& =21_{\{0<k\}} \int q^{k-1}(1-q)^{n-k} \nu_{1}(\mathrm{~d} q),  \tag{7.12}\\
c_{\mathrm{r}}(\vec{\gamma}) & =21_{\{k<n\}} \int q^{k}(1-q)^{n-k-1} \nu_{\mathrm{r}}(\mathrm{~d} q)
\end{array}\right\} \quad \text { where } k:=\left|\left\{i: \gamma_{i}=+1\right\}\right| .
$$

We modify our reference web by setting

$$
\begin{equation*}
\mathcal{W}_{0}^{\prime}:=\lim _{\Delta_{m} \uparrow C^{*}} \operatorname{switch}_{\Delta_{m}}\left(\mathcal{W}_{0}\right) \quad \text { where } \quad C^{*}:=C_{1}(+1, \ldots,+1) \cup C_{\mathrm{r}}(-1, \ldots,-1) \tag{7.13}
\end{equation*}
$$

By Proposition 3.6 (i), $\ell_{1}$ and $\ell_{\mathrm{r}}$ are also the intersection local time measures for the modified reference web $\mathcal{W}_{0}^{\prime}$. Since $\mathcal{W}_{0}^{\prime}$ is a.s. uniquely determined by $\mathcal{W}_{0}$ and the set $C^{*}$, and since conditional on $\mathcal{W}_{0}$ and $C^{*}$, the sets $C(\vec{\gamma})$ with $\gamma \neq(-1, \ldots,-1),(+1, \ldots,+1)$ are independent Poisson point sets with intensity $c_{1}(\vec{\gamma}) \ell_{1}+c_{\mathrm{r}}(\vec{\gamma}) \ell_{\mathrm{r}}$, by Theorem 4.6, we can define a Brownian net $\mathcal{N}_{n}$ with set of separation points $S_{n}$ by (recall (7.5))

$$
\begin{equation*}
\mathcal{N}_{n}:=\lim _{\Delta_{m} \uparrow S_{n}} \operatorname{hop}_{\Delta_{m}}\left(\mathcal{W}_{0}^{\prime}\right) \quad \text { where } \quad S_{n}:=\bigcup_{\vec{\gamma} \in \Gamma_{n}} C(\vec{\gamma}) . \tag{7.14}
\end{equation*}
$$

By Theorem 4.6, conditional on $\mathcal{N}_{n}$, the random variables $\left(\operatorname{sign}_{\mathcal{W}_{0}^{\prime}}(z)\right)_{z \in S_{n}}$ are i.i.d. with

$$
\begin{equation*}
\mathbb{P}\left[\operatorname{sign}_{\mathcal{W}_{0}^{\prime}}(z)=+1 \mid \mathcal{N}_{n}\right]=\frac{\sum_{\vec{\gamma} \in \Gamma_{n}} c_{\mathrm{r}}(\vec{\gamma})}{\sum_{\vec{\gamma} \in \Gamma_{n}}\left(c_{1}(\vec{\gamma})+c_{\mathrm{r}}(\vec{\gamma})\right)} . \tag{7.15}
\end{equation*}
$$

Using this and the independence of the Poisson point sets $\left\{C(\vec{\gamma}): \vec{\gamma} \in \Gamma_{n}\right\}$, it is straightforward to check from (3.21) and (7.12) that conditional on $\mathcal{N}_{n}$, the random variables in (7.9) are i.i.d. with law $\Lambda_{n}$ defined in (7.6). By Lemma 6.17.

$$
\begin{equation*}
\mathcal{W}_{i}=\lim _{\Delta_{m} \uparrow C^{i}} \operatorname{switch}_{\Delta_{m}}\left(\mathcal{W}_{0}^{\prime}\right) \quad \text { where } \quad C^{i}:=\bigcup_{\substack{\vec{\gamma} \in \Gamma_{n} \\ \gamma_{i}=+1}} C_{\mathbf{l}}(\vec{\gamma}) \cup \bigcup_{\substack{\vec{\gamma} \in \Gamma_{n} \\ \gamma_{i}=-1}} C_{\mathrm{r}}(\vec{\gamma}) . \tag{7.16}
\end{equation*}
$$

Therefore, by Proposition 4.5, we see that (7.10) holds.
The speed of $\mathcal{W}_{0}^{\prime}$ is given by

$$
\begin{equation*}
\beta_{0}+c_{1}(+1, \ldots,+1)-c_{\mathrm{r}}(-1, \ldots,-1) \tag{7.17}
\end{equation*}
$$

and therefore the left speed of $\mathcal{N}_{n}$ is given by

$$
\begin{align*}
\beta_{0} & +c_{1}(+1, \ldots,+1)-c_{\mathrm{r}}(-1, \ldots,-1)-\sum_{\vec{\gamma} \in \Gamma_{n}} c_{\mathrm{r}}(\vec{\gamma}) \\
& =\beta-2 \nu_{1}([0,1])+2 \nu_{\mathrm{r}}([0,1])+c_{1}(+1, \ldots,+1)-\sum_{\vec{\gamma} \neq(+1, \ldots,+1)} c_{\mathrm{r}}(\vec{\gamma}) \\
& =\beta-2 \int \nu_{1}(\mathrm{~d} q)\left(1-q^{n-1}\right)+2 \int \nu_{\mathrm{r}}(\mathrm{~d} q)\left(1-\sum_{k=0}^{n-1}\binom{n}{k} q^{k}(1-q)^{n-k-1}\right) \\
& =\beta-2 \int(1-q) \nu_{1}(\mathrm{~d} q) \sum_{k=0}^{n-2} q^{k}-2 \int q \nu_{\mathrm{r}}(\mathrm{~d} q) \sum_{k=0}^{n-2} q^{k}=\beta-2 \int \nu(\mathrm{~d} q) \sum_{k=0}^{n-2} q^{k}=\beta_{-}(n), \tag{7.18}
\end{align*}
$$

where we have used $(3.20,(3.21), 2.3)$ and the fact that

$$
\begin{equation*}
\sum_{k=0}^{n-1}\binom{n}{k} q^{k}(1-q)^{n-k-1}=(1-q)^{-1}\left(1-q^{n}\right)=\sum_{k=0}^{n-1} q^{k} \tag{7.19}
\end{equation*}
$$

which is true even for $q=1$, even though the intermediate step is not defined in this case. The calculation for $\beta_{+}(n)$ is completely analogous.

We are left with the task to prove 7.8 . The inclusion $\mathcal{N}_{n}(\mathcal{D}) \supset \mathcal{H}_{\mathcal{T}}\left(\mathcal{W}_{1} \cup \cdots \cup \mathcal{W}_{n}\right)(\mathcal{D})$ follows from the fact that $\mathcal{N}_{n}$ is closed under hopping at deterministic times, see SS08, Lemma 8.3]. To prove the converse inclusion, by the compactness of $\mathcal{N}_{n}$, it suffices to prove that for each $t_{1}<\cdots<t_{m}$ with $t_{1}, \ldots, t_{m} \in \mathcal{T}, z=\left(x, t_{0}\right) \in \mathcal{D}$ with $t_{0}<t_{1}$, and $\pi \in \mathcal{N}_{n}(z)$, we can find $\pi^{\prime} \in \mathcal{H}_{\mathcal{T}}\left(\mathcal{W}_{1} \cup \cdots \cup \mathcal{W}_{n}\right)$ starting from $z$ such that $\pi\left(t_{i}\right)=\pi^{\prime}\left(t_{i}\right)$ for $i=1, \ldots, m$. By the finite graph representation (in particular, by Corollary 6.6) and the fact that for each separation point $z$ of $\mathcal{N}_{n}$, there exists $1 \leq i, j \leq n$ such that $\operatorname{sign}_{\mathcal{W}_{i}}(z)=-1$ and $\operatorname{sign}_{\mathcal{W}_{j}}(z)=+1$, we can find a $\pi^{\prime \prime}$ starting at $z$ and satisfying $\pi\left(t_{i}\right)=\pi^{\prime \prime}\left(t_{i}\right)$ for $i=1, \ldots, m$ that is obtained by concatenating finitey many paths in $\mathcal{W}_{1}, \ldots, \mathcal{W}_{n}$ at separation points of $\mathcal{N}_{n}$. By the fact that $\mathcal{T}$ is dense and the structure of separation points (see Proposition 4.3), we can modify $\pi^{\prime \prime}$ a bit such that the concatenation takes place at times in $\mathcal{T}$.

Proof of Theorem 3.9 Let $\mathcal{U}_{\langle k\rangle}^{i}(i=1, \ldots, n)$ be the discrete webs in $\sqrt[7.2]{ }$, and let $\mathcal{V}_{\langle k\rangle}$ be the discrete net defined by

$$
\begin{equation*}
\mathcal{V}_{\langle k\rangle}:=\left\{\pi: \pi(t+1)-\pi(t) \in\left\{\alpha_{(p(t), t)}^{\langle k\rangle 1}, \ldots, \alpha_{(p(t), t)}^{\langle k\rangle}\right\} \forall t \geq \sigma_{\pi}\right\} \tag{7.20}
\end{equation*}
$$

By Theorem 6.11, $\mathcal{V}_{\langle k\rangle}$, diffusively rescaled, converges to a Brownian net with left and right speeds given by

$$
\begin{align*}
\beta_{-}(n) & =\lim _{k \rightarrow \infty} \varepsilon_{k}^{-1} \mathbb{E}\left[\alpha_{z}^{\langle k\rangle} \wedge \cdots \wedge \alpha_{z}^{\langle k\rangle}\right]=\lim _{k \rightarrow \infty} \varepsilon_{k}^{-1} \int_{n-2} \mu_{k}(\mathrm{~d} q)\left(q^{n}-\left(1-q^{n}\right)\right) \\
& =\lim _{k \rightarrow \infty} \varepsilon_{k}^{-1} \int \mu_{k}(\mathrm{~d} q)\left((2 q-1)-2 q(1-q) \sum_{k=0}^{n-2} q^{k}\right)=\beta-2 \int \nu(\mathrm{~d} q) \sum_{k=0}^{n-2} q^{k},  \tag{7.21}\\
\beta_{+}(n) & =\lim _{k \rightarrow \infty} \varepsilon_{k}^{-1} \mathbb{E}\left[\alpha_{z}^{\langle k\rangle 1} \vee \cdots \vee \alpha_{z}^{\langle k\rangle n}\right]=\beta+2 \int \nu(\mathrm{~d} q) \sum_{k=0}^{m-2}(1-q)^{k},
\end{align*}
$$

where we have used 1.7 . Let $S_{\langle k\rangle}$ be the set of separation points of $\mathcal{V}_{\langle k\rangle}$ and set $\vec{\alpha}_{z}^{\langle k\rangle}:=$ $\left(\alpha_{z}^{\langle k\rangle 1}, \ldots, \alpha_{z}^{\langle k\rangle 1}\right)$. Then, conditional on $\mathcal{V}_{\langle k\rangle}$, the random variables $\left(\vec{\alpha}_{z}^{\langle k\rangle}\right)_{z \in S_{\langle k\rangle}}$ are i.i.d. with

$$
\begin{equation*}
\mathbb{P}\left[\vec{\alpha}_{z}^{\langle k\rangle}=\vec{\gamma} \mid \mathcal{V}_{\langle k\rangle}\right]=\frac{1}{Z_{k}} \int \mu_{k}(\mathrm{~d} q) q^{l}(1-q)^{n-l}, \quad \text { where } \quad l:=\left|\left\{i: \gamma^{i}=+1\right\}\right| \tag{7.22}
\end{equation*}
$$

and $Z_{k}$ is a normalization constant. Using (1.7), it is easy to check that this conditional law converges as $k \rightarrow \infty$ to the law in (7.6).

For each $i=1, \ldots, n$, the pairs $\left(\mathcal{V}_{\langle k\rangle}, \mathcal{U}_{\langle k\rangle}^{i}\right)$ are distributed as the discrete nets and webs in Theorem 6.15, so by that theorem, Lemma 6.7 (a) and the remarks below it, we can couple our random variables in such a way that

$$
\begin{equation*}
S_{\varepsilon_{k}}\left(\mathcal{V}_{\langle k\rangle}, \mathcal{U}_{\langle k\rangle}^{1}, \ldots, \mathcal{U}_{\langle k\rangle}^{n}\right) \underset{n \rightarrow \infty}{\longrightarrow}\left(\mathcal{N}, \mathcal{W}_{1}, \ldots, \mathcal{W}_{n}\right) \tag{7.23}
\end{equation*}
$$

where $\mathcal{N}$ is a Brownian net with left and right speeds $\beta_{-}(n), \beta_{+}(n), \mathcal{W}_{1}, \ldots, \mathcal{W}_{n}$ are Brownian webs with drift $\beta$, such that each separation point of $\mathcal{N}$ is of type $(1,2)$ in each $\mathcal{W}_{i}$ and

$$
\begin{equation*}
\mathcal{W}_{i}:=\left\{\pi \in \mathcal{N}: \operatorname{sign}_{\pi}(z)=\operatorname{sign}_{\mathcal{W}_{i}}(z) \forall z \in S \text { s.t. } \pi \text { enters } z\right\} \quad(i=1, \ldots, n), \tag{7.24}
\end{equation*}
$$

where $S$ is the set of separation points of $\mathcal{N}$. Much in the same way as in the proof of property (ii) of Theorem 6.15, we find that conditional on $\mathcal{N}$, the random variables

$$
\begin{equation*}
\left(\operatorname{sign}_{\mathcal{W}_{1}}(z), \ldots, \operatorname{sign}_{\mathcal{W}_{n}}(z)\right)_{z \in S} \tag{7.25}
\end{equation*}
$$

are i.i.d. with common law as in 7.6). By Lemma 7.1, this proves the convergence of the moment measures in (7.2) to those in (7.1) and hence, by [Daw91, Thm. 3.2.9], the convergence in (3.27).

### 7.2 Proof of the marking constructions of Howitt-Warren flows

In this section, we prove our main results, Theorems 3.7 and 4.7 on the construction of Howitt-Warren flows inside a Brownian web and net. It turns out that we already have most ingredients of the proofs. The main point that still needs to be settled is to verify that our construction agrees with the original definition of Howitt-Warren flows based on $n$-point motions and the Howitt-Warren martingale problem.

Proposition 7.2 (Identification of $n$-point motions) Let $\beta_{0} \in \mathbb{R}$ and let $\nu_{1}, \nu_{\mathrm{r}}$ be finite measures on $[0,1]$. Let $\left(\mathcal{W}_{0}, \mathcal{M}\right)$ be a marked reference web as in Theorem 3.7 and conditional on $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, let $\mathcal{W}_{1}, \ldots, \mathcal{W}_{n}$ be $n$ independent sample webs constructed as in (3.19). For each deterministic $z \in \mathbb{R}^{2}$, let $\pi_{z}^{i}$ denote the a.s. unique element of $\mathcal{W}_{i}(z)$. Then, for each $\vec{x} \in \mathbb{R}^{n}$ and $s \in \mathbb{R}$, the process

$$
\begin{equation*}
\left(\pi_{\left(x_{1}, s\right)}^{1}(s+t), \ldots, \pi_{\left(x_{n}, s\right)}^{n}(s+t)\right)_{t \geq 0} \tag{7.26}
\end{equation*}
$$

solves the Howitt-Warren martingale problem with drift $\beta$ and characteristic measure $\nu$ given by (3.20) and (3.21).

Proof. Instead of attempting a direct proof we will use discrete approximation. It is easy to verify that Theorem 3.9 implies the convergence of the $n$-point motions of diffusively rescaled discrete Howitt-Warren flows to the $n$-point motions of the quenched law $\mathbb{Q}$, while by Proposition A.5, the same discrete $n$-point motions converge to a solution of the Howitt-Warren martingale problem. The proposition then follows.
Proof of Theorem 3.7 We start by checking that the random kernels $K_{s, t}^{+}$defined in 3.22 form a stochastic flow of kernels on $\mathbb{R}$ as in Definition 2.1. Indeed, Property (i) follows from the fact that

$$
\begin{align*}
& \int_{\mathbb{R}} K_{s, t}^{+}(x, \mathrm{~d} y) K_{t, u}^{+}(y, \mathrm{~d} z)=\int_{\mathbb{R}} \mathbb{P}\left[\pi_{(x, s)}^{+}(t) \in \mathrm{d} y \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right] \mathbb{P}\left[\pi_{(y, t)}^{+}(u) \in \mathrm{d} z \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right] \\
& \quad=\int_{\mathbb{R}} \mathbb{P}\left[\pi_{(x, s)}^{+}(t) \in \mathrm{d} y \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right] \mathbb{P}\left[\pi_{(y, t)}^{+}(u) \in \mathrm{d} z \mid\left(\mathcal{W}_{0}, \mathcal{M}\right), \pi_{(x, s)}^{+}(t)=y\right]  \tag{7.27}\\
& \quad=\int_{\mathbb{R}} \mathbb{P}\left[\pi_{\left(\pi_{(x, s)}^{+}(t), t\right)}^{+}(u) \in \mathrm{d} z \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right]=\int_{\mathbb{R}} \mathbb{P}\left[\pi_{(x, s)}^{+}(u) \in \mathrm{d} z \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right] \text { a.s., }
\end{align*}
$$

where we have used that $\pi_{(x, s)}^{+}(t)$ and $\pi_{(y, t)}^{+}(u)$ are conditionally independent given $\left(\mathcal{W}_{0}, \mathcal{M}\right)$ and $\pi_{\left(\pi_{(x, s)}^{+}(t), t\right)}^{+}(u)=\pi_{(x, s)}^{+}(u)$ a.s., which follows from the fact that for deterministic $t$, a.s.
every point in $\mathbb{R} \times\{t\}$ is of type $(0,1),(0,2)$ or $(1,1)$ (see Proposition 3.3). Property (ii) of Definition 2.1 follows from the fact that the restrictions of $\left(\mathcal{W}, \mathcal{M}, \mathcal{W}_{0}\right)$ to disjoint time intervals are independent, which follows from the analogue property for a single Brownian web which is proved by discrete approximation. Property (iii), finally, is obvious from the translation invariance of our definitions. Since $K_{s, t}^{\uparrow}(x, \cdot)=K_{s, t}^{+}(x, \cdot)$ a.s. for deterministic $s \leq t$ and $x \in \mathbb{R}$, the same conclusions can be drawn for $K_{s, t}^{\uparrow}$.

To identify $\left(K_{s, t}^{+}\right)_{s \leq t}$ (and likewise $\left(K_{s, t}^{\uparrow}\right)_{s \leq t}$ ) as a Howitt-Warren flow with drift $\beta$ and characteristic measure $\nu$, therefore, it suffices to check that for each deterministic $\vec{x} \in \mathbb{R}^{n}$ and $s \leq t$, one has (compare (2.2))

$$
\begin{equation*}
\mathbb{E}\left[K_{s, t}^{+}\left(x_{1}, \cdot\right) \cdots K_{s, t}^{+}\left(x_{n}, \cdot\right)\right]=\mathbb{P}\left[\vec{X}_{t-s}^{\vec{x}} \in \cdot\right], \tag{7.28}
\end{equation*}
$$

where $\vec{X} \vec{x}$ is a solution the Howitt-Warren martingale problem with drift $\beta$ and characteristic measure $\nu$, started in $\vec{X}_{0}^{\vec{x}}=\vec{x}$. Since

$$
\begin{equation*}
\mathbb{E}\left[K_{s, t}^{+}\left(x_{1}, A_{1}\right) \cdots K_{s, t}^{+}\left(x_{n}, A_{n}\right)\right]=\mathbb{P}\left[\pi_{\left(x_{1}, s\right)}^{1}(t) \in A_{1}, \ldots, \pi_{\left(x_{n}, s\right)}^{n}(t) \in A_{n}\right] \tag{7.29}
\end{equation*}
$$

where $\pi_{\left(x_{1}, s\right)}^{1}, \ldots, \pi_{\left(x_{1}, s\right)}^{n}$ are as in Proposition 7.2 , our claim follows from that result. The fact that $\left(\mathcal{W}_{0}, \mathcal{M}, \mathcal{W}\right)$ and $\left(\mathcal{W}, \mathcal{M}, \mathcal{W}_{0}\right)$ are equally distributed if $\nu_{1}=\nu_{\mathrm{r}}$ follows from the somewhat stronger Proposition 7.3 below.

Proposition 7.3 (Exchangeability of reference web) Let $\left(\mathcal{W}_{0}, \mathcal{M}\right)$ be a marked reference web as in Theorem 3.7 and conditional on $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, let $\left(\mathcal{W}_{1}, \mathcal{W}_{2}, \ldots\right)$ be independent sample webs constructed as in (3.19). Assume that $\nu_{1}=\nu_{\mathrm{r}}$. Then the sequence of Brownian webs $\left(\mathcal{W}_{0}, \mathcal{W}_{1}, \mathcal{W}_{2}, \ldots\right)$ is exchangeable.

Proof. In the set-up of Lemma 7.1, we will show that if $\nu_{1}=\nu_{\mathrm{r}}$, then the joint law of $\left(\mathcal{W}_{0}, \ldots, \mathcal{W}_{n}\right)$ is equal to the law of $\left(\mathcal{W}_{1}, \ldots, \mathcal{W}_{n+1}\right)$, which is clearly exchangeable. To see this, let $C$ be defined as in the proof of Lemma 7.1 and in analogy with (7.14), set

$$
\begin{equation*}
\mathcal{N}_{n}^{\prime}:=\lim _{\Delta_{m} \uparrow C} \operatorname{hop}_{\Delta_{m}}\left(\mathcal{W}_{0}\right) \tag{7.30}
\end{equation*}
$$

Then $\mathcal{N}_{n}^{\prime}$ is a Brownian net with set of separation points $C$. For $i=0, \ldots, n$ and $z \in C$, set $\alpha_{z}^{i}:=\operatorname{sign}_{\mathcal{V}_{i}}(z)$ and as in (7.9) let $\vec{\alpha}_{z}:=\left(\alpha_{z}^{1}, \ldots, \alpha_{z}^{n}\right)$. In a similar way as in the proof of Lemma 7.1, we check that conditional on $\mathcal{N}_{n}^{\prime}$, the random variables $\left(\alpha_{z}^{0}, \vec{\alpha}_{z}\right)_{z \in C}$ are i.i.d. with

$$
\begin{aligned}
& \mathbb{P}\left[\left(\alpha_{z}^{0}, \vec{\alpha}_{z}\right)=\left(\gamma_{0}, \vec{\gamma}\right) \mid \mathcal{N}_{n}^{\prime}\right]=\frac{c_{1}(\vec{\gamma})}{\sum_{\vec{\gamma} \neq(-1, \ldots,-1)} c_{1}(\vec{\gamma})+\sum_{\vec{\gamma} \neq(+1, \ldots,+1)} c_{\mathrm{r}}(\vec{\gamma})} \quad \text { if } \gamma_{0}=-1, \\
& \mathbb{P}\left[\left(\alpha_{z}^{0}, \vec{\alpha}_{z}\right)=\left(\gamma_{0}, \vec{\gamma}\right) \mid \mathcal{N}_{n}^{\prime}\right]=\frac{c_{\mathrm{r}}(\vec{\gamma})}{\sum_{\vec{\gamma} \neq(-1, \ldots,-1)} c_{1}(\vec{\gamma})+\sum_{\vec{\gamma} \neq(+1, \ldots,+1)} c_{\mathrm{r}}(\vec{\gamma})} \quad \text { if } \gamma_{0}=+1,
\end{aligned}
$$

where $c_{1}(\vec{\gamma}), c_{\mathrm{r}}(\vec{\gamma})$ are defined in (7.12). In particular, if $\nu_{1}=\nu_{\mathrm{r}}$, then

$$
\begin{equation*}
c_{1}(\vec{\gamma})=c(-1, \vec{\gamma}) \quad \text { and } \quad c_{\mathrm{r}}(\vec{\gamma})=c(+1, \vec{\gamma}), \tag{7.31}
\end{equation*}
$$

where we define

$$
c\left(\gamma_{0}, \vec{\gamma}\right):=21_{\{0<k<n+1\}} \int q^{k-1}(1-q)^{(n+1)-k-1} \nu(\mathrm{~d} q) \quad \text { with } \quad k:=\left|\left\{i: 0 \leq i \leq n, \gamma_{i}=+1\right\}\right| .
$$

From this, it is easy to check that $\mathcal{N}_{n}^{\prime}$ has left and right speeds $\beta_{-}(n+1), \beta_{+}(n+1)$. By Lemma 7.1, since $\mathcal{W}_{0}, \ldots, \mathcal{W}_{n}$ can be constructed inside $\mathcal{N}_{n}^{\prime}$ as in 7.10), it follows that $\left(\mathcal{W}_{0}, \ldots, \mathcal{W}_{n}\right)$ is equally distributed with $\left(\mathcal{W}_{1}, \ldots, \mathcal{W}_{n+1}\right)$.
The next lemma implies Theorem 4.7.
Lemma 7.4 (Limit of moment measures) Let $\beta_{0} \in \mathbb{R}$, let $\nu_{1}, \nu_{\mathrm{r}}$ be finite measures on $[0,1]$, let $\beta$ and $\nu$ be given by (3.20) and (3.21), and assume the left and right speeds $\beta_{-}, \beta_{+}$ defined in (2.12) satisfy $-\infty<\beta_{-}$, $\beta_{+}<\infty$. Let $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, with $\mathcal{M}=\left\{\left(z, \omega_{z}\right): z \in M\right\}$, be a marked reference web as in Theorem 3.7 and conditional on $(\mathcal{W}, \mathcal{M})$, let $\left(\mathcal{W}_{i}\right)_{i \geq 1}$ be independent sample webs constructed as in (3.19). Then there exists a Brownian net $\mathcal{N}_{\infty}$, which is determined a.s. uniquely by $\left(\mathcal{W}_{0}, \mathcal{M}\right)$ and has left and right speeds $\beta_{-}, \beta_{+}$, such that for any deterministic countable set $\mathcal{D} \subset \mathbb{R}^{2}$ and countable dense set of times $\mathcal{T} \subset \mathbb{R}$,

$$
\begin{equation*}
\overline{\mathcal{N}_{\infty}(\mathcal{D})}=\overline{\mathcal{H}_{\mathcal{T}}\left(\bigcup_{i \geq 1} \mathcal{W}_{i}\right)(\mathcal{D})} \text { a.s. } \tag{7.32}
\end{equation*}
$$

Let $S_{\infty}$ be the set of separation points of $\mathcal{N}_{\infty}$. Then $S_{\infty} \subset M$ and conditional on $\mathcal{N}_{\infty}$, the collection of random variables $\omega=\left(\omega_{z}\right)_{z \in S_{\infty}}$ is i.i.d. with law $\bar{\nu}$ defined in Theorem 4.7. A.s., each $z \in S_{\infty}$ is of type $(1,2)$ in each $\mathcal{W}_{i}(i \geq 1)$. Conditional on $\left(\mathcal{N}_{\infty}, \omega\right)$, the random variables $\left(\alpha_{z}^{i}\right)_{z \in S_{\infty}}^{i \geq 1}$ defined by $\alpha_{z}^{i}:=\operatorname{sign}_{\mathcal{W}_{i}}(z)$ are independent with $\mathbb{P}\left[\alpha_{z}^{i}=+1 \mid\left(\mathcal{N}_{\infty}, \omega\right)\right]=\omega_{z}$, and one has

$$
\begin{equation*}
\mathcal{W}_{i}=\left\{\pi \in \mathcal{N}_{\infty}: \operatorname{sign}_{\pi}(z)=\alpha_{z}^{i} \forall z \in S_{\infty} \text { s.t. } \pi \text { enters } z\right\} \quad(i \geq 1) \tag{7.33}
\end{equation*}
$$

Proof. This is very similar to the proof of Lemma 7.1 so we will only sketch the main line of proof. By the assumption that the speeds $\beta_{-}, \beta_{+}$are finite, conditional on $\mathcal{W}_{0}$, the set $M$ is a Poisson point set with intensity $c_{1} \ell_{1}+c_{\mathrm{r}} \ell_{\mathrm{r}}$, where $c_{1}, c_{\mathrm{r}}<\infty$ are given by

$$
\begin{equation*}
c_{1}:=2 \int q^{-1} \nu_{1}(\mathrm{~d} q) \quad \text { and } \quad c_{\mathrm{r}}:=2 \int(1-q)^{-1} \nu_{\mathrm{r}}(\mathrm{~d} q) . \tag{7.34}
\end{equation*}
$$

In analogy with (7.13), we set

$$
\begin{equation*}
\mathcal{W}_{0}^{\prime}:=\lim _{\Delta_{m} \uparrow C^{*}}\left(\mathcal{W}_{0}\right) \quad \text { with } \quad C^{*}:=\left\{z \in M_{1}: \omega_{z}=1\right\} \cup\left\{z \in M_{\mathrm{r}}: \omega_{z}=0\right\} \tag{7.35}
\end{equation*}
$$

where $M_{1}, M_{\mathrm{r}}$ denote the restrictions of $M$ to the sets of points of type $(1,2)_{1}$ and $(1,2)_{\mathrm{r}}$ in $\mathcal{W}_{0}$, respectively. Then, conditional on $\mathcal{W}_{0}^{\prime}$, the set $\left\{\left(z, \omega_{z}\right): z \in M \backslash C^{*}\right\}$ is a Poisson point set on $\mathbb{R}^{2} \times(0,1)$ with intensity

$$
\begin{equation*}
\ell_{1}(\mathrm{~d} z) \otimes 2 q^{-1} 1_{\{q<1\}} \nu_{1}(\mathrm{~d} q)+\ell_{\mathrm{r}}(\mathrm{~d} z) \otimes 2(1-q)^{-1} 1_{\{0<q\}} \nu_{\mathrm{r}}(\mathrm{~d} q) . \tag{7.36}
\end{equation*}
$$

Next, in analogy with 7.14, we define a Brownian net $\mathcal{N}_{\infty}$ with set of separation points $S_{\infty}:=M \backslash C^{*}$ by $\mathcal{N}_{\infty}:=\lim _{\Delta_{m} \uparrow S_{\infty}} \operatorname{hop}_{\Delta_{m}}\left(\mathcal{W}_{0}^{\prime}\right)$. Then Proposition 4.5 implies (7.33), while a calculation similar to (7.18) shows that the left and right speeds of $\mathcal{N}_{\infty}$ are the constants $\beta_{-}, \beta_{+}$from (2.12).

By Theorem 4.6 and (7.36), conditional on $\mathcal{N}_{\infty}$, the random variables $\left(\operatorname{sign}_{\mathcal{W}_{0}^{\prime}}(z)\right)_{z \in M \backslash C^{*}}$ are i.i.d., where $\mathbb{P}\left[\operatorname{sign}_{\mathcal{W}_{0}^{\prime}}(z)=+1 \mid \mathcal{N}_{\infty}\right]=c_{\mathrm{r}}^{\prime} /\left(c_{1}^{\prime}+c_{\mathrm{r}}^{\prime}\right)$ and $c_{1}^{\prime}:=2 \int q^{-1} 1_{\{q<1\}} \nu_{1}(\mathrm{~d} q), c_{\mathrm{r}}^{\prime}:=$ $2 \int(1-q)^{-1} 1_{\{0<q\}} \nu_{\mathrm{r}}(\mathrm{d} q)$. Using this, (3.21) and 7.36), we see that conditional on $\mathcal{N}_{\infty}$, the collection of random variables $\omega=\left(\omega_{z}\right)_{z \in M}$ is i.i.d. with common law $\bar{\nu}$ from Theorem 4.7. Since $\left(\mathcal{N}_{\infty}, \omega\right)$ is determined a.s. by $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, and since conditional on $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, the random
variables $\left(\alpha_{z}^{i}\right)_{z \in S_{\infty}}^{i \geq 1}$ are independent with $\mathbb{P}\left[\alpha_{z}^{i}=+1 \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right]=\omega_{z}$, the same statement holds for the conditional law given $\left(\mathcal{N}_{\infty}, \omega\right)$. The proof of $(\overline{7.32})$, finally, is completely analogous to the proof of formula (7.8) from Lemma 7.1.

Proof of Theorem 4.7. In Lemma 7.4. $\left(\mathcal{N}_{\infty}, \omega\right)$ is determined a.s. by $\left(\mathcal{W}_{0}, \mathcal{M}\right)$. Moreover, by (7.33), the conditional law $\mathbb{P}\left[\mathcal{W}_{i} \in \cdot \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right]$ is a function of $\left(\mathcal{N}_{\infty}, \omega\right)$ only. It follows that $\mathbb{P}\left[\mathcal{W}_{i} \in \cdot \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right]=\mathbb{P}\left[\mathcal{W}_{i} \in \cdot \mid\left(\mathcal{N}_{\infty}, \omega\right)\right]$.

### 7.3 Some immediate consequences of our construction

Proof of Proposition 3.8 (b)-(d). Let $\mathbb{Q}$ be the Howitt-Warren quenched law defined in (3.24). Then $\mathbb{Q}$ is a random probability law on the space $\mathcal{K}(\Pi)$ of compact subsets of the space of paths. We will be interested in a.s. properties of $\mathbb{Q}$ that hold for almost every realization of $\left(\mathcal{W}_{0}, \mathcal{M}\right)$. Let $\mathcal{W}$ be a random variable with law $\mathbb{Q}$. Since averaged over the law of $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, the set $\mathcal{W}$ is distributed as a Brownian web, it follows that for a.e. realization of $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, the compact set $\mathcal{W} \subset \Pi$ will satify all the a.s. properties of a Brownian web, such as the classification of special points. In particular, we can define the special paths $\pi_{z}^{+}$and $\pi_{z}^{\uparrow}$ for each $z \in \mathbb{R}^{2}$. Then

$$
\begin{equation*}
K_{s, t}^{+}(x, A):=\mathbb{Q}\left[\pi_{(x, s)}^{+}(t) \in A\right] \quad \text { and } \quad K_{s, t}^{\uparrow}(x, A):=\mathbb{Q}\left[\pi_{(x, s)}^{\uparrow}(t) \in A\right] \tag{7.37}
\end{equation*}
$$

are well-defined for every $s \leq t, x \in \mathbb{R}, A \in \mathcal{B}(\mathbb{R})$.
To prove part (b), it then suffices to note that for every $s \leq t_{n}, t_{n} \rightarrow t, x \in \mathbb{R}$ and continuous $f: \mathbb{R} \rightarrow \mathbb{R}$,

$$
\begin{equation*}
\int K_{s, t_{n}}^{+}(x, \mathrm{~d} y) f(y)=\mathbb{Q}\left[f\left(\pi_{(x, s)}^{+}\left(t_{n}\right)\right)\right] \underset{n \rightarrow \infty}{\longrightarrow} \mathbb{Q}\left[f\left(\pi_{(x, s)}^{+}(t)\right)\right]=\int K_{s, t}^{+}(x, \mathrm{~d} y) f(y) \tag{7.38}
\end{equation*}
$$

where we also use the symbol $\mathbb{Q}$ to denote expectation with respect to the probability law $\mathbb{Q}$, and we have used the continuity of $t \mapsto \pi_{(x, s)}^{+}$. The same proof works for $K_{s, t}^{\dagger}$.

The proof of part (c) is similar, where this time, for any $s<t, \mathbb{R} \ni x_{n} \downarrow x$ and $A \in \mathcal{B}(\mathbb{R})$,

$$
\begin{equation*}
K_{s, t}^{+}\left(x_{n}, A\right)=\mathbb{Q}\left[\pi_{\left(x_{n}, s\right)}^{+}(t) \in A\right] \rightarrow \mathbb{Q}\left[\pi_{(x, s)}^{+}(t) \in A\right]=K_{s, t}^{+}\left(x_{n}, A\right), \tag{7.39}
\end{equation*}
$$

where we have used that by [SS08, Lemma 3.4 (a)], under $\mathbb{Q}$, there is a random $m$ such that $\pi_{\left(x_{n}, s\right)}^{+}(t)=\pi_{(x, s)}^{+}(t)$ for all $n \geq m$. The existence of left limits follows in the same way.

To prove part (d), we observe that for all $s \leq t \leq u, x \in \mathbb{R}$ and $A \in \mathcal{B}(\mathbb{R})$,

$$
\begin{align*}
\mathbb{Q} & {\left[\pi_{(x, s)}^{\uparrow}(u) \in A\right]=\int \mathbb{Q}\left[\pi_{(x, s)}^{\uparrow}(u) \in A \mid \pi_{(x, s)}^{\uparrow}(t)=y\right] \mathbb{Q}\left[\pi_{(x, s)}^{\uparrow}(t) \in \mathrm{d} y\right] } \\
& =\int \mathbb{Q}\left[\pi_{(y, t)}^{\uparrow}(u) \in A \mid \pi_{(x, s)}^{\uparrow}(t)=y\right] \mathbb{Q}\left[\pi_{(x, s)}^{\uparrow}(t) \in \mathrm{d} y\right]  \tag{7.40}\\
& =\int \mathbb{Q}\left[\pi_{(y, t)}^{\uparrow}(u) \in A\right] \mathbb{Q}\left[\pi_{(x, s)}^{\uparrow}(t) \in \mathrm{d} y\right]
\end{align*}
$$

where we have conditioned on the value of $\pi_{(x, s)}^{\uparrow}(t)$, used the fact that $\pi_{(y, t)}^{\uparrow}$ is the continuation of any incoming path at $(y, t)$, and in the last step we have used that under the law $\mathbb{Q}$, for any $y \in \mathbb{R}$ that may depend on the marked reference web $\left(\mathcal{W}_{0}, \mathcal{M}\right)$ but not on the sample web $\mathcal{W}$, the path $\pi_{(y, t)}^{\uparrow}$ is independent of $\pi_{(x, s)}^{\uparrow}(t)$. To prove this independence, for any $t_{1} \leq t_{2}$,
let $\left.\mathcal{W}\right|_{t_{1}} ^{t_{2}}$ denote the restriction of $\mathcal{W}$ to the time interval $\left[t_{1}, t_{2}\right]$, i.e., $\left.\mathcal{W}\right|_{t_{1}} ^{t_{2}}:=\left\{\left.\pi\right|_{t_{1}} ^{t_{2}}: \pi \in \mathcal{W}\right\}$ where $\left.\pi\right|_{t_{1}} ^{t_{2}}:=\left\{(\pi(u), u): u \in\left[t_{1}, t_{2}\right] \cap\left[\sigma_{\pi}, \infty\right]\right\}$ is the restriction of a path $\pi$ to $\left[t_{1}, t_{2}\right]$. It follows from the marking construction in Theorem 3.7 that for all $t_{1} \leq t_{2} \leq t_{3},\left.\mathcal{W}\right|_{t_{1}} ^{t_{2}}$ and $\left.\mathcal{W}\right|_{t_{2}} ^{t_{3}}$ are independent under $\mathbb{Q}$. Since $\pi_{(y, t)}^{\uparrow}$ is a function of $\left.\mathcal{W}\right|_{t-\varepsilon} ^{\infty}$ for each $\varepsilon>0$, we conclude that $\pi_{(x, s)}^{\uparrow}(t-\varepsilon)$ is independent of $\pi_{(y, t)}^{\uparrow}$ under $\mathbb{Q}$ for each $\varepsilon>0$. Since $\pi_{(x, s)}^{\uparrow}(t)=\lim _{\varepsilon \rightarrow 0} \pi_{(x, s)}^{\uparrow}(t-\varepsilon)$, it follows that $\pi_{(x, s)}^{\uparrow}(t)$ is independent of $\pi_{(y, t)}^{\uparrow}$ under $\mathbb{Q}$.

The proof of Proposition 3.8 (a) needs a bit of preparation. We start by proving the statement for the Arratia flow.

Proposition 7.5 (Measurability of special paths) There exists a measurable function $\mathcal{K}(\Pi) \times \mathbb{R}^{2} \ni(\mathcal{A}, z) \mapsto \pi_{z}^{+}(\mathcal{A}) \in \Pi$ such that if $\mathcal{W}$ is a Brownian web, then almost surely, for all $z \in \mathbb{R}^{2}$ the path $\pi_{z}^{+}(\mathcal{W})=\pi_{z}^{+}$is the special path in $\mathcal{W}(z)$ defined below Proposition 3.3. An analogue statement holds for $\pi_{z}^{\uparrow}$.

Proof. It suffices to define $\pi_{z}^{+}(\mathcal{A})$ on the measurable set of all $\mathcal{A} \in \mathcal{K}(\Pi)$ such that $\mathcal{A}(z)$ contains a single path $\pi_{z}$ for each $z \in \mathbb{Q}^{2}$. For any $r \in \mathbb{R}$, define $\lfloor r\rfloor_{n} \uparrow r$ by $\lfloor r\rfloor_{n}:=\sup \left\{r^{\prime} \in\right.$ $\left.\mathbb{Z} / n: r^{\prime}<r\right\}$ and similarly, set $\lceil r\rceil_{n}:=\inf \left\{r^{\prime} \in \mathbb{Z} / n: r^{\prime}>r\right\}$. Then $(\mathcal{A},(x, t)) \mapsto \pi_{\left(\lceil x\rceil_{n},\lfloor t\rfloor_{m}\right)}$ is a measurable function. If $\mathcal{W}$ is a Brownian web, then applying Lemma 7.6 to the dual web $\hat{\mathcal{W}}$ we see that for each $t \in \mathbb{R}$, there exists at most one $x \in \mathbb{Q}$ such that $\mathcal{W}(x, t)$ contains more than one path. It follows that for each $(x, t) \in \mathbb{R}^{2}$, there is at most one $n$ for which the limit $\lim _{m \rightarrow \infty} \pi_{\left(\lceil x\rceil_{n},\lfloor t\rfloor_{m}\right)}$ does not exist, hence the double limit

$$
\begin{equation*}
\pi_{(x, t)}^{+}=\lim _{n \rightarrow \infty} \lim _{m \rightarrow \infty} \pi_{\left(\lceil x\rceil_{n},[t]_{m}\right)} \tag{7.41}
\end{equation*}
$$

is well-defined and gives the right-most path in $\mathcal{W}(x, t)$. Since pointwise limits of measurable functions are meaurable, restricting ourselves to a suitable measurable subset of $\mathcal{K}(\Pi)$, we see that $\pi_{z}^{+}$depends measurably jointly on $z$ and the Brownian web $\mathcal{W}$.

To also prove the statement for $\pi_{z}^{\uparrow}$, we note that the dual $\hat{\mathcal{W}}$ of a Brownian web $\mathcal{W}$ is a measurable function of $\mathcal{W}$ and that by what we have just proved, both the left-most and rightmost dual paths $\hat{\pi}_{z}^{-}$and $\hat{\pi}_{z}^{+}$depend measurably jointly on $z$ and $\hat{\mathcal{W}}$. For any $z=(x, t) \in \mathbb{R}^{2}$, set $\tau_{z}:=\inf \left\{s: s>0, \hat{\pi}_{z}^{-}(t-s)=\hat{\pi}_{z}^{+}(t-s)\right\}$ and

$$
z^{\prime}:=\left(\frac{1}{2}\left(\pi_{z}^{-}\left(t-\frac{1}{2} \tau_{z}\right)+\pi_{z}^{+}\left(t-\frac{1}{2} \tau_{z}\right)\right), t-\frac{1}{2} \tau_{z}\right) .
$$

Then $z^{\prime}$ depends measurably jointly on $z$ and $\mathcal{W}$ and $\pi_{z}^{\uparrow}$ is the restriction of $\pi_{z^{\prime}}^{+}$to $[t, \infty)$.
Lemma 7.6 (No simultaneous incoming paths) Let $\mathcal{W}$ be a Brownian web and let $x, y \in$ $\mathbb{R}, x \neq y$ be deterministic positions. Then a.s., there exist no time $t \in \mathbb{R}$ such that there exists paths $\pi, \pi^{\prime} \in \mathcal{W}$ with $\sigma_{\pi}, \sigma_{\pi^{\prime}}<t, \pi(t)=x, \pi^{\prime}(t)=y$.

Proof. By [SS08, Lemma 3.4 (b)] it suffices to prove the statement for paths $\pi, \pi^{\prime}$ started at deterministic points. The statement then follows from the fact that two-dimensional Brownian motion a.s. does not hit deterministic points.

Proposition 7.7 (Measurability of quenched laws on path space) Let $\mathbb{Q}_{z}^{+}$and $\mathbb{Q}_{z}^{\dagger}$ be the quenched laws on path space defined in (4.18). Then $\mathbb{R}^{2} \times \Omega \ni(z, \omega) \mapsto \mathbb{Q}_{z}^{+}(\omega) \in \mathcal{M}_{1}(\Pi)$ is a measurable map. An analogue statement holds for $\mathbb{Q}_{z}^{\dagger}$.

Proof. The quenched law $\mathbb{Q}$ is a random variable taking values in $\mathcal{M}_{1}(\mathcal{K}(\Pi))$, i.e., a measurable map $\Omega \ni \omega \mapsto \mathbb{Q}(\omega) \ni \mathcal{M}_{1}(\mathcal{K}(\Pi))$, where $(\omega, \mathcal{F}, \mathbb{P})$ is our underlying probability space. Since $\mathbb{Q}_{z}^{+}=\mathbb{Q} \circ \pi_{z}^{+-1}$, where $(\mathcal{A}, z) \mapsto \pi_{z}^{+}(\mathcal{A})$ is the measurable map from Proposition 7.5. the statement follows from Lemma C. 3 in the appendix. The same argument applies to $\mathbb{Q}_{z}^{1}$.

Proof of Proposition 3.8 (a). Define a continuous map $\Pi \times \mathbb{R} \ni(\pi, t) \mapsto \psi_{t}(\pi) \in \mathbb{R}$ by $\psi_{t}(\pi):=\pi\left(\sigma_{\pi} \vee t\right)$. Then, since

$$
\begin{equation*}
K_{s, t}^{+}(x, \cdot)=\mathbb{Q}_{(x, s)}^{+} \circ \psi_{t}^{-1}=\mathbb{Q} \circ\left(\psi_{t} \circ \pi_{(x, s)}^{+}\right)^{-1} \quad(s, t, x \in \mathbb{R}, s \leq t) \tag{7.42}
\end{equation*}
$$

the statement follows from Lemma C. 3 in the appendix. The same argument applies to $K_{s, t}^{\uparrow}(x, \cdot)$.

Proof of Proposition 2.3. Immediate from Proposition 3.8 (b) and (d).
Proposition 2.4 is a direct consequence of the following proposition, which is formulated on the level of quenched laws on the space of webs.

Proposition 7.8 (Scaling of quenched laws) Let $\mathbb{Q}$ be a Howitt-Warren quenched law with drift $\beta$ and characteristic measure $\nu$. Define scaling maps $S_{a}: R_{\mathrm{c}}^{2} \rightarrow R_{\mathrm{c}}^{2}(a>0)$ as in (3.25) and let $T_{a}: R_{\mathrm{c}}^{2} \rightarrow R_{\mathrm{c}}^{2}(a \in \mathbb{R})$ be defined by $T_{a}(x, t):=(x+a t, t)$. Then:
(a) For each $a>0, S_{a}(\mathbb{Q})$ is a Howitt-Warren quenched law with drift $a^{-1} \beta$ and characteristic measure $a^{-1} \nu$.
(b) For each $a \in \mathbb{R}, T_{a}(\mathbb{Q})$ is a Howitt-Warren quenched law with drift $\beta+a$ and characteristic measure $\nu$.

Proof. To prove part (a) set, choose $\varepsilon_{k}$ and $\mu_{k}$ such that 1.7) holds and set $\varepsilon_{k}^{\prime}:=a \varepsilon_{k}$. Then, by Theorem 3.9. $S_{\varepsilon_{k}}\left(\mathbf{Q}_{\langle k\rangle}\right)$ converges weakly in law to $\mathbb{Q}$ while, since the $\mu_{k}$ satisfy 1.7) with $\varepsilon_{k}$ replaced by $\varepsilon_{k}^{\prime}$ and $\beta$ and $\nu$ replaced by $a^{-1} \beta$ and $a^{-1} \nu$, respectively, $S_{\varepsilon_{k}^{\prime}}\left(\mathbf{Q}_{\langle k\rangle}\right)$ converges weakly in law to a Howitt-Warren quenched law with this drift and characteristic measure. Obviously, $S_{\varepsilon_{k}^{\prime}}\left(\mathbf{Q}_{\langle k\rangle}\right)=S_{a}\left(S_{\varepsilon_{k}}\left(\mathbf{Q}_{\langle k\rangle}\right)\right)$ also converges to $S_{a}(\mathbb{Q})$, so the latter is a Howitt-Warren quenched law with drift $a^{-1} \beta$ and characteristic measure $a^{-1} \nu$.

To prove part (b), let $\left(\mathcal{W}_{0}, \mathcal{M}, \mathcal{W}\right)$ be a marked reference Brownian web and sample Brownian web as in Theorem 3.7. Then $T_{a}\left(\mathcal{W}_{0}\right)$ is a Brownian web with drift $\beta_{0}+a$. It follows from Proposition 3.4 that $T_{a}(\ell)$ is the intersection local time measure of $T_{a}\left(\mathcal{W}_{0}\right)$ and its dual. It follows that conditional on $\mathcal{W}_{0}$, the set $T_{a}(\mathcal{M}):=\left\{\left(T_{a}(z), \omega_{z}\right):\left(z, \omega_{z}\right) \in \mathcal{M}\right\}$ is a Poisson point set with intensity as in 3.16, with $\ell_{1}$ and $\ell_{\mathrm{r}}$ replaced by $T_{a}\left(\ell_{1}\right)$ and $T_{a}\left(\ell_{\mathrm{r}}\right)$. Since $T_{a}(\mathcal{W})$ is constructed from $T_{a}\left(\mathcal{W}_{0}\right)$ and $T_{a}(\mathcal{M})$ in the same way as $\mathcal{W}$ is constructed from $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, in particular, it follows that $\mathbb{P}\left[T_{a}(\mathcal{W}) \in \cdot \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right]$ is a Howitt-Warren quenched law with drift $\beta+a$ and characteristic measure $\nu$.

Proof of Proposition 2.4. Immediate from Proposition 7.8.

## 8 Support properties

In this section, we will first prove Theorem 4.8 on the characterization of Brownian half-nets, establish its connection to Howitt-Warren flows and prove some of its basic properties. We will then prove Theorem 4.9 on the image set of the support of the Howitt-Warren quenched law $\mathbb{Q}$, from which Theorems 2.5 and 2.7 on the support properties of Howitt-Warren processes follow immediately.

### 8.1 Generalized Brownian nets

Proof of Theorem 4.8. Let $\gamma_{n} \in \mathcal{H}_{-}$and $\gamma_{n} \rightarrow \gamma$ in $\Pi$. If $\gamma \notin \mathcal{H}_{-}$so that it crosses some $\pi \in \mathcal{W}$ from left to right, i.e., $\gamma(s)<\pi(s)$ and $\pi(t)<\gamma(t)$ for some $s<t$, then $\gamma_{n}$ crosses $\pi$ from left to right for all $n$ large, a contradiction. Therefore $\mathcal{H}_{-}$is a.s. closed.

The two characterizations of $\mathcal{H}_{-}$in Theorem 4.8 (i) and (ii) are equivalent, because if $\gamma \in \Pi$ crosses some $\pi \in \mathcal{W}$ from left to right, then by the non-crossing property of paths in $\mathcal{W}$ and $\hat{\mathcal{W}}, \gamma$ must also cross some $\hat{\pi} \in \hat{\mathcal{W}}$ from left to right, and the same is true if $\mathcal{W}$ and $\hat{\mathcal{W}}$ are interchanged.

It only remains to show that for each deterministic $z=(x, t) \in \mathbb{R}^{2}$, if $\pi_{z}$ denotes the a.s. unique path in $\mathcal{W}$ starting from $z$, then $\pi_{z}$ is the maximal element in $\mathcal{H}_{-}(z)$. Certainly $\pi_{z} \in$ $\mathcal{H}_{-}(z)$. Note that $z$ is a.s. of type $(0,1)$ in $\mathcal{W}$ by Prop. 3.3, and hence for any positive sequence $\varepsilon_{n} \downarrow 0, \pi_{\left(x+\varepsilon_{n}, t\right)} \rightarrow \pi_{z}$ as $n \rightarrow \infty$. If $\gamma \in \mathcal{H}_{-}(z)$, then it cannot cross $\pi_{\left(x+\varepsilon_{n}, t\right)} \in \mathcal{W}\left(\left(x+\varepsilon_{n}, t\right)\right)$ from left to right. Therefore $\gamma \leq \pi_{\left(x+\varepsilon_{n}, t\right)}$ on $[t, \infty)$ for all $n \in \mathbb{N}$, which implies that $\gamma \leq \pi_{z}$ on $[t, \infty)$. Therefore $\pi_{z}$ is the maximal element in $\mathcal{H}_{-}(z)$.

For any $-\infty \leq \beta_{-} \leq \beta_{+} \leq \infty$ with $\beta_{-}<\infty$ and $-\infty<\beta_{+}$, we define a generalized Brownian net with speeds $\beta_{-}, \beta_{+}$to be a Brownian net with these speeds if $\beta_{-}, \beta_{+}$are both finite, a Brownian half-net with these speeds if one of $\beta_{-}, \beta_{+}$is infinite, and the space of all paths $\Pi$ if both speeds $\beta_{-}, \beta_{+}$are infinite.

Consider a reference Brownian web $\mathcal{W}_{0}$ and set of marked points $\mathcal{M}$ as in Theorem 3.7 and conditional on $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, construct an i.i.d. sequence $\mathcal{W}_{1}, \mathcal{W}_{2}, \ldots$ of sample Brownian webs as in (7.3). For each $n \geq 1$, let $\mathcal{N}_{n}$ denote the Brownian net containing $\mathcal{W}_{1}, \ldots, \mathcal{W}_{n}$ introduced in Lemma 7.1. Recall that $\mathcal{N}_{n}$ has left and right speeds $\beta_{-}(n), \beta_{+}(n)$ given by (2.3), which converge as $n \rightarrow \infty$ to the speeds $\beta_{-}, \beta_{+}$given by (2.12).

Lemma 8.1 (Generalized Brownian net associated with Howitt-Warren flow) Let $\beta_{0} \in \mathbb{R}$ and let $\nu_{1}, \nu_{\mathrm{r}}$ be finite measures on $[0,1]$. Let $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, with $\mathcal{M}=\left\{\left(z, \omega_{z}\right): z \in M\right\}$, be a marked reference web as in Theorem 3.7 and conditional on $(\mathcal{W}, \mathcal{M})$, let $\left(\mathcal{W}_{i}\right)_{i \geq 1}$ be independent sample webs constructed as in (3.19). Then there exists a generalized Brownian net $\mathcal{N}_{\infty}$ with left and right speeds given by (2.12), which is a.s. uniquely defined by $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, such that for any deterministic countable set $\mathcal{D} \subset \mathbb{R}^{2}$ and countable dense set of times $\mathcal{T} \subset \mathbb{R}$

$$
\begin{equation*}
\overline{\mathcal{N}_{\infty}(\mathcal{D})}=\overline{\mathcal{H}_{\mathcal{T}}\left(\bigcup_{i \geq 1} \mathcal{W}_{i}\right)(\mathcal{D})} \tag{8.1}
\end{equation*}
$$

Proof. We treat the cases when $\mathcal{N}_{\infty}$ is a Brownian net, a Brownian halfnet or the space of all paths $\Pi$ separately. If the speeds $\beta_{-}, \beta_{+}$from (2.12) are both finite, then the statements follow from Lemma 7.4 .

If only one of the speeds $\beta_{-}, \beta_{+}$from $(2.12)$ is finite, then by symmetry, we may without loss of generality assume that $-\infty<\beta_{-}$and $\beta_{+}=\infty$. We claim that without loss of generality,
we may further assume that $\nu_{\mathrm{r}}=0$. To see this, let $M_{1}$ and $M_{\mathrm{r}}$ be the restrictions of $M$ to the sets of points of type $(1,2)_{1}$ and $(1,2)_{\mathrm{r}}$ in $\mathcal{W}_{0}$, respectively. As a first step, we reduce our problem to the case that $\nu_{1}(\{1\})=0=\nu_{\mathrm{r}}(\{0\})$. If this is not yet the case, then define $\mathcal{W}_{0}^{\prime}$ and $C^{*}$ as in 7.35) and replace $\mathcal{W}_{0}$ by $\mathcal{W}_{0}^{\prime}$ and $M$ by $M \backslash C^{*}$.

Next, we observe that conditional on $\mathcal{W}_{0}$, the set $M_{\mathrm{r}}$ is a Poisson point set with intensity $c_{\mathrm{r}} \ell_{\mathrm{r}}$, where $c_{\mathrm{r}}:=2 \int(1-q)^{-1} \nu_{\mathrm{r}}(\mathrm{d} q)<\infty$ by our assumption that $-\infty<\beta_{-}$. Let

$$
\begin{equation*}
\mathcal{W}_{0}^{\prime}:=\lim _{\Delta_{m} \uparrow M_{\mathrm{r}}} \operatorname{switch}_{\Delta_{m}}\left(\mathcal{W}_{0}\right) \tag{8.2}
\end{equation*}
$$

Using Proposition 3.6, it is not hard to see that conditional on $\mathcal{W}_{0}^{\prime}$, the set $\left\{\left(z, \omega_{z}\right): z \in M\right\}$ is a Poisson point set on $\mathbb{R}^{2} \times(0,1]$ with intensity

$$
\begin{equation*}
\ell_{1}(\mathrm{~d} z) \otimes\left(21_{\{0<q\}} q^{-1} \nu_{1}(\mathrm{~d} q)+2(1-q)^{-1} \nu_{\mathrm{r}}(\mathrm{~d} q)\right)=\ell_{1}(\mathrm{~d} z) \otimes 2 q^{-1}(1-q)^{-1} 1_{\{0<q\}} \nu(\mathrm{d} q), \tag{8.3}
\end{equation*}
$$

where we have used (3.21) and the fact that $\nu_{1}(\{1\})=0=\nu_{\mathrm{r}}(\{0\})$. By Proposition 3.6 and Lemma 6.17, conditional on $\left(\mathcal{W}_{0}^{\prime}, \omega\right)$, the random variables $\left(\alpha_{z}^{i}\right)_{z \in M}^{i \geq 1}$ defined by $\alpha_{z}^{i}:=\operatorname{sign}_{\mathcal{W}_{i}}(z)$ are independent with $\mathbb{P}\left[\alpha_{z}^{i}=+1 \mid\left(\mathcal{W}_{0}^{\prime}, \omega\right)\right]=\omega_{z}$, and

$$
\begin{equation*}
\mathcal{W}_{i}=\lim _{\Delta_{m} \uparrow A_{i} \cup B_{i}} \operatorname{switch}_{\Delta_{m}}\left(\mathcal{W}_{0}^{\prime}\right) \quad(i \geq 1) \tag{8.4}
\end{equation*}
$$

where $A_{i}=\left\{z \in M: \alpha_{z}^{i} \neq \operatorname{sign}_{\mathcal{W}_{0}^{\prime}}(z)\right\}$ and $B_{i}$ is an independent Poisson point set with intensity $2 \nu_{1}(\{0\})$. Thus, replacing our reference web $\mathcal{W}_{0}$ by $\mathcal{W}_{0}^{\prime}$, we have reduced our problem to the case that $\nu_{\mathrm{r}}=0$ and $\nu_{1}(\mathrm{~d} q)=(1-q)^{-1} \nu(\mathrm{~d} q)$.

In light of this, assume from now on that $\nu_{\mathrm{r}}=0$. Then (3.20) tells us that $\beta=\beta_{0}+$ $2 \nu_{1}([0,1])=\beta_{0}+2 \int(1-q)^{-1} \nu(\mathrm{~d} q)$, hence $\beta_{0}=\beta_{-}$, the left speed from $(2.12)$. Let $\mathcal{N}_{\infty}$ be the Brownian halfnet with left Brownian web $\mathcal{W}_{0}$. Let $\mathcal{N}_{n}(n \geq 1)$ be the Brownian nets defined in Lemma 7.1. We recall from formulas (7.13) and (7.14) in the proof of that lemma that $\mathcal{N}_{n}$ is constructed by switching and then allowing hopping at subsets of the set $C:=\bigcup_{i}\left(A_{i} \cup B_{i}\right)$. Since points in $A_{i} \cup B_{i}$ are of type $(1,2)_{1}$ in $\mathcal{W}_{0}$, it follows that paths in $\mathcal{N}_{n}$ cannot cross paths in $\mathcal{W}_{0}$ from right to left, hence $\mathcal{N}_{n} \subset \mathcal{N}_{\infty}$.

Let $\left(\mathcal{W}_{n}^{\mathrm{l}}, \mathcal{W}_{n}^{\mathrm{r}}\right)$ be the left-right Brownian web associated with $\mathcal{N}_{n}$. Let $z=(x, s) \in \mathbb{R}^{2}$ be deterministic and let $l_{z}^{n}, r_{z}^{n}$ denote the a.s. unique elements of $\mathcal{W}_{n}^{1}(z), \mathcal{W}_{n}^{\mathrm{r}}(z)$, respectively. Since $\mathcal{N}_{n} \subset \mathcal{N}_{n+1}$, one has $l_{z}^{n} \downarrow l_{z}^{\infty}$ and $r_{z}^{n} \uparrow r_{z}^{\infty}$ for some functions $l_{z}^{\infty}:[s, \infty) \rightarrow[-\infty, \infty)$ and $r_{z}^{\infty}:[s, \infty) \rightarrow(-\infty, \infty]$. Since $\mathcal{N}_{n} \subset \mathcal{N}_{\infty}$ for each $n \geq 1$, we have $\pi_{z}^{0} \leq l_{z}^{\infty}$, where $\pi_{z}^{0}$ denotes the a.s. unique element of $\mathcal{W}_{0}$ starting at $z$. Since $\pi_{z}^{0}(t)$ and $l_{z}^{\infty}(t)$ are normally distributed with the same mean for each $t>s$, we must have $\pi_{z}^{0}=l_{z}^{\infty}$. On the other hand, since the $r_{z}^{n}$ are Brownian motions with drifts tending to infinity we must have $r_{z}^{\infty}(t)=\infty$ for all $t>s$.

We are now ready to prove 8.1. Since $\mathcal{W}_{i} \subset \mathcal{N}_{n} \subset \mathcal{N}_{\infty}$ for each $i \leq n$ and since by [SS08, Lemma 8.3], a Brownian net is closed under hopping at deterministic times, we see that $\mathcal{H}_{\mathcal{T}}\left(\bigcup_{i \geq 1} \mathcal{W}_{i}\right)=\bigcup_{n \geq 1} \mathcal{H}_{\mathcal{T}}\left(\mathcal{W}_{1} \cup \cdots \cup \mathcal{W}_{n}\right) \subset \bigcup_{n \geq 1} \mathcal{N}_{n} \subset \mathcal{N}_{\infty}$ and therefore $\overline{\mathcal{H}_{\mathcal{T}}\left(\bigcup_{i \geq 1} \mathcal{W}_{i}\right)(\mathcal{D})} \subset$ $\overline{\mathcal{N}_{\infty}(\mathcal{D})}$. To prove the other inclusion, we first observe that $\overline{\mathcal{N}_{n}(\mathcal{D})}=\overline{\mathcal{H}_{\mathcal{T}}\left(\mathcal{W}_{1} \cup \cdots \cup \mathcal{W}_{n}\right)(\mathcal{D})} \subset$ $\overline{\mathcal{H}_{\mathcal{T}}\left(\bigcup_{i \geq 1} \mathcal{W}_{i}\right)(\mathcal{D})}$ for all $n \geq 1$, hence $\bigcup_{n \geq 1} \mathcal{N}_{n}(\mathcal{D}) \subset \overline{\mathcal{H}_{\mathcal{T}}\left(\bigcup_{i \geq 1} \mathcal{W}_{i}\right)(\mathcal{D})}$. Since $\mathcal{N}_{1} \subset \mathcal{N}_{2} \subset \ldots$ and since by Theorem4.1, each $\mathcal{N}_{n}$ is closed under hopping at intersection times, it follows that also $\bigcup_{n \geq 1} \mathcal{N}_{n}$ is closed under hopping at intersection times, i.e., $\mathcal{H}_{\text {int }}\left(\bigcup_{n \geq 1} \mathcal{N}_{n}\right)=\bigcup_{n \geq 1} \mathcal{N}_{n}$. In view of this, we have $\overline{\mathcal{H}_{\text {int }}\left(\bigcup_{n \geq 1} \mathcal{N}_{n}\right)(\mathcal{D})}=\overline{\bigcup_{n \geq 1} \mathcal{N}_{n}(\mathcal{D})} \subset \overline{\mathcal{H}_{\mathcal{T}}\left(\bigcup_{i \geq 1} \mathcal{W}_{i}\right)(\mathcal{D})}$, so it suffices to show that $\mathcal{N}_{\infty}(\mathcal{D}) \subset \overline{\mathcal{H}_{\text {int }}\left(\bigcup_{n \geq 1} \mathcal{N}_{n}\right)(\mathcal{D})}$. By Lemma 8.4 below, it suffices to show that each path $\pi \in \mathcal{N}_{\infty}(\mathcal{D})$ with $-\infty<\pi<\infty$ on $\left[\sigma_{\pi}, \infty\right)$ can be approximated by paths in $\overline{\mathcal{H}_{\text {int }}\left(\bigcup_{n \geq 1} \mathcal{N}_{n}\right)(\mathcal{D})}$.


Figure 13: Approximation of paths in generalized Brownian nets. On the left: a Brownian halfnet with finite left speed. On the right: the case when both the left and right speeds are infinite.

Let $\pi$ be such a path, $\varepsilon>0$ and $T<\infty$. Let $z_{0}=\left(x_{0}, t_{0}\right)$ denote the starting point of $\pi$ and inductively choose times $t_{k}(k \geq 1)$ and paths $l_{k} \in \mathcal{W}_{0}\left(\pi\left(t_{k}\right), t_{k}\right)(k \geq 0)$ such that $t_{k}=\inf \left\{t \geq t_{k-1}: \pi(t)-l_{k-1}(t) \geq \varepsilon\right\}$. (See Figure 13.) Since paths in $\mathcal{N}_{\infty}$ do not cross paths in $\mathcal{W}_{0}$ from right to left, we can moreover choose the $l_{k}$ such that $l_{k} \leq \pi$ for each $k \geq 0$. The continuity of $\pi$ and the equicontinuity of $\mathcal{W}_{0}$ imply that $t_{m} \geq T$ for some $m \geq 1$. Let $\mathcal{D}^{\prime} \subset \mathbb{R}^{2}$ be a deterministic countable dense set. Choose $z_{k}=\left(x_{k}, t_{k}^{\prime}\right) \in \mathcal{D}^{\prime}(k \geq 1)$ such that $t_{k}<t_{k}^{\prime}<t_{k+1}, x_{k}<l_{k-1}\left(t_{k}^{\prime}\right)$, and $\pi-l_{k-1}<2 \varepsilon$ on $\left[t_{k-1}, t_{k}^{\prime}\right]$ and choose $n_{k}$ such that the right-most path $r_{z_{k}}^{n_{k}}$ in $\mathcal{N}_{n_{k}}\left(z_{k}\right)$ crosses $l_{k}$ before time $t_{k+1}$ and before $\pi-l_{k-1}$ exceeds $2 \varepsilon$. Then the concatenation of the paths $l_{0}, r_{z_{1}}^{n_{1}}, l_{1}, r_{z_{2}}^{n_{2}}, \ldots, r_{z_{m}}^{n_{m}}, l_{m}$ approximates the path $\pi$ on $\left[t_{0}, T\right]$ within distance $2 \varepsilon$. We are not quite done yet, however, since the paths $l_{k} \in \mathcal{W}_{0}$ are not elements of $\bigcup_{n \geq 1} \mathcal{N}_{n}$.

To finish the argument, we will show that each path in $\mathcal{W}_{0}$ can be approximated by paths in $\bigcup_{n \geq 1} \mathcal{N}_{n}$ and in case $z_{0}$ is a deterministic point, that $l_{0}$ can be approximated by paths in $\bigcup_{n \geq 1} \mathcal{N}_{n}\left(z_{0}\right)$. Indeed, this follows from the fact, proved above, that for each deterministic $z \in \mathbb{R}^{2}$, the left-most path $l_{z}^{n}$ in $\mathcal{N}_{n}(z)$ converges to the a.s. unique path in $\mathcal{W}_{0}(z)$, and the fact that $\mathcal{W}_{0}=\overline{\mathcal{W}_{0}\left(\mathcal{D}^{\prime}\right)}$. Replacing the paths $l_{0}, \ldots, l_{m}$ by sufficiently close approximating paths $l_{0}^{\prime}, \ldots, l_{m}^{\prime} \in \bigcup_{n>1} \mathcal{N}_{n}$, with $l_{0}^{\prime}$ also starting at $z_{0}$, we see that these approximating paths are crossed by the paths $r_{z_{1}}^{n_{1}}, \ldots, r_{z_{m}}^{n_{m}}$ and hence $\pi$ can be approximated by a concatenation of $l_{0}^{\prime}, r_{z_{1}}^{n_{1}}, l_{1}^{\prime}, r_{z_{2}}^{n_{2}}, \ldots, r_{z_{m}}^{n_{m}}, l_{m}^{\prime}$. This completes the proof for case $-\infty<\beta_{-}$and $\beta_{+}=\infty$, where
we have identified $N_{\infty}$ as a Brownian half-net.
The proof for the case $\beta_{-}=-\infty$ and $\beta_{+}=\infty$ is similar, but easier. In this case, for an arbitrary path $\pi \in \Pi$ with starting point $z_{0}=\left(x_{0}, t_{0}\right) \in \mathcal{D}$, we inductively choose times $t_{k}(k \geq 1)$ and paths $\pi_{k} \in \mathcal{W}_{1}\left(\pi\left(t_{k}\right), t_{k}\right)(k \geq 0)$ such that $t_{k}=\inf \left\{t \geq t_{k-1}\right.$ : $\left.\left|\pi(t)-\pi_{k-1}(t)\right| \geq \varepsilon\right\}$. If $\pi_{k}\left(t_{k}\right)<\pi\left(t_{k}\right)$, then we use a right-most path $r_{z_{k}}^{n_{k}}$ of the Brownian net $\mathcal{N}_{n}:=\overline{\mathcal{H}_{\mathcal{T}}\left(\mathcal{W}_{1} \cup \cdots \cup \mathcal{W}_{n}\right)}$ to connect $\pi_{k}$ to $\pi_{k+1}$ and if $\pi\left(t_{k}\right)<\pi_{k}\left(t_{k}\right)$, then we use a left-most path $l_{z_{k}}^{n_{k}}$ of $\mathcal{N}_{n}$.

The proof of Lemma 8.1 has a useful corollary.
Corollary 8.2 (Paths going to infinity) Let $\mathcal{N}$ be a generalized Brownian net with infinite right speed. Then, for each deterministic $z=(x, t) \in \mathbb{R}^{2}$, there a.s. exist paths $r_{z}^{n} \in \mathcal{N}(z)$ with $r_{z}^{n}(u) \uparrow \infty$ for all $u>t$.

Proof. If $\mathcal{N}$ is a Brownian half-net, then we may take for $r_{z}^{n}$ the a.s. unique right-most paths in the Brownian nets $\mathcal{N}_{n}$ as in the proof of Lemma 8.1. If $\mathcal{N}=\Pi$ the statement is trivial.

Lemma 8.3 (Hopping with left-most paths) Let $\mathcal{N}$ be a generalized Brownian net with left speed $-\infty<\beta_{-}$and let $\mathcal{W}^{1}$ be its associated left Brownian web. Then, a.s. for each $\pi \in \mathcal{N}$ and $l \in \mathcal{W}^{1}$ with $\pi(t)<l(t)$ at $t:=\sigma_{\pi} \vee \sigma_{l}$, the following statements hold:
(i) The path $\pi^{\prime}$ defined by $\sigma_{\pi^{\prime}}:=\sigma_{\pi}, \pi^{\prime}:=\pi$ on $\left[\sigma_{\pi}, t\right]$ and $\pi^{\prime}:=\pi \wedge l$ on $[t, \infty)$ satisfies $\pi^{\prime} \in \mathcal{N}$.
(ii) The path $\pi^{\prime \prime}$ defined by $\sigma_{\pi^{\prime \prime}}:=\sigma_{l}, \pi^{\prime \prime}:=l$ on $\left[\sigma_{l}, t\right]$ and $\pi^{\prime \prime}:=\pi \vee l$ on $[t, \infty)$ satisfies $\pi^{\prime \prime} \in \mathcal{N}$.

Proof. Set $\tau:=\inf \{u \geq t: \pi(u)=l(u)\}$. Since paths in $\mathcal{N}$ cannot cross paths in $\mathcal{W}^{l}$ from right to left, one has $\pi^{\prime}=\pi$ on $\left[\sigma_{\pi}, \tau\right]$ and $\pi^{\prime}=l$ on $[\tau, \infty]$, and also, $\pi^{\prime \prime}=l$ on $\left[\sigma_{l}, \tau\right]$ and $\pi^{\prime \prime}=\pi$ on $[\tau, \infty]$. Therefore, if $\mathcal{N}$ is a Brownian net, then the statements follow from the fact that $\mathcal{N}$ is closed under hopping at intersection times (see Theorem 4.1). We do not know if Brownian half-nets are closed under hopping at intersection times, so in this case we check directly that the paths $\pi^{\prime}$ and $\pi^{\prime \prime}$ do not cross paths in $\mathcal{W}^{l}$ from right to left and hence (by Theorem 4.8) are elements of $\mathcal{N}$. If $l^{\prime} \in \mathcal{W}^{1}$ satisfies $l^{\prime}(s)<\pi^{\prime}(s)$ for some $s \in\left[\sigma_{\pi}, \infty\right)$, then $l^{\prime} \leq \pi^{\prime}$ on $[s, \tau]$ by the fact that $\pi \in \mathcal{N}$ while $l^{\prime} \leq \pi^{\prime}$ on $[\tau, \infty)$ by the fact that $l^{\prime}$ cannot cross $l$. If $l^{\prime} \in \mathcal{W}^{1}$ satisfies $l^{\prime}(s)<\pi^{\prime \prime}(s)$ for some $s \in\left[\sigma_{l}, \tau\right]$, then $l^{\prime} \leq l \leq \pi^{\prime \prime}$ on $[s, \infty)$ by the fact that $l^{\prime}$ cannot cross $l$, while if $l^{\prime} \in \mathcal{W}^{1}$ satisfies $l^{\prime}(s)<\pi^{\prime \prime}(s)$ for some $s \in[\tau, \infty]$, then $l^{\prime} \leq \pi \leq \pi^{\prime \prime}$ on $[s, \infty)$ by the fact that $l^{\prime}$ cannot cross $\pi$ from left to right.

Lemma 8.4 (Finite paths) Let $\mathcal{N}$ be a generalized Brownian net and let $\mathcal{N}_{\text {fin }}:=\{\pi \in \mathcal{N}$ : $-\infty<\pi<\infty$ on $\left.\left[\sigma_{\pi}, \infty\right)\right\}$. Then a.s. for each $z \in \mathbb{R}^{2}, \overline{\mathcal{N}_{\text {fin }}(z)}=\mathcal{N}(z)$.

Proof. We need to show that for each $\pi \in \mathcal{N}(z)$ there exist $\pi_{n} \in \mathcal{N}_{\text {fin }}(z)$ with $\pi_{n} \rightarrow \pi$. In case both the left and right speeds of $\mathcal{N}$ are infinite, and hence $\mathcal{N}=\Pi$, we may simply take $\pi_{n}^{\prime}:=-n \vee\left(n \wedge \pi_{n}\right)$. Otherwise, by symmetry, we may without loss of generality assume that the left speed of $\mathcal{N}$ is finite. Let $\mathcal{W}^{1}$ be the left Brownian web. Write $z=(x, t)$ and for $n$ large enough such that $-n \leq x \leq n$ and $t \leq n$, choose $l_{n}^{-}, l_{n}^{+} \in \mathcal{W}^{1}$ with $\sigma_{l_{n}^{-}}=\sigma_{l_{n}^{+}}:=t$ such
that $l_{n}^{-} \leq-n$ and $n \leq l_{n}^{+}$on $[t, n]$ and set $\pi_{n}:=l_{n}^{-} \vee\left(l_{n}^{+} \wedge \pi\right)$. Then $\pi_{n} \in \mathcal{N}$ by Lemma 8.3. $-\infty<\pi_{n}<\infty$ on $[t, \infty)$, and $\pi_{n} \rightarrow \pi$.

In Lemmas $7.1,7.4$ and 8.1 , we have seen objects of the form $\overline{\mathcal{N}(\mathcal{D})}$ where $\mathcal{N}$ is a generalized Brownian net and $\mathcal{D} \subset \mathbb{R}^{2}$ is a deterministic countable set. Naively, one might guess that $\overline{\mathcal{N}(\mathcal{D})}=\mathcal{N}(\overline{\mathcal{D}})$, where $\overline{\mathcal{D}}$ denotes the closure of $\mathcal{D}$ in $R_{\mathrm{c}}^{2}$. It turns out, however, that this is not always true. In particular, it may happen that $(*,-\infty) \in \overline{\mathcal{D}}$ but $\mathcal{N}(*,-\infty) \not \subset \overline{\mathcal{N}(\mathcal{D})}$. Our next result shows that this is indeed all that can go wrong. It is not very difficult, but rather tedious, to give a precise decription of $\overline{\mathcal{N}(\mathcal{D})}(*,-\infty)$ in terms of the shape of $\mathcal{D}$ near $(*,-\infty)$. Since we will not need such a precise description, we will settle for the following lemma.

Lemma 8.5 (Closure of paths started from a countable set) Let $\mathcal{N}$ be a generalized Brownian net, let $\mathcal{D} \subset \mathbb{R}^{2}$ be a deterministic countable set, let $\overline{\mathcal{N}(\mathcal{D})}$ denote the closure of $\mathcal{N}(\mathcal{D})$ in $\Pi$ and let $\overline{\mathcal{D}}$ denote the closure of $\mathcal{D}$ in $R_{\mathrm{c}}^{2}$. Then

$$
\begin{equation*}
\mathcal{N}(\overline{\mathcal{D}} \backslash\{(*,-\infty)\}) \subset \overline{\mathcal{N}(\mathcal{D})} \subset \mathcal{N}(\overline{\mathcal{D}}) \quad \text { a.s. } \tag{8.5}
\end{equation*}
$$

If $\mathcal{D}$ is dense in $\mathbb{R}^{2}$, then moreover $\mathcal{N}(*,-\infty) \subset \overline{\mathcal{N}(\mathcal{D})}$.
Proof. Since $\mathcal{N}$ is closed and since convergence of paths implies convergence of their starting points, the inclusion $\overline{\mathcal{N}(\mathcal{D})} \subset \mathcal{N}(\overline{\mathcal{D}})$ is trivial. As a next step, we will show that $\mathcal{N}_{\text {fin }}(\hat{\mathcal{D}}) \subset$ $\overline{\mathcal{N}(\mathcal{D})}$, where $\hat{\mathcal{D}}$ denotes the closure of $\mathcal{D}$ in $\mathbb{R}^{2}$ and $\mathcal{N}_{\text {fin }}$ is defined in Lemma 8.4. If $\mathcal{N}$ is a Brownian net, then the statement follows from [SS08, Lemma 8.1]. If $\mathcal{N}=\Pi$, the statement is trivial, so it suffices to treat the case when $\mathcal{N}$ is a Brownian half-net with, say, a finite left speed. We need to show that any path in $\mathcal{N}_{\text {fin }}(\hat{\mathcal{D}})$ can be approximated by paths in $\mathcal{N}(\mathcal{D})$. We use the steering argument from the proof of Lemma 8.1 (see Figure 13), where in this case, for any $z_{0}=\left(x_{0}, t_{0}\right) \in \hat{\mathcal{D}}$, we choose for $l_{0}$ the left-most path in $\mathcal{W}_{0}\left(z_{0}\right)$, which by the fact that paths in $\mathcal{N}$ cannot cross paths in $\mathcal{W}_{0}$ from right to left is guaranteed to stay on the left of any path $\pi \in \mathcal{N}$ starting from $z_{0}$. We need to show that $l_{0}$ can be approximated by paths in $\bigcup_{n \geq 1} \mathcal{N}_{n}$ that moreover start in $\mathcal{D}$. Since in the proof of Lemma 8.1 it has already been shown that any path in $\mathcal{W}_{0}(\mathcal{D})$ can be approximated by paths in $\bigcup_{n \geq 1} \mathcal{N}_{n}(\mathcal{D})$, by a diagonal argument, it suffices to show that $l_{0}$ can be approximated by paths in $\mathcal{W}_{0}(\mathcal{D})$. By the structure of special points in a Brownian web, such an approximation is possible unless there exists a dual path $\hat{l} \in \hat{\mathcal{W}}_{0}$ entering $z_{0}$ and an $\varepsilon>0$ such that $\left\{z=(x, t):\left|z-z_{0}\right|<\varepsilon, x \leq \hat{l}(t)\right\} \cap \mathcal{D}=\emptyset$. But by Lemma 8.6 below, such a dual path a.s. does not exist for any $z_{0} \in \hat{\mathcal{D}}$.

This completes the proof that $\mathcal{N}_{\text {fin }}(\hat{\mathcal{D}}) \subset \overline{\mathcal{N}(\mathcal{D})}$. By Lemma 8.4, it follows that also $\mathcal{N}(\hat{\mathcal{D}}) \subset \overline{\mathcal{N}(\mathcal{D})}$. Obviously, the trivial path starting at time $(*, \infty)$ is an element of $\overline{\mathcal{N}(\mathcal{D})}$ if and only if there exists $\mathcal{D} \ni z_{n} \rightarrow(*, \infty)$. To complete the proof that $\mathcal{N}(\overline{\mathcal{D}} \backslash\{(*,-\infty)\}) \subset \overline{\mathcal{N}(\mathcal{D})}$, by symmetry, it therefore suffices to show that if $\pi \in \mathcal{N}(\overline{\mathcal{D}})$ starts at some $z \in\{-\infty\} \times \mathbb{R}$, then $\pi$ can be approximated by $\pi_{n} \in \mathcal{N}(\hat{\mathcal{D}})$. Let $\mathcal{D} \ni z_{n} \rightarrow z$ and in case either the left or right speed of $\mathcal{N}$ is finite, let $\pi_{n}^{\prime}$ be a.s. unique path in the left or right Brownian web starting from $z_{n}$. Then by Lemma 8.3, $\pi_{n}:=\pi_{n}^{\prime} \vee \pi \in \mathcal{N}$ and $\pi_{n} \rightarrow \pi$. If both speeds are infinite, we may use the same argument with $\pi_{n}^{\prime}$ replaced by a constant path. This completes the proof of (8.5).

To prove that also $\mathcal{N}(*,-\infty) \subset \overline{\mathcal{N}(\mathcal{D})}$ if $\mathcal{D}$ is dense in $\mathbb{R}^{2}$, by $(8.5)$, it suffices to prove that each path $\pi \in \mathcal{N}(*,-\infty)$ can be approximated by paths $\pi_{n} \in \mathcal{N}\left(R_{\mathrm{c}}^{2} \backslash\{(*,-\infty)\}\right)$. In view of this, taking for $\pi_{n}$ the restriction of $\pi$ to $[-n, \infty]$ completes the proof of the lemma.

A statement very similar to the lemma below has been demonstrated in the proof of [SS08, Lemma 8.1].

Lemma 8.6 (The Brownian web does not skim closed sets) Let $\mathcal{W}$ be a Brownian web and let $K \subset \mathbb{R}^{2}$ be a deterministic closed set. Then a.s. there exist no $z_{0}=\left(x_{0}, t_{0}\right) \in K$ and $\pi \in \mathcal{W}$ entering $z_{0}$, such that $\left\{z=(x, t) \in \mathbb{R}^{2}:\left|z-z_{0}\right|<\varepsilon, x<\pi(t)\right\} \cap K=\emptyset$ for some $0<\varepsilon<t-\sigma_{\pi}$.

Proof. It suffice to prove the statement for paths $\pi$ started from a deterministic point and for deterministic $\varepsilon>0$. By cutting $K$ into countably any pieces of diameter at most $\varepsilon / 4$ and using translation invariance, we can reduce the problem to the following statement: let $\pi$ be a Brownian motion started at time zero in the origin, let $K$ be a deterministic closed subset of $\mathbb{R}^{2}$, let $\varepsilon>0$ and let $U:=\{(x, t): t \geq \varepsilon, x<\pi(t)\}$. Then the event that $U \cap K=\emptyset$ but $\bar{U} \cap K \neq \emptyset$ has probability zero. To see that this is the case, set $\pi^{\varepsilon}(t):=\pi(t)-\pi(\varepsilon)(t \geq \varepsilon)$. Then, conditional on the path $\left(\pi^{\varepsilon}(t)\right)_{t \geq \varepsilon}$, there is at most one value of $\pi(\varepsilon)$ for which the event $\{U \cap K=\emptyset, \bar{U} \cap K \neq \emptyset\}$ occurs. Since $\pi(\varepsilon)$ is independent of $\left(\pi^{\varepsilon}(t)\right)_{t \geq \varepsilon}$ and normally distributed, we see that the conditional probability of $\{U \cap K=\emptyset, \bar{U} \cap K \neq \emptyset\}$ given $\left(\pi^{\varepsilon}(t)\right)_{t \geq \varepsilon}$ is zero, hence integrating over the distribution of $\left(\pi^{\varepsilon}(t)\right)_{t \geq \varepsilon}$ yields the desired result.

Let $\mathcal{N}$ be a generalized Brownian net with left and right speeds $\beta_{-}, \beta_{+}$(which may be infinite). Then, generalizing (4.6) and (6.41), for any closed subset $A \subset \mathbb{R}$, setting

$$
\begin{equation*}
\xi_{t}^{A}:=\{\pi(t): \pi \in \mathcal{N}(A \times\{0\})\} \quad(t \geq 0) \tag{8.6}
\end{equation*}
$$

defines a Markov process taking values in the space of closed subsets of $\mathbb{R}$, which we call the branching-coalescing point set with left and right speeds $\beta_{-}, \beta_{+}$.

Lemma 8.7 (Edge of a branching-coalescing point set) Let $A$ be a deterministic nonempty closed subset of the real line and let $\left(\xi_{t}^{A}\right)_{t \geq 0}$ be the branching-coalescing point set with left and right speeds $\beta_{-}, \beta_{+}$defined in 8.6). Set $r_{t}:=\sup \left(\xi_{t}^{A}\right)(t \geq 0)$. Then:
(a) If $\beta_{+}<\infty$ and $r_{0}<\infty$, then $\left(r_{t}\right)_{t \geq 0}$ is a Brownian motion with drift $\beta_{+}$. If $\beta_{+}<\infty$ and $r_{0}=\infty$, then $r_{t}=\infty$ for all $t \geq 0$.
(b) If $\beta_{+}=\infty$, then $r_{t}=\infty$ for all $t>0$.
(c) If $\beta_{-}=-\infty$ and $\beta_{+}<\infty$, then $\xi_{t}^{A}=\left(-\infty, r_{t}\right] \cap \mathbb{R}$ for all $t>0$.
(d) If $\beta_{-}=-\infty$ and $\beta_{+}=\infty$, then $\xi_{t}^{A}=\mathbb{R}$ for all $t>0$.

Proof. To prove part (a), let $\mathcal{W}^{\mathrm{r}}$ be the right Brownian web associated with $\mathcal{N}$, which is to be interpreted as the right Brownian web in case both speeds are finite. If $\sup (A)<\infty$, then let $r$ be the a.s. unique path in $\mathcal{W}^{\mathrm{r}}$ started from $r_{0}=\sup (A)$. Now $r$ is a Brownian motion with drift $\beta_{+}, r \in \mathcal{N}$, and $\pi \leq r$ for any path in $\mathcal{N}$ started from $A \times\{0\}$ by the fact that paths in $\mathcal{N}$ cannot cross paths in $\mathcal{W}^{\mathrm{r}}$ (see Theorem4.8(i) and [SS08, Prop. 1.8]). If $\sup (A)=\infty$, then choose $x_{n} \in A$ with $x_{n} \uparrow \infty$ and let $r_{n}$ be the a.s. unique paths in $\mathcal{W}^{\mathrm{r}}$ started from ( $x_{n}, 0$ ). Then $\mathbb{P}\left[\inf \left\{r^{n}(t): 0 \leq t \leq T\right\} \leq N\right] \rightarrow 0$ as $n \uparrow \infty$ for each $N, T<\infty$, hence $\sup _{n} r_{t}^{n}=\infty$ for all $t \geq 0$ a.s.

Part (b) is an immediate consequence of Corollary 8.2.
To prove part (c), we first consider the case that $\sup (A)<\infty$. Let $\mathcal{W}^{\mathrm{r}}$ be the right Brownian web of $\mathcal{N}$ and let $r$ be the a.s. unique path in $\mathcal{W}^{r}$ started from $r_{0}=\sup (A)$. It has been shown in the proof of part (a) that $\xi_{t}^{A} \subset\left(-\infty, r_{t}\right]$ for all $t \geq 0$. To prove the other inclusion, it suffices to show that for each $r^{\prime} \in \mathcal{W}^{\mathrm{r}}$ started from $\left\{(x, t): t \geq 0, x<r_{t}\right\}$ and $\varepsilon>0$ there exists a path $\pi \in \mathcal{N}$ started at time zero from $\pi(0)=\sup (A)$ such that $\pi=r^{\prime}$
on $\left[\sigma_{r^{\prime}}+\varepsilon, \infty\right)$. Let $\mathcal{D} \subset \mathbb{R}^{2}$ be a deterministic countable dense set. By Corollary 8.2 , we can find some path $\pi^{\prime} \in \mathcal{N}(\mathcal{D})$ that starts on the right of $r$ and crosses both $r$ and $r^{\prime}$ before time $\sigma_{r^{\prime}}+\varepsilon$. Let $\pi$ be the concatenation of $r, \pi$ and $r^{\prime}$. Then $\pi \in \mathcal{N}$ by Lemma 8.3, $\pi(0)=\sup (A)$, and $\pi=r^{\prime}$ on $\left[\sigma_{r^{\prime}}+\varepsilon, \infty\right)$. The proof in case $\sup (A)=\infty$ is similar, where now instead of $r$ we use a sequence of paths $r^{n}$ started from points $\left(x_{n}, 0\right)$ with $A \ni x_{n} \uparrow \infty$.

Part (d), finally, is trivial since in this case $\mathcal{N}=\Pi$.

### 8.2 Support properties of Howitt-Warren flows and quenched laws

We are now ready to prove Theorems 2.5 and 2.7 on the support properties of Howitt-Warren processes and Theorem 4.9 on the support of quenched laws. We start by preparing for the proof of the latter. In line with notation introduced in Section 4.5, we set $\mathbb{Q}_{z}^{+}:=\mathbb{P}\left[\pi_{z}^{+} \in\right.$ $\left.\cdot \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right]$, where $\left(\mathcal{W}_{0}, \mathcal{M}\right)$ is the marked reference Brownian web as in Theorem 3.7 and $\pi_{z}^{+}$ denotes the right-most path in the sample Brownian web $\mathcal{W}$ started from $z$.

Lemma 8.8 (Support of quenched laws on the space of paths) Conditional on a marked reference Brownian web $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, let $\left(\mathcal{W}_{i}\right)_{i \geq 1}$ be an i.i.d. sequence of sample Brownian webs as in (7.3), and let $\mathcal{N}_{\infty}$ be the generalized Brownian net associated with $\left(\mathcal{W}_{0}, \mathcal{M}\right)$ defined in Lemma 8.1. Then, for any deterministic $z \in \mathbb{R}^{2}$,

$$
\begin{equation*}
\operatorname{supp}\left(\mathbb{Q}_{z}^{+}\right)=\mathcal{N}_{\infty}(z) \quad \text { a.s. } \tag{8.7}
\end{equation*}
$$

Proof. Since $\mathbb{Q}_{z}^{+}$is the conditional law, given the marked reference Brownian web $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, of the a.s. unique path in $\mathcal{W}_{1}$ starting at $z$, and since $\mathcal{W}_{1} \subset \mathcal{N}_{\infty}$ a.s., the inclusion $\operatorname{supp}\left(\mathbb{Q}_{z}^{+}\right) \subset$ $\mathcal{N}_{\infty}(z)$ is trivial.

To prove the other inclusion, by Lemma 8.1, it suffices to show that

$$
\begin{equation*}
\mathcal{H}_{\mathcal{T}}\left(\mathcal{W}_{1} \cup \cdots \cup \mathcal{W}_{n}\right)(z) \subset \operatorname{supp}\left(\mathbb{Q}_{z}^{+}\right) \quad \text { a.s. } \tag{8.8}
\end{equation*}
$$

for each $n \geq 1$. Fix $1 \leq i_{0}, \ldots, i_{m} \leq n$ and $t_{1}, \ldots, t_{m} \in \mathcal{T}$ with $t_{1}<\cdots<t_{m}$, and set $t_{0}:=-\infty, t_{m+1}:=+\infty$. Let $\mathcal{W}$ be the concatenation of $\mathcal{W}_{i_{0}}, \ldots, \mathcal{W}_{i_{m}}$ on the time intervals $\left[t_{0}, t_{1}\right], \ldots,\left[t_{m}, t_{m+1}\right]$, i.e., $\mathcal{W}$ is the set of all paths $\pi$ such that for each $k=0, \ldots, m$ there is a $\pi^{\prime} \in \mathcal{W}_{i_{k}}$ with $\pi=\pi^{\prime}$ on $\left[\sigma_{\pi}, \infty\right] \cap\left[t_{k}, t_{k+1}\right]$. Since conditional on the marked reference web $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, restrictions of a sample Brownian web $\mathcal{W}_{i}$ to disjoint space-time regions are independent, we see that the conditional distribution of $\mathcal{W}$ equals that of the $\mathcal{W}_{i}$ 's. In particular, conditional on $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, the a.s. unique path in $\mathcal{W}(z)$ is distributed with law $\mathbb{Q}_{z}^{+}$. Since $t_{1}, \ldots, t_{m} \in \mathcal{T}$ are arbitrary, this proves (8.8).

Proof of Theorem 4.9. Without loss of generality, we may assume that $\mu$ is a probability measure. In the set-up of Lemma 8.8, let $\left(Z_{j}\right)_{j \geq 1}$ be an i.i.d. sequence of $\mathbb{R}^{2}$-valued random variables with law $\mu$, independent of the marked reference Brownian web $\left(\mathcal{W}_{0}, \mathcal{M}\right)$ and sequence of sample Brownian webs $\left(\mathcal{W}_{i}\right)_{i \geq 1}$. Then, conditional on $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, for each $i, j \geq 1$, the random variable $\pi_{Z_{j}}^{+i}$ has law $\int \mu(\mathrm{d} z) \mathbb{Q}_{z}^{+}$, where $\pi_{z}^{+i}$ denotes the rightmost path in $\mathcal{W}_{i}$ started at $z$. Since $\mathcal{W}_{i} \subset \mathcal{N}_{\infty}$ and $Z_{j} \in \operatorname{supp}(\mu)$ a.s., we see that $\pi_{Z_{j}}^{+i} \in \mathcal{N}_{\infty}(\operatorname{supp}(\mu))$ a.s. and hence $\operatorname{supp}\left(\int \mu(\mathrm{d} z) \mathbb{Q}_{z}^{+}\right) \subset \overline{\mathcal{N}_{\infty}(\operatorname{supp}(\mu))}$ a.s.

To prove the other inclusion, set $\mathcal{D}:=\left\{Z_{j}: j \geq 0\right\}$. Since $\pi_{Z_{j}}^{+i} \in \operatorname{supp}\left(\int \mu(\mathrm{~d} z) \mathbb{Q}_{z}^{+}\right)$a.s. for each $i, j \geq 1$, and conditional on $\left(\mathcal{W}_{0}, \mathcal{M},\left(Z_{j}\right)_{j \geq 1}\right)$, the random variable $\pi_{Z_{j}}^{+i}$ has law $\mathbb{Q}_{Z_{j}}^{+}$, we see that $\operatorname{supp}\left(\mathbb{Q}_{Z_{j}}^{+}\right) \subset \operatorname{supp}\left(\int \mu(\mathrm{d} z) \mathbb{Q}_{z}^{+}\right)$a.s. for each $j \geq 1$. By Lemma 8.8, it follows
that $\mathcal{N}_{\infty}(\mathcal{D}) \subset \operatorname{supp}\left(\int \mu(\mathrm{d} z) \mathbb{Q}_{z}^{+}\right)$a.s., which by Lemma 8.5 implies that $\mathcal{N}_{\infty}(\operatorname{supp}(\mu)) \subset$ $\operatorname{supp}\left(\int \mu(\mathrm{d} z) \mathbb{Q}_{z}^{+}\right)$a.s. (Recall that $\operatorname{supp}(\mu)$ is the support of $\mu$ in $\mathbb{R}^{2}$, not $R_{\mathrm{c}}^{2}$, which is why Lemma 8.5 can be applied here.)

Proposition 8.9 (Support of Howitt-Warren process) Let $\left(\mathcal{W}_{0}, \mathcal{M}\right)$ be a marked reference Brownian web, let $\mathcal{N}_{\infty}$ be the generalized Brownian net associated with $\left(\mathcal{W}_{0}, \mathcal{M}\right)$ defined in Lemma 8.1, and for each closed $A \subset \mathbb{R}$, let $\left(\xi_{t}^{A}\right)_{t \geq 0}$ be the branching-coalescing point set associated with $\mathcal{N}_{\infty}$ defined in (8.6). Then, for each deterministic finite measure $\rho_{0}$ on $\mathbb{R}$, alomst surely

$$
\begin{equation*}
\operatorname{supp}\left(\rho_{t}\right)=\xi_{t}^{\operatorname{supp}\left(\rho_{0}\right)} \quad(t \geq 0) \tag{8.9}
\end{equation*}
$$

where $\left(\rho_{t}\right)_{t \geq 0}$ is the Howitt-Warren process defined as in 2.1) with $K_{0, t}=K_{0, t}^{+}$or $K_{0, t}^{\uparrow}$, where $\left(K_{s, t}^{+}\right)_{s \leq t}$ and $\left(K_{s, t}^{\uparrow}\right)_{s \leq t}$ are the versions of the Howitt-Warren flow defined in Theorem 3.7.

Proof. In line with notation introduced in (6.10), let $\Pi\left(\Sigma_{0}\right):=\left\{\pi \in \Pi: \sigma_{\pi}=0\right\}$. For each $t \geq 0$, define a continuous map $\psi_{t}: \Pi\left(\Sigma_{0}\right) \rightarrow[-\infty, \infty]$ by $\psi_{t}(\pi):=\pi(t)$. Recall that if $E, F$ are Polish spaces, $\mu$ is a finite measure on $E, f: E \rightarrow F$ is a continuous function, and $f(\mu)$ denotes the image of $\mu$ under $f$, then $\operatorname{supp}(f(\mu))=\overline{f(\operatorname{supp}(\mu))}$. Then by Theorem 4.9.

$$
\begin{align*}
& \operatorname{supp}\left(\rho_{t}\right)=\operatorname{supp}\left(\psi_{t}\left(\int \rho_{0}(\mathrm{~d} x) \mathbb{Q}_{(x, 0)}^{+}\right)\right)=\overline{\psi_{t}\left(\operatorname{supp}\left(\int \rho_{0}(\mathrm{~d} x) \mathbb{Q}_{(x, 0)}^{+}\right)\right)}  \tag{8.10}\\
& \quad=\overline{\psi_{t}\left(\mathcal{N}_{\infty}\left(\operatorname{supp}\left(\rho_{0}\right) \times\{0\}\right)\right)}=\xi_{t}^{\operatorname{supp}\left(\rho_{0}\right)} \quad(t \geq 0)
\end{align*}
$$

By the remarks above Theorem 4.9. replacing $\mathbb{Q}_{(x, 0)}^{+}$with $\mathbb{Q}_{(x, 0)}^{\uparrow}$ makes no difference.
Proof of Theorems 2.5 and 2.7. Immediate from Lemma 8.7 and Proposition 8.9.

## 9 Atomic or non-atomic

In this section, we use our construction of the Howitt-Warren flows in Theorems 3.7 and 4.7 to prove Theorem 2.8 on the atomicness/non-atomicness of the Howitt-Warren processes. Parts (a), (b) and (c) are proved in Sections 9.1, 9.2 and 9.3 , respectively.

### 9.1 Atomicness at deterministic times

To prove Theorem 2.8 (a) on the atomicness of any Howitt-Warren process $\left(\rho_{t}\right)_{t \geq 0}$ at deterministic times, we need the following lemma.

Lemma 9.1 (Coincidence of points entered by a path) Let $\left(\mathcal{W}, \mathcal{W}^{\prime}\right)$ be a pair of sticky Brownian webs and for $t \in \mathbb{R}$, let $I(t):=\left\{\pi(t): \pi \in \mathcal{W}, \sigma_{\pi}<t\right\}$ and let $I^{\prime}(t)$ be defined similarly with $\mathcal{W}$ replaced by $\mathcal{W}^{\prime}$. Then for each deterministic $t \in \mathbb{R}$, a.s. $I(t)=I^{\prime}(t)$.

Proof. For each $s<t$, set $I(s, t):=\left\{\pi(t): \pi \in \mathcal{W}, \sigma_{\pi}=s\right\}$. Let $\mathcal{T} \subset \mathbb{R}$ be a deterministic countable dense set. It suffices to show that for each $s, t \in \mathcal{T}$ with $s<t$ one has $I(s, t) \subset I^{\prime}(t)$. As in the proof of Theorem 3.5 (in Section 6.6), we may without loss of generality assume that $\mathcal{W}$ and $\mathcal{W}^{\prime}$ are embedded in a Brownian net $\mathcal{N}$ with associated left-right Brownian web $\left(\mathcal{W}^{\mathrm{l}}, \mathcal{W}^{\mathrm{r}}\right)$ and set of separation points $S$ in such a way that

$$
\begin{align*}
\mathcal{W} & =\left\{\pi \in \mathcal{N}: \operatorname{sign}_{\pi}(z)=\alpha_{z} \forall z \in S \text { s.t. } \pi \text { enters } z\right\} \\
\mathcal{W}^{\prime} & =\left\{\pi \in \mathcal{N}: \operatorname{sign}_{\pi}(z)=-\alpha_{z} \forall z \in S \text { s.t. } \pi \text { enters } z\right\} \tag{9.1}
\end{align*}
$$

where conditional on $\mathcal{N}$, the $\left(\alpha_{z}\right)_{z \in S}$ are i.i.d. $\{-1,+1\}$-valued random variables. Then $I(s, t) \subset\left\{\pi(t): \pi \in \mathcal{N}: \sigma_{\pi}=s\right\}$, hence by the structure of special points in the Brownian net at deterministic times (see Proposition 6.4), for each $z \in I(s, t)$ there exist $l \in \mathcal{W}^{1}$ and $r \in \mathcal{W}^{\mathrm{r}}$ such that $l \sim_{\text {in }}^{z} r$. Since paths in $\mathcal{N}$ are contained between left-most and right-most paths [SS08, Prop. 1.8], any path in $\mathcal{N}$ started between $l$ and $r$ must enter $z$. In particular (using [SS08, Prop. 3.6] to show that there are points that lie strictly between $l$ and $r$ ), it follows that there are paths in $\mathcal{W}^{\prime}$ that enter $z$.
Proof of Theorem 2.8 (a). Let $\left(\mathcal{W}_{0}, \mathcal{M}, \mathcal{W}\right)$ be a reference web, a set of marked points, and a sample web as constructed in Theorem 3.7. For $t \in \mathbb{R}$, let $I_{0}(t):=\left\{\pi(t): \pi \in \mathcal{W}_{0}, \sigma_{\pi}<t\right\}$. By Theorems 3.7 and 3.5, $\left(\mathcal{W}_{0}, \mathcal{W}\right)$ is a pair of sticky Brownian webs, while by Theorem 3.7, we have the representation

$$
\begin{equation*}
\rho_{t}=\int \rho_{0}(\mathrm{~d} x) \mathbb{P}\left[\pi_{(x, 0)}^{\uparrow}(t) \in \cdot \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right] . \tag{9.2}
\end{equation*}
$$

By Lemma 9.1. this implies that a.s. $\left\{\pi_{(x, 0)}^{\uparrow}(t)\right\}_{x \in \mathbb{R}} \subset I_{0}(t)$, which is a countable set. Therefore $\rho_{t}$ must be atomic with atoms in $I_{0}(t)$.

The proof of Theorem 2.8 (a) shows that at deterministic times, $\rho_{t}$ is concentrated on the countable set $I_{0}(t)$ of points where there is an incoming path of the reference Brownian web. The following lemma, which will be needed in the proof of Theorem 9.6 below, tells us which points in $I_{0}(t)$ carry positive mass.

Lemma 9.2 (Position of atoms at deterministic times) Let $\rho_{0}$ be a deterministic finite measure on $\mathbb{R}$, let $\left(\rho_{t}\right)_{t \geq 0}$ be a Howitt-Warren process constructed as in 9.2) and let $I_{0}(t):=$ $\left\{\pi(t): \pi \in \mathcal{W}_{0}, \sigma_{\pi}<t\right\}$. Then for each deterministic $t>0$, one has

$$
\begin{equation*}
\left\{x \in \mathbb{R}: \rho_{t}(\{x\})>0\right\}=I_{0}(t) \cap \operatorname{supp}\left(\rho_{t}\right) \quad \text { a.s. } \tag{9.3}
\end{equation*}
$$

Proof. The inclusion $\subset$ follows from our proof of Theorem 2.8 (a). If both the left and right speeds of the Howitt-Warren flow are finite, then by Proposition 2.6 (d) and Theorem 2.7 (a), $\operatorname{supp}\left(\rho_{t}\right)$ consists of isolated points, hence $\left\{x \in \mathbb{R}: \rho_{t}(\{x\})>0\right\}=\operatorname{supp}\left(\rho_{t}\right)$, which by the inclusion $\subset$ implies that we have in fact equality in (9.3).

If at least one of the speeds of the Howitt-Warren flow is infinite, then by Theorem 2.7 (b) and $(\mathrm{c})$, either $\operatorname{supp}\left(\rho_{t}\right)=\mathbb{R}$ or $\operatorname{supp}\left(\rho_{t}\right)$ is a halfline. To prove the inclusion $\supset$ in (9.3) in this case, let $x$ be some point in $I_{0}(t) \cap \operatorname{supp}\left(\rho_{t}\right)$. We will treat the cases that $x$ lies in the interior or on the boundary of $\operatorname{supp}\left(\rho_{t}\right)$ separately.

Assume for the moment that $x$ lies in the interior of $\operatorname{supp}\left(\rho_{t}\right)$, and let $\left(\mathcal{W}_{0}, \mathcal{M}, \mathcal{W}\right)$ be the reference web, set of marked points, and sample web used for constructing $\rho$. By Lemma 9.1, there exists a.s. an incoming path $\pi \in \mathcal{W}$ at $x$, so by the structure of special points of the Brownian web (Proposition 3.3), there exist $\hat{\pi}_{-}, \hat{\pi}_{+} \in \hat{\mathcal{W}}(x, t)$ such that $\tau:=\sup \{s<t$ : $\left.\hat{\pi}_{-}(s)=\hat{\pi}_{+}(s)\right\}<t$ and $\hat{\pi}_{-}<\hat{\pi}_{+}$on $(\tau, t)$. Choose deterministic $t_{n} \uparrow t$. Then

$$
\begin{equation*}
\lim _{n \rightarrow \infty} \mathbb{P}\left[\tau<t_{n} \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right]=1 \quad \text { a.s. } \tag{9.4}
\end{equation*}
$$

so by the fact that $x$ lies in the interior of $\operatorname{supp}\left(\rho_{t}\right)$, we can choose $n$ sufficiently large, depending on $\mathcal{W}_{0}$ and $\mathcal{M}$ but not on $\mathcal{W}$, such that

$$
\begin{equation*}
\mathbb{P}\left[\tau<t_{n},\left(\hat{\pi}_{-}\left(t_{n}\right), \hat{\pi}_{+}\left(t_{n}\right)\right) \subset \operatorname{supp}\left(\rho_{t_{n}}\right) \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right]>0 \tag{9.5}
\end{equation*}
$$

Conditional on $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, the restrictions of the paths $\hat{\pi}_{-}, \hat{\pi}_{+}$to the time interval $\left[t_{n}, t\right]$ are independent of $\rho_{t_{n}}$, and, since every path in $\mathcal{W}$ started from $\left(\hat{\pi}_{-}\left(t_{n}\right), \hat{\pi}_{+}\left(t_{n}\right)\right) \times\left\{t_{n}\right\}$ must pass through $x$,

$$
\begin{equation*}
\rho_{t}(\{x\}) \geq \int \mathbb{E}\left[\rho_{t_{n}}\left(\left(\hat{\pi}_{-}\left(t_{n}\right), \hat{\pi}_{+}\left(t_{n}\right)\right)\right) \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right] \tag{9.6}
\end{equation*}
$$

Applying (9.6) with $n$ large enough such that (9.5) holds shows that $\rho_{t}(\{x\})>0$.
If $x$ lies on the boundary of $\operatorname{supp}\left(\rho_{t}\right)$, then by symmetry, it suffices to consider the case that $\beta_{-}=-\infty, \beta_{+}<\infty$ and $\operatorname{supp}\left(\rho_{t}\right)=(-\infty, x]$. In this case, as in the proof of Lemma 8.1, we can without loss of generality assume that $\left(\mathcal{W}_{0}, \mathcal{M}, \mathcal{W}\right)$ are constructed as in Theorem 3.7 with $\nu_{1}=0$. In this case, $\mathcal{W}_{0}$ is the right Brownian web associated with the Brownian halfnet $\mathcal{N}_{\infty}$ from Lemma 8.1, hence by Proposition 8.9, $\operatorname{supp}\left(\rho_{t}\right)=\left(-\infty, \pi^{0}(t)\right](t \geq 0)$ where $\pi^{0}$ is the a.s. unique path in $\mathcal{W}_{0}$ starting at time zero from $\sup \left(\operatorname{supp}\left(\rho_{0}\right)\right)$. Define dual paths $\hat{\pi}_{-}, \hat{\pi}_{+}$and their first meeting time $\tau$ as before. Since $\left(\mathcal{W}, \mathcal{W}_{0}\right)$ is a left-right Brownian web, by Proposition 6.4, any path $\pi \in \mathcal{W}$ entering $x$ must satisfy $\pi(s)=\pi^{0}(s)$ for some $\sigma_{\pi}<s<t$, and hence, since paths in $\mathcal{W}$ cannot cross paths in $\mathcal{W}_{0}$ from left to right, $\tau^{\prime}:=\sup \left\{s<t: \pi^{0}(s)<\hat{\pi}_{-}(s)\right\}<t$ a.s. Then, for deterministic $t_{n} \uparrow t$,

$$
\begin{equation*}
\lim _{n \rightarrow \infty} \mathbb{P}\left[\tau \vee \tau^{\prime}<t_{n} \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right]=1 \quad \text { a.s. } \tag{9.7}
\end{equation*}
$$

so we can choose $n$ large enough such that this conditional probability is positive. Since $\left(\mathcal{W}, \mathcal{W}_{0}\right)$ is a left-right Brownian web, by Proposition 6.4, $\hat{\pi}_{-}\left(t_{n}\right)<\pi^{0}\left(t_{n}\right)$ a.s. on the event $\tau \vee \tau^{\prime}<t_{n}$. Now the argument proceeds as in the case when $x$ lies in the interior of $\operatorname{supp}\left(\rho_{t}\right)$, with the interval $\left(\hat{\pi}_{-}\left(t_{n}\right), \hat{\pi}_{+}\left(t_{n}\right)\right)$ replaced by $\left(\hat{\pi}_{-}\left(t_{n}\right), \pi^{0}\left(t_{n}\right)\right)$.

### 9.2 Non-atomicness at random times for non-erosion flows

The proof of Theorem 2.8 (b) is a bit involved. When $\nu(0,1)>0$ and $\beta_{+}-\beta_{-}<\infty$, the Howitt-Warren process $\left(\rho_{t}\right)_{t \geq 0}$ can be constructed from a Brownain net as in Theorem4.7. In particular, at each separation point of $\mathcal{N}$, mass splits binarily. We can then use the fact that separation points are dense in space and time to split each atom in $\rho_{0}$ infinitely often to reach a random time $t$ when $\rho_{t}$ contains no atoms. This is closely related to, and in fact implies, the fact that the branching-coalescing point set $\xi_{t}$, whose properties are listed in Prop. 2.6, admit random times when $\xi_{t}$ has no isolated points. When $\beta_{+}-\beta_{-}=\infty$, it turns out that as long as $\nu(0,1)>0$, the picture of binary splitting of mass is still valid, although more work will be required. To avoid repetition, we will prove Theorem 2.8 (b) directly under the general assumption $\nu(0,1)>0$, which we assume from now on. Who wants to see a similar, but simpler proof should read the proof that there exist random times when $\xi_{t}$ has no isolated points in [SSS09, Prop. 3.14].

Let $\left(\mathcal{W}_{0}, \mathcal{M}, \mathcal{W}\right)$ be a reference web, a set of marked points, and a sample web with quenched law $\mathbb{Q}$ as in Theorem 3.7, where we take $\nu_{1}=\nu_{\mathrm{r}}=\nu$. We will use the version of Howitt-Warren flow $\left(K_{s, t}^{\uparrow}\right)_{s<t}$ in Theorem 3.7 to represent the Howitt-Warren process

$$
\begin{equation*}
\rho_{t}=\int \rho_{s}(\mathrm{~d} x) K_{s, t}^{\uparrow}(x, \cdot)=\int \rho_{s}(\mathrm{~d} x) \mathbb{P}\left[\pi_{(x, s)}^{\uparrow}(t) \in \cdot \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right]=\int \rho_{s}(\mathrm{~d} x) \mathbb{Q}\left[\pi_{(x, s)}^{\uparrow}(t) \in \cdot\right] \tag{9.8}
\end{equation*}
$$

where $K^{\uparrow}$ satisfies Prop. 2.3 (i)', which implies that (9.8) holds a.s. for all $s<t$.
The heuristics outlined before require two ingredients. First we need to establish the existence of bottlenecks, i.e., points in space-time where most of the mass in $\rho_{0}$ must enter.

Secondly, at such bottlenecks, mass is split binarily. When the flow can be embedded in a Brownian net, separation points provide the bottlenecks. In general, choose $\alpha \in(0,1)$ such that $\nu([\alpha, 1-\alpha])>0$. Then

$$
\begin{equation*}
A_{1}^{\alpha}:=\left\{z:\left(z, \omega_{z}\right) \in \mathcal{M}, \omega_{z} \in[\alpha, 1-\alpha], z \text { is of type }(1,2)_{1} \text { in } \mathcal{W}_{0}\right\} \tag{9.9}
\end{equation*}
$$

will provide the bottlenecks we need, and we will show that an atom of size 1 entering any $z \in A_{1}^{\alpha}$ will be split into atoms of size no larger than $1-\alpha / 4$.

Lemma 9.3 (Bottlenecks) Let $\rho_{0}=\delta_{0}$, the delta mass at 0 . Then a.s. w.r.t. $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, for any $\delta>0$, we can find a point $z=(x, u) \in A_{1}^{\alpha}$ with $u \in(0, \delta)$ such that $\rho_{u}(\{x\}) \geq 1-\delta$.

Proof. By Theorems 3.7 and $3.5,\left(\mathcal{W}_{0}, \mathcal{W}\right)$ is a pair of sticky Brownian webs with drift $\beta$ and coupling parameter $\kappa=2 \nu([0,1])$. In particular, by Propositions 7.2 and 7.3 , if $\pi^{0}$ resp. $\pi$ denotes the a.s. unique path in $\mathcal{W}_{0}$ resp. $\mathcal{W}$ starting from the origin, then $\left(\pi^{0}, \pi\right)$ solves the Howitt-Warren martingale problem, which is equivalent to (2.6) and (2.9). Without loss of generality, we will assume $\beta=0$.

Fix any $\varepsilon>0$. For each $0 \leq s \leq t$, consider the following subset of $(1,2)_{1}$ points along $\pi^{0}$ :

$$
I_{[s, t]}^{\varepsilon}:=\left\{(y, v): s \leq v \leq t, y=\pi^{0}(v)=\hat{\pi}_{\left(\pi^{0}(u)+\varepsilon, u\right)}^{0}(v) \text { for some } u \in(v, t]\right\}
$$

where $\hat{\pi}_{\left(\pi^{0}(u)+\varepsilon, u\right)}^{0}$ is any path in the dual reference web $\hat{\mathcal{W}}_{0}$ starting from $\left(\pi^{0}(u)+\varepsilon, u\right)$. It was shown in [NRS10, Lemma 7.2] that $I_{[s, t]}^{\varepsilon}$ contains exactly the points of intersection on the time interval $[s, t]$ between $\pi^{0}$ and a backward Brownian motion $\hat{\pi}^{0}$ starting at $\left(\pi^{0}(t)+\varepsilon, t\right)$, which is Skorohod reflected between $\left(\pi^{0}(r)\right)_{0 \leq r \leq t}$ and $\left(\pi^{0}(r)+\varepsilon\right)_{0 \leq r \leq t}$. In particular, after reversing time and centering, the triple

$$
\left(\vec{Z}_{r}\right)_{0 \leq r \leq t}:=\left(\pi^{0}(t-r)-\pi^{0}(t), \hat{\pi}^{0}(t-r)-\pi^{0}(t), \ell_{1}\left(I_{[t-r, t]}^{\varepsilon}\right)\right)_{0 \leq r \leq t}
$$

is a strong Markov process, where $\ell_{1}\left(I_{[t-r, t]}^{\varepsilon}\right)$ is the intersection local time measure of the set $I_{[t-r, t]}^{\varepsilon}$ and is a finite continuous increasing process. Also note that $\left(\pi^{0}(t-r)-\pi^{0}(r)\right)_{0 \leq r \leq t}$ is distributed as a standard Brownian motion.

By the definition of $A_{1}^{\alpha}$ and our construction of $\mathcal{M}$ in Theorem 3.7, conditional on $\mathcal{W}_{0}$, $A_{1}^{\alpha} \cap I_{[0, t]}^{\varepsilon}$ is a Poisson subset of $I_{[0, t]}^{\varepsilon}$ with intensity measure $21_{\left\{z \in I_{[0, t]}^{\varepsilon}\right\}}^{\}} \ell_{1}(\mathrm{~d} z) \int_{[\alpha, 1-\alpha]} q^{-1} \nu(\mathrm{~d} q)$. In particular, $A_{1}^{\alpha} \cap I_{[0, t]}^{\varepsilon}$ is a.s. a finite set. If $A_{1}^{\alpha} \cap I_{[0, t]}^{\varepsilon} \neq \emptyset$, then let $\tau^{\varepsilon, t}:=\sup \{r \in[0, t]$ : $\left.\left(\pi^{0}(r), r\right) \in A_{1}^{\alpha} \cap I_{[0, t]}^{\varepsilon}\right\}$, and denote the corresponding point in $A_{1}^{\alpha} \cap I_{[0, t]}^{\varepsilon}$ by $z^{\varepsilon, t}=\left(x^{\varepsilon, t}, \tau^{\varepsilon, t}\right)$. If $A_{1}^{\alpha} \cap I_{[0, t]}^{\varepsilon}=\emptyset$, then we set $\tau^{\varepsilon, t}=0$ and $z^{\varepsilon, t}=(0,0)$. Therefore

$$
\mathbb{P}\left[\tau^{\varepsilon, t}>0 \mid \mathcal{W}_{0}\right]=1-e^{-c_{\alpha, \nu} \ell_{1}\left(\left[I_{[0, t]}^{\varepsilon}\right)\right.}, \quad \text { where } c_{\alpha, \nu}:=2 \int_{[\alpha, 1-\alpha]} q^{-1} \nu(\mathrm{~d} q)
$$

Now observe that we can construct $\tau^{\varepsilon, t}$ by setting $t-\tau^{\varepsilon, t}:=t \wedge \inf \left\{r \geq 0: c_{\alpha, \nu} \ell_{1}\left(I_{[t-r, t]}^{\varepsilon}\right) \geq L\right\}$, where $L$ is an independent mean 1 exponential random variable. In particular, conditional on $L, t-\tau^{\varepsilon, t}$ is a stopping time for the process $\vec{Z}$. Therefore, conditional on $\tau^{\varepsilon, t}$ and $\left(\vec{Z}_{s}\right)_{0 \leq s \leq t-\tau^{\varepsilon, t}}$, the law of $\left(\pi^{0}(s), \pi(s)\right)_{0 \leq s \leq \tau^{\varepsilon, t}}$ remains the same as before because conditional on $\tau^{\varepsilon, t},\left(\pi_{s}^{0}\right)_{0 \leq s \leq \tau^{\varepsilon, t}}$ is a Brownian motion independent of $\left(\vec{Z}_{s}\right)_{0 \leq s \leq t-\tau^{\varepsilon, t}}$, and $\pi$ is constructed
by switching among the truncated paths $\left\{(\gamma(s))_{s_{0} \leq s \leq \tau^{\varepsilon, t}}: \gamma \in \mathcal{W}_{0}, \sigma_{\gamma}=s_{0}<\tau^{\varepsilon, t}\right\}$ using independent Poisson processes, which are all independent of $\left(\vec{Z}_{s}\right)_{0 \leq s \leq t-\tau, t}$. Therefore

$$
\begin{aligned}
\mathbb{P}\left[\tau^{\varepsilon, t}>0, \pi^{0}\left(\tau^{\varepsilon, t}\right)=\pi\left(\tau^{\varepsilon, t}\right)\right] & =\mathbb{E}\left[1_{\left\{\tau^{\varepsilon, t}>0\right\}} \mathbb{P}\left[\pi^{0}\left(\tau^{\varepsilon, t}\right)=\pi\left(\tau^{\varepsilon, t}\right) \mid \tau^{\varepsilon, t},\left(\vec{Z}_{s}\right)_{0 \leq s \leq t-\tau^{\varepsilon, t}}\right]\right] \\
& =\mathbb{E}\left[1_{\left\{\tau^{\varepsilon, t}>0\right\}} \phi\left(\tau^{\varepsilon, t}\right)\right],
\end{aligned}
$$

where $\phi(r):=\mathbb{P}\left[\pi^{0}(r)=\pi(r)\right]$. Note that $\left|\pi^{0}(s)-\pi(s)\right|$ is distributed as a Brownian motion starting at 0 and sticky reflected at 0 . It was shown in [SSS09, Lemma A.2] that $\phi(r)=$ $\mathbb{P}\left[\left|\pi^{0}(r)-\pi(r)\right|=0\right]$ is monotone in $r$ and increases to 1 as $r \downarrow 0$. Therefore

$$
\begin{equation*}
\mathbb{P}\left[\tau^{\varepsilon, t}>0, \pi^{0}\left(\tau^{\varepsilon, t}\right)=\pi\left(\tau^{\varepsilon, t}\right)\right] \geq \phi(t)\left(1-\mathbb{E}\left[e^{-c_{\alpha, \nu} \ell_{1}\left(I_{[0, t]}^{\varepsilon}\right]}\right]\right) \tag{9.10}
\end{equation*}
$$

Note that almost surely,

$$
\ell_{1}\left(I_{[0, t]}^{\varepsilon}\right) \underset{\varepsilon \rightarrow 0}{\longrightarrow} \ell_{1}\left(\left\{\left(\pi^{0}(s), s\right): s \in[0, t]\right\}\right),
$$

where the right-hand side is easily seen to be a.s. infinite using the fact that $\ell_{1}\left(I_{[0, t]}^{\varepsilon}\right) \geq$ $\sum_{i=1}^{\left\lfloor t \varepsilon^{-2}\right\rfloor} \ell_{1}\left(I_{\left[(i-1) \varepsilon^{2}, i \varepsilon^{2}\right]}^{\varepsilon}\right)$, and $\varepsilon^{-1} \ell_{1}\left(I_{\left[(i-1) \varepsilon^{2}, i \varepsilon^{2}\right]}^{\varepsilon}\right), i \in \mathbb{N}$, are i.i.d. with the same distribution as $\ell_{1}\left(I_{[0,1]}^{1}\right)$. Therefore 9.10 implies

$$
\lim _{t \downarrow 0} \lim _{\varepsilon \downarrow 0} \mathbb{P}\left[\tau^{\varepsilon, t}>0, \pi^{0}\left(\tau^{\varepsilon, t}\right)=\pi\left(\tau^{\varepsilon, t}\right)\right]=1 .
$$

Since

$$
\begin{aligned}
\mathbb{P}\left[\tau^{\varepsilon, t}>0, \pi^{0}\left(\tau^{\varepsilon, t}\right)=\pi\left(\tau^{\varepsilon, t}\right)\right] & =\mathbb{E}\left[1_{\left\{\tau^{\varepsilon, t}>0\right\}} \mathbb{P}\left[\pi^{0}\left(\tau^{\varepsilon, t}\right)=\pi\left(\tau^{\varepsilon, t}\right) \mid \tau^{\varepsilon, t}, \mathcal{W}_{0}, \mathcal{M}\right]\right] \\
& =\mathbb{E}\left[1_{\left\{\tau^{\varepsilon, t}>0\right\}} \rho_{\tau^{\varepsilon, t}}\left(\left\{\pi^{0}\left(\tau^{\varepsilon, t}\right)\right\}\right)\right],
\end{aligned}
$$

and $\rho_{\tau^{\varepsilon, t}}\left(\left\{\pi^{0}\left(\tau^{\varepsilon, t}\right)\right\}\right) \in[0,1]$, it follows that for any $\delta \in(0,1)$, we must have

$$
\lim _{t \downarrow 0} \lim _{\varepsilon \downarrow 0} \mathbb{P}\left[\tau^{\varepsilon, t}>0, \rho_{\tau^{\varepsilon, t}}\left(\left\{\pi^{0}\left(\tau^{\varepsilon, t}\right)\right\}\right) \geq 1-\delta\right]=1 .
$$

Therefore for any sequence $\left(\varepsilon_{n}, t_{n}\right) \rightarrow(0,0)$, the event $\left\{\tau^{\varepsilon_{n}, t_{n}}>0, \rho_{\tau^{\varepsilon_{n}, t_{n}}}\left(\left\{\pi^{0}\left(\tau^{\varepsilon_{n}, t_{n}}\right)\right\}\right) \geq\right.$ $1-\delta\}$ must occur a.s. infinitely often. The lemma then follows.

Next we show that mass entering any $z \in A_{1}^{\alpha}$ must be split into smaller atoms.
Lemma 9.4 (Splitting of mass) Almost surely w.r.t. $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, for each $z=(x, u) \in A_{1}^{\alpha}$ and for any $\varepsilon, \delta>0$, we can find $h>0$ such that if $\left(\rho_{t}\right)_{t \geq u}$ is defined as in 9.8) with $\rho_{u}=\delta_{x}$, the delta mass at $x$, then for all $t \in(u, u+h)$, we have $\rho_{t}([x-\varepsilon, x+\varepsilon]) \geq 1-\delta$ and $\left|\rho_{t}\right|_{\infty}:=\sup _{y \in \mathbb{R}} \rho_{t}(\{y\}) \leq 1-\alpha / 4$.

Proof. We shall embed the reference and sample webs $\mathcal{W}_{0}$ and $\mathcal{W}$ in a Brownian net whose separation points contain all of $A_{1}^{\alpha}$. Recall from Theorem 3.7 that $\mathcal{W}$ is constructed from $\left(\mathcal{W}_{0}, \mathcal{M}\right)$ by switching paths in $\mathcal{W}_{0}$ at a countable collection of $(1,2)$ points $A \cup B$, where $A$ is a random subset of the marked points in $\mathcal{M}$, with $z$ included in $A$ independently for each $\left(z, \omega_{z}\right) \in \mathcal{M}$ with probability $\omega_{z}$ if $z$ is of type $(1,2)_{1}$ in $\mathcal{W}_{0}$, and with probability $1-\omega_{z}$ if $z$ is of type $(1,2)_{\mathrm{r}}$; while $B$ is an independent Poisson point set on (1,2) points with intensity
measure $2 \nu(\{0\}) \ell_{1}+2 \nu(\{1\}) \ell_{\mathrm{r}}$. Note that conditional on $\mathcal{W}_{0}, A \cup B \cup A_{1}^{\alpha}$ is a Poisson point process on $(1,2)$ points of $\mathcal{W}_{0}$ with intensity measure

$$
2 \nu([0,1]) \ell_{\mathrm{r}}(\mathrm{~d} z)+2\left(\nu([0, \alpha))+\int_{[\alpha, 1-\alpha]} q^{-1} \nu(\mathrm{~d} q)+\nu((1-\alpha, 1])\right) \ell_{1}(\mathrm{~d} z)
$$

By Theorem 4.6, if we allow hopping at all points in $A \cup B \cup A_{1}^{\alpha}$, then we obtain a Brownian net $\mathcal{N}$ with left and right speeds

$$
\beta_{-}=\beta-2 \nu([0,1]) \quad \text { and } \quad \beta_{+}=\beta+2\left(\nu([0, \alpha))+\int_{[\alpha, 1-\alpha]} q^{-1} \nu(\mathrm{~d} q)+\nu((1-\alpha, 1])\right) .
$$

Clearly $\mathcal{W}_{0}$ and $\mathcal{W}$ are subsets of $\mathcal{N}$, and all points in $A_{1}^{\alpha}$ are separation points of $\mathcal{N}$. In particular, a.s. each point in $A_{1}^{\alpha}$ is of type $(1,2)$ in $\mathcal{W}$. Given $z=(x, u) \in A_{1}^{\alpha}$, let $\pi_{z}^{-}$resp. $\pi_{z}^{+}$ denote the left resp. right of the two outgoing paths in $\mathcal{W}$ at $z$, and recall from after Prop. 3.3 that $\pi_{z}^{\uparrow}$ picks from $\pi_{z}^{+}$and $\pi_{z}^{-}$the natural continuation of any paths in $\mathcal{W}$ entering $z$. By 9.8), given $\rho_{u}=\delta_{x}$, we have $\rho_{t}(\cdot)=\mathbb{P}\left[\pi_{z}^{\uparrow}(t) \in \cdot \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right]$ for all $t>u$. Since $\pi_{z}^{\uparrow}$ is a.s. a continuous path starting at $z$, for any $\varepsilon, \delta>0$, we can choose $h>0$ sufficiently small such that $\rho_{t}([x-\varepsilon, x+\varepsilon])=\mathbb{P}\left[\pi_{z}^{\uparrow}(t) \in[x-\varepsilon, x+\varepsilon] \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right] \geq 1-\delta$ for all $t \in[u, u+h]$. This establishes the first part of the lemma.

Since $z \in A_{1}^{\alpha}$, we have $\left(z, \omega_{z}\right) \in \mathcal{M}$ for some $\omega_{z} \in[\alpha, 1-\alpha]$. By our construction of $\mathcal{W}$ in Theorem 3.7 and the natural coupling between sticky Brownian webs and Brownian nets given in Lemma 6.16, we note that conditional on $\left(\mathcal{W}_{0}, \mathcal{M}\right)$ and $A \cup B \backslash\{z\}, \mathcal{W}$ is uniquely determined except for the orientation of paths in $\mathcal{W}$ entering $z$, which is then resolved by an independent random variable $\alpha_{z}$ with $\mathbb{P}\left[\alpha_{z}=1\right]=\omega_{z}$ and $\mathbb{P}\left[\alpha_{z}=-1\right]=1-\omega_{z}$. More precisely, we set $\operatorname{sign}_{\mathcal{W}}(z)=\alpha_{z}$, and $\pi_{z}^{\uparrow}=\pi_{z}^{+}$when $\alpha_{z}=1$ and $\pi_{z}^{\uparrow}=\pi_{z}^{-}$when $\alpha_{z}=-1$. Therefore with $\rho_{u}=\delta_{x}$, we have for all $t>u$,

$$
\begin{equation*}
\rho_{t}(\cdot)=\omega_{z} \mathbb{P}\left[\pi_{z}^{+}(t) \in \cdot \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right]+\left(1-\omega_{z}\right) \mathbb{P}\left[\pi_{z}^{-}(t) \in \cdot \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right] . \tag{9.11}
\end{equation*}
$$

Almost surely, $\pi_{z}^{-}<\pi_{z}^{+}$on $(u, u+h)$ for some $h>0$. Therefore for $h>0$ sufficiently small,

$$
\begin{equation*}
\mathbb{P}\left(\pi_{z}^{-}(t)<\pi_{z}^{+}(t) \text { for all } t \in(u, u+h) \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right)>\frac{1}{2} \tag{9.12}
\end{equation*}
$$

 $\rho_{t}(\{y\})>1-\alpha / 4$ for some $t \in(u, u+h)$ and $y \in \mathbb{R}$, then by 9.11), we must have

$$
\mathbb{P}\left[\pi_{z}^{+}(t) \neq y \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right] \leq \frac{\alpha}{4 \omega_{z}} \leq \frac{1}{4} \quad \text { and } \quad \mathbb{P}\left[\pi_{z}^{-}(t) \neq y \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right] \leq \frac{\alpha}{4\left(1-\omega_{z}\right)} \leq \frac{1}{4}
$$

since $\omega_{z} \in[\alpha, 1-\alpha]$. This implies that

$$
\mathbb{P}\left(\pi_{z}^{-}(t)<\pi_{z}^{+}(t) \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right) \leq \mathbb{P}\left[\pi_{z}^{+}(t) \neq y \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right]+\mathbb{P}\left[\pi_{z}^{-}(t) \neq y \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right] \leq \frac{1}{2}
$$

contradicting (9.12). This completes the proof of Lemma 9.4 .
Lemmas 9.3 and 9.4 immediately imply the following.
Lemma 9.5 Let $\rho_{0}=\delta_{x}$ for some $x \in \mathbb{R}$, and let $\alpha \in(0,1)$ satisfy $\nu([\alpha, 1-\alpha])>0$. Then for any $\varepsilon, \delta, h>0$, a.s. there exists $(u, v) \subset(0, h)$ such that $\rho_{t}([x-\varepsilon, x+\varepsilon]) \geq 1-\delta$ for all $t \in[0, v]$ and $\left|\rho_{t}\right|_{\infty}:=\sup _{y \in \mathbb{R}} \rho_{t}(\{y\}) \leq 1-\alpha / 5$ for all $t \in[u, v]$.

Now we can prove the existence of random times when $\rho_{t}$ is non-atomic, if $\nu(0,1)>0$.
Theorem 2.8 (b). It suffices to show that for each interval $(u, v)$ with $0<u<v \in \mathbb{Q}$, a.s. there exists $t \in(u, v)$ such that $\rho_{t}$ has no atoms. Denote $\lambda:=\sup _{y \in \mathbb{R}} \rho_{u}(\{y\})$, and let $\rho_{t}$ be defined using $\left(K_{s, t}^{\uparrow}\right)_{s<t}$ as in 9.8 .

By Theorem 2.8 (a), which we have already established, $\rho_{u}$ is a.s. atomic. Since $\rho_{0}$ is assumed to be a finite measure (for infinite $\rho_{0}$, see the remark after Theorem 2.9), for any $\varepsilon>0$, we can find a finite set of atoms of $\rho_{u}$ at $\left\{x_{1}, \cdots, x_{k}\right\}$ with $\lambda=\rho_{u}\left(\left\{x_{1}\right\}\right) \geq \rho_{u}\left(\left\{x_{2}\right\}\right) \geq \cdots$ such that $\rho_{u}\left(\mathbb{R} \backslash\left\{x_{1}, \cdots, x_{k}\right\}\right) \leq \varepsilon \lambda$. For $t \geq u$, let

$$
\rho_{t}^{\mathrm{a}}:=\int_{x \notin\left\{x_{1}, \cdots, x_{k}\right\}} K_{u, t}^{\uparrow}(x, \cdot) \rho_{u}(\mathrm{~d} x), \quad \text { and } \quad \rho_{t}^{(i)}:=\rho_{u}\left(\left\{x_{i}\right\}\right) K_{u, t}^{\uparrow}\left(x_{i}, \cdot\right) \quad \text { for } 1 \leq i \leq k .
$$

We can find $\delta_{1}, \cdots, \delta_{k}>0$ such that

$$
\left[x_{i}-\delta_{i}, x_{i}+\delta_{i}\right] \cap\left[x_{j}-\delta_{j}, x_{j}+\delta_{j}\right]=\emptyset \quad \text { for all } i \neq j .
$$

By Theorem 2.9, for each $1 \leq i \leq k,\left(\rho_{t}^{(i)}\right)_{t \geq u}$ has a.s. continuous sample path in $\mathcal{M}_{\mathrm{g}}(\mathbb{R})$. Therefore we can choose $h_{1}>0$ sufficiently small such that

$$
\sum_{i=1}^{k} \rho_{t}^{(i)}\left(\mathbb{R} \backslash\left[x_{i}-\delta_{i}, x_{i}+\delta_{i}\right]\right) \leq \varepsilon \lambda \quad \text { for all } t \in\left[u, u+h_{1}\right],
$$

and by Lemma 9.5. we can choose $\left(u_{1}, v_{1}\right) \subset\left(u, u+h_{1}\right)$ with $u_{1}, v_{1} \in \mathbb{Q}$ such that

$$
\sup _{y \in \mathbb{R}} \rho_{t}^{(1)}(\{y\}) \leq(1-\alpha / 5) \lambda \quad \text { for all } t \in\left[u_{1}, v_{1}\right]
$$

Our construction guarantees that for $t \in\left[u_{1}, v_{1}\right]$, the atom of $\rho_{u}$ at $x_{1}$ is split into atoms of size no larger than $(1-\alpha / 5+2 \varepsilon) \lambda$, and the remaining atoms of $\rho_{u}$ at $\left\{x_{2}, \cdots, x_{k}\right\}$ can only gain altogether a mass of at most $2 \varepsilon \lambda$, and there is no merging and formation of new atoms with size larger than $2 \varepsilon \lambda$. By choosing $\varepsilon>0$ sufficiently small and repeating a finite number of times (say $m$ times) the same argument, we can split the remaining atoms of $\rho_{u}$ at $\left\{x_{1}, \cdots, x_{k}\right\}$ and find $u_{m}, v_{m} \in \mathbb{Q}$ such that $\sup _{y \in \mathbb{R}} \rho_{t}(\{y\}) \leq(1-\alpha / 6) \lambda$ for all $t \in\left[u_{m}, v_{m}\right]$. By repeating the whole argument above, for each $n \in \mathbb{N}$, we can inductively find $u^{(n)}, v^{(n)} \in \mathbb{Q}$ with $u<u^{(n-1)}<u^{(n)}<v^{(n)}<v^{(n-1)}<v$, such that $\sup _{y \in \mathbb{R}} \rho_{t}(\{y\}) \leq 1 / n$ for all $t \in$ $\left[u^{(n)}, v^{(n)}\right]$. Any $t \in \bigcap_{n \in \mathbb{N}}\left[u^{(n)}, v^{(n)}\right] \neq \emptyset$ then gives a time when $\rho_{t}$ contains no atoms. In fact, $\bigcap_{n \in \mathbb{N}}\left[u^{(n)}, v^{(n)}\right]$ contains a single point since $\rho_{t}$ is a.s. atomic at deterministic times by Theorem 2.8 (a).

### 9.3 Atomicness at all times for erosion flows

In this subsection, we prove Theorem 2.8 (c). In fact, we will prove the following stronger result. Note that below, when we apply Theorem 3.7, we deviate from the canonical choice $\nu_{1}=\nu_{\mathrm{r}}=\nu$.

Theorem 9.6 (Atomicness of erosion flows) Let $\rho=\left(\rho_{t}\right)_{t \geq 0}$ be a Howitt-Warren process with drift $\beta \in \mathbb{R}$ and characteristic measure of the form $\nu=c_{0} \delta_{0}+c_{1} \delta_{1}$, with $c_{0}, c_{1} \geq 0$, started in some deterministic, finite nonzero measure $\rho_{0}$ on $\mathbb{R}$. Let $\rho$ be constructed as $\rho_{t}=$
$\int \rho_{0}(\mathrm{~d} x) K_{0, t}^{\uparrow}(x, \cdot)$, where $\left(K_{s, t}^{\uparrow}\right)_{s \leq t}$ is the Howitt-Warren flow constructed as in Theorem 3. $\uparrow$ using a reference Brownian web $\mathcal{W}_{0}$ with drift $\beta_{0}=\beta-2 c_{0}+2 c_{1}, \nu_{1}=c_{0} \delta_{0}$ and $\nu_{\mathrm{r}}=c_{1} \delta_{1}$. Then a.s., $\rho_{t}$ is purely atomic at each $t>0$ and

$$
\begin{align*}
& \left\{(x, t) \in \mathbb{R}^{2}: t>0, \rho_{t}(\{x\})>0\right\}  \tag{9.13}\\
& \quad=\left\{(x, t) \in \mathbb{R}^{2}: t>0, x \in \operatorname{supp}\left(\rho_{t}\right), \exists \pi \in \mathcal{W}_{0} \text { s.t. } \sigma_{\pi}<t \text { and } \pi(t)=x\right\}
\end{align*}
$$

Remark. Note that if $c_{0}$ and $c_{1}$ are both strictly positive, then by Theorem $2.7(\mathrm{c}), \operatorname{supp}\left(\rho_{t}\right)=$ $\mathbb{R}$, and hence the right-hand side of $(9.13$ is just the set

$$
\begin{equation*}
\left\{(\pi(t), t): t>0, \pi \in \mathcal{W}_{0}, \sigma_{\pi}<t\right\} \tag{9.14}
\end{equation*}
$$

We state as an open problem whether in this case, $\left.\mathcal{W}_{0}\right|_{0} ^{\infty}$, the restriction of $\mathcal{W}_{0}$ to the time interval $[0, \infty]$, can a.s. be uniquely reconstructed from $\left(\rho_{t}\right)_{t \geq 0}$. In fact, it seems likely that $\left.\mathcal{W}_{0}\right|_{0} ^{\infty}$ consists of all paths $\pi \in \Pi$ starting at $\sigma_{\pi} \geq 0$ such that $\rho_{t}(\{\pi(t)\})$ is locally uniformly bounded away from zero on $\left(\sigma_{\pi}, \infty\right)$. Note that $\left.\mathcal{W}_{0}\right|_{0} ^{\infty}$ cannot be reconstructed from the set in (9.14), since switching the orientation of finitely many points of type $(1,2)$ does not change this set.

Recall that for an erosion flow, $\mathbb{Q}_{z}^{\uparrow}:=\mathbb{P}\left[\pi_{z}^{\uparrow} \in \cdot \mid \mathcal{W}_{0}\right]$ as defined in 4.18 is the quenched law of a Markov process in a random environment. Theorem 9.6 says that conditional on $\mathcal{W}_{0}$, this Markov process has the property that at all fixed times $t$ (that may depend on $\mathcal{W}_{0}$ but not on $\mathcal{W}$ ), the motion is located in the countable set $I_{0}(t)$ of points where there is an incoming path from $\mathcal{W}_{0}$. This type of behavior is reminiscent of the FIN diffusion defined in [FIN02], which is concentrated on a random countable set at each deterministic time.

Our proof of Theorem 9.6 is based on the following lemma, which controls the speed of mass loss of an erosion flow along a path $\pi^{0} \in \mathcal{W}_{0}$. The proof of this lemma is somewhat involved. The intuitive idea behind it is that, for erosion flows, mass must dissipate continuously, which is contrary to the case $\nu(0,1)>0$ where mass undergoes binary splitting. However, a crude estimate on the loss of mass from $\pi^{0}$ will show that all mass is lost instantly. Indeed, we conjecture that $\left(\rho_{t}\left(\pi^{0}(t)\right)\right)_{t \geq 0}$ as a function of time has locally unbounded variation, which means that in each positive time interval an infinite amount of mass leaves and rejoins $\pi^{0}$. Nevertheless, formula 9.15 below shows that the decrease of this process is, in a sense, Hölder continuous for any exponent $\gamma<1 / 2$.

Lemma 9.7 (Atomic mass along a reference Brownian web path) In the set-up of Theorem 9.6, assume that $\rho_{0}=\delta_{0}$ and let $\pi^{0}$ be the a.s. unique path starting from the origin in the reference web $\mathcal{W}_{0}$. Then a.s. with respect to $\mathcal{W}_{0}$, for each $\gamma \in(0,1 / 2)$, there exists a constant $0<C_{\gamma, \mathcal{W}_{0}}<\infty$ depending on $\gamma$ and $\mathcal{W}_{0}$, such that

$$
\begin{equation*}
\rho_{0}\left(\pi^{0}(0)\right)-\rho_{t}\left(\pi^{0}(t)\right) \leq C_{\gamma, \mathcal{W}_{0}} t^{\gamma} \quad \text { for all } 0 \leq t \leq 1 \tag{9.15}
\end{equation*}
$$

Moreover, for any $\delta>0$, there exists a deterministic constant $0<C_{\delta, \gamma}<\infty$, such that the random constant $C_{\gamma, \mathcal{W}_{0}}$ satisfies

$$
\begin{equation*}
\mathbb{P}\left[C_{\gamma, \mathcal{W}_{0}}>u\right] \leq \frac{C_{\delta, \gamma}}{u^{\delta}} \wedge 1 \quad \text { for all } u>0 \tag{9.16}
\end{equation*}
$$

Proof. To prove (9.15), we will apply a one-sided version of Kolmogorov's moment criterion, Theorem E.1, to the process $X_{t}:=\rho_{t}\left(\pi^{0}(t)\right)$. Note that we cannot expect the gain of mass at $\pi^{0}$ to be continuous due to the merging of atoms. Therefore the standard version of Kolmogorov's moment criterion is not applicable.

Let $\pi$ denote the a.s. unique path starting from the origin in the sample web $\mathcal{W}$. Then by construction,

$$
X_{s}=\mathbb{Q}\left[\pi(s)=\pi^{0}(s)\right]=\mathbb{P}\left[\pi(s)=\pi^{0}(s) \mid \mathcal{W}_{0}\right] \quad \text { a.s. for all } s \geq 0
$$

For any $0 \leq s \leq t$,

$$
\begin{aligned}
X_{s}-X_{t} & =\mathbb{P}\left[\pi(s)=\pi^{0}(s) \mid \mathcal{W}_{0}\right]-\mathbb{P}\left[\pi(t)=\pi^{0}(t) \mid \mathcal{W}_{0}\right] \\
& \leq \mathbb{P}\left[\pi(s)=\pi^{0}(s), \pi(t) \neq \pi^{0}(t) \mid \mathcal{W}_{0}\right]
\end{aligned}
$$

Let $\left(\mathcal{W}_{i}\right)_{i \in \mathbb{N}}$ be i.i.d. copies of the sample web $\mathcal{W}$ conditional on $\mathcal{W}_{0}$, with $\mathcal{W}_{i}(0,0)=\left\{\pi^{i}\right\}$. Then for any $k \in \mathbb{N}$,

$$
\begin{equation*}
\mathbb{E}\left[\left(\left(X_{s}-X_{t}\right)^{+}\right)^{k}\right] \leq \mathbb{P}\left[\pi^{i}(s)=\pi^{0}(s), \pi^{i}(t) \neq \pi^{0}(t) \text { for all } 1 \leq i \leq k\right] \tag{9.17}
\end{equation*}
$$

Since $\left(\mathcal{W}^{i}\right)_{i \in \mathbb{N}}$ are constructed from $\mathcal{W}_{0}$ by independent Poisson marking and switching of $(1,2)$ points of $\mathcal{W}_{0}$ with intensity measure $2 c_{0} \ell_{1}+2 c_{1} \ell_{\mathrm{r}}$, if we allow paths in $\mathcal{W}_{0}$ to hop at all such marked points, then by Theorem 4.6, we obtain a Brownian net $\mathcal{N}$ with left and right speeds $\beta_{-}=\beta_{0}-2 k c_{1}, \beta_{+}=\beta_{0}+2 k c_{0}$. By Proposition 4.5, $\mathcal{W}_{0}, \cdots, \mathcal{W}_{k}$ are all contained in $\mathcal{N}$. Since the Poisson marked $(1,2)$ points of $\mathcal{W}_{0}$ are a.s. distinct for different $\mathcal{W}_{i}$, at each separation point $z$ of $\mathcal{N}$, we must have $\operatorname{sign}_{z}\left(\mathcal{W}_{0}\right)=\operatorname{sign}_{z}\left(\mathcal{W}_{i}\right)$ for all but one $i \in\{1, \cdots, k\}$. Note that at this place, we make essential use of the fact that we have an erosion flow; in fact, this is the only place in the proof of Theorem 9.6 where we will use this. Let $N_{s, t}$ be the number of $s, t$-relevant separation points along $\pi^{0}$ on the time interval $(s, t)$. Then the event in the RHS of 9.17) can only occur if $N_{s, t} \geq k$.

We will next bound $\mathbb{P}\left[N_{s, t} \geq k\right]$ and show that

$$
\mathbb{E}\left[\left(\left(X_{s}-X_{t}\right)^{+}\right)^{k}\right] \leq \mathbb{P}\left[N_{s, t} \geq k\right] \leq C_{k}(t-s)^{\frac{k}{2}} \quad(0 \leq s \leq t \leq 1)
$$

for some $C_{k}$ depending only on $k$. We can then apply the one-sided Kolmogorov's moment criterion, Theorem E.1, to deduce (9.15).

Let $\left(\mathcal{W}^{1}, \mathcal{W}^{\mathrm{r}}\right)$ be the left-right Brownian web associated with $\mathcal{N}$, and let $\left(\hat{\mathcal{W}}^{1}, \hat{\mathcal{W}}^{\mathrm{r}}\right)$ be the corresponding dual left-right Brownian web. Note that $\left(\mathcal{W}^{\mathrm{l}}, \mathcal{W}_{0}\right)$ and $\left(\mathcal{W}_{0}, \mathcal{W}^{\mathrm{r}}\right)$ each form a left-right Brownian web. For any deterministic $0 \leq s<t \leq 1$, the $s, t$-relevant separation points along $\pi^{0}$ which are of type $(1,2)_{1}$ in $\mathcal{W}_{0}$ can be constructed by first following $\hat{r}_{1} \in \hat{\mathcal{W}}^{\mathrm{r}}\left(\pi^{0}(t), t\right)$, starting on the right of $\pi^{0}$, until the first time $\tau_{1}$ when $\hat{r}_{1}$ crosses $\pi^{0}$ from right to left. This gives the first $s, t$-relevant separation point $\left(\pi^{0}\left(\tau_{1}\right), \tau_{1}\right)$ of type $(1,2)_{1}$ along $\pi^{0}$. We then repeat the above procedure by following $\hat{r}_{2} \in \hat{\mathcal{W}}^{\mathrm{r}}\left(\pi^{0}\left(\tau_{1}\right), \tau_{1}\right)$ until the first time $\tau_{2}$ when $\hat{r}_{2}$ crosses $\pi^{0}$ from right to left. Iterating this procedure until time $s$ exhausts all $s, t$ relevant separation points of type $(1,2)_{1}$ along $\pi^{0}$, the total number of which will be denoted by $N_{s, t}^{1}$ and is a.s. finite by Proposition 6.1 (b). By [SSS09, Lemma 2.2] and its proof, if we define $\hat{r}=\hat{r}_{1}$ on $\left[\tau_{1}, t\right], \hat{r}=\hat{r}_{2}$ on $\left[\tau_{2}, \tau_{1}\right], \ldots$, then $\hat{r}$ is distributed as a Brownian motion $B^{1}$ with drift $-\beta_{+}$starting from $\left(\pi^{0}(t), t\right)$ running backward in time and Skorohod reflected to the right of $\pi^{0}$. More precisely, $(\hat{r}(t-v))_{v \geq 0}$ solves the Skorohod equation

$$
\begin{array}{lr}
\mathrm{d} \hat{r}(t-v)=\mathrm{d} B^{1}(v)+\mathrm{d} \Delta^{1}(v), & 0 \leq v \leq t-s, \\
\mathrm{~d} \hat{r}(t-v)=\mathrm{d} B^{1}(v), & t-s \leq v, \tag{9.18}
\end{array}
$$

where $\Delta^{1}(v)$ is an increasing process with $\int_{0}^{t-s} 1_{\left\{\hat{r}(t-v) \neq \pi^{0}(t-v)\right\}} \mathrm{d}^{1}(v)=0$, and $\hat{r}$ is subject to the constraint $\hat{r}(v) \geq \pi^{0}(v)$ for all $0 \leq v \leq t-s$. Furthermore, by [SSS09, Lemma 2.2], conditional on $\pi^{0}$, the $s, t$-relevant separation points of type $(0,1)_{1}$ along $\pi^{0}$ is distributed as a Poisson point process along $\pi^{0}$ with intensity measure $2 k c_{0} \mathrm{~d} \Delta^{1}(v)$ on the projected time interval $[s, t]$, where $2 k c_{0}$ is the difference between the drifts of $\mathcal{W}_{0}$ and $\mathcal{W}_{r}$, and its appearance in the intensity measure can be deduced from the fact that Brownian nets of different left-right speeds are related by changing the drift and performing diffusive rescaling. In particular, $N_{s, t}^{1}$ is distributed as a Poisson random variable with mean $2 k c_{0} \Delta^{1}(t-s)$. Therefore for any $k_{1} \geq 0$,

$$
\begin{equation*}
\mathbb{P}\left[N_{s, t}^{1} \geq k_{\mathrm{l}}\right]=\sum_{j=k_{1}}^{\infty} \frac{1}{j!} \mathbb{E}\left[e^{-2 k c_{0} \Delta^{1}(t-s)}\left(2 k c_{0} \Delta^{1}(t-s)\right)^{j}\right] \leq \sum_{j=k_{1}}^{\infty} \frac{1}{j!} \mathbb{E}\left[\left(2 k c_{0} \Delta^{1}(t-s)\right)^{j}\right] . \tag{9.19}
\end{equation*}
$$

The Skorohod equation (9.18) admits a pathwise unique solution (see e.g. KKS91, Sec. 3.6.C]) with $\Delta^{1}(t-s)=-\inf _{0 \leq v \leq t-s}\left(B^{\mathrm{l}}(v)-\pi^{0}(t-v)\right)$. By the independence of $B^{1}$ and $\pi^{0}, B^{\mathrm{l}}(v)-$ $\pi^{0}(t-v)$ is distributed as $-\sqrt{2} \tilde{B}^{1}(v)-2 k c_{0} v$ for a standard Brownian motion $\tilde{B}^{1}$ starting from 0 . Therefore

$$
\Delta^{1}(t-s)=\sqrt{t-s} \sup _{0 \leq v \leq 1}\left(\sqrt{2} \tilde{B}^{1}(v)+2 k c_{0} v \sqrt{t-s}\right) \leq \sqrt{t-s}\left(\sqrt{2} \sup _{0 \leq v \leq 1} \tilde{B}^{1}(v)+2 k c_{0} \sqrt{t-s}\right) .
$$

Since $\sup _{0 \leq v \leq 1} \tilde{B}^{1}(v)$ is equally distributed with $\left|\tilde{B}^{1}(1)\right|$, which has Gaussian tail distribution, substituting the above bound into (9.19) yields

$$
\mathbb{P}\left[N_{s, t}^{1} \geq k_{1}\right] \leq C_{k, k_{1}}(t-s)^{k_{1} / 2} \quad \text { for all } 0 \leq t-s \leq 1
$$

for some $C_{k, k_{1}}$ depending only on $k$ and $k_{1}$ and $c_{0}$. If we let $N_{s, t}^{\mathrm{r}}$ denote the number of $s, t$-relevant separation points along $\pi^{0}$ which are of type $(1,2)_{\mathrm{r}}$ in $\mathcal{W}_{0}$, then similarly, $N_{s, t}^{\mathrm{r}}$ is distributed as a Poisson random variable with mean $2 k c_{1} \Delta^{\mathrm{r}}(t-s)$, where $\Delta^{\mathrm{r}}(t-s)=$ $\sup _{0 \leq v<t-s}\left(B^{\mathrm{r}}(v)-\pi^{0}(t-v)\right)$ for an independent Brownian motion $B^{\mathrm{r}}$ with drift $-\beta_{-}$starting from $\pi^{\overline{0}}(t)$. Also, for any $k_{\mathrm{r}} \geq 0$,

$$
\mathbb{P}\left[N_{s, t}^{\mathrm{r}} \geq k_{\mathrm{r}}\right] \leq C_{k, k_{\mathrm{r}}}(t-s)^{k_{\mathrm{r}} / 2} \quad \text { for all } 0 \leq t-s \leq 1
$$

for some $C_{k, k_{\mathrm{r}}}$ depending only on $k$ and $k_{\mathrm{r}}$ and $c_{1}$.
Observe that if we impose the partial order $\prec$ on $\mathcal{C}_{0}([0, t-s])$, the space of continuous functions with value 0 at 0 , where $f, g \in \mathcal{C}_{0}([0, t-s])$ satisfies $f \prec g$ if $f(v) \leq g(v)$ for all $v \in[0, t-s]$, then conditional on $B^{\mathrm{l}}$ and $B^{\mathrm{r}}, \Delta^{\mathrm{l}}(t-s)$ is increasing in $\left(\pi^{0}(t-v)-\pi^{0}(t)\right)_{v \in[0, t-s]}$, while $\Delta^{\mathrm{r}}(t-s)$ is decreasing in $\left(\pi^{0}(t-v)-\pi^{0}(t)\right)_{v \in[0, t-s]}$. This implies that for any $k_{1}, k_{\mathrm{r}} \geq 0$, $\mathbb{P}\left[N_{s, t}^{\mathrm{l}} \geq k_{1} \mid \pi^{0}, B^{\mathrm{l}}\right]$ is increasing in $\left(\pi^{0}(t-v)-\pi^{0}(t)\right)_{v \in[0, t-s]}$, while $\mathbb{P}\left[N_{s, t}^{\mathrm{r}} \geq k_{\mathrm{r}} \mid \pi^{0}, B^{\mathrm{r}}\right]$ is decreasing in $\left(\pi^{0}(t-v)-\pi^{0}(t)\right)_{v \in[0, t-s]}$, and hence the same holds for $\mathbb{P}\left[N_{s, t}^{1} \geq k_{1} \mid \pi^{0}\right]$ and $\mathbb{P}\left[N_{s, t}^{\mathrm{r}} \geq k_{\mathrm{r}} \mid \pi^{0}\right]$. Since the Brownian motion $\left(\pi^{0}(t-v)-\pi^{0}(t)\right)_{v \in[0, t-s]}$ satisfies the FKG inequality w.r.t. the partial order $\prec$ (see e.g. Bar05]), and the events $\left\{N_{s, t}^{1} \geq k_{1}\right\}$ and $\left\{N_{s, t}^{\mathrm{r}} \geq\right.$ $\left.k_{\mathrm{r}}\right\}$ are independent conditional on $\pi^{0}$, we have

$$
\begin{aligned}
\mathbb{P}\left[N_{s, t}^{1} \geq k_{1}, N_{s, t}^{\mathrm{r}} \geq k_{\mathrm{r}}\right] & =\mathbb{E}\left[\mathbb{P}\left[N_{s, t}^{1} \geq k_{1} \mid \pi^{0}\right] \mathbb{P}\left[N_{s, t}^{\mathrm{r}} \geq k_{\mathrm{r}} \mid \pi^{0}\right]\right] \\
& \leq \mathbb{P}\left[N_{s, t}^{\mathrm{r}} \geq k_{1}\right] \mathbb{P}\left[N_{s, t}^{\mathrm{r}} \geq k_{\mathrm{r}}\right] \leq C_{k, k_{1}} C_{k, k_{\mathrm{r}}} t^{\frac{k_{1}+k_{\mathrm{r}}}{2}} .
\end{aligned}
$$

Since $N_{s, t}=N_{s, t}^{1}+N_{s, t}^{\mathrm{r}}$, substituting this bound into (9.17) then gives

$$
\begin{equation*}
\mathbb{E}\left[\left(\left(X_{s}-X_{t}\right)^{+}\right)^{k}\right] \leq \mathbb{P}\left[N_{s, t} \geq k\right] \leq \sum_{i=0}^{k} \mathbb{P}\left[N_{s, t}^{1} \geq i, N_{s, t}^{\mathrm{r}} \geq k-i\right] \leq C_{k}(t-s)^{\frac{k}{2}} \tag{9.20}
\end{equation*}
$$

for some $C_{k}$ depending only on $k$.
We can now apply the one-sided Kolmogorov's moment criterion, Theorem E.1, which implies that for any $\gamma \in(0,1 / 2)$ and $\delta>0$, if we choose $k$ sufficiently large in 9.20) such that $(k / 2-1) / k>\gamma$ and $k / 2-1-\gamma k>\delta$, then there exists a random constant $C_{\gamma, \mathcal{W}_{0}}$ such that

$$
\begin{equation*}
X_{0}-X_{t}=\rho_{0}\left(\pi^{0}(0)\right)-\rho_{t}\left(\pi^{0}(t)\right) \leq C_{\gamma, \mathcal{W}_{0}} t^{\gamma} \quad \text { for all } t \in \mathbb{Q}_{2} \cap[0,1] \tag{9.21}
\end{equation*}
$$

where the distribution of $C_{\gamma, \mathcal{W}_{0}}$ satisfies 9.16).
To extend (9.21) from dyadic rational $t \in \mathbb{Q}_{2} \cap[0,1]$ to all $t \in[0,1]$, we note that $\rho_{t}$ is a.s. continuous in $t$, which implies that for any $\varepsilon>0$ and any $t_{n} \in \mathbb{Q}_{2} \cap[0,1]$ with $t_{n} \rightarrow t$, we have

$$
\rho_{t}\left(\left[\pi^{0}(t)-\varepsilon, \pi^{0}(t)+\varepsilon\right]\right) \geq \limsup _{n \rightarrow \infty} \rho_{t_{n}}\left(\left[\pi^{0}(t)-\varepsilon, \pi^{0}(t)+\varepsilon\right]\right) \geq \limsup _{n \rightarrow \infty} \rho_{t_{n}}\left(\pi^{0}\left(t_{n}\right)\right)
$$

by the a.s. continuity of the path $\pi^{0}$. Since $\varepsilon>0$ can be arbitrarily small, the above bound is also true for $\varepsilon=0$, which by (9.21) implies (9.15).

Proof of Theorem 9.6. By the third remark after Theorem [2.9, it suffices to consider the case $\rho_{0}$ is a probability measure on $\mathbb{R}$. Let $I_{0}(t):=\left\{\pi(t): \pi \in \mathcal{W}_{0}, \sigma_{\pi}<t\right\}$. We start by proving that

$$
\begin{equation*}
\mathbb{P}\left[\rho_{t}\left(I_{0}(t)\right)=1 \forall t>0\right]=1 \tag{9.22}
\end{equation*}
$$

The proof of Theorem 2.8 (a) in Section 9.1 showed that

$$
\begin{equation*}
\mathbb{P}\left[\rho_{t}\left(I_{0}(t)\right)=1\right]=1 \quad(t>0) \tag{9.23}
\end{equation*}
$$

Fix some deterministic $s>0$. For each $x \in \mathbb{R}$, let $\rho^{(x)}$ denote the Howitt-Warren process started at time $s$ from $\rho_{s}^{(x)}:=\delta_{x}$, constructed from the same reference web $\mathcal{W}_{0}$ as $\rho$. Then $\rho_{t}=\sum_{x \in I_{0}(s)} a_{x} \rho_{s}^{(x)}(t \geq s)$, where $a_{x}:=\rho_{s}(\{x\})$. If $x$ is deterministic or if $x \in I_{0}(s)$, then there is an a.s. unique path in $\mathcal{W}_{0}((x, s))$; let $\pi_{(x, s)}^{0}$ denote this path. Then, for each deterministic $\varepsilon \in(0,1)$ and $u>s$,

$$
\begin{aligned}
\mathbb{P} & \left.\inf _{t \in[s, u]} \rho_{t}\left(I_{0}(t)\right) \leq 1-\varepsilon\right] \leq \mathbb{P}\left[\inf _{t \in[s, u]} \sum_{x \in I_{0}(s)} \rho_{t}\left(\left\{\pi_{(x, s)}^{0}(t)\right\}\right) \leq 1-\varepsilon\right] \\
& \leq \mathbb{P}\left[\inf _{t \in[s, u]} \sum_{x \in I_{0}(s)} a_{x} \rho_{t}^{(x)}\left(\left\{\pi_{(x, s)}^{0}(t)\right\}\right) \leq 1-\varepsilon\right] \\
& \leq \mathbb{P}\left[\sum_{x \in I_{0}(s)} a_{x} \inf _{t \in[s, u]} \rho_{t}^{(x)}\left(\left\{\pi_{(x, s)}^{0}(t)\right\}\right) \leq 1-\varepsilon\right] \\
& =\mathbb{P}\left[\sum_{x \in I_{0}(s)} a_{x} \sup _{t \in[s, u]}\left(1-\rho_{t}^{(x)}\left(\left\{\pi_{(x, s)}^{0}(t)\right\}\right)\right) \geq \varepsilon\right] \\
& \leq \varepsilon^{-6} \mathbb{E}\left[\left(\sum_{x \in I_{0}(s)} a_{x} \sup _{t \in[s, u]}\left(1-\rho_{t}^{(x)}\left(\left\{\pi_{(x, s)}^{0}(t)\right\}\right)\right)\right)^{6}\right] \\
& \leq \varepsilon^{-6} \mathbb{E}\left[\sum_{x \in I_{0}(s)} a_{x}\left(\sup _{t \in[s, u]}\left(1-\rho_{t}^{(x)}\left(\left\{\pi_{(x, s)}^{0}(t)\right\}\right)\right)\right)^{6}\right] \\
& =\varepsilon^{-6} \mathbb{E}\left[\left(\sup _{t \in[s, u]}\left(1-\rho_{t}^{(0)}\left(\left\{\pi_{(0, s)}^{0}(t)\right\}\right)\right)\right)^{6}\right]
\end{aligned}
$$

where in the last inequality we applied the Hölder inequality w.r.t. the probability law given by the $\left(a_{x}\right)_{x \in I_{0}(s)}$, and in the last equality we used the spatial translation invariance of $\mathcal{W}_{0}$.

By Lemma 9.7 with $\gamma=1 / 3$,

$$
\begin{equation*}
\mathbb{E}\left[\sup _{t \in[s, u]}\left(1-\rho_{t}^{(0)}\left(\pi_{(0, s)}^{0}(t)\right)\right)^{6}\right] \leq \mathbb{E}\left[C_{\gamma, \mathcal{W}_{0}}^{6}(u-s)^{2}\right]=C(u-s)^{2} \tag{9.24}
\end{equation*}
$$

for some finite $C>0$, since $C_{\gamma, \mathcal{W}_{0}}$ has finite moments of all orders by 9.16]. Therefore, by our previous calculation, uniformly for all deterministic $0<s<u$,

$$
\begin{equation*}
\mathbb{P}\left[\inf _{t \in[s, u]} \rho_{t}\left(I_{0}(t)\right) \leq 1-\varepsilon\right] \leq C \varepsilon^{-6}(u-s)^{2} \tag{9.25}
\end{equation*}
$$

It follows that for each $n \geq 1$,

$$
\begin{equation*}
\mathbb{P}\left[\inf _{t \in\left[2^{-n}, 1\right]} \rho_{t}\left(I_{0}(t)\right) \leq 1-\varepsilon\right] \leq \sum_{k=2}^{2^{n}} \mathbb{P}\left[\inf _{t \in\left[(k-1) 2^{-n}, k 2^{-n}\right]} \rho_{t}\left(I_{0}(t)\right) \leq 1-\varepsilon\right] \leq C \varepsilon^{-6} 2^{-n} \tag{9.26}
\end{equation*}
$$

Letting first $n \rightarrow \infty$ and then $\varepsilon \rightarrow 0$ shows that that a.s. $\rho_{t}\left(I_{0}(t)\right)=1$ for all $t \in(0,1]$, and similarly for all $t>0$. This completes the proof of (9.22). In particular, this shows that almost surely, $\rho_{t}$ is atomic for all $t>0$.

To complete the proof, we need to show that almost surely, for all $t>0$,

$$
\begin{equation*}
\left\{x \in \mathbb{R}: \rho_{t}(\{x\})>0\right\}=I_{0}(t) \cap \operatorname{supp}\left(\rho_{t}\right) \tag{9.27}
\end{equation*}
$$

The inclusion $\subset$ follows from 9.22 . The inclusion $\supset$ is trivial for the Arratia flow, for which the characteristic measure $\nu=0$. To prove it for erosion flows with $\nu=c_{0} \delta_{0}+c_{1} \delta_{1}$ $\left(c_{0}+c_{1}>0\right)$, by Lemma 8.7 and Proposition 8.9, we only need to consider the cases that either $\operatorname{supp}\left(\rho_{t}\right)=\mathbb{R}$ for all $t>0$, or $\operatorname{supp}\left(\rho_{t}\right)$ is a halfline whose moving boundary is a path in $\mathcal{W}_{0}$. Let $\mathcal{T} \subset(0, \infty)$ be a deterministic countable dense set. Then for each $t>0$ and $x \in I_{0}(t) \cap \operatorname{supp}\left(\rho_{t}\right)$, we can find a time $s \in \mathcal{T}$ and $y \in I_{0}(s) \cap \operatorname{supp}\left(\rho_{s}\right)$ such that $\pi_{(y, s)}^{0}(t)=x$, where $\pi_{(y, s)}^{0}$ denotes the a.s. unique path in $\mathcal{W}_{0}$ starting from $(y, s)$. Since by Lemma 9.2, we have $\left\{x \in \mathbb{R}: \rho_{s}(\{x\})>0\right\} \supset I_{0}(s) \cap \operatorname{supp}\left(\rho_{s}\right)$, it suffices to prove that if $\pi^{0}$ is a path in $\mathcal{W}_{0}$ and $\rho_{s}\left(\left\{\pi^{0}(s)\right\}\right)>0$ for some $s \in \mathcal{T}$, then $\rho_{t}\left(\left\{\pi^{0}(t)\right\}\right)>0$ for all $t \geq s$.

By translation invariance, and the fact that $I_{0}(s)$ is independent of the restriction of $\mathcal{W}_{0}$ to $[s, \infty]$, it suffices to prove that almost surely, if $\pi^{0}$ is the unique path in $\mathcal{W}_{0}$ starting from the origin and $\rho_{0}=\delta_{0}$, then $\rho_{t}\left(\left\{\pi^{0}(t)\right\}\right)>0$ for all $t \geq 0$. To see this, set

$$
\begin{equation*}
\psi_{s, t}:=K_{s, t}^{\uparrow}\left(\pi^{0}(s),\left\{\pi^{0}(t)\right\}\right) \quad(0 \leq s \leq t) \tag{9.28}
\end{equation*}
$$

and observe that $\psi_{s, t} \psi_{t, u} \leq \psi_{s, u}(0 \leq s \leq t \leq u)$. It follows from (9.24) that

$$
\begin{equation*}
\mathbb{P}\left[\inf _{t \in[s, u]} \psi_{s, t} \leq 1 / 2\right] \leq C(u-s)^{2} \tag{9.29}
\end{equation*}
$$

for some finite $C>0$, and therefore

$$
\begin{equation*}
\mathbb{P}\left[\inf _{t \in[0,1]} \psi_{0, t}=0\right] \leq \sum_{k=1}^{2^{n}} \mathbb{P}\left[\psi_{(k-1) 2^{-n}, k 2^{-n}} \leq 1 / 2\right] \leq C 2^{-n} \tag{9.30}
\end{equation*}
$$

Letting $n \rightarrow \infty$ shows that a.s. $\rho_{t}\left(\left\{\pi^{0}(t)\right\}\right)=\psi_{0, t}>0$ for all $t \in[0,1]$ and similarly for all $t>0$.

Proof of Theorem 2.8 (c). Immediate from Theorem 9.6.

## 10 Infinite starting mass and discrete approximation

In this section, we prove Theorems 2.92 .10 , which will be based on our construction of the Howitt-Warren flows in Theorems 3.7 and 4.7.

### 10.1 Proof of Theorem 2.9

We first prove part (b), assuming $\beta_{+}-\beta_{-}<\infty$. By Theorem 4.7, we have the following representation for a Howitt-Warren process with drift $\beta$ and characteristic measure $\nu$ :

$$
\begin{equation*}
\rho_{t}=\int \rho_{s}(\mathrm{~d} x) K_{s, t}^{+}(x, \cdot)=\int \rho_{s}(\mathrm{~d} x) \mathbb{P}\left[\pi_{(x, s)}^{+}(t) \in \cdot \mid(\mathcal{N}, \omega)\right] \quad \text { for all } s<t \tag{10.1}
\end{equation*}
$$

where $\mathcal{N}$ is the Brownian net with left and right speeds $\beta_{-}$resp. $\beta_{+}, \omega:=\left(\omega_{z}\right)_{z \in S}$ are i.i.d. marks attached to the separation points $S$ of $\mathcal{N}$, and $\pi_{(x, s)}^{+}$is the rightmost path starting from $(x, s)$ in the sample web $\mathcal{W}$ conditional on the random environment $(\mathcal{N}, \omega)$. By construction, $\mathcal{W} \subset \mathcal{N}$ a.s., therefore for any $L>0, K_{0, t}^{+}(x,[-L, L])=0$ if $|x|$ is sufficiently large. This implies that if $\rho_{0} \in \mathcal{M}_{\text {loc }}$, then a.s. $\rho_{t} \in \mathcal{M}_{\text {loc }}$ for all $t \geq 0$.

Let $s_{n}, t_{n}, t, \rho_{s_{n}}^{\langle n\rangle}$ and $\rho_{0}$ be as in Theorem 2.9. By the representation 10.1, proving (2.18) with vague convergence in $\mathcal{M}_{\text {loc }}(\mathbb{R})$ amounts to showing that, for all $f \in C_{c}(\mathbb{R})$,

$$
\begin{equation*}
\int \rho_{s_{n}}^{\langle n\rangle}(\mathrm{d} x) \mathbb{E}\left[f\left(\pi_{\left(x, s_{n}\right)}^{+}\left(t_{n}\right)\right) \mid(\mathcal{N}, \omega)\right] \underset{n \rightarrow \infty}{\longrightarrow} \int \rho_{0}(\mathrm{~d} x) \mathbb{E}\left[f\left(\pi_{(x, 0)}^{+}(t)\right) \mid(\mathcal{N}, \omega)\right] \quad \text { a.s. } \tag{10.2}
\end{equation*}
$$

Note that the sample web $\mathcal{W}$ constructed in Theorem 4.7 and from which we draw $\pi_{(,,))}^{+}$, is distributed as a Brownian web with drift $\beta$. Therefore for each $(x, s) \in \mathbb{R}^{2}$, a.s. $(x, s)$ is of type $(0,1)$ in $\mathcal{W}$, and hence $\pi_{\left(x_{n}, s_{n}\right)}^{+} \rightarrow \pi_{(x, s)}^{+}$for any $\left(x_{n}, s_{n}\right) \rightarrow(x, s)$. By Fubini, a.s. w.r.t. $(\mathcal{N}, \omega)$, there exists a set $A_{(\mathcal{N}, \omega)} \subset \mathbb{R}$ with full $\rho_{0}$ measure such that for all $x \in A_{(\mathcal{N}, \omega)},(x, 0)$ is of type $(0,1)$ in $\mathcal{W}$ a.s. w.r.t. the quenched law $\mathbb{Q}:=\mathbb{P}(\mathcal{W} \in \cdot \mid(\mathcal{N}, \omega)$. This implies that

$$
\phi_{(\mathcal{N}, \omega)}(x, s ; t):=\mathbb{E}\left[f\left(\pi_{(x, s)}^{+}(t)\right) \mid(\mathcal{N}, \omega)\right]
$$

is continuous at all $x \in A_{(\mathcal{N}, \omega)}, s=0$ and $t \geq 0$. Since $\mathcal{W} \subset \mathcal{N}$ a.s. and $f$ has compact support, there exists $L>0$ such that

$$
\begin{equation*}
\phi_{(\mathcal{N}, \omega)}\left(x, s_{n} ; t_{n}\right)=\mathbb{E}\left[f\left(\pi_{\left(x, s_{n}\right)}^{+}\left(t_{n}\right)\right) \mid(\mathcal{N}, \omega)\right]=0 \quad \text { for all }|x|>L, n \in \mathbb{N} \tag{10.3}
\end{equation*}
$$

Choose $-L_{1}<-L$ and $L_{2}>L$ to be points of continuity of $\rho_{0}$. Note that restricted to $\left(-L_{1}, L_{2}\right), \rho_{s_{n}}^{\langle n\rangle}$ converges weakly to $\rho_{0}$. The convergence in 10.2 then follows from the continuous mapping theorem for weak convergence. The almost sure path continuity of $\left(\rho_{t}\right)_{t \geq 0}$ follows from 10.2 by setting $s_{n}=0$ and $\rho_{s_{n}}^{\langle n\rangle}=\rho_{0}$. This proves part (b).

We now prove part (a). When $\beta_{+}-\beta_{-}=\infty, 10.3$ may fail and we need to control in any finite region the inflow of measure from arbitrarily far away. By Theorem 3.7, we have the representation:

$$
\begin{equation*}
\rho_{t}=\int \rho_{s}(\mathrm{~d} x) K_{s, t}^{+}(x, \cdot)=\int \rho_{s}(\mathrm{~d} x) \mathbb{P}\left[\pi_{(x, s)}^{+}(t) \in \cdot \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right] \quad \text { for all } s<t . \tag{10.4}
\end{equation*}
$$

Here $\left(\mathcal{W}_{0}, \mathcal{M}\right)$ are the reference web and its associated set of marked points as in Theorem 3.7, where for definitiveness we can choose $\nu_{1}=\nu_{\mathrm{r}}=\nu$, so that $\mathcal{W}_{0}$ and the sample web $\mathcal{W}$ are
both Brownian webs with drift $\beta$. In particular, $\pi_{(x, 0)}^{+}$is distributed as a Brownian motion with drift $\beta$ starting from $x$ at time 0 . Without loss of generality, assume $\beta=0$. Since $\rho_{0} \in \mathcal{M}_{g}(\mathbb{R})$, it is then easy to check that for any $c>0$,

$$
\mathbb{E}\left[\int e^{-c y^{2}} \rho_{t}(\mathrm{~d} y)\right]=\frac{1}{\sqrt{2 \pi t}} \int \rho_{0}(\mathrm{~d} x) \int e^{-c y^{2}} e^{-\frac{(y-x)^{2}}{2 t}} \mathrm{~d} y \leq C_{1} \int \rho_{0}(\mathrm{~d} x) e^{-C_{2} x^{2}}<\infty
$$

for some $C_{1}, C_{2}>0$, which implies (2.17) and that $\rho_{t} \in \mathcal{M}_{g}(\mathbb{R})$ almost surely.
By the representation (10.4), proving (2.18) with convergence in $\mathcal{M}_{g}(\mathbb{R})$ amounts to showing that, for all $c>0$ and all bounded continuous function $f: \mathbb{R} \rightarrow \mathbb{R}$, we have

$$
\begin{align*}
& \int e^{-\varepsilon x^{2}} \rho_{s_{n}}^{\left\langle n_{n}\right.} \\
&(\mathrm{d} x) e^{\varepsilon x^{2}} \mathbb{E}\left[f\left(\pi_{\left(x, s_{n}\right)}^{+}\left(t_{n}\right)\right) e^{-c \pi_{\left(x, s_{n}\right)}^{+}\left(t_{n}\right)^{2}} \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right]  \tag{10.5}\\
& \int e^{-\varepsilon x^{2}} \rho_{0}(\mathrm{~d} x) e^{\varepsilon x^{2}} \mathbb{E}\left[f\left(\pi_{(x, 0)}^{+}(t)\right) e^{-c \pi_{(x, 0)}^{+}(t)^{2}} \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right] \quad \text { a.s. }
\end{align*}
$$

where $\varepsilon>0$ is chosen small. Denote

$$
\begin{equation*}
\phi\left(\mathcal{W}_{0}, \mathcal{M}\right)(x, s ; t):=e^{\varepsilon x^{2}} \mathbb{E}\left[f\left(\pi_{(x, s)}^{+}(t)\right) e^{-c \pi_{(x, s)}^{+}(t)^{2}} \mid\left(\mathcal{W}_{0}, \mathcal{M}\right)\right] . \tag{10.6}
\end{equation*}
$$

As before, a.s. w.r.t. $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, there exists $A_{\left(\mathcal{W}_{0}, \mathcal{M}\right)} \subset \mathbb{R}$ with full $\rho_{0}$ measure such that $\phi_{\left(\mathcal{W}_{0}, \mathcal{M}\right)}(x, s ; t)$ is continuous at all $x \in A_{\left(\mathcal{W}_{0}, \mathcal{M}\right)}, s=0$ and $t \geq 0$. Our assumption $\rho_{s_{n}}^{\langle n\rangle} \rightarrow \rho_{0}$ in $\mathcal{M}_{\mathrm{g}}(\mathbb{R})$ implies that $e^{-\varepsilon x^{2}} \rho_{s_{n}}^{\langle n\rangle}(\mathrm{d} x)$ converges weakly to $e^{-\varepsilon x^{2}} \rho_{0}(\mathrm{~d} x)$. Therefore 10.5 follows from the continuous mapping theorem for weak convergence, provided we show that a.s. w.r.t. $\left(\mathcal{W}_{0}, \mathcal{M}\right)$,

$$
\begin{equation*}
\sup _{y \in \mathbb{R}, 0 \leq u \leq v \leq t}\left|\phi_{\left(\mathcal{W}_{0}, \mathcal{M}\right)}(y, u ; v)\right|<\infty \quad \text { for all } t>0 \tag{10.7}
\end{equation*}
$$

so that we can apply the bounded convergence theorem. We verify (10.7) by Borel-Cantelli.
Without loss of generality, assume $|f|_{\infty}=1$. Fix $t>0$. For each $m \in \mathbb{Z}$, we have

$$
\begin{equation*}
\mathbb{E}\left[\sup _{\substack{y \in m, m+1] \\ 0 \leq u \leq v \leq t}}\left|\phi_{\left(\mathcal{W}_{0}, \mathcal{M}\right)}(y, u ; v)\right|\right] \leq C e^{3 \varepsilon m^{2}} \mathbb{E}\left[\sup _{\substack{y \in m, m+1] \\ 0 \leq u \leq v \leq t}} e^{-c \pi_{(y, u)}^{+}(v)^{2}}\right]<\infty . \tag{10.8}
\end{equation*}
$$

Now consider $m \geq L$ for some fixed large $L$. By the coalescing property of paths in the sample web $\mathcal{W}$, if $\pi_{(3 m / 4,0)}^{+} \in \mathcal{W}$ starting from $(3 m / 4,0)$ stays within $[m / 2, m]$ on the time interval $[0, t]$, then $\inf _{y \in[m, m+1], 0 \leq u \leq v \leq t} \pi_{(y, u)}^{+}(v) \geq m / 2$. Therefore

$$
\mathbb{E}\left[\sup _{\substack{y \in[m, m+m] \\ 0 \leq u \leq v \leq t}} e^{-c \pi_{(y, u)}^{+}(v)^{2}}\right] \leq \mathbb{P}\left(\sup _{0 \leq s \leq t}\left|\pi_{(3 m / 4,0)}^{+}(s)-\frac{3 m}{4}\right| \geq \frac{m}{4}\right)+e^{-\frac{c m^{2}}{4}} \leq C e^{-\frac{m^{2}}{32 t}}+e^{-\frac{c m^{2}}{4}},
$$

and hence

$$
\begin{equation*}
\mathbb{E}\left[\sup _{\substack{y \in m, m+1] \\ 0 \leq u \leq v \leq t}}\left|\phi_{\left(\mathcal{W}_{0}, \mathcal{M}\right)}(y, u ; v)\right|\right] \leq C e^{-\left(\frac{1}{32 t}-3 \varepsilon\right) m^{2}}+C e^{-\left(\frac{c}{4}-3 \varepsilon\right) m^{2}} \leq C_{1} e^{-C_{2} m^{2}} \tag{10.9}
\end{equation*}
$$

for some $C_{1}, C_{2}>0$ depending only on $t$ and $c$ if we choose $\varepsilon>0$ sufficiently small. Thus

$$
\mathbb{P}\left(\sup _{\substack{y \in[[, m+1] \\ 0 \leq u \leq v \leq t}}\left|\phi_{\left(\mathcal{W}_{0}, \mathcal{M}\right)}(y, u ; v)\right|>1\right) \leq C_{1} e^{-C_{2} m^{2}}
$$

By Borel-Cantelli, a.s. w.r.t. $\left(\mathcal{W}_{0}, \mathcal{M}\right)$, there exists a random $N_{+}>L$ sufficiently large such that $\sup _{y \geq N_{+}, 0 \leq u \leq v \leq t}\left|\phi_{\left(\mathcal{W}_{0}, \mathcal{M}\right)}(y, u ; v)\right| \leq 1$. Similarly, a.s. there exists $N_{-}<-L$ such that $\sup _{y \leq N_{-}, 0 \leq u \leq v \leq t}\left|\phi_{\left(\mathcal{W}_{0}, \mathcal{M}\right)}(y, u ; v)\right| \leq 1$. Combined with 10.8), this implies 10.7), and hence (10.5). The almost sure path continuity of $\left(\rho_{t}\right)_{t \geq 0}$ in $\mathcal{M}_{\mathrm{g}}(\mathbb{R})$ follows from (10.5) by setting $s_{n}=0$ and $\rho_{s_{n}}^{\langle n\rangle}=\rho_{0}$.

### 10.2 Proof of Theorem 2.10

The proof is similar to that of Theorem 2.9. The complication lies again with infinite $\bar{\rho}_{t}^{\langle k\rangle}$ and $\rho_{t}$. Without loss of generality, assume that the Howitt-Warren process $\left(\rho_{t}\right)_{t \geq 0}$ has drift $\beta=0$. First we note that there exists a countable family of bounded continuous functions $\left\{f_{n}\right\}_{n \in \mathbb{N}}$ such that a sequence of finite measures $\xi_{k} \in \mathcal{M}(\mathbb{R})$ converges weakly to $\xi \in \mathcal{M}(\mathbb{R})$ if and only if $\int f_{n}(x) \xi_{k}(\mathrm{~d} x) \rightarrow \int f_{n}(x) \xi(\mathrm{d} x)$ for all $n \in \mathbb{N}$, see e.g. Res87, Proof of Prop. 3.17]. Since $\xi_{k} \rightarrow \xi$ in $\mathcal{M}_{\mathrm{g}}(\mathbb{R})$ is equivalent to weak convergence of $e^{-c x^{2}} \xi_{k}(\mathrm{~d} x)$ to $e^{-c x^{2}} \xi(\mathrm{~d} x)$ for all $c>0$, to prove the weak convergence in 2.20 on path space with uniform topology, it suffices to show that for any finite sets $K \subset(0, \infty)$ and $\Lambda \subset \mathbb{N}$, we have

$$
\begin{equation*}
\left(F_{c, n}^{\langle k\rangle}(t):=\int e^{-c x^{2}} f_{n}(x) \bar{\rho}_{t}^{-k\rangle}(\mathrm{d} x)\right)_{c \in K, n \in \Lambda} \underset{n \rightarrow \infty}{\Longrightarrow}\left(F_{c, n}(t):=\int e^{-c x^{2}} f_{n}(x) \rho_{t}(\mathrm{~d} x)\right)_{c \in K, n \in \Lambda}, \tag{10.10}
\end{equation*}
$$

where $\Rightarrow$ denotes weak convergence of $\mathcal{C}([0, T], \mathbb{R})^{|K|+|\Lambda|}$-valued random variables.
For any $c>0$ and $n \in \mathbb{N}, F_{c, n} \in \mathcal{C}([0, T], \mathbb{R})$ a.s. by Theorem 2.9. By similar reasoning, $\bar{\rho}_{t}^{\langle k\rangle}$ has a.s. continuous sample path in $\mathcal{M}_{\mathrm{g}}(\mathbb{R})$ and hence $F_{c, n}^{\langle k\rangle} \in \mathcal{C}([0, T], \mathbb{R})$. Since $\bar{\rho}_{0}^{\langle k\rangle}$ converges weakly to $\rho_{0}$ as $\mathcal{M}_{\mathrm{g}}(\mathbb{R})$-valued random variables, the Skorohod representation theorem for weak convergence (see e.g. [Bil99, Theorem 6.7]) allows a coupling between $\left(\bar{\rho}_{0}^{\langle k\rangle}\right)_{k \in \mathbb{N}}$ and $\rho_{0}$ such that $\bar{\rho}_{0}^{\langle k\rangle} \rightarrow \rho_{0}$ in $\mathcal{M}_{\mathrm{g}}(\mathbb{R})$ almost surely. Therefore we may assume that $\left(\bar{\rho}_{0}^{\langle k\rangle}\right)_{n \in \mathbb{N}}$ and $\rho_{0}$ are deterministic and $\bar{\rho}_{0}^{\langle k\rangle} \rightarrow \rho_{0}$.

Recall from (3.3) and Section 3.5 the discrete quenched law $\mathbf{Q}_{\langle k\rangle}$ associated with a discrete Howitt-Warren flow with characteristic measure $\mu_{k}$, and recall from (3.25) and Theorem 3.9 the diffusive scaling map $S_{\varepsilon}$ and its action on a quenched law $\mathbf{Q}$. We have the representation

$$
\bar{\rho}_{t}^{\langle k\rangle}=\int \bar{\rho}_{0}(\mathrm{~d} x) S_{\varepsilon_{k}}\left(\mathbf{Q}_{\langle k\rangle}\right)\left[\pi_{(x, 0)}^{\langle k\rangle}(t) \in \cdot\right]
$$

where for $(x, s) \in S_{\varepsilon_{k}}\left(\mathbb{Z}_{\text {even }}^{2}\right), \pi_{(x, s)}^{\langle k\rangle}$ is the unique path starting from $(x, s)$ in a discrete sample web $\mathcal{W}^{\langle k\rangle}$ with quenched law $S_{\varepsilon_{k}}\left(\mathbf{Q}_{\langle k\rangle}\right)$. Similarly, by Theorem 3.7.

$$
\rho_{t}=\int \rho_{0}(\mathrm{~d} x) \mathbb{Q}\left[\pi_{(x, 0)}^{+}(t) \in \cdot\right],
$$

with the quenched law $\mathbb{Q}$ defined as in (3.24). For any $L>0$, we can then write

$$
\begin{aligned}
F_{c, n}^{\langle k\rangle}(t) & =F_{c, n}^{\langle k\rangle,[-L, L]}(t)+F_{c, n}^{\langle k\rangle,[-L, L]^{c}}(t) \\
F_{c, n}(t) & =F_{c, n}^{[-L, L]}(t)+F_{c, n}^{[-L, L]^{c}}(t),
\end{aligned}
$$

where for any $I \subset \mathbb{R}$,

$$
\begin{aligned}
F_{c, n}^{\langle k\rangle, I}(t) & =\int_{I} \bar{\rho}_{0}(\mathrm{~d} x) S_{\varepsilon_{k}}\left(\mathbf{Q}_{\langle k\rangle}\right)\left[f_{n}\left(\pi_{(x, 0)}^{\langle k\rangle}(t)\right) e^{-c \pi_{(x, 0)}^{(k\rangle}(t)^{2}}\right], \\
F_{c, n}^{I}(t) & =\int_{I} \rho_{0}(\mathrm{~d} x) \mathbb{Q}\left[f_{n}\left(\pi_{(x, 0)}^{+}(t)\right) e^{-c \pi_{(x, 0)}^{+}(t)^{2}}\right] .
\end{aligned}
$$

To prove 10.10 , it suffices to show that for any $\varepsilon>0, c>0$ and $n \in \mathbb{N}$, we can choose $L$ large such that

$$
\begin{gather*}
\limsup _{k \rightarrow \infty} \mathbb{E}\left[\left|F_{c, n}^{\langle k\rangle,[-L, L]^{c}}\right|_{\infty}\right] \leq \varepsilon  \tag{10.11}\\
\mathbb{E}\left[\left|F_{c, n}^{[-L, L]^{c}}\right|_{\infty}\right] \leq \varepsilon \tag{10.12}
\end{gather*}
$$

where $|\cdot|_{\infty}$ denotes the supremum norm on $\mathcal{C}([0, T], \mathbb{R})$, and furthermore,

$$
\begin{equation*}
\left(F_{c, n}^{\langle k\rangle,[-L, L]}\right)_{c \in K, n \in \Lambda} \underset{n \rightarrow \infty}{\Longrightarrow}\left(F_{c, n}^{[-L, L]}\right)_{c \in K, n \in \Lambda} \tag{10.13}
\end{equation*}
$$

with $\Rightarrow$ denoting weak convergence of $\mathcal{C}([0, T], \mathbb{R})^{|K|+|\Lambda|}$-valued random variables.
Fix $0<\varepsilon<\inf K$ and define

$$
\begin{align*}
\phi_{\mathbf{Q}_{\langle k\rangle}}^{\langle k\rangle, c, n}(x, t) & :=e^{\varepsilon x^{2}}\left(S_{\varepsilon_{k}} \mathbf{Q}_{\langle k\rangle}\right)\left[f_{n}\left(\pi_{(x, 0)}^{\langle k\rangle}(t)\right) e^{-c \pi_{(x, 0)}^{\langle k\rangle}(t)^{2}}\right]  \tag{10.14}\\
\phi_{\mathbb{Q}}^{c, n}(x, t) & :=e^{\varepsilon x^{2}} \mathbb{Q}\left[f_{n}\left(\pi_{(x, 0)}^{+}(t)\right) e^{-c \pi_{(x, 0)}^{+}(t)^{2}}\right]
\end{align*}
$$

Then

$$
\left|F_{c, n}^{[-L, L]^{c}}(\cdot)\right|_{\infty}=\left|\int_{[-L, L]^{c}} e^{-\varepsilon x^{2}} \rho_{0}(\mathrm{~d} x) \phi_{\mathbb{Q}}^{c, n}(x, \cdot)\right|_{\infty} \leq \int_{[-L, L]^{c}} e^{-\varepsilon x^{2}} \rho_{0}(\mathrm{~d} x)\left|\phi_{\mathbb{Q}}^{c, n}(x, \cdot)\right|_{\infty}
$$

Note that $\phi_{\mathbb{Q}}^{c, n}(x, t)$ is exactly $\phi_{\left(\mathcal{W}_{0}, \mathcal{M}\right)}(x, s ; t)$ in 10.6 if we set $s=0$ and $f=f_{n}$. Since $e^{-\varepsilon x^{2}} \rho_{0}(\mathrm{~d} x)$ is a finite measure, 10.12 then follows from 10.9 . Note that 10.9 is based on Brownian motion estimates, and analogues of 10.9 for $\phi_{\mathbf{Q}_{\langle k\rangle}}^{\langle k\rangle}, c, n$ can be established using random walk estimates, which are furthermore uniform in $k \in \mathbb{N}$. Such a uniform estimate and the fact that $e^{-\varepsilon x^{2}} \bar{\rho}_{0}(\mathrm{~d} x)$ converges weakly to $e^{-\varepsilon x^{2}} \rho_{0}(\mathrm{~d} x)$ then imply 10.11 . We omit the details here.

Lastly, we prove 10.13 , where we may assume that $-L$ and $L$ are points of continuity of $\rho_{0}(\mathrm{~d} x)$ so that restricted to $[-L, L], e^{-\varepsilon x^{2}} \bar{\rho}_{0}(\mathrm{~d} x)$ converges weakly to $e^{-\varepsilon x^{2}} \rho_{0}(\mathrm{~d} x)$. By Skorohod representation, the weak convergence of $S_{\varepsilon_{k}} \mathbf{Q}_{\langle k\rangle}$ to $\mathbb{Q}$ in Theorem 3.9 can be turned into a.s. convergence in $\mathcal{M}_{1}(\mathcal{K}(\Pi))$ via a suitable coupling, which we now assume. Then the discrete sample web $\mathcal{W}^{\langle k\rangle}$ with law $S_{\varepsilon_{k}} \mathbf{Q}_{\langle k\rangle}$ converges weakly to the sample web $\mathcal{W}$ with law $\mathbb{Q}$, where the convergence can again assumed to be a.s. in $\mathcal{K}(\Pi)$ by Skorohod representation. Then for each $c>0$ and $k, n \in \mathbb{N}$,

$$
\left|F_{c, n}^{\langle k\rangle,[-L, L]}-F_{c, n}^{[-L, L]}\right|_{\infty}=\left|\int_{-L}^{L} e^{-\varepsilon x^{2}} \bar{\rho}_{0}^{\langle k\rangle}(\mathrm{d} x) \phi_{\mathbf{Q}_{\langle k\rangle}^{\langle k\rangle}, c, n}(x, \cdot)-\int_{-L}^{L} e^{-\varepsilon x^{2}} \rho_{0}(\mathrm{~d} x) \phi_{\mathbb{Q}}^{c, n}(x, \cdot)\right|_{\infty}
$$

First we claim that for a.e. $x \in[-L, L]$ w.r.t. the measure $e^{-\varepsilon x^{2}} \rho_{0}(\mathrm{~d} x)$, if $x_{n} \rightarrow x$ for some sequence $x_{n} \in \varepsilon_{k} \mathbb{Z}_{\text {even }}$, then for each $c>0$ and $n \in \mathbb{N}$,

$$
\begin{equation*}
\left|\phi_{\mathbf{Q}_{\langle k\rangle}}^{\langle k\rangle, c, n}\left(x_{n}, \cdot\right)-\phi_{\mathbb{Q}}^{c, n}(x, \cdot)\right|_{\infty} \underset{n \rightarrow \infty}{\longrightarrow} 0 . \tag{10.15}
\end{equation*}
$$

Indeed, since the law of $\mathcal{W}$ averaged over the randomness of $\mathbb{Q}$ is that of a Brownian web, each deterministic $(x, 0)$ is a.s. of type $(0,1)$ in $\mathcal{W}$, and by Fubini, the same is true for $e^{-\varepsilon x^{2}} \rho_{0}$ a.e. $x$. Therefore $\mathcal{W}^{\langle k\rangle} \rightarrow \mathcal{W}$ a.s. in $\mathcal{K}(\Pi)$ implies that $\pi_{\left(x_{k}, 0\right)}^{\langle k\rangle} \rightarrow \pi_{(x, 0)}^{+}$in $\mathcal{C}([0, T], \mathbb{R})$, which
when plugged into the definitions in 10.14 then implies 10.15 . If we regard $\phi_{\mathbf{Q}_{\langle k\rangle}}^{\langle k\rangle, c, n}(x, \cdot)$ and $\phi_{\mathbb{Q}}^{c, n}(x, \cdot)$ as mappings from $\mathbb{R}$ to $\mathcal{C}([0, T], \mathbb{R})$ and note that $\left|\phi_{\mathbf{Q}_{(k\rangle}}^{\langle k\rangle, c, n}(x, \cdot)\right|_{\infty}$ and $\left|\phi_{\mathbb{Q}}^{c, n}(x, \cdot)\right|_{\infty}$ are bounded uniformly in $\mathbf{Q}_{\langle k\rangle}, \mathbb{Q}$ and $x \in[-L, L]$, then the continuous mapping theorem for weak convergence implies that $\left|F_{c, n}^{\langle k),[-L, L]}-F_{c, n}^{[-L, L]}\right|_{\infty} \rightarrow 0$ a.s. for each $c \in K$ and $n \in \Lambda$, which then implies (10.13).

## 11 Ergodic properties

In this section, we prove Theorem 2.11 2.12 on homogeneous invariant laws for Howitt-Warren processes. By the observation that $\rho_{t}$ in (2.1) depends linearly on the initial condition $\rho_{0}$, the Howitt-Warren process falls in the class of linear systems, the theory of which for processes on $\mathbb{Z}^{d}$ was developed by Liggett and Spitzer, see e.g. [LS81] and Lig05, Chap. IX]. We will adapt the theory of linear systems to our continuum setting. The main tools are duality, second moment calculations, and coupling, which will be developed in successive subsections. Duality is used to give a simple construction of the family of ergodic homogeneous invariant laws. Second moment calculations determine spatial correlations for the homogeneous invariant laws, and are used to prove the uniform integrability of the Howitt-Warren process $\left(\rho_{t}\right)_{t \geq 0}$ over time, as well as to show that certain spatial ergodic properties of the initial measure $\rho_{0}$ are preserved by the dynamics even in the limit $t \rightarrow \infty$. The last point will be crucial for proving convergence to homogeneous invariant laws, which is based on coupling arguments. Most of our arguments are adapted from [LS81 and Lig05, Chap. IX], to which we will refer many of the details. The main difference in our proof lies in the second moment calculations of Lemma 11.6, for which we need to devise a different and perhaps more robust approach than the one used in LS81, Lig05.

### 11.1 Dual smoothing process

Similar to the linear systems on $\mathbb{Z}^{d}$ studied in LS81, Lig05, the Howitt-Warren process $\left(\rho_{t}\right)_{t \geq 0}$ is dual to a function-valued smoothing process with random kernels. Analogous to the construction of the Howitt-Warren process from the Howitt-Warren flow $\left(K_{s, t}^{+}\right)_{s<t}$ as in (2.1), we can define a function-valued dual process $\left(\zeta_{t}\right)_{t \geq 0}$ by

$$
\begin{equation*}
\zeta_{t}(x)=\int K_{-t, 0}^{+}(x, \mathrm{~d} y) \zeta_{0}(y)=\int \mathbb{Q}\left[\zeta_{0}\left(\pi_{(x,-t)}^{+}(0)\right)\right] \quad \text { for all } x \in \mathbb{R}, \tag{11.1}
\end{equation*}
$$

where $\mathbb{Q}$ is the quenched law of a sample web $\mathcal{W}$ defined as in Theorem 3.7, and $\pi_{(x, s)}^{+}$is the a.s. unique rightmost path in $\mathcal{W}$ starting from $(x, s)$. A natural state space for $\left(\zeta_{t}\right)_{t \geq 0}$ is $D_{b}(\mathbb{R})$, the space of bounded cádlág functions on $\mathbb{R}$. Note that $\pi_{(x,-t)}^{+}(0)$ is cádlág in $x$. With this observation, it is easy to see that if $\zeta_{0} \in D_{b}(\mathbb{R})$, then $\zeta_{t} \in D_{b}(\mathbb{R})$ for all $t>0$.

We have the following duality relation between $\left(\rho_{t}\right)_{t \geq 0}$ and $\left(\zeta_{t}\right)_{t \geq 0}$.
Lemma 11.1 (Duality) Let $\rho_{0} \in \mathcal{M}_{\mathrm{loc}}(\mathbb{R})$, or $\rho_{0} \in \mathcal{M}_{\mathrm{g}}(\mathbb{R})$ if $\beta_{+}-\beta_{-}=\infty$ in 2.12). Let $\zeta_{0} \in D_{b}(\mathbb{R})$. Assume that either $\rho_{0}$ is a finite measure or $\zeta_{0}$ has bounded support. Then for all $t \geq 0$,

$$
\begin{equation*}
\int \zeta_{0}(x) \rho_{t}(\mathrm{~d} x) \stackrel{\text { dist }}{=} \int \zeta_{t}(x) \rho_{0}(\mathrm{~d} x) . \tag{11.2}
\end{equation*}
$$

Proof. Follows from the definition of $\rho_{t}$ and $\zeta_{t}$, and the equality in law between $K_{0, t}^{+}$and $K_{-t, 0}^{+}$.

The advantage of working with the smoothing process $\left(\zeta_{t}\right)_{t \geq 0}$ is that there is a natural martingale associated with it.

Lemma 11.2 (Extinction vs uniform integrability) Let $\zeta_{0} \in D_{b}(\mathbb{R})$ have bounded support. Then $[\zeta]_{t}:=\int \zeta_{t}(x) \mathrm{d} x$ is a martingale which a.s. has a limit $[\zeta]_{\infty}$ as $t \rightarrow \infty$. Furthermore, either $[\zeta]_{\infty}=0$ a.s. for all $\zeta_{0} \in D_{b}(\mathbb{R})$ with bounded support, or $\left([\zeta]_{t}\right)_{t \geq 0}$ is uniformly integrable for all $\zeta_{0} \in D_{b}(\mathbb{R})$ with bounded support. We say the finite smoothing process $\zeta$ dies out in the first case, and survives in the second case.

Proof. By separating $\zeta_{0}$ into its positive and negative parts and by the linear dependence of $\zeta_{t}$ on $\zeta_{0}$, we may assume $\zeta_{0} \geq 0$. Note that for $0 \leq s<t$,

$$
\begin{align*}
t & =\int \zeta_{t}(x) \mathrm{d} x=\iint K_{-t, 0}^{+}(x, \mathrm{~d} y) \zeta_{0}(y) \mathrm{d} x  \tag{11.3}\\
& =\iiint K_{-t,-s}^{+}(x, \mathrm{~d} z) K_{-s, 0}^{+}(z, \mathrm{~d} y) \zeta_{0}(y) \mathrm{d} x=\iint \zeta_{s}(z) K_{-t,-s}^{+}(x, \mathrm{~d} z) \mathrm{d} x
\end{align*}
$$

By the independence of $K_{-t,-s}^{+}$and $\left(K_{-u, 0}^{+}\right)_{0 \leq u \leq s}, K_{-t,-s}^{+}$is independent of $\left(\zeta_{u}\right)_{0 \leq u \leq s}$. By the translation invariance in law of $K_{-t,-s}^{+}(x, \cdot)$ in $x$, we note that $\int \mathbb{E}\left[K_{-t,-s}^{+}(x, \cdot)\right] \mathrm{d} x$ is simply the Lebesgue measure. Therefore

$$
\mathbb{E}\left[[\zeta]_{t} \mid\left([\zeta]_{u}\right)_{0 \leq u \leq s}\right]=\iint \zeta_{s}(z) \mathbb{E}\left[K_{-t,-s}^{+}(x, \mathrm{~d} z)\right] \mathrm{d} x=\int \zeta_{s}(z) \mathrm{d} z=[\zeta]_{s}
$$

which proves the martingale property of $[\zeta]_{t}$. Since $[\zeta]_{t}$ is furthermore non-negative, a.s. it has a limit $[\zeta]_{\infty}$.

The dichotomy between a.s. extinction and uniform integrability of $\left([\zeta]_{t}\right)_{t \geq 0}$ follows from a similar argument as in the proof of Lig05, Thm. IX.2.4.(a)]. Given $\zeta_{0}(x)=1_{[0,1]}(x)$, let $\lambda:=\mathbb{E}\left[[\zeta]_{\infty}\right]$. First we claim that for any $\zeta_{0} \in D_{b}(\mathbb{R})$ with bounded support, we have $\left.\mathbb{E}[\zeta \zeta]_{\infty}\right]=\lambda[\zeta]_{0}$. By the linear dependence of $\zeta_{\infty}$ on $\zeta_{0}$ and the translation invariance of the kernels $\left(K_{s, t}^{+}\right)_{s \leq t}$, the claim holds for all $\zeta_{0}$ which are linear combinations of characteristic functions of finite intervals. Since all $\zeta_{0} \in D_{b}(\mathbb{R})$ with bounded support can be approximated from above and below by such functions, and $[\zeta]_{\infty}$ depends monotonically on $\zeta_{0}$, the claim holds for all $\zeta_{0} \in D_{b}(\mathbb{R})$ with bounded support. The dichotomy amounts to showing either $\lambda=0$ or $\lambda=1$.

Note that the RHS of 11.3 can be interpreted as $\left[\tilde{\zeta^{1}}\right]_{t-s}:=\int \tilde{\zeta}_{t-s}(x) \mathrm{d} x$ for a smoothing process $\tilde{\zeta}$ defined from the time-shifted kernels $\left(K_{-r-s,-s}^{+}\right)_{r \geq 0}$ with initial condition $\tilde{\zeta}_{0}=\zeta_{s}$. In particular, a.s. $[\tilde{\zeta}]_{t-s}$ tends to a limit $[\tilde{\zeta}]_{\infty}$ as $t \rightarrow \infty$. Letting $t \rightarrow \infty$ in (11.3) then gives $[\zeta]_{\infty}=[\tilde{\zeta}]_{\infty}$. Therefore by Jensen's inequality,

$$
\mathbb{E}\left[e^{-[\zeta]_{\infty}}\right]=\mathbb{E}\left[e^{-[\tilde{\zeta}]_{\infty}}\right] \geq \mathbb{E}\left[e^{\left.-\mathbb{E}[\tilde{\zeta}]_{\infty} \mid \tilde{\zeta}_{0}=\zeta_{s}\right]}\right]=\mathbb{E}\left[e^{-\lambda[\zeta]_{s}}\right]
$$

where we now take the limit $s \rightarrow \infty$ and obtain by the bounded convergence theorem that

$$
\begin{equation*}
\mathbb{E}\left[e^{\left.-[\zeta]_{\infty}\right]} \geq \mathbb{E}\left[e^{-\lambda[\zeta]_{\infty}}\right]\right. \tag{11.4}
\end{equation*}
$$

Since $\zeta_{0} \geq 0$ by assumption in (11.3), we have $[\zeta]_{\infty} \geq 0$. Assume further that $[\zeta]_{0}>0$. If $[\zeta]_{\infty}=0$ a.s., then $\lambda=0$. If $[\zeta]_{\infty}>0$ with positive probability, then because $\lambda \in[0,1]$, 11.4) can only hold if $\lambda=1$.

Lemmas 11.1 and 11.2 imply the weak convergence of the Howitt-Warren process $\rho_{t}$ with initial condition $\rho_{0}(\mathrm{~d} x)=c \mathrm{~d} x$ to a homogeneous invariant law. Recall the set of invariant laws $\mathcal{I}$ and $\mathcal{T}$ from Theorem 2.11.

Lemma 11.3 (Construction of homogeneous invariant laws) Assume that $\rho_{0}(\mathrm{~d} x)=$ $c \mathrm{~d} x$ for some $c \geq 0$. Then there exists $\Lambda_{c} \in \mathcal{I} \cap \mathcal{T}$ such that

$$
\begin{equation*}
\mathcal{L}\left(\rho_{t}\right) \underset{t \rightarrow \infty}{\Longrightarrow} \Lambda_{c}, \tag{11.5}
\end{equation*}
$$

where $\mathcal{L}(\cdot)$ denotes law and $\Rightarrow$ denotes weak convergence of probability laws on $\mathcal{M}_{\mathrm{loc}}(\mathbb{R})$. If the finite smoothing process $\zeta$ survives, then $\int \rho([0,1]) \Lambda_{c}(\mathrm{~d} \rho)=c$; otherwise $\Lambda_{c}=\delta_{0}$, the delta measure concentrated in the zero measure on $\mathbb{R}$. Furthermore, $\Lambda_{c}(\mathrm{~d}(c \rho))=\Lambda_{1}(\mathrm{~d} \rho)$.

Proof. Since $\rho_{0}(\mathrm{~d} x)=c \mathrm{~d} x$, by the translation invariance of $\left(K_{0, t}^{+}(x, \cdot)\right)_{x \in \mathbb{R}}$ in space, we have $\mathbb{E}\left[\rho_{t}(\mathrm{~d} x)\right]=c \mathrm{~d} x$ for all $t \geq 0$. In particular, for any bounded interval $I \subset \mathbb{R},\left(\rho_{t}(I)\right)_{t \geq 0}$ is a tight family of random variables, which implies that $\left(\rho_{t}\right)_{t \geq 0}$ is a tight family of $\mathcal{M}_{\text {loc }}(\mathbb{R})$ valued random variables (see e.g. Res87, Lemma 3.20]). In fact $\left(\rho_{t}\right)_{t \geq 0}$ is also a tight family of $\mathcal{M}_{\mathrm{g}}(\mathbb{R})$-valued random variables (recall 2.15$)$. This follows from the additional observation that for any $a>0,\left(\int e^{-a x^{2}} \rho_{t}(\mathrm{~d} x)\right)_{t \geq 0}$ is a tight family of real-valued random variables, because

$$
\sup _{t \geq 0} \mathbb{E}\left[\int e^{-a x^{2}} \rho_{t}(\mathrm{~d} x)\right]=c \int e^{-a x^{2}} \mathrm{~d} x<\infty
$$

Let $\zeta_{t}$ be the dual smoothing process with initial condition $\zeta_{0} \in C_{c}(\mathbb{R})$, the space of continuous functions with compact support. If $\rho_{t}$ converges weakly to a $\mathcal{M}_{\mathrm{g}}(\mathbb{R})$-valued random variable $\rho^{*}$ along a subsequence $t_{n} \uparrow \infty$, then by Lemmas 11.1 and 11.2 , we must have equality in distribution between $\int \zeta_{0}(x) \rho^{*}(\mathrm{~d} x)$ and $c[\zeta]_{\infty}$. Since the law of $[\zeta]_{\infty}$ does not depend on $t_{n} \uparrow \infty$, and $\zeta_{0}$ can be any function in $C_{c}(\mathbb{R})$, the law of $\rho^{*}$ is uniquely determined (see e.g. [Res87, Prop. 3.19]). Together with tightness, this implies that $\rho_{t} \Longrightarrow \rho_{t \rightarrow \infty}^{\Longrightarrow} \rho^{*}$ as $\mathcal{M}_{\mathrm{g}}(\mathbb{R})$-valued random variables, and we denote $\Lambda_{c}:=\mathcal{L}\left(\rho^{*}\right)$. The fact that $\Lambda_{c} \in \mathcal{I}$ then follows from the Feller property of $\left(\rho_{t}\right)_{t \geq 0}$ implied by Theorem 2.9 , and clearly $\Lambda_{c} \in \mathcal{T}$. Since $\mathbb{E}\left[\rho^{*}([0,1])\right]=c \mathbb{E}\left[[\zeta]_{\infty}\right]$ with $\zeta_{0}=1_{[0,1]}$, the dichotomy between $\int \rho([0,1]) \Lambda_{c}(\mathrm{~d} \rho)=c$ and $\Lambda_{c}=\delta_{0}$ follows from Lemma 11.2. The scaling relation between $\Lambda_{1}$ and $\Lambda_{c}$ is trivial.

When the characteristic measure $\nu$ for the Howitt-Warren flow is not zero so that the flow is not purely coalescing, the possibility of $\Lambda_{c}=\delta_{0}$ in Lemma 11.3 can be ruled out by showing the uniform integrability of $\rho_{t}([0,1])$ in $t \geq 0$. This can be accomplished by the second moment calculation in Lemma 11.6 below. In any event, we can deduce the extremality of $\Lambda_{c}$ in $\mathcal{I} \cap \mathcal{T}$ using Lemmas 11.1 and 11.3 .

Lemma 11.4 (Extremality of $\Lambda_{c}$ ) For all $c \geq 0$, we have $\Lambda_{c} \in(\mathcal{I} \cap \mathcal{T})_{\mathrm{e}}$.
Proof. The proof is the same as that of [Lig05, Lemma IX.2.9]. We include it here for the reader's convenience. Assume that $\Lambda_{c}=\alpha \mu_{1}+(1-\alpha) \mu_{2}$ for some $\alpha \in(0,1)$ and $\mu_{1}, \mu_{2} \in \mathcal{I} \cap \mathcal{T}$
with $\int \rho([0,1]) \mu_{i}(\mathrm{~d} \rho)=c_{i}$, where $c=\alpha c_{1}+(1-\alpha) c_{2}$. Then for any $\zeta_{0} \in \mathcal{C}_{\mathrm{c}}(\mathbb{R})$ and $i=1,2$,

$$
\begin{align*}
\int \mathbb{E}\left[e^{-\int \zeta_{0}(x) \rho_{0}(\mathrm{~d} x)}\right] \mu_{i}\left(\mathrm{~d} \rho_{0}\right) & =\int \mathbb{E}\left[e^{-\int \zeta_{0}(x) \rho_{t}(\mathrm{~d} x)}\right] \mu_{i}\left(\mathrm{~d} \rho_{0}\right)=\int \mathbb{E}\left[e^{-\int \zeta_{t}(x) \rho_{0}(\mathrm{~d} x)}\right] \mu_{i}\left(\mathrm{~d} \rho_{0}\right) \\
& \geq \mathbb{E}\left[e^{-\iint \zeta_{t}(x) \rho_{0}(\mathrm{~d} x) \mu_{i}\left(\mathrm{~d} \rho_{0}\right)}\right]=\mathbb{E}\left[e^{-c_{i} \int \zeta_{t}(x) \mathrm{d} x}\right]  \tag{11.6}\\
& =\mathbb{E}\left[\left.e^{-\frac{c_{i}}{c} \int \zeta_{0}(x) \rho_{t}(\mathrm{~d} x)} \right\rvert\, \rho_{0} \equiv c\right] \underset{t \rightarrow \infty}{\longrightarrow} \int \mathbb{E}\left[e^{-\frac{c_{i}}{c} \int \zeta_{0}(x) \rho(\mathrm{d} x)}\right] \Lambda_{c}(\mathrm{~d} \rho),
\end{align*}
$$

where we used the fact that $\mu_{i} \in \mathcal{I} \cap \mathcal{T}$, duality, Jensen's inequality, and Lemma 11.3. Denote $\phi(a)=\int \mathbb{E}\left[e^{-a \int \zeta_{0}(x) \rho(\mathrm{d} x)}\right] \Lambda_{c}(\mathrm{~d} \rho)$. Since $\left.\Lambda_{c}=\alpha \mu_{1}+(1-\alpha) \mu_{2}, 11.6\right)$ implies that

$$
\begin{equation*}
\phi(1) \geq \alpha \phi\left(\frac{c_{1}}{c}\right)+(1-\alpha) \phi\left(\frac{c_{2}}{c}\right) . \tag{11.7}
\end{equation*}
$$

If $\Lambda_{c}=\delta_{0}$, then the extremality of $\Lambda_{c}$ is trivial; otherwise we can find $\zeta_{0} \in \mathcal{C}_{\mathbf{c}}(\mathbb{R})$ such that $\phi$ is strictly convex, which implies equality in (11.7) and hence $c_{1}=c_{2}=c$. Therefore we have equality in 11.7) for all $\zeta_{0} \in \mathcal{C}_{\mathbf{c}}(\mathbb{R})$, and we can then deduce from (11.6) that

$$
\int \mathbb{E}\left[e^{-\int \zeta_{0}(x) \rho_{0}(\mathrm{~d} x)}\right] \mu_{i}\left(\mathrm{~d} \rho_{0}\right)=\int \mathbb{E}\left[e^{-\int \zeta_{0}(x) \rho(\mathrm{d} x)}\right] \Lambda_{c}(\mathrm{~d} \rho),
$$

which implies that $\mu_{i}=\Lambda_{c}$.
We remark that Lemmas 11.3 and 11.4 can also be deduced from the convergence to invariant laws proved below using coupling. However, the proof by duality illustrates a useful tool.

### 11.2 Second moment calculations

Following [LS81] and Lig05], we first introduce for each $c>0$ the subset of translation invariant probability laws $\mathcal{T}_{c} \subset \mathcal{T}$, where $\Gamma \in \mathcal{T}$ is in $\mathcal{T}_{c}$ if and only if $\int \rho([0,1]) \Gamma(\mathrm{d} \rho)=c$, $\int \rho([0,1])^{2} \Gamma(\mathrm{~d} \rho)<\infty$, and

$$
\begin{equation*}
\int\left(\int \frac{1}{\sqrt{2 \pi t}} e^{-\frac{x^{2}}{2 t}} \rho(\mathrm{~d} x)-c\right)^{2} \Gamma(\mathrm{~d} \rho) \underset{t \rightarrow \infty}{\longrightarrow} 0 \tag{11.8}
\end{equation*}
$$

For a Howitt-Warren process with initial law $\mathcal{L}\left(\rho_{0}\right) \in \mathcal{T}_{c}$, we can perform second moment calculations for $\rho_{t}$ as $t \rightarrow \infty$ (see Lemma 11.6). Furthermore, if $\mathcal{L}\left(\rho_{0}\right) \in \mathcal{T}_{c}$, then any weak limit of $\rho_{t}$ as $t \rightarrow \infty$ is also in $\mathcal{T}_{c}$ (see Corollary 11.9), which will be crucial for proving the convergence of $\mathcal{L}\left(\rho_{t}\right)$ to the homogeneous ergodic law $\Lambda_{c}$. First we note that

Lemma 11.5 If $\Gamma \in \mathcal{T}_{\mathrm{e}}, \int \rho([0,1]) \Gamma(\mathrm{d} \rho)=c>0$ and $\int \rho([0,1])^{2} \Gamma(\mathrm{~d} \rho)<\infty$, then $\Gamma \in \mathcal{T}_{c}$. Conversely, any $\Gamma \in \mathcal{T}_{c}$ is a mixture of laws in $\mathcal{T}_{\mathrm{e}}$ satisfying the conditions above.

Proof. Our assumption implies that

$$
\begin{equation*}
\int\left(\frac{\rho([-L, L])}{2 L}-c\right)^{2} \Gamma(\mathrm{~d} \rho) \underset{L \rightarrow \infty}{\longrightarrow} 0 \tag{11.9}
\end{equation*}
$$

by the $L_{2}$ ergodic theorem. By the layercake representation,

$$
\begin{aligned}
\int_{\mathbb{R}} \frac{1}{\sqrt{2 \pi t}} e^{-\frac{x^{2}}{2 t}} \rho(\mathrm{~d} x)-c & =\frac{1}{\sqrt{2 \pi}} \int_{\mathbb{R}} e^{-\frac{y^{2}}{2}}\left(\frac{\rho(\sqrt{t} \mathrm{~d} y)}{\sqrt{t}}-c \mathrm{~d} y\right) \\
& =\frac{1}{\sqrt{2 \pi}} \int_{\mathbb{R}} \int_{0}^{1} 1_{\left\{z<e^{-\frac{y^{2}}{2}}\right\}} \mathrm{d} z\left(\frac{\rho(\sqrt{t} \mathrm{~d} y)}{\sqrt{t}}-c \mathrm{~d} y\right) \\
& =\frac{1}{\sqrt{2 \pi}} \int_{0}^{1}\left(\frac{\rho([-\sqrt{-2 t \ln z}, \sqrt{-2 t \ln z]})}{2 \sqrt{-2 t \ln z}}-c\right) 2 \sqrt{-2 \ln z} \mathrm{~d} z
\end{aligned}
$$

where we note that $\frac{2 \sqrt{-2 \ln z}}{\sqrt{2 \pi}} \mathrm{~d} z$ is a probability distribution on $[0,1]$ independent of $t$. Substituting this representation into the left hand side of (11.8), applying the Hölder inequality with respect to $\frac{2 \sqrt{-2 \ln z}}{\sqrt{2 \pi}} \mathrm{~d} z$, and applying 11.9 then proves 11.8. A more general argument using Bochner's theorem can be found in the proof of [Lig73, Theorem 5.6] or Lig05, Corollary II.8.20]. The second statement in Lemma 11.5 follows from the ergodic decomposition of $\Gamma \in \mathcal{T}_{c}$.

Lemma 11.6 (Second moment calculation) Let $\left(\rho_{t}\right)_{t \geq 0}$ be a Howitt-Warren process with drift $\beta \in \mathbb{R}$ and characteristic measure $\nu \neq 0$. If $\mathcal{L}\left(\rho_{0}\right) \in \mathcal{T}_{1}$, then for all $\phi, \psi \in C_{c}(\mathbb{R})$, we have

$$
\begin{equation*}
\lim _{t \rightarrow \infty} \mathbb{E}\left[\int \phi(x) \rho_{t}(\mathrm{~d} x) \int \psi(y) \rho_{t}(\mathrm{~d} y)\right]=\int \phi(x) \mathrm{d} x \int \psi(y) \mathrm{d} y+\frac{\int \phi(x) \psi(x) \mathrm{d} x}{2 \nu([0,1])} \tag{11.10}
\end{equation*}
$$

Proof. We may assume $\psi, \psi \in C_{c}(\mathbb{R})$ are non-negative. Since such functions can be approximated from above and below by finite linear combinations of indicator functions of finite intervals, it suffices to prove 11.10) for $\phi=1_{I_{1}}$ and $\psi=1_{I_{2}}$ for some finite intervals $I_{1}$ and $I_{2}$. Since $2 \phi \psi=(\phi+\psi)^{2}-\phi^{2}-\psi^{2}$, it suffices to consider only $I_{1}=I_{2}$, and we may even take $I_{1}=[0,1]$, so that (11.10) reduces to showing

$$
\begin{equation*}
\lim _{t \rightarrow \infty} \mathbb{E}\left[\rho_{t}([0,1])^{2}\right]=1+\frac{1}{2 \nu([0,1])} \tag{11.11}
\end{equation*}
$$

By Theorem 3.7, we have the representation

$$
\begin{equation*}
\mathbb{E}\left[\rho_{t}([0,1])^{2}\right]=\mathbb{E}\left[\iint \mathbb{Q}^{\otimes 2}\left(\left(\pi_{(x, 0)}^{1,+}(t), \pi_{(y, 0)}^{2,+}(t)\right) \in[0,1]^{2}\right) \rho_{0}(\mathrm{~d} x) \rho_{0}(\mathrm{~d} y)\right] \tag{11.12}
\end{equation*}
$$

where $\mathbb{Q}^{\otimes 2}$ denotes the 2-fold product measure, and $\pi_{(\cdot, \cdot)}^{1,+}$ resp. $\pi_{(\cdot,)}^{2,+}$ are rightmost elements in two independent sample webs $\mathcal{W}^{1}$ resp. $\mathcal{W}^{2}$, both with quenched law $\mathbb{Q}$. With respect to $\mathbb{E}\left[\mathbb{Q}^{\otimes 2}\right],\left(\pi^{1,+}(x, 0), \pi_{(y, 0)}^{2,+}\right)$ is the two-point motion of the Howitt-Warren flow with drift $\beta$ and characteristic measure $\nu$, and hence solves the Howitt-Warren martingale problem under conditions 2.6 and 2.9. In particular, $R_{t}:=\pi_{(y, 0)}^{2,+}(t)-\pi_{(x, 0)}^{1,+}(t)$ is an autonomous Brownian motion with stickiness at the origin, and conditional on $\left(R_{t}\right)_{t \geq 0}, S_{t}:=\pi_{(x, 0)}^{1,+}(t)+\pi_{(y, 0)}^{2,+}(t)$ is distributed as a time change of an independent Brownian motion with drift $2 \beta$. We leave the verification of this statement as an exercise to the reader. A similar statement for a pair of Brownian motions satisfying (4.1) can be found in [SS08, Lemma 2.2].

Let $\bar{\rho}_{0}^{\otimes 2}:=\mathbb{E}\left[\rho_{0}^{\otimes 2}\right]$, and let $\hat{\rho}_{0}^{\otimes 2}$ denote the image measure of $\bar{\rho}_{0}^{\otimes 2}$ under the change of coordinates $(x, y) \rightarrow(r, s):=(y-x, x+y)$. Then by Lemma 11.7 below, $\hat{\rho}_{0}^{\otimes 2}(\mathrm{~d} r \mathrm{~d} s)=\alpha(\mathrm{d} r) \mathrm{d} s$. Therefore we can rewrite 11.12) as

$$
\begin{align*}
\mathbb{E}\left[\rho_{t}([0,1])^{2}\right] & =\iint \mathbb{P}_{(r, s)}\left(\left|R_{t}\right| \leq 1,\left|R_{t}\right| \leq S_{t} \leq 2-\left|R_{t}\right|\right) \alpha(\mathrm{d} r) \mathrm{d} s \\
& =2 \int \mathbb{E}_{r}\left[\left(1-\left|R_{t}\right|\right) 1_{\left\{\left|R_{t}\right| \leq 1\right\}}\right] \alpha(\mathrm{d} r) \\
& =2 \int_{0}^{1} \int_{\mathbb{R}} \mathbb{P}_{r}\left(\left|R_{t}\right| \leq a\right) \alpha(\mathrm{d} r) \mathrm{d} a, \tag{11.13}
\end{align*}
$$

where $\mathbb{P}_{(r, s)}$ denotes probability for $\left(R_{t}, S_{t}\right)$ starting at $(r, s)$, and we have used the fact that conditioned on $\left(R_{t}\right)_{t \geq 0},\left(S_{t}\right)_{t \geq 0}$ with differential initial conditions can be coupled together simply by translation.

In (11.13), let $f_{t, a}(r):=\mathbb{P}_{r}\left(\left|R_{t}\right| \leq a\right)$, which is even, and strictly decreasing on $[0, \infty)$. The latter follows from the fact that $\left|R_{t}\right|$ is a reflected Brownian motion with stickiness at the origin, and there is a natural coupling through coalescence for $\left|R_{t}\right|$ starting at differential initial conditions on $[0, \infty)$. By the layercake representation, we may rewrite 11.13) as

$$
\begin{align*}
\mathbb{E}\left[\rho_{t}([0,1])^{2}\right] & =2 \int_{0}^{1} \int_{\mathbb{R}} \int_{0}^{\infty} 1_{\left\{y \leq f_{t, a}(r)\right\}} \mathrm{d} y \alpha(\mathrm{~d} r) \mathrm{d} a \\
& =2 \int_{0}^{1} \int_{0}^{\infty} \alpha\left(\left[-f_{t, a}^{-1}(y), f_{t, a}^{-1}(y)\right]\right) \mathrm{d} y \mathrm{~d} a \\
& =2 \int_{0}^{1} \int_{0}^{\infty}-\frac{\alpha([-r, r])}{r} r \mathrm{~d} f_{t, a}(r) \mathrm{d} a \tag{11.14}
\end{align*}
$$

Note that $-r \mathrm{~d} f_{t, a}(r)$ is a finite measure on $(0, \infty)$ with total mass $\int_{0}^{\infty} f_{t, a}(r) \mathrm{d} r$, and for any $u>0$, integrating by parts gives

$$
\int_{0}^{u}-r \mathrm{~d} f_{t, a}(r)=-u f_{t, a}(u)+\int_{0}^{u} f_{t, a}(r) \mathrm{d} r \underset{t \rightarrow \infty}{\longrightarrow} 0
$$

since $f_{t, a}(r)$ is decreasing on $[0, \infty)$ and $f_{t, a}(0) \rightarrow 0$ as $t \rightarrow \infty$ by basic properties of $\left|R_{t}\right|$. Therefore the sequence of measures $-r \mathrm{~d} f_{t, a}(r)$ shifts its mass to $\infty$ as $t \rightarrow \infty$. Since $\lim _{r \rightarrow \infty} \frac{\alpha([-r, r])}{r}=1$ by Lemma 11.7 below, we deduce from 11.13) that

$$
\begin{equation*}
\lim _{t \rightarrow \infty} \mathbb{E}\left[\rho_{t}([0,1])^{2}\right]=2 \int_{0}^{1} \lim _{t \rightarrow \infty} \int_{0}^{\infty} f_{t, a}(r) \mathrm{d} r \mathrm{~d} a \tag{11.15}
\end{equation*}
$$

Since $\mathrm{d} r+\frac{1}{4 \nu([0,1])} \delta_{0}(r)$ is an invariant measure for $\left|R_{t}\right|$ by Lemma 11.8 , for any $t \geq 0$,

$$
\frac{1}{4 \nu([0,1])} f_{t, a}(0)+\int_{0}^{\infty} f_{t, a}(r) \mathrm{d} r=\frac{1}{4 \nu([0,1])} \mathbb{P}_{0}\left(\left|R_{t}\right| \leq a\right)+\int_{0}^{\infty} \mathbb{P}_{r}\left(\left|R_{t}\right| \leq a\right)=\frac{1}{4 \nu([0,1])}+a
$$

Since $f_{t, a}(0) \rightarrow 0$ as $t \rightarrow \infty$, we obtain $\lim _{t \rightarrow \infty} \int_{0}^{\infty} f_{t, a}(r) \mathrm{d} r=\frac{1}{4 \nu([0,1])}+a$. Substituting this into (11.15) then gives (11.11).

Lemma 11.7 (Ergodicity of the second moment measure) Let $\mathcal{L}(\rho) \in \mathcal{T}_{1}$. Let $\bar{\rho}^{\otimes 2}:=$ $\mathbb{E}\left[\rho^{\otimes 2}\right]$ denote the second moment measure of $\rho$, and let $\hat{\rho}^{\otimes 2}$ denote the image measure of $\bar{\rho}^{\otimes 2}$ under the change of coordinates $(x, y) \rightarrow(r, s):=(y-x, x+y)$. Then $\hat{\rho}^{\otimes 2}(\mathrm{~d} r \mathrm{~d} s)=\alpha(\mathrm{d} r) \mathrm{d} s$, where $\alpha(A)=\hat{\rho}^{\otimes 2}(A \times[0,1])$ for all $A \in \mathcal{B}(\mathbb{R})$, and $\lim _{L \rightarrow \infty} \frac{1}{L} \alpha([-L, L])=1$.

Proof. By translation invariance in law of $\rho, \bar{\rho}^{\otimes 2}(\cdot)=\bar{\rho}^{\otimes 2}(\cdot+(a, a))$ for all $a \in \mathbb{R}$. Therefore $\hat{\rho}^{\otimes 2}(\mathrm{~d} r \mathrm{~d} s)$ is translation invariant in $s$, and hence $\hat{\rho}^{\otimes 2}$ admits the desired factorization. Therefore

$$
\alpha([-L, L])=\int 1_{\{[-L, L]\}}(r) 1_{\{[0,1]\}}(s) \hat{\rho}^{\otimes 2}(\mathrm{~d} r \mathrm{~d} s)=\int 1_{\{[-L, L]\}}(r) 1_{\{[0,1]\}}(r+s) \hat{\rho}^{\otimes 2}(\mathrm{~d} r \mathrm{~d} s) .
$$

Transforming back into the variables $(x, y)$ and the measure $\bar{\rho}^{\otimes 2}$ then gives

$$
\alpha([-L, L])=\bar{\rho}^{\otimes 2}(\{(x, y): 0 \leq y \leq 1 / 2,|x-y| \leq L\})
$$

which is bounded between $\bar{\rho}^{\otimes}([-L+1, L-1] \times[0,1 / 2])$ and $\bar{\rho}^{\otimes}([-L-1, L+1] \times[0,1 / 2])$. Note that

$$
\begin{align*}
\frac{1}{L} \bar{\rho}^{\otimes}([-L+1, L-1] \times[0,1 / 2]) & =\frac{1}{L} \mathbb{E}[\rho([0,1 / 2]) \rho([-L+1, L-1])] \\
& \xrightarrow[L \rightarrow \infty]{\longrightarrow} 2 \mathbb{E}[\rho([0,1 / 2])]=1, \tag{11.16}
\end{align*}
$$

provided that $\frac{1}{2 L} \rho([-L+1, L-1]) \rightarrow 1$ in $L_{2}$. Indeed, by Lemma 11.5 , there exists a probability measure $\gamma(\mathrm{d} \Lambda)$ on $\mathcal{T}_{1} \cap \mathcal{T}_{\mathrm{e}}$ such that $\mathcal{L}(\rho)=\int_{\mathcal{T}_{1} \cap \mathcal{T}_{\mathrm{e}}} \Lambda \gamma(\mathrm{d} \Lambda)$. Therefore

$$
\begin{aligned}
& \mathbb{E}\left[\left(\frac{1}{2 L} \rho([-L+1, L-1])-1\right)^{2}\right] \\
= & \int_{\mathcal{T}_{1} \cap \mathcal{T}_{\mathrm{e}}} \int\left(\frac{1}{2 L} \rho([-L+1, L-1])-1\right)^{2} \Lambda(\mathrm{~d} \rho) \gamma(\mathrm{d} \Lambda) \underset{L \rightarrow \infty}{\longrightarrow} 0,
\end{aligned}
$$

since the integrand w.r.t. $\gamma(\mathrm{d} \Lambda)$ tends to $0 \gamma$ a.s. by the $L_{2}$ ergodic theorem applied to $\Lambda$, and is dominated uniformly in $L$ by $2+2 \int \rho([0,1])^{2} \Lambda(\mathrm{~d} \rho)$, which is integrable by the assumption that $\mathbb{E}\left[\rho([0,1])^{2}\right]<\infty$. This proves $\left.\sqrt{11.16}\right)$, and the same can be proved for $\bar{\rho}^{\otimes}([-L-1, L+$ $1] \times[0,1 / 2])$, the upper bound on $\alpha[-L, L]$. Therefore $\lim _{L \rightarrow \infty} \frac{1}{L} \alpha([-L, L])=1$.

Lemma 11.8 (Invariant measure for the two point motion) Let $\left(X_{t}, Y_{t}\right)$ be two coupled Brownian motions solving the Howitt-Warren martingale problem under conditions (2.6) and 2.9). Let $R_{t}=Y_{t}-X_{t}$. Then $\mathrm{d} r+\frac{1}{4 \nu(0,1])} \delta_{0}(r)$ on $[0, \infty)$ is an invariant measure for $\left|R_{t}\right|$, and $\mathrm{d} x \mathrm{~d} y+\frac{1}{2 \nu([0,1])} \delta_{x}(y) \mathrm{d} x$ on $\mathbb{R}^{2}$ is an invariant measure for $\left(X_{t}, Y_{t}\right)$.

Proof. Note that $\left|R_{t}\right|$ is uniquely characterized in law by the following two properties: (1) $\left|R_{t}\right|-4 \nu([0,1]) \int_{0}^{t} 1_{\left\{\left|R_{s}\right|=0\right\}} \mathrm{d} s$ is a martingale; (2) $\langle | R_{t}\left|,\left|R_{t}\right|\right\rangle=2 \int_{0}^{t} 1_{\left\{\left|R_{s}\right| \neq 0\right\}} \mathrm{d} s$. These two properties are clearly satisfied by the solution of the following SDE

$$
\begin{equation*}
\mathrm{d} Z_{t}=1_{\left\{Z_{t} \neq 0\right\}} \sqrt{2} \mathrm{~d} B_{t}+1_{\left\{Z_{t}=0\right\}} 4 \nu([0,1]) \mathrm{d} t \tag{11.17}
\end{equation*}
$$

where $B_{t}$ is a standard Brownian motion, and $Z_{t}$ is constrained to be non-negative. For the existence and pathwise uniqueness of the solution to this SDE, see e.g. [SS08, Lemma 2.2]. The solution of 11.17) generates a Feller semigroup $\left(S_{t}\right)_{t \geq 0}$ on the Banach space $C_{0}([0, \infty)$ ), the space of continuous functions on $[0, \infty)$ which vanish at $\infty$ and equipped with the supremum norm. By Itô's formula, the generator $L$ for $S_{t}$ is given by

$$
\begin{equation*}
L f(x)=1_{\{x \neq 0\}} f^{\prime \prime}(x)+1_{\{x=0\}} 4 \nu([0,1]) f^{\prime}(x) . \tag{11.18}
\end{equation*}
$$

Let $\mathcal{D} \subset C_{0}([0, \infty))$ denote the domain of $L$. If $f \in \mathcal{D}$, then $L f \in C_{0}([0, \infty))$. In particular, we must have $f^{\prime \prime}(0)=4 \nu([0,1]) f^{\prime}(0)$ and $f^{\prime \prime} \in C_{0}([0, \infty))$. Together with $f \in C_{0}([0, \infty))$, this also implies that $f^{\prime} \in C_{0}([0, \infty))$. Conversely, if $f \in D:=\left\{f \in C_{0}([0, \infty)): f^{\prime}, f^{\prime \prime} \in\right.$ $C_{0}([0, \infty))$ and $\left.f^{\prime \prime}(0)=4 \nu([0,1]) f^{\prime}(0)\right\}$, then it is not difficult to see from Itô's formula that $f \in \mathcal{D}$. Therefore $\mathcal{D}=D$. If we denote $\mu(\mathrm{d} x)=\mathrm{d} x+\frac{1}{4 \nu([0,1])} \delta_{0}(x)$, then we have

$$
\begin{equation*}
\int_{0}^{\infty} L f(x) \mu(\mathrm{d} x)=\int_{0}^{\infty} f^{\prime \prime}(x) \mathrm{d} x+f^{\prime}(0)=0 \quad \text { for all } f \in \mathcal{D} . \tag{11.19}
\end{equation*}
$$

From (11.19), we can deduce that $\mu$ is an invariant measure for $\left(\left|R_{t}\right|\right)_{t \geq 0}$. Indeed, for any $f \in \mathcal{D} \cap C_{c}([0, \infty))$ with compact support and for any $t>0$, we have

$$
\begin{align*}
& \int_{0}^{\infty} S_{t} f(x) \mu(\mathrm{d} x)-\int_{0}^{\infty} f(x) \mu(\mathrm{d} x)=\int_{0}^{\infty}\left(S_{t} f(x)-f(x)\right) \mu(\mathrm{d} x) \\
= & \int_{0}^{\infty} \int_{0}^{t} \frac{\mathrm{~d}}{\mathrm{~d} s} S_{s} f(x) \mathrm{d} s \mu(\mathrm{~d} x)=\int_{0}^{\infty} \int_{0}^{t} L S_{s} f(x) \mathrm{d} s \mu(\mathrm{~d} x) \\
= & \int_{0}^{t} \int_{0}^{\infty} L S_{s} f(x) \mu(\mathrm{d} x) \mathrm{d} s=0, \tag{11.20}
\end{align*}
$$

where we have used Fubini based on the fact that $L S_{s} f(x)=S_{s} L f(x)$ is decaying superexponentially in $x$ because $L f$ has compact support and $\left|R_{t}\right|$ is distributed as a Brownian motion on $(0, \infty)$; and in 11.20, we have applied 11.19) using the fact that $f \in \mathcal{D}$ implies $S_{t} f \in \mathcal{D}$. Since 11.20 holds for all $f \in \mathcal{D} \cap C_{c}([0, \infty))$, which is a measure determining class, $\mu$ is an invariant measure.

The symmetry of $R_{t}$ implies that $\tilde{\mu}(\mathrm{d} r):=\mathrm{d} r+\frac{1}{2 \nu([0,1])} \delta_{0}(r)$ on $\mathbb{R}$ is an invariant measure for $R_{t}$, and the translation invariance of $\left(X_{t}, Y_{t}\right)$ along the diagonal implies that $\tilde{\mu}(\mathrm{d} r) \mathrm{d} s$ is an invariant measure for $\left(R_{t}, S_{t}\right)$ with $R_{t}=Y_{t}-X_{t}$ and $S_{t}=Y_{t}+X_{t}$. A change of coordinates then verifies that $\mathrm{d} x \mathrm{~d} y+\frac{1}{2 \nu([0,1])} \delta_{x}(y) \mathrm{d} x$ is an invariant measure for $\left(X_{t}, Y_{t}\right)$.

Remark. In LS81 and Lig05, the analogue of Lemma 11.6 is proved by treating the two-point motion as a perturbation of two independent one-point motions. This approach requires exact calculations involving the two-point motion and is not clear how to implement in the continuous space setting. Our approach reduces the task to first identifying the invariant measure for the two-point motion, which when integrated over the test function $\phi(x) \psi(y)$ gives the RHS of (11.10), and then using qualitative properties of the two-point motion together with the ergodicity of the initial condition to remove the dependence on the initial condition.

The following corollary of Lemma 11.6 is the analogue of [LS81, Lemma (5.3)(b)] for our model, and will be crucial in proving convergence to the homogeneous invariant laws.

Corollary 11.9 (Preservation of $\mathcal{T}_{1}$ ) Assume the same conditions as in Lemma 11.6. Then the law of any subsequential weak limit of $\left(\rho_{t}\right)_{t \geq 0}$ is also in $\mathcal{T}_{1}$.
Proof. Let $\rho^{*}$ be the weak limit of $\rho_{t_{n}}$ along a subsequence $t_{n} \uparrow \infty$. Clearly $\mathcal{L}\left(\rho^{*}\right) \in \mathcal{T}$. Lemma 11.6 implies the uniform integrability of $\left(\rho_{t_{n}}([0,1])\right)_{n \in \mathbb{N}}$, and hence $\mathbb{E}\left[\rho^{*}([0,1])\right]=1$ since $\mathbb{E}\left[\rho_{t}([0,1])\right]=1$ for all $t \geq 0$. By Fatou's lemma, Lemma 11.6 also implies that for all non-negative $\phi, \psi \in C_{c}(\mathbb{R})$,

$$
\begin{equation*}
\mathbb{E}\left[\int \phi(x) \rho^{*}(\mathrm{~d} x) \int \psi(y) \rho^{*}(\mathrm{~d} y)\right] \leq \int \phi(x) \mathrm{d} x \int \psi(y) \mathrm{d} y+\frac{\int \phi(x) \psi(x) \mathrm{d} x}{2 \nu([0,1])} \tag{11.21}
\end{equation*}
$$

Approximating $\frac{1}{\sqrt{2 \pi t}} e^{-\frac{x^{2}}{2 t}}$ from below by functions with compact support, we then have

$$
\mathbb{E}\left[\left(\int \frac{1}{\sqrt{2 \pi t}} e^{-\frac{x^{2}}{2 t}} \rho^{*}(\mathrm{~d} x)\right)^{2}\right] \leq 1+\frac{1}{4 \nu([0,1]) \sqrt{\pi t}} .
$$

Together with $\mathbb{E}\left[\int \frac{1}{\sqrt{2 \pi t}} e^{-\frac{x^{2}}{2 t}} \rho^{*}(\mathrm{~d} x)\right]=1$, this implies 11.8 with $c=1$ and $\Gamma=\mathcal{L}\left(\rho^{*}\right)$. Therefore $\mathcal{L}\left(\rho^{*}\right) \in \mathcal{T}_{1}$.

### 11.3 Coupling and convergence

The definition of the Howitt-Warren process $\left(\rho_{t}\right)_{t \geq 0}$ from the kernels $\left(K_{s, t}^{+}\right)_{s<t}$ constructed in Theorem 3.7 gives a natural coupling between $\left(\rho_{t}\right)_{t \geq 0}$ with different initial conditions. This coupling is monotone in the sense that if $\rho_{0}^{1}(A) \geq \rho_{0}^{2}(A)$ for all $A \in \mathcal{B}(\mathbb{R})$, which we denote by $\rho_{0}^{1} \succ \rho_{0}^{2}$, then $\rho_{t}^{1} \succ \rho_{t}^{2}$ a.s. for all $t>0$. Through this coupling, we will prove the weak convergence of $\rho_{t}$ to a mixture of homogeneous invariant laws under suitable assumption on $\mathcal{L}\left(\rho_{0}\right)$. The first observation is the following.

Lemma 11.10 (Coupled Howitt-Warren processs) Let $\left(\rho_{t}^{1}\right)_{t \geq 0}$ and $\left(\rho_{t}^{2}\right)_{t \geq 0}$ be HowittWarren processes with drift $\beta$ and characteristic measure $\nu$, defined from the same HowittWarren flow $\left(K_{s, t}^{+}\right)_{s<t}$ as in (2.1). Assume that $\mathcal{L}\left(\rho_{0}^{1}\right), \mathcal{L}\left(\rho_{0}^{2}\right) \in \mathcal{T}$ and $\mathbb{E}\left[\rho_{0}^{1}([0,1])\right]<\infty$, $\mathbb{E}\left[\rho_{0}^{2}([0,1])\right]<\infty$. Then any weak limit point $\left(\rho^{1 *}, \rho^{2 *}\right)$ of $\left(\rho_{t}^{1}, \rho_{t}^{2}\right)_{t \geq 0}$ as $t \rightarrow \infty$ satisfies $\mathbb{P}\left(\rho^{1 *} \succ \rho^{2 *}\right.$ or $\left.\rho^{2 *} \succ \rho^{1 *}\right)=1$.

Proof. Let $\left(\rho_{t}^{1}-\rho_{t}^{2}\right)=\left(\rho_{t}^{1}-\rho_{t}^{2}\right)^{+}-\left(\rho_{t}^{1}-\rho_{t}^{2}\right)^{-}$denote the Jordan decomposition of $\rho_{t}^{1}-\rho_{t}^{2}$, and let $\left|\rho_{t}^{1}-\rho_{t}^{2}\right|:=\left(\rho_{t}^{1}-\rho_{t}^{2}\right)^{+}+\left(\rho_{t}^{1}-\rho_{t}^{2}\right)^{-}$denote the total variation measure of $\rho_{t}^{1}-\rho_{t}^{2}$. Recall the quenched law $\mathbb{Q}$ from (3.24). For any $0<s<t$, almost surely we have

$$
\left|\rho_{t}^{1}-\rho_{t}^{2}\right|=\left|\int\left(\rho_{s}^{1}-\rho_{s}^{2}\right)(\mathrm{d} x) \mathbb{Q}\left[\pi_{(x, s)}^{+}(t) \in \cdot\right]\right| \prec \int\left|\rho_{s}^{1}-\rho_{s}^{2}\right|(\mathrm{d} x) \mathbb{Q}\left[\pi_{(x, s)}^{+}(t) \in \cdot\right] .
$$

Therefore

$$
\begin{aligned}
\mathbb{E}\left[\left|\rho_{t}^{1}-\rho_{t}^{2}\right|([0,1])\right] & \leq \mathbb{E}\left[\int\left|\rho_{s}^{1}-\rho_{s}^{2}\right|(\mathrm{d} x) \mathbb{Q}\left[\pi_{(x, s)}^{+}(t) \in[0,1]\right]\right] \\
& =\int \mathbb{E}\left[\left|\rho_{s}^{1}-\rho_{s}^{2}\right|\right](\mathrm{d} x) \mathbb{E}\left[\mathbb{Q}\left[\pi_{(x, s)}^{+}(t) \in[0,1]\right]\right]=\mathbb{E}\left[\left|\rho_{s}^{1}-\rho_{s}^{2}\right|([0,1])\right],
\end{aligned}
$$

where we have used the independence between $\rho_{s}^{1}, \rho_{s}^{2}$ and $\left(\mathbb{Q}\left[\pi_{(x, s)}^{+}(t) \in[0,1]\right]\right)_{x \in \mathbb{R}}$, the translation invariance of $\mathbb{E}\left|\rho_{s}^{1}-\rho_{s}^{2}\right|$, and the fact that $\pi_{(x, s)}^{+}$is distributed as a standard Brownian motion starting at $x$ at time $s$ under the law $\mathbb{E} \mathbb{Q}$. Therefore $\mathbb{E}\left[\left|\rho_{t}^{1}-\rho_{t}^{2}\right|([0,1])\right]$ decreases monotonically to a non-negative limit as $t \uparrow \infty$.

Note that the Lemma follows once we show that for any $\varepsilon>0$ and any $\phi, \psi \in C_{c}(\mathbb{R})$ with $0 \leq \phi, \psi \leq 1$, we have

$$
\begin{equation*}
\mathbb{P}\left(\int \phi(x) \rho_{t}^{1}(\mathrm{~d} x)-\int \phi(x) \rho_{t}^{2}(\mathrm{~d} x)>\varepsilon \text { and } \int \psi(x) \rho_{t}^{2}(\mathrm{~d} x)-\int \psi(x) \rho_{t}^{1}(\mathrm{~d} x)>\varepsilon\right) \underset{t \rightarrow \infty}{\longrightarrow} 0 \tag{11.22}
\end{equation*}
$$

Suppose that (11.22) fails so that for some $\varepsilon>0$ and $\phi, \psi \in C_{c}(\mathbb{R})$ with $0 \leq \phi, \psi \leq 1$, the probability in 11.22 is bounded uniformly from below by $\delta>0$ along a sequence $t_{i} \uparrow \infty$.

Choose $L>0$ large such that $\phi$ and $\psi$ vanish outside $[-L, L]$. Given $\rho_{t_{i}}^{1}$ and $\rho_{t_{i}}^{2}$ satisfying the conditions in the probability in 11.22 , we have

$$
\varepsilon<\int \phi(x)\left(\rho_{t_{i}}^{1}-\rho_{t_{i}}^{2}\right)(\mathrm{d} x) \leq \int \phi(x)\left(\rho_{t_{i}}^{1}-\rho_{t_{i}}^{2}\right)^{+}(\mathrm{d} x) \leq\left(\rho_{t_{i}}^{1}-\rho_{t_{i}}^{2}\right)^{+}([-L, L])
$$

and similarly $\left(\rho_{t_{i}}^{1}-\rho_{t_{i}}^{2}\right)^{-}([-L, L])>\varepsilon$. For such a realization of $\rho_{t_{i}}^{1}$ and $\rho_{t_{i}}^{2}$,

$$
\begin{aligned}
& \left|\rho_{t_{i}+1}^{1}-\rho_{t_{i}+1}^{2}\right|([0,1]) \\
\leq & \int\left|\rho_{t_{i}}^{1}-\rho_{t_{i}}^{2}\right|(\mathrm{d} x) \mathbb{Q}\left[\pi_{\left(x, t_{i}\right)}^{+}\left(t_{i}+1\right) \in[0,1]\right]-2 \varepsilon \mathbb{Q}\left[\pi_{\left(-L, t_{i}\right)}^{+}\left(t_{i}+1\right)\right. \\
= & \left.\pi_{\left(L, t_{i}\right)}^{+}\left(t_{i}+1\right) \in[0,1]\right],
\end{aligned}
$$

where we observed that the mass assigned by $\left(\rho_{t_{i}}^{1}-\rho_{t_{i}}^{2}\right)^{+}$and $\left(\rho_{t_{i}}^{1}-\rho_{t_{i}}^{2}\right)^{-}$to $[-L, L]$ are carried by $\left(\pi_{\left(x, t_{i}\right)}^{+}\left(t_{i}+1\right)\right)_{x \in \mathbb{R}}$ to the same point in $[0,1]$ when $\pi_{\left(-L, t_{i}\right)}^{+}\left(t_{i}+1\right)=\pi_{\left(L, t_{i}\right)}^{+}\left(t_{i}+1\right) \in[0,1]$. Therefore

$$
\begin{equation*}
\mathbb{E}\left[\left|\rho_{t_{i}+1}^{1}-\rho_{t_{i}+1}^{2}\right|([0,1])\right] \leq \mathbb{E}\left[\left|\rho_{t_{i}}^{1}-\rho_{t_{i}}^{2}\right|([0,1])\right]-2 \varepsilon \delta h, \tag{11.23}
\end{equation*}
$$

where $h=\mathbb{E} \mathbb{Q}\left[\pi_{\left(-L, t_{i}\right)}^{+}\left(t_{i}+1\right)=\pi_{\left(L, t_{i}\right)}^{+}\left(t_{i}+1\right) \in[0,1]\right]>0$ is independent of $t_{i}$. Since 11.23) holds for all $t_{i}$, this contradicts the fact that $\mathbb{E}\left[\left|\rho_{t}^{1}-\rho_{t}^{2}\right|([0,1])\right]$ decreases monotonically to a non-negative limit as $t \uparrow \infty$.

Lemma 11.11 (Convergence to $\Lambda_{c}$ ) Let $\left(\rho_{t}\right)_{t \geq 0}$ be a Howitt-Warren process with drift $\beta$ and characteristic measure $\nu \neq 0$. If $\mathcal{L}\left(\rho_{0}\right) \in \mathcal{T}_{\mathrm{e}}$ and $\mathbb{E}\left[\rho_{0}([0,1])\right]=c<\infty$, then $\rho_{t}$ converges weakly to $\Lambda_{c}$, which was defined in Lemma 11.3. If $\mathcal{L}\left(\rho_{0}\right) \in \mathcal{T}_{\mathrm{e}}$ and $\mathbb{E}\left[\rho_{0}([0,1])\right]=\infty$, then $\rho_{t}$ has no weak limit which is supported on $\mathcal{M}_{\mathrm{loc}}(\mathbb{R})$.

Proof. Without loss of generality, assume $c=1$. Let $\left(\rho_{t}^{2}\right)_{t \geq 0}$ be a Howitt-Warren process with initial condition $\rho_{0}^{2}$ such that $\mathcal{L}\left(\rho_{0}^{2}\right)=\Lambda_{1}$, and let $\rho_{t}$ and $\rho_{t}^{2}$ be defined from the same Howitt-Warren flow $\left(K_{s, t}^{+}\right)_{s<t}$. As in the proof of Lemma 11.3, we note that $\left(\rho_{t}\right)_{t \geq 0}$ (resp. $\left.\left(\rho_{t}, \rho_{t}^{2}\right)_{t \geq 0}\right)$ is a tight family of $\mathcal{M}_{\mathrm{g}}(\mathbb{R})$ (resp. $\left.\mathcal{M g}_{\mathrm{g}}(\mathbb{R})^{2}\right)$ valued random variable. Therefore any subsequential weak limit $\Gamma \in \mathcal{M}_{1}\left(\mathcal{M}_{\mathrm{g}}(\mathbb{R})\right)$ of $\mathcal{L}\left(\rho_{t}\right)$ as $t \rightarrow \infty$ can be realized as the marginal law of the first component of a random couple $\left(\rho^{*}, \rho^{2 *}\right) \in \mathcal{M}_{\mathrm{g}}(\mathbb{R})^{2}$, which arises as a subsequential weak limit of $\left(\rho_{t}, \rho_{t}^{2}\right)_{t \geq 0}$. By Lemma 11.10, $\mathbb{P}\left(\rho^{*} \succ \rho^{2 *}\right.$ or $\left.\rho^{2 *} \succ \rho^{*}\right)=1$. Therefore for any rational $a<b$, by the translation invariance in law of $\left(\rho^{*}-\rho^{2 *}\right) 1_{\left\{\rho^{*} \succ \rho^{2 *}\right\}}$ and $\left(\rho^{2 *}-\rho^{*}\right) 1_{\left\{\rho^{2 *} \succ \rho^{*}\right\}}$, we have

$$
\begin{align*}
& \mathbb{E}\left|\rho^{*}([a, b))-\rho^{2 *}([a, b))\right|=E\left|\int \frac{b-a}{\sqrt{2 \pi t}} e^{-\frac{x^{2}}{2 t}}\left(\rho^{*}-\rho^{2 *}\right)(\mathrm{d} x)\right| \\
\leq & (b-a) \mathbb{E}\left|\int \frac{1}{\sqrt{2 \pi t}} e^{-\frac{x^{2}}{2 t}} \rho^{*}(\mathrm{~d} x)-1\right|+(b-a) \mathbb{E}\left|\int \frac{1}{\sqrt{2 \pi t}} e^{-\frac{x^{2}}{2 t}} \rho^{2 *}(\mathrm{~d} x)-1\right| . \tag{11.24}
\end{align*}
$$

If we first restrict ourselves to the case $\mathbb{E}\left[\rho_{0}([0,1])^{2}\right]<\infty$, then $\mathcal{L}\left(\rho_{0}\right) \in \mathcal{T}_{1}$ by Lemma 11.5 , and by Corollary 11.9, $\mathcal{L}\left(\rho^{*}\right) \in \mathcal{T}_{1}$ and $\mathcal{L}\left(\rho^{2 *}\right)=\Lambda_{1} \in \mathcal{T}_{1}$. Definition of $\mathcal{T}_{1}$ implies that both terms in (11.24) vanish as $t \rightarrow \infty$, and hence $\rho^{*}([a, b))=\rho^{2 *}([a, b))$ a.s. for all rational $a<b$. Since $\{[a, b)\}_{a<b \in \mathbb{Q}}$ is measure determining, we have $\rho^{*}=\rho^{2 *}$ a.s., and hence $\mathcal{L}\left(\rho_{t}\right)$ converges weakly to $\Lambda_{1}$.

If $\mathbb{E}\left[\rho_{0}([0,1])^{2}\right]=\infty$, then we can approximate $\rho_{0}$ by $\left(\rho_{0}^{n}\right)_{n \in \mathbb{N}}$ with $\mathcal{L}\left(\rho_{0}^{n}\right) \in \mathcal{T}_{\text {e }}$ such that $\mathbb{E}\left[\rho_{0}^{n}([0,1])^{2}\right]<\infty$ and $\rho_{0}^{n}$ increases monotonically to $\rho_{0}$ almost surely. For instance, given $\rho_{0}$,
we can sample a uniform random variable $U$ on $[0,1]$ and then define $\rho_{0}^{n}$ on $[U+k, U+k+1)$ for each $k \in \mathbb{Z}$ by $\rho_{0}^{n}=\rho_{0}$ on $[U+k, U+k+1)$ if $\rho_{0}([U+k, U+k+1)) \leq n$, and set $\rho_{0}^{n}=0$ on $[U+k, U+k+1)$ otherwise. Then $\mathbb{E}\left[\rho_{0}^{n}([0,1])\right]=1-\varepsilon_{n}$ for some $\varepsilon_{n} \downarrow 0$. Our argument above shows that $\mathcal{L}\left(\rho_{t}^{n}\right)$ converges weakly to $\Lambda_{1-\varepsilon_{n}}$. Since $\rho_{0} \succ \rho_{0}^{n}$ a.s. for all $n \in \mathbb{N}$, any weak limit point $\Gamma$ of $\mathcal{L}\left(\rho_{t}\right)_{t \geq 0}$ stochastically dominates $\Lambda_{c}$ for all $c<1$. Since $\int \rho([0,1]) \Gamma(\mathrm{d} \rho) \leq 1$ by Fatou, we must have $\Gamma=\Lambda_{1}$.

If $\mathcal{L}\left(\rho_{0}\right) \in \mathcal{T}_{\mathrm{e}}$ and $\mathbb{E}\left[\rho_{0}([0,1])\right]=\infty$, then by the same argument as above, any weak limit point of $\rho_{t}$ stochastically dominates $\Lambda_{c}$ for all $c>0$, which is not possible for an $\mathcal{M}_{\mathrm{loc}}(\mathbb{R})$ valued random variable since $\Lambda_{1}$ is not concentrated on the zero measure by our assumption $\nu \neq 0$.

From Lemma 11.11, we can deduce that
Lemma 11.12 (Extremal measures in $\mathcal{I} \cap \mathcal{T}$ ) For the Howitt-Warren processs with drift $\beta$ and characteristic measure $\nu \neq 0$, we have $(\mathcal{I} \cap \mathcal{T})_{\mathrm{e}}=\left\{\Lambda_{c}: c \geq 0\right\}$.

Proof. If $\mathcal{L}\left(\rho_{0}\right) \in(\mathcal{I} \cap \mathcal{T})_{\mathrm{e}}$, then $\mathcal{L}\left(\rho_{0}\right)$ can be decomposed into measures in $\mathcal{T}_{\mathrm{e}}$ with different mean densities, which by Lemma 11.11 converges to mixtures of $\left(\Lambda_{c}\right)_{c \geq 0}$. Therefore by the extremality of $\rho_{0}$, we must have $\mathcal{L}\left(\rho_{0}\right)=\Lambda_{c}$ for some $0 \leq c<\infty$, and hence $(\mathcal{I} \cap \mathcal{T})_{\mathrm{e}} \subset\left\{\Lambda_{c}\right.$ : $c \geq 0\}$. The converse $\left\{\Lambda_{c}: c \geq 0\right\} \subset(\mathcal{I} \cap \mathcal{T})_{\mathrm{e}}$ has been established in Lemma 11.4.

### 11.4 Proof of Theorems $2.11 \sqrt{2.12}$

Proof of Theorem 2.11. Part (a) follows from Lemma 11.12, where the scaling relation $\Lambda_{c}(\mathrm{~d}(c \rho))=\Lambda_{1}(\mathrm{~d} \rho)$ is trivial, while 2.21) and 2.22 follow from Lemma 11.3 and Lemma 11.6 applied to $\rho_{0}(\mathrm{~d} x)=\mathrm{d} x$. Parts (b) and (c) follow from Lemma 11.11 and 11.6, while part (d) follows from spatial ergodic decomposition and Lemma 11.11 .

Proof of Theorem 2.12. Part (a) follows from Theorem 2.7 (a) and Proposition 2.6 (c). Part (b) follows from Theorem 2.8 (a).

## A The Howitt-Warren martingale problem

Howitt and Warren HW09a, Thm 2.1] formulated a martingale problem for a class of sticky Brownian motions on $\mathbb{R}$, for which they showed that for each deterministic initial state $\vec{x} \in \mathbb{R}^{n}$, there exists a unique solution in distribution to their martingale problem. Moreover, they showed that the family of all solutions to their martingale problem forms a consistent Feller family HW09a, Prop. 8.1], which defines a stochastic flow of kernels we call a Howitt-Warren flow. In this appendix, we show that our formulation of the Howitt-Warren martingale problem in Definition 2.2 is equivalent to Howitt and Warren's original formulation in HW09a. The advantage of our formulation is that we use a much simpler set of test functions, which somewhat simplies the proof of the convergence of the $n$-point motions of discrete Howitt-Warren flows to their continuous counterparts. This convergence result is formulated in Proposition A.5, and is used to verify that the flows we construct in Theorem 3.7 are indeed Howitt-Warren flows. A similar convergence result for the $n$-point motions of a continuous time version of the discrete Howitt-Warren flows was established previously in HW09a. We will also give some new parametrizations of Howitt-Warren martingale problems in Lemma A.2.

## A. 1 Different formulations

Let us first recall the original formulation of the Howitt-Warren martingale problem from HW09a, and then state two lemmas that show how one can go from their formulation to ours in Definition 2.2 and vice versa. The proof of these lemmas will be given in the next subsection.

Recall that if $Y$ is a continuous semimartingale, then there exist a unique continuous process $Y^{\mathcal{C}}$ with bounded variation such that $Y-Y^{\mathcal{C}}$ is a martingale. The process $Y^{\mathcal{C}}$ is called the compensator of $Y$. Now if $Y_{1}$ and $Y_{2}$ are continuous, square integrable semimartingales, then by definition, the covariance process $\left\langle Y_{1}, Y_{2}\right\rangle$ of $Y_{1}$ and $Y_{2}$ is the compensator of ( $Y_{1}-$ $\left.Y_{1}^{\mathcal{C}}\right)\left(Y_{2}-Y_{2}^{\mathcal{C}}\right)$, i.e., $\left\langle Y_{1}, Y_{2}\right\rangle$ is the unique continuous process of bounded variation such that

$$
\begin{equation*}
t \mapsto\left(Y_{1}(t)-Y_{1}^{\mathcal{C}}(t)\right)\left(Y_{2}(t)-Y_{2}^{\mathcal{C}}(t)\right)-\left\langle Y_{1}, Y_{2}\right\rangle(t) \tag{A.1}
\end{equation*}
$$

is a martingale. We generalize our definition of the Howitt-Warren martingale problem as follows.

Remark A. 1 (Initial states with infinite second moments) The solutions to a HowittWarren martingale problem (for given $\beta, \nu$ and $n$ ) form a Feller process. Therefore, if $\mathbb{P}_{\vec{x}}$ denotes the law of the solution of the Howitt-Warren martingale problem with initial state $\vec{x}$, and $\rho$ is any probability law on $\mathbb{R}^{n}$, then $\int \rho(\mathrm{d} \vec{x}) \mathbb{P}_{\vec{x}}$ is the law of some Markov process in $\mathbb{R}^{n}$. Generalizing Definition 2.2, we may call such a process $\vec{X}$ the solution to the Howitt-Warren martingale problem with initial law $\rho$, even though $\vec{X}$ is not square integrable if $\rho$ does not have a finite second moment.

We now turn our attention to the original formulation of Howitt and Warren's martingale problem in HW09a. Recall that our formulation of the Howitt-Warren martingale in Definition 2.2 is based on the constants $\left(\beta_{+}(m)\right)_{m \geq 1}$ defined in (2.3). Instead, Howitt and Warren's formulation of their martingale problem is based on real constants $(\theta(k, l))_{k, l \geq 0}$ satisfying

$$
\begin{array}{rlrl}
\text { (i) } & \theta(k, l) & \geq 0 & \\
\text { (ii) } & \theta(k, l) & =\theta(k+1, l)+\theta(k, l+1) &  \tag{A.2}\\
(k, l \geq 0)
\end{array}
$$

The $\theta(k, l)$ 's are related to the $\beta_{+}(m)$ 's by

$$
\begin{equation*}
\beta_{+}(m)=\theta(0,0)-2 \theta(0, m) \quad(m \geq 1) \tag{A.3}
\end{equation*}
$$

while their relation to the constant $\beta$ and measure $\nu$ is described by

$$
\begin{align*}
& \theta(k, l)=\int \nu(\mathrm{d} q) q^{k-1}(1-q)^{l-1} \quad(k, l \geq 1)  \tag{A.4}\\
& \theta(1,0)-\theta(0,1)=\beta
\end{align*}
$$

Note that we have now three ways to parametrize Howitt-Warren martingale problems: we may use the pair $(\beta, \nu)$, the constants $\left(\beta_{+}(m)\right)_{m \geq 1}$, or the constants $(\theta(k, l))_{k, l \geq 0}$. The next lemma shows how to go from one parametrization to another.

## Lemma A. 2 (Different parametrizations)

(a) Let $(\theta(k, l))_{k, l \geq 0}$ be real constants satisfying (A.2). Then there exists a unique $\beta \in \mathbb{R}$ and a finite measure $\nu$ on $[0,1]$ such that A.4 holds.
(b) Let $\beta \in \mathbb{R}$ and let $\nu$ be a finite measure on $[0,1]$. Then there exists a function $\theta: \mathbb{N}^{2} \rightarrow \mathbb{R}$ satisfying (A.2) such that (A.4) holds. Any other $\theta^{\prime}$ satisfies (A.2) and (A.4) if and only if

$$
\begin{equation*}
\theta^{\prime}(k, l)=\theta(k, l)+c\left(1_{\{k=0\}}+1_{\{l=0\}}\right) \quad(k, l \geq 0) \tag{A.5}
\end{equation*}
$$

for some $c \in \mathbb{R}$, and we say that $\theta$ and $\theta^{\prime}$ are equivalent.
(c) Let $\beta \in \mathbb{R}$ and let $\nu$ be a finite measure on $[0,1]$. Let $(\theta(k, l))_{k, l \geq 0}$ be real constants satisfying (A.2), and let $\left(\beta_{+}(m)\right)_{m \geq 1}$ be real constants. Then of the relations (2.3), A.3), and (A.4), any two imply the third one.

By definition, a weak total order on $\{1, \ldots, n\}$ is a relation $\prec$ such that
(i) $i \prec i$,
(ii) $\quad i \prec j \prec k$ implies $i \prec k$,
(iii) there exist no $i, j$ with $i \nprec j$ and $j \nprec i$.

Each weak total order $\prec$ on $\{1, \ldots, n\}$ defines a nonempty cell $C_{\prec} \subset \mathbb{R}^{n}$ by

$$
\begin{equation*}
C_{\prec}:=\left\{\vec{x} \in \mathbb{R}^{n}: x_{i} \leq x_{j} \text { if and only if } i \prec j\right\} . \tag{A.7}
\end{equation*}
$$

We note that cells belonging to different weak total orders are disjoint, and that the union of all such cells is $\mathbb{R}^{n}$. For example:

$$
\begin{equation*}
\left\{\vec{x}: x_{1}<x_{3}<x_{2}\right\}, \quad\left\{\vec{x}: x_{2}=x_{3}<x_{1}\right\}, \quad \text { and } \quad\left\{\vec{x}: x_{1}=x_{2}=x_{3}\right\} \tag{A.8}
\end{equation*}
$$

are three of the thirteen cells that make up $\mathbb{R}^{3}$. Let $L_{n}$ be the linear space consisting of all continuous real functions on $\mathbb{R}^{n}$ that are piecewise linear on each cell $C_{\prec}$, i.e.,

$$
\begin{align*}
L_{n}:=\{f: f & \text { is a continuous function } f: \mathbb{R}^{n} \rightarrow \mathbb{R} \text { such that for each weak total } \\
& \text { order } \left.\prec \text { there exists a linear function } l: \mathbb{R}^{n} \rightarrow \mathbb{R} \text { with } f=l \text { on } C_{\prec}\right\} . \tag{A.9}
\end{align*}
$$

For each $\vec{x} \in \mathbb{R}^{n}$, let us define

$$
\begin{equation*}
\operatorname{Ran}(\vec{x}):=\bigcup_{i=1}^{n}\left\{x_{i}\right\} \tag{A.10}
\end{equation*}
$$

and for each $x \in \operatorname{Ran}(\vec{x})$, let us write

$$
\begin{equation*}
J_{x}:=\left\{i \in\{1, \ldots, n\}: x_{i}=x\right\} \tag{A.11}
\end{equation*}
$$

For disjoint $I, J \subset\{1, \ldots, n\}$ let us define a vector $\vec{v}_{I, J} \in \mathbb{R}^{n}$ by

$$
v_{I, J}(i):=\left\{\begin{align*}
1 & \text { if } i \in I  \tag{A.12}\\
-1 & \text { if } i \in J \\
0 & \text { otherwise }
\end{align*}\right.
$$

For any $\vec{v} \in \mathbb{R}^{n}$, let $\nabla_{\vec{v}}$ denote the one-sided derivative

$$
\begin{equation*}
\nabla_{\vec{v}} f(\vec{x}):=\lim _{\varepsilon \downarrow 0} \varepsilon^{-1}(f(\vec{x}+\varepsilon \vec{v})-f(\vec{x})) \tag{A.13}
\end{equation*}
$$

Let $(\theta(k, l))_{k, l \geq 0}$ be real constants satisfying A.2p. Then, by definition, $\mathcal{A}_{n}^{\theta}$ is the linear operator acting on functions in $L_{n}$ defined by

$$
\begin{equation*}
\mathcal{A}_{n}^{\theta} f(\vec{x}):=\sum_{x \in \operatorname{Ran}(\vec{x})} \sum_{I \subset J_{x}} \theta\left(|I|,\left|J_{x} \backslash I\right|\right) \nabla_{\vec{v}_{I, J_{x} \backslash I}} f(\vec{x}) . \tag{A.14}
\end{equation*}
$$

The original formulation of the Howitt-Warren martingale problem in HW09a differs from our formulation in that formula (2.7) is replaced by the requirement that for each $f \in L_{n}$

$$
\begin{equation*}
f(\vec{X}(t))-\int_{0}^{t} \mathcal{A}_{n}^{\theta} f(\vec{X}(s)) \mathrm{d} s \tag{A.15}
\end{equation*}
$$

is a martingale with respect to the filtration generated by $\vec{X}$. To see that this is equivalent to the formulation in Definition 2.2, we need the following lemma, the proof of which is not entirely trivial.

Lemma A. 3 (Action of operator on basis vectors) Let $f_{\Delta}, g_{\Delta}$ be defined as in (2.5). Then:
(a) The functions

$$
\begin{equation*}
\left\{f_{\Delta}: \emptyset \neq \Delta \subset\{1, \ldots, n\}\right\} \tag{A.16}
\end{equation*}
$$

form a basis for the space $L_{n}$.
(b) Let $(\theta(k, l))_{k, l \geq 0}$ be real constants satisfying A.2) and let $\left(\beta_{+}(m)\right)_{m \geq 1}$ be given by A.3). Then for each nonempty $\Delta \subset\{1, \ldots, n\}$, one has

$$
\begin{equation*}
\mathcal{A}_{n}^{\theta} f_{\Delta}(\vec{x})=\beta_{+}\left(g_{\Delta}(\vec{x})\right) \quad\left(\vec{x} \in \mathbb{R}^{n}\right) . \tag{A.17}
\end{equation*}
$$

(c) If $\theta$ and $\theta^{\prime}$ satisfy A.2) and are equivalent in the sense of A.5), then $\mathcal{A}_{n}^{\theta}=\mathcal{A}_{n}^{\theta^{\prime}}$.

## A. 2 Proof of the equivalence of formulations

To prepare for the proof of Lemma A.2, we start with the following lemma.
Lemma A. 4 (Moments defining a measure) Let $(\phi(k, l))_{k, l \geq 0}$ be real constants such that

$$
\begin{align*}
& \text { (i) } \quad \phi(k, l) \geq 0  \tag{A.18}\\
& \text { (ii) } \quad \phi(k, l)=\phi(k+1, l)+\phi(k, l+1)
\end{align*}
$$

for all $k, l \geq 0$. Then there exists a unique finite measure $\nu$ on $[0,1]$ such that

$$
\begin{equation*}
\phi(k, l)=\int \nu(\mathrm{d} q) q^{k}(1-q)^{l} \quad(k, l \geq 0) \tag{A.19}
\end{equation*}
$$

Proof. Let $-\Delta$ be the operator, acting on sequences of real constants $\left(a_{k}\right)_{k \geq 0}$ as $((-\Delta) a)_{k}:=$ $a_{k}-a_{k+1}$. Setting $a_{k}:=\phi(k, 0)$, we observe that $((-\Delta) a)_{k}=\phi(k, 1)(k \geq 0)$ and more generally $\left((-\Delta)^{l} a\right)_{k}=\phi(k, l) \geq 0(k, l \geq 0)$. This qualifies $\left(a_{n}\right)_{n \in \mathbb{N}}$ as a completely monotone sequence, which by [Fel66, Theorem VII.3.2] can be represented as $a_{k}=\int \nu(\mathrm{d} q) q^{k}$ for some finite measure $\nu$ on $[0,1]$. Using A.18) (ii), this implies A.19).
Proof of Lemma A. 2 Part (a) is a straightforward consequence of Lemma A. 4 To prove part (b), note that by (A.4), $\nu$ uniquely determines $\theta(k, l)$ for $k, l \geq 1$, which is easily seen to
satisfy A.2 for $k, l \geq 1$. Once $\theta(1,0)$ and $\theta(0,1)$ are chosen, $\theta(k, 0)$ and $\theta(0, l)$ for $k, l \geq 0$ are uniquely determined from the recursion relation (A.2) (ii). Since $\theta(1,0)-\theta(0,1)=\beta$, it follows that $\theta$ is uniquely determined up to the equivalence defined in A .5 ).

To prove part (c), we observe that (2.3) and A.4, together with A.2) (ii), imply that

$$
\begin{equation*}
\beta_{+}(1)=\beta=\theta(1,0)-\theta(0,1)=\theta(0,0)-2 \theta(0,1) \tag{A.20}
\end{equation*}
$$

and

$$
\begin{aligned}
\beta_{+}(m) & =\beta+2 \int \nu(\mathrm{~d} q) \sum_{k=1}^{m-1}(1-q)^{k-1}=\theta(1,0)-\theta(0,1)+2 \sum_{k=1}^{m-1} \theta(1, k) \\
& =\theta(1,0)-\theta(0,1)+2(\theta(0,1)-\theta(0, m))=\theta(0,0)-2 \theta(0, m) \quad(m \geq 2)
\end{aligned}
$$

This shows that (2.3) and (A.4) imply A.3). Running the argument backward, we also see that (A.4) and A.3) imply (2.3). Finally, (2.3) and A.3 imply that

$$
\begin{align*}
\theta(0,0)-2 \theta(0,1) & =\beta \\
\theta(0,0)-2 \theta(0, m) & =\beta+2 \int q^{-1}\left(1-(1-q)^{m-1}\right) \nu(\mathrm{d} q) \quad(m \geq 2) \tag{A.21}
\end{align*}
$$

from which it is not hard to derive A.4 using A.2 (ii).
Proof of Lemma A.3 As a first step towards proving part (a), we start by proving that the functions $\left\{f_{\Delta}: \emptyset \neq \Delta \subset\{1, \ldots, n\}\right\}$ are linearly independent. Consider the set $\{0,1\}^{n} \subset \mathbb{R}^{n}$. For each $A \subset\{1, \ldots, n\}$, define $g_{A}:\{0,1\}^{n} \rightarrow \mathbb{R}$ by

$$
g_{A}(\vec{x}):= \begin{cases}1-f_{A}(\vec{x}) & \text { if } A \neq \emptyset  \tag{A.22}\\ 1 & \text { if } A=\emptyset\end{cases}
$$

It is not hard to see that $g_{A} g_{B}=g_{A \cup B}$ and that the functions $\left\{g_{A}: A \subset\{1, \ldots, n\}\right\}$ separate points. Therefore, by the Stone-Weierstrass theorem, they span the space of all real functions on $\{0,1\}^{n}$. Since this space has dimension $2^{n}$ and since $\left\{g_{A}: A \subset\{1, \ldots, n\}\right\}$ has $2^{n}$ elements, we conclude that the $g_{A}$ 's are linearly independent and hence the same is true for the $f_{\Delta}$ 's.

We next prove that the $f_{\Delta}$ 's span $L_{n}$. Obviously $f_{\Delta} \in L_{n}$ for each $\emptyset \neq \Delta \subset\{1, \ldots, n\}$. Therefore, since the $f_{\Delta}$ 's are linearly independent and since $\left\{f_{\Delta}: \emptyset \neq \Delta \subset\{1, \ldots, n\}\right\}$ has $2^{n}-1$ elements, it suffices to show that $\operatorname{dim}\left(L_{n}\right) \leq 2^{n}-1$. We proceed by induction. It is easy to check that $L_{1}$ is the space of all linear functions from $\mathbb{R}$ to $\mathbb{R}$, which has dimension one. Now assume that $\operatorname{dim}\left(L_{n}\right) \leq 2^{n}-1$. We claim that $\operatorname{dim}\left(L_{n+1}\right) \leq 2^{n+1}-1$. Each function $f \in L_{n+1}$ can be uniquely written as

$$
\begin{equation*}
f(\vec{x})=\sum_{i=1}^{n+1} c_{i}(\vec{x}) x_{i} \tag{A.23}
\end{equation*}
$$

where the functions $c_{1}, \ldots, c_{n+1}$ are piecewise constant on each cell $C_{\prec}$. In fact, since functions in $L_{n}$ are continuous, we must have that the function $c_{n+1}$ depends only on the relative order of $x_{n+1}$ with respect to the first $n$ coordinates and does not change if we interchange the order of two other coordinates $x_{j}, x_{k}$ with $j, k \leq n$. In particular, for each $A \subset\{1, \ldots, n\}$, if we set

$$
\begin{equation*}
U_{A}:=\left\{\vec{x} \in \mathbb{R}^{n+1}: x_{i}<x_{n+1} \forall i \in A, x_{i}>x_{n+1} \forall i \in\{1, \ldots, n\} \backslash A\right\} \tag{A.24}
\end{equation*}
$$

then

$$
\begin{equation*}
c_{n+1}(\vec{x})=l_{A} \quad\left(x \in U_{A}\right) \tag{A.25}
\end{equation*}
$$

for some constant $l_{A} \in \mathbb{R}$. Let $l$ be the linear map defined by

$$
l(f):=\left(l_{A}(f)\right)_{A \subset\{1, \ldots, n\}} \quad\left(f \in L_{n}\right)
$$

Then $\operatorname{Ker}(l)$ consists of all functions in $L_{n+1}$ that do not depend on the variable $x_{n+1}$, hence $\operatorname{Ker}(l) \subset L_{n}$. It follows that

$$
\operatorname{dim}\left(L_{n+1}\right)=\operatorname{dim}(\operatorname{Ker}(l))+\operatorname{dim}(\operatorname{Ran}(l)) \leq\left(2^{n}-1\right)+2^{n}=2^{n+1}-1
$$

as claimed.
To prove part (b) of the lemma, we need to calculate

$$
\begin{equation*}
\mathcal{A}_{n}^{\theta} f_{\Delta}(\vec{x})=\sum_{x \in \operatorname{Ran}(\vec{x})} \sum_{I \subset J_{x}} \theta\left(|I|,\left|J_{x} \backslash I\right|\right) \nabla_{\vec{v}_{I, J_{x} \backslash I}} f_{\Delta}(\vec{x}) \tag{A.26}
\end{equation*}
$$

Let us define

$$
\begin{align*}
H(\vec{x}) & :=J_{f_{\Delta}(\vec{x})}=\left\{i \in\{1, \ldots, n\}: x_{i}=f_{\Delta}(\vec{x})\right\}  \tag{A.27}\\
G(\vec{x}) & :=H(\vec{x}) \cap \Delta=\left\{i \in \Delta: x_{i}=f_{\Delta}(\vec{x})\right\}
\end{align*}
$$

Recalling A.12 and A.13, we see that

$$
\nabla_{\vec{v}_{I, J}} f_{\Delta}(\vec{x})=\left\{\begin{align*}
+1 & \text { if } I \cap G(\vec{x}) \neq \emptyset  \tag{A.28}\\
-1 & \text { if } J \supset G(\vec{x}), \\
0 & \text { otherwise. }
\end{align*}\right.
$$

Inserting this into A.26 we see that

$$
\begin{equation*}
\mathcal{A}_{n}^{\theta} f_{\Delta}(\vec{x})=\sum_{I \subset H(\vec{x})} \theta(|I|,|H(\vec{x}) \backslash I|)\left(1_{\{I \cap G(\vec{x}) \neq \emptyset\}}-1_{\{I \cap G(\vec{x})=\emptyset\}}\right), \tag{A.29}
\end{equation*}
$$

where we have used that for $I \subset H(\vec{x})$, one has $(H(\vec{x}) \backslash I) \supset G(\vec{x})$ if and only if $I \cap G(\vec{x})=\emptyset$.
We claim that A.29 can be rewritten as

$$
\begin{equation*}
\mathcal{A}_{n}^{\theta} f_{\Delta}(\vec{x})=\sum_{I \subset G(\vec{x})} \theta(|I|,|G(\vec{x}) \backslash I|)\left(1_{\{I \cap G(\vec{x}) \neq \emptyset\}}-1_{\{I \cap G(\vec{x})=\emptyset\}}\right) \tag{A.30}
\end{equation*}
$$

To see this, note that if $H^{\prime}$ is a set such that $G(\vec{x}) \subset H^{\prime} \subset H(\vec{x})$ and $H^{\prime}$ contains one element less than $H(\vec{x})$, then since it does not make a difference for the sign of a term in A.29) whether we include this element in $I$ or in $H(\vec{x}) \backslash I$, we have

$$
\begin{align*}
\mathcal{A}_{n}^{\theta} f_{\Delta}(\vec{x}) & =\sum_{I \subset H^{\prime}}\left(\theta\left(|I|+1,\left|H^{\prime} \backslash I\right|\right)+\theta\left(|I|,\left|H^{\prime} \backslash I\right|+1\right)\right)\left(1_{\{I \cap G(\vec{x}) \neq \emptyset\}}-1_{\{I \cap G(\vec{x})=\emptyset\}}\right) \\
& =\sum_{I \subset H^{\prime}} \theta\left(|I|,\left|H^{\prime} \backslash I\right|\right)\left(1_{\{I \cap G(\vec{x}) \neq \emptyset\}}-1_{\{I \cap G(\vec{x})=\emptyset\}}\right) \tag{A.31}
\end{align*}
$$

where we have used (A.2) (ii). Continuing this process of removing points from $H(\vec{x})$ we arrive at A.30).

We may rewrite A.30 as

$$
\begin{align*}
\mathcal{A}_{n}^{\theta} f_{\Delta}(\vec{x}) & =\sum_{I \subset G(\vec{x})} \theta(|I|,|G(\vec{x}) \backslash I|)\left(1-2 \cdot 1_{\{I \cap G(\vec{x})=\emptyset\}}\right) \\
& =\left(\sum_{I \subset G(\vec{x})} \theta(|I|,|G(\vec{x}) \backslash I|)\right)-2 \theta(0,|G(\vec{x})|) \tag{A.32}
\end{align*}
$$

The same sort of argument as in A.31 shows that the first term on the right-hand side of A.32) equals $\theta(0,0)$ and hence, recalling A.3) and the fact that $|G(\vec{x})|=g_{\Delta}(\vec{x})$ (see (2.5), we arrive at A.17).

Part (c) is a trivial consequence of parts (a) and (b) and the fact that if $\theta$ and $\theta^{\prime}$ are equivalent in the sense of A.5), then they define the same $\left(\beta_{+}(m)\right)_{m \geq 1}$ through A.3).

## A. 3 Convergence of discrete n-point motions

In this section we prove that if $\mu_{k}$ are probability measures on $[0,1]$ satisfying (1.7), then the diffusively rescaled discrete $n$-point motions associated with the $\mu_{k}$ converge in law to the Markov process defined by the Howitt-Warren martingale problem with drift $\beta$ and characteristic measure $\nu$. To formulate this precisely, fix $\mu_{k}$ satisfying $\sqrt[1.7]{ }$, let $\vec{X}^{\langle k\rangle}$ be discrete $n$-point motions associated with the $\mu_{k}$, started in deterministic initial states $\vec{x}^{(k\rangle}$ and linearly interpolated between integer times, and let $\vec{Y}^{\langle k\rangle}$ defined by

$$
\begin{equation*}
Y_{i}^{\langle k\rangle}(t):=\varepsilon_{k} X_{i}^{\langle k\rangle}\left(t / \varepsilon_{k}^{2}\right) \quad(i=1, \ldots, n, t \geq 0) \tag{A.33}
\end{equation*}
$$

denote the process $\vec{X}^{\langle k\rangle}$, diffusively rescaled with $\varepsilon_{k}$. Let $\mathcal{C}_{\mathbb{R}^{n}}[0, \infty)$ denote the space of continuous functions from $[0, \infty)$ to $\mathbb{R}^{n}$, equipped with the topology of local uniform convergence. Then, in analogy with [HW09a, Thm. 8.1], we have the following result.

Proposition A. 5 (Convergence of the $n$-point motions) Assume that the initial states satisfy

$$
\begin{equation*}
\varepsilon_{k} \vec{x}^{\langle k\rangle} \underset{k \rightarrow \infty}{\longrightarrow} \vec{x} \tag{A.34}
\end{equation*}
$$

for some $\vec{x} \in \mathbb{R}^{n}$. Then

$$
\begin{equation*}
\mathbb{P}\left[\left(\vec{Y}^{\langle k\rangle}(t)\right)_{t \geq 0} \in \cdot\right] \underset{k \rightarrow \infty}{\Longrightarrow} \mathbb{P}\left[(\vec{X}(t))_{t \geq 0} \in \cdot\right] \tag{A.35}
\end{equation*}
$$

where $\Rightarrow$ denotes weak convergence of probability laws on $\mathcal{C}_{\mathbb{R}^{n}}[0, \infty)$ and $\vec{X}$ is the unique solution of the Howitt-Warren martingale problem with drift $\beta$ and characteristic measure $\nu$, started in the initial state $\vec{X}_{0}=\vec{x}$.

Remark. There is an analoguous statement for random initial states, see Remark A.1.
We will actually prove a somewhat stronger statement than the convergence in A.35), since we will show that the intersection times of the rescaled discrete process also converge to those of the limiting process. For technical reasons, it will be convenient to interpolate in a piecewise constant, rather than in a linear way. Therefore, we set (compare A.33))

$$
\begin{equation*}
\mathrm{Y}_{i}^{\langle k\rangle}(t):=\varepsilon_{k} X_{i}^{\langle k\rangle}\left(\left\lfloor t / \varepsilon_{k}^{2}\right\rfloor\right) \quad(i=1, \ldots, n, t \geq 0) \tag{A.36}
\end{equation*}
$$

We view $\overrightarrow{\mathrm{Y}}{ }^{\langle k\rangle}$ as a process with paths in $\mathcal{D}_{\mathbb{R}^{n}}[0, \infty)$, the space of cádlág functions from $[0, \infty)$ to $\mathbb{R}^{n}$, equipped with the Skorohod topology. Letting $\vec{Y}^{\langle k\rangle}, \overrightarrow{\mathrm{Y}}^{\langle k\rangle}$ denote the linearly interpolated and piecewise constant processes, respectively, we have

$$
\begin{equation*}
\sup _{t \geq 0}\left|Y_{i}^{\langle k\rangle}(t)-Y_{i}^{\langle k\rangle}(t)\right|=\varepsilon_{k} \underset{k \rightarrow \infty}{\longrightarrow} 0 \tag{A.37}
\end{equation*}
$$

From this, it is easy to see that Proposition A.5 is implied by the following, somewhat stronger result.

Proposition A. 6 (Convergence including intersection times) Let $\vec{X}^{\langle k\rangle}$ be discrete $n$ point motions associated with probability measures $\mu_{k}$ satisfying (1.7), started from initial states $\vec{x}^{\langle k\rangle}$, and let $\overrightarrow{\mathrm{Y}}^{\langle k\rangle}$ denote $\vec{X}^{\langle k\rangle}$ diffusively rescaled as in A.36). Let $\vec{X}$ be the unique solution of the Howitt-Warren martingale problem with drift $\beta$ and characteristic measure $\nu$, started in $\vec{X}_{0}=\vec{x}$. Define $n \times n$ matrix valued processes $Z^{\langle k\rangle}$ and $Z$ by

$$
\begin{align*}
& \text { (i) } Z_{i j}^{\langle k\rangle}(t):=\int_{0}^{t} 1_{\left\{\mathrm{Y}_{i}^{\langle k\rangle}(s)=\mathrm{Y}_{j}^{\langle k\rangle}(s)\right\}} \mathrm{d} s,  \tag{A.38}\\
& \text { (ii) } Z_{i j}(t):=\int_{0}^{t} 1_{\left\{X_{i}(s)=X_{j}(s)\right\}} \mathrm{d} s .
\end{align*}
$$

Then, assuming that the initial states satisfy

$$
\begin{equation*}
\varepsilon_{k} \vec{x}^{\langle k\rangle} \underset{k \rightarrow \infty}{\longrightarrow} \vec{x} \tag{A.39}
\end{equation*}
$$

one has

$$
\begin{equation*}
\mathbb{P}\left[\left(\overrightarrow{\mathrm{Y}}^{\langle k\rangle}(t), Z^{\langle k\rangle}(t)\right)_{t \geq 0} \in \cdot\right] \underset{k \rightarrow \infty}{\Longrightarrow} \mathbb{P}\left[(\vec{X}(t), Z(t))_{t \geq 0} \in \cdot\right], \tag{A.40}
\end{equation*}
$$

where $\Rightarrow$ denotes weak convergence of probability laws on path space.
Proof. When $\vec{X}^{\langle k\rangle}$ is the $n$-point motion of a continuous time version of the discrete HowittWarren flow, the same result has been proved by Howitt and Warren in HW09a, Prop. 6.3] (for tightness in their case, see the remarks above their formula (6.13).) Our proof copies their proof in many places, except that we use a different argument to get convergence of the compensators of $f_{\Delta}\left(\overrightarrow{\mathrm{Y}}^{\langle k\rangle}\right)$ and we have also simplified their proof somewhat due to our reformulation of their martingale problem.

Let $P^{\langle k\rangle}$ be the transition kernel from $\mathbb{Z}^{n}$ to $\mathbb{Z}^{n}$ defined by

$$
\begin{equation*}
P^{\langle k\rangle}(x, y):=\prod_{x \in \operatorname{Ran}(\vec{x})} \int \mu_{k}(\mathrm{~d} q) \prod_{i \in J_{x}}\left(1_{\left\{y_{i}=x_{i}+1\right\}} q+1_{\left\{y_{i}=x_{i}-1\right\}}(1-q)\right) \quad\left(x, y \in \mathbb{Z}^{n}\right), \tag{A.41}
\end{equation*}
$$

where $\operatorname{Ran}(\vec{x})$ and $J_{x}$ are defined in A.10) and A.11. We adopt the notation

$$
\begin{equation*}
P^{\langle k\rangle} f(x):=\sum_{y \in \mathbb{Z}^{n}} P^{\langle k\rangle}(x, y) f(y) \quad\left(x \in \mathbb{Z}^{n}, f: \mathbb{Z}^{n} \rightarrow \mathbb{R}\right), \tag{A.42}
\end{equation*}
$$

whenever the infinite sum is well-defined.
We observe that $\vec{X}^{\langle k\rangle}$ is a Markov chain with transition kernel $P^{\langle k\rangle}$. Since we start $\vec{X}^{\langle k\rangle}$ in an initial state $\vec{X}^{\langle k\rangle}(0)=x^{\langle k\rangle} \in\left(\mathbb{Z}_{\text {even }}\right)^{n}$, because of the nature of the transition mechanism, we have $\vec{X}^{\langle k\rangle}(t) \in\left(\mathbb{Z}_{\text {even }}\right)^{n}$ at even times and $\vec{X}^{\langle k\rangle}(t) \in\left(\mathbb{Z}_{\text {odd }}\right)^{n}$ at odd times.

For $\emptyset \neq \Delta \subset\{1, \ldots, n\}$, let $f_{\Delta}, g_{\Delta}$ be the functions defined in (2.5). By standard theory, for each $\emptyset \neq \Delta \subset\{1, \ldots, n\}$, the discrete-time process

$$
\begin{equation*}
f_{\Delta}\left(\vec{X}^{\langle k\rangle}(t)\right)-\sum_{s=0}^{t-1}\left(P^{\langle k\rangle} f_{\Delta}\left(\vec{X}^{\langle k\rangle}(s)\right)-f_{\Delta}\left(\vec{X}^{\langle k\rangle}(s)\right)\right) \tag{A.43}
\end{equation*}
$$

is a martingale with respect to the filtration generated by $\vec{X}^{\langle k\rangle}$. We observe that if either $x \in\left(\mathbb{Z}_{\text {even }}\right)^{n}$ or $x \in\left(\mathbb{Z}_{\text {odd }}\right)^{n}$, then under the transition kernel $P^{\langle k\rangle}$ the maximum $f_{\Delta}(x)=$ $\max _{i \in \Delta} x_{i}$ moves down by one with probability $\mu_{k}(\mathrm{~d} q)(1-q)^{g_{\Delta}(x)}$ and up by one with the remaining probability, hence

$$
\begin{align*}
P^{\langle k\rangle} f_{\Delta}(x)-f_{\Delta}(x) & =\int \mu_{k}(\mathrm{~d} q)\left(1-2(1-q)^{g_{\Delta}(x)}\right)  \tag{A.44}\\
& =\beta_{k}\left(g_{\Delta}(x)\right) \quad\left(x \in\left(\mathbb{Z}_{\text {even }}\right)^{n} \text { or }\left(\mathbb{Z}_{\text {even }}\right)^{n}\right),
\end{align*}
$$

where we have introduced the notation

$$
\begin{equation*}
\beta_{k}(m):=\int \mu_{k}(\mathrm{~d} q)\left(1-2(1-q)^{m}\right) \quad(m \geq 1) \tag{A.45}
\end{equation*}
$$

Setting $\beta_{k}:=\beta_{k}(1)=\int \mu_{k}(\mathrm{~d} q)(2 q-1)$, by standard theory, one may moreover check that

$$
\begin{equation*}
\left(X_{i}^{\langle k\rangle}(t)-\beta_{k} t\right)\left(X_{j}^{\langle k\rangle}(t)-\beta_{k} t\right)-\sum_{s=0}^{t-1} \Gamma_{i j}^{\langle k\rangle}\left(\vec{X}^{\langle k\rangle}(s)\right) \tag{A.46}
\end{equation*}
$$

is a martingale, where

$$
\begin{align*}
\Gamma_{i j}^{\langle k\rangle}(\vec{x}) & :=\sum_{y \in \mathbb{Z}} P^{\langle k\rangle}(x, y)\left(y_{i}-x_{i}\right)\left(y_{j}-x_{j}\right)-\beta_{k}^{2} \\
& =\left\{\begin{array}{cl}
1-\beta_{k}^{2} & \text { if } i=j, \\
\int \mu_{k}(\mathrm{~d} q)(1-4 q(1-q))-\beta_{k}^{2} & \text { if } i \neq j, x_{i}=x_{j}, \\
0 & \text { otherwise }
\end{array}\right.  \tag{A.47}\\
& =\left(1-\beta_{k}^{2}-2\left(\beta_{k}(2)-\beta_{k}(1)\right) 1_{\{i \neq j\}}\right) 1_{\left\{x_{i}=x_{j}\right\}} .
\end{align*}
$$

For the process $\overrightarrow{\mathrm{Y}}^{\langle k\rangle}$ defined in A.36, our arguments so far show that for each $\emptyset \neq \Delta \subset$ $\{1, \ldots, n\}$,

$$
\begin{equation*}
f_{\Delta}\left(\overrightarrow{\mathrm{Y}}^{\langle k\rangle}(t)\right)-\varepsilon_{k}^{-1} \int_{0}^{\lfloor t\rfloor_{k}} \beta_{k}\left(g_{\Delta}\left(\overrightarrow{\mathrm{Y}}^{\langle k\rangle}(s)\right)\right) \mathrm{d} s \tag{A.48}
\end{equation*}
$$

is a martingale, where $\lfloor t\rfloor_{k}:=\varepsilon_{k}^{2}\left\lfloor t / \varepsilon_{k}^{2}\right\rfloor$ denotes the time $t$ rounded downwards to the next time in $\varepsilon_{k}^{2} \mathbb{N}$. Moreover, for each $1 \leq i, j \leq n$, the process

$$
\begin{align*}
& \left(\mathrm{Y}_{i}^{\langle k\rangle}(t)-\varepsilon_{k}^{-1} \beta_{k} t\right)\left(\mathrm{Y}_{j}^{\langle k\rangle}(t)-\varepsilon_{k}^{-1} \beta_{k} t\right) \\
& \quad-\left(1-\beta_{k}^{2}-2\left(\beta_{k}(2)-\beta_{k}(1)\right) 1_{\{i \neq j\}}\right) \int_{0}^{\lfloor t\rfloor_{k}} 1_{\left\{\mathrm{Y}_{i}^{\langle k\rangle}(s)=\mathrm{Y}_{i}^{\langle k\rangle}(s)\right\}} \mathrm{d} s \tag{A.49}
\end{align*}
$$

is a martingale with respect to the filtration generated by $\overrightarrow{\mathrm{Y}}^{\langle k\rangle}$. It follows from our assumption (1.7) (see also (2.4)) that

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \varepsilon_{k}^{-1} \beta_{k}(m)=\beta_{+}(m) \quad(m \geq 1) \tag{A.50}
\end{equation*}
$$

Standard results (Donsker's invariance principle) tell us that for $k \rightarrow \infty$, each component $\mathrm{Y}_{i}^{\langle k\rangle}$ of the rescaled process converges weakly in law, on the space $\mathcal{D}_{\mathbb{R}}[0, \infty)$, to a Brownian motion with drift $\beta=\beta_{+}(1)$. This implies that the laws of the processes $\overrightarrow{\mathrm{Y}}^{\langle k\rangle}$ (viewed as probability laws on $\mathcal{D}_{\mathbb{R}^{n}}[0, \infty)$ ) are tight. Let $Z^{\langle k\rangle}$ be the matrix valued processes defined in A. 38 (i). Since the slope of each $Z_{i j}^{\langle k\rangle}$ is between zero and one, tightness for these processes is immediate.

By going to a subsequence if necessary, we may assume that the joint processes ( $\overrightarrow{\mathrm{Y}}^{\langle k\rangle}, Z^{\langle k\rangle}$ ) converges weakly in law, and by Skorohod's representation theorem (see e.g. Bil99, Theorem 6.7]), we can couple the ( $\left.\overrightarrow{\mathrm{Y}}^{\langle k\rangle}, Z^{\langle k\rangle}\right)$ 's such that the convergence is almost sure. Let $(\vec{X}, \tilde{Z})$ denote the limiting process. Then, taking the limit in A.49), using A.50), we see that

$$
\begin{equation*}
\left(X_{i}(t)-\beta t\right)\left(X_{j}(t)-\beta t\right)-\tilde{Z}_{i j}(t) \tag{A.51}
\end{equation*}
$$

is a martingale, hence

$$
\begin{equation*}
\left\langle X_{j}, X_{j}\right\rangle(t)=\tilde{Z}_{i j}(t)=\lim _{k \rightarrow \infty} \int_{0}^{t} 1_{\left\{\mathrm{Y}_{i}^{(k\rangle}(s)=\mathrm{Y}_{j}^{\langle k\rangle}(s)\right\}} \mathrm{d} s \quad \forall t \geq 0,1 \leq i, j \leq n \quad \text { a.s. } \tag{A.52}
\end{equation*}
$$

Since, for given $t>0$, the function $w \mapsto \int_{0}^{t} 1_{\left\{w_{i}(s)=w_{j}(s)\right\}} \mathrm{d} s$ is upper semicontinuous with respect to the topology on $\mathcal{D}_{\mathbb{R}^{n}}[0, \infty)$, formula A.52) implies that

$$
\begin{equation*}
\left\langle X_{i}, X_{j}\right\rangle(t) \leq \int_{0}^{t} 1_{\left\{X_{i}(s)=X_{j}(s)\right\}} \mathrm{d} s \quad(t \geq 0,1 \leq i, j \leq n) . \tag{A.53}
\end{equation*}
$$

To prove also the other inequality in (A.53), we use an argument due to Howitt and Warren (see the proof of formula (6.9) in HW09a). For any real square integrable semimartingale $W$, one can define a 'local time' $L(x, t)$ such that

$$
\begin{equation*}
\int_{0}^{t} f(W(s)) \mathrm{d}\langle W, W\rangle(s)=\int_{\mathbb{R}} f(x) L(x, t) \mathrm{d} x \tag{A.54}
\end{equation*}
$$

(See [BY81, formula (3)].) Applying this to the semimartingale $X_{i}-X_{j}$ and the function $f=1_{\{0\}}$, we find that

$$
\begin{equation*}
\int_{0}^{t} 1_{\left\{X_{i}(s)=X_{j}(s)\right\}} \mathrm{d}\left\langle X_{i}-X_{j}, X_{i}-X_{j}\right\rangle(s)=\int_{\mathbb{R}} 1_{\{0\}}(x) L(x, t) \mathrm{d} x=0 . \tag{A.55}
\end{equation*}
$$

Since $X_{i}, X_{j}$ are Brownian motions, we have

$$
\begin{align*}
\left\langle X_{i}-X_{j}, X_{i}-X_{j}\right\rangle(t) & =\left\langle X_{i}, X_{i}\right\rangle(t)+\left\langle X_{j}, X_{j}\right\rangle(t)-2\left\langle X_{i}, X_{j}\right\rangle(t)  \tag{A.56}\\
& =2 t-2\left\langle X_{i}, X_{j}\right\rangle(t) .
\end{align*}
$$

Inserting this into A.55 yields

$$
\begin{equation*}
\int_{0}^{t} 1_{\left\{X_{i}(s)=X_{j}(s)\right\}} \mathrm{d} s=\int_{0}^{t} 1_{\left\{X_{i}(s)=X_{j}(s)\right\}} \mathrm{d}\left\langle X_{i}, X_{j}\right\rangle(s) \tag{A.57}
\end{equation*}
$$

On the other hand, A.53) implies that

$$
\begin{equation*}
\int_{0}^{t} 1_{\left\{X_{i}(s) \neq X_{j}(s)\right\}} \mathrm{d}\left\langle X_{i}, X_{j}\right\rangle(s)=0 . \tag{A.58}
\end{equation*}
$$

Combining this with A.57) yields

$$
\begin{equation*}
\int_{0}^{t} 1_{\left\{X_{i}(s)=X_{j}(s)\right\}} \mathrm{d} s=\left\langle X_{i}, X_{j}\right\rangle(t) \quad(t \geq 0,1 \leq i, j \leq n) \tag{A.59}
\end{equation*}
$$

as claimed.
By A.52) and A.59), we conclude that

$$
\begin{equation*}
\int_{0}^{t} 1_{\left\{\mathrm{Y}_{i}^{\langle k\rangle}(s)=\mathrm{Y}_{j}^{\langle k\rangle}(s)\right\}} \mathrm{d} s \underset{k \rightarrow \infty}{\longrightarrow} \int_{0}^{t} 1_{\left\{X_{i}(s)=X_{j}(s)\right\}} \mathrm{d} s \quad \forall t \geq 0,1 \leq i, j \leq n \quad \text { a.s. } \tag{A.60}
\end{equation*}
$$

The lower semicontinuity of the map $w \mapsto \int_{0}^{t} 1_{\left\{w_{i}(s)<w_{j}(s)\right\}} \mathrm{d} s$ implies that

$$
\begin{equation*}
\liminf _{k \rightarrow \infty} \int_{0}^{t} 1_{\left\{\mathrm{Y}_{i}^{\langle k\rangle}(s)<\mathrm{Y}_{j}^{\langle k\rangle}(s)\right\}} \mathrm{d} s \geq \int_{0}^{t} 1_{\left\{X_{i}(s)<X_{j}(s)\right\}} \mathrm{d} s \quad \forall t \geq 0,1 \leq i, j \leq n \quad \text { a.s. } \tag{A.61}
\end{equation*}
$$

Combining this with A.60 we see that a.s., for all $t \geq 0$ and $1 \leq i, j \leq n$,

$$
\begin{align*}
& \limsup _{k \rightarrow \infty} \int_{0}^{t} 1_{\left\{\mathrm{Y}_{i}^{\langle k\rangle}(s)>\mathrm{Y}_{j}^{\langle k\rangle}(s)\right\}} \mathrm{d} s \\
& \quad \leq 1-\lim _{k \rightarrow \infty} \int_{0}^{t} 1_{\left\{\mathrm{Y}_{i}^{\langle k\rangle}(s)=\mathrm{Y}_{j}^{\langle k\rangle}(s)\right\}} \mathrm{d} s-\liminf _{k \rightarrow \infty} \int_{0}^{t} 1_{\left\{\mathrm{Y}_{i}^{\langle k\rangle}(s)<\mathrm{Y}_{j}^{\langle k\rangle}(s)\right\}} \mathrm{d} s  \tag{A.62}\\
& \quad \leq 1-\int_{0}^{t} 1_{\left\{X_{i}(s)=X_{j}(s)\right\}} \mathrm{d} s-\int_{0}^{t} 1_{\left\{X_{i}(s)<X_{j}(s)\right\}} \mathrm{d} s=\int_{0}^{t} 1_{\left\{X_{i}(s)>X_{j}(s)\right\}} \mathrm{d} s
\end{align*}
$$

which together with A.61) shows that

$$
\begin{equation*}
\int_{0}^{t} 1_{\left\{\mathrm{Y}_{i}^{\langle k\rangle}(s)<\mathrm{Y}_{j}^{\langle k\rangle}(s)\right\}} \mathrm{d} s \underset{k \rightarrow \infty}{\longrightarrow} \int_{0}^{t} 1_{\left\{X_{i}(s)<X_{j}(s)\right\}} \mathrm{d} s \quad \forall t \geq 0,1 \leq i, j \leq n \quad \text { a.s. } \tag{A.63}
\end{equation*}
$$

By Lemma A. 7 below, this implies that

$$
\begin{equation*}
\int_{0}^{t}\left|1_{\left\{\mathrm{Y}_{i}^{\langle k\rangle}(s)<\mathrm{Y}_{j}^{\langle k\rangle}(s)\right\}}-1_{\left\{X_{i}(s)<X_{j}(s)\right\}}\right| \mathrm{d} s \underset{k \rightarrow \infty}{\longrightarrow} 0 \quad \forall t \geq 0,1 \leq i, j \leq n \quad \text { a.s. }, \tag{A.64}
\end{equation*}
$$

which in turn implies that a.s., for each $t \geq 0$ and weak total order $\prec$ on $\{1, \ldots, n\}$, one has

$$
\begin{equation*}
\int_{0}^{t}\left|1_{\left\{\mathrm{Y}^{\langle k\rangle}(s) \in C \prec\right\}}-1_{\{X(s) \in C<\}}\right| \mathrm{d} s \underset{k \rightarrow \infty}{\longrightarrow} 0 \tag{A.65}
\end{equation*}
$$

where $C_{\prec}$ is the cell defined in A.7). Since $g_{\Delta}(x)$ depends only on the relative order of the coordinates $x_{1}, \ldots, x_{n}$, formulas A.65 and A.50 imply that for each $\emptyset \neq \Delta \subset\{1, \ldots, n\}$,

$$
\begin{equation*}
\varepsilon_{k}^{-1} \int_{0}^{\lfloor t\rfloor_{k}} \beta_{k}\left(g_{\Delta}\left(\overrightarrow{\mathrm{Y}}^{\langle k\rangle}(s)\right)\right) \mathrm{d} s \underset{k \rightarrow \infty}{\longrightarrow} \int_{0}^{t} \beta_{+}\left(g_{\Delta}(\vec{X}(s))\right) \mathrm{d} s \quad \forall t \geq 0 \quad \text { a.s. } \tag{A.66}
\end{equation*}
$$

Taking the limit $k \rightarrow \infty$ in A.48 using A.66) and the fact that $g_{\Delta}$ is a bounded function (indeed, $1 \leq g_{\Delta}(x) \leq|\Delta|$ ), we find that for each $\emptyset \neq \Delta \subset\{1, \ldots, n\}$,

$$
\begin{equation*}
f_{\Delta}(\vec{X}(t))-\int_{0}^{t} \beta_{+}\left(g_{\Delta}(\vec{X}(s))\right) \mathrm{d} s \tag{A.67}
\end{equation*}
$$

is a martingale with respect to the filtration generated by $\vec{X}$. Together with A.59 this shows that $\vec{X}$ solves the Howitt-Warren martingale problem, completing our proof.

Lemma A. 7 (Convergence of integrals) Let $T>0$ and let $\rho$ be a finite measure on $[0, T]$.
(a) Let $f_{k}, f$ be Borel measurable real functions on $[0, T]$ such that $\sup _{k}\left\|f_{k}\right\|<\infty$, where $\|\cdot\|$ denotes the supremumnorm. Assume that

$$
\begin{equation*}
\int_{[0, t]} \rho(\mathrm{d} s) f_{k}(s) \underset{k \rightarrow \infty}{\longrightarrow} \int_{[0, t]} \rho(\mathrm{d} s) f(s) \quad(0 \leq t \leq T) \tag{A.68}
\end{equation*}
$$

Then

$$
\begin{equation*}
\int_{A} \rho(\mathrm{~d} s) f_{k}(s) \underset{k \rightarrow \infty}{\longrightarrow} \int_{A} \rho(\mathrm{~d} s) f(s) \tag{A.69}
\end{equation*}
$$

for each Borel measurable $A \subset[0, T]$.
(b) Let $A_{k}, A \subset[0, T]$ be Borel measurable. Assume that

$$
\begin{equation*}
\int_{[0, t]} \rho(\mathrm{d} s) 1_{A_{k}}(s) \underset{k \rightarrow \infty}{\longrightarrow} \int_{[0, t]} \rho(\mathrm{d} s) 1_{A}(s) \quad(0 \leq t \leq T) \tag{A.70}
\end{equation*}
$$

Then

$$
\begin{equation*}
\int_{[0, T]} \rho(\mathrm{d} s)\left|1_{A_{k}}(s)-1_{A}(s)\right| \underset{k \rightarrow \infty}{\longrightarrow} 0 \tag{A.71}
\end{equation*}
$$

Proof. To prove part (a), let $\mathcal{G}$ be the set of Borel measurable subsets $A \subset[0, T]$ for which A.69 holds. It is clear that $A, B \in \mathcal{G}, A \supset B$ implies that $A \backslash B \in \mathcal{G}$. We claim that moreover, if $A_{n} \in \mathcal{G}$ satisfy $A_{n} \uparrow A$ for some $A \subset[0, T]$, then $A \in \mathcal{G}$. To see this, write

$$
\begin{align*}
& \left|\int_{A} \rho(\mathrm{~d} s) f_{k}(s)-\int_{A} \rho(\mathrm{~d} s) f(s)\right| \\
& \quad \leq\left|\int_{A} \rho(\mathrm{~d} s) f_{k}(s)-\int_{A_{n}} \rho(\mathrm{~d} s) f_{k}(s)\right|+\left|\int_{A_{n}} \rho(\mathrm{~d} s) f(s)-\int_{A} \rho(\mathrm{~d} s) f(s)\right| \\
& \quad+\left|\int_{A_{n}} \rho(\mathrm{~d} s) f_{k}(s)-\int_{A_{n}} \rho(\mathrm{~d} s) f(s)\right|  \tag{A.72}\\
& \quad \leq 2 \rho\left(A \backslash A_{n}\right) \sup _{m}\left\|f_{m}\right\|+\left|\int_{A_{n}} \rho(\mathrm{~d} s) f_{k}(s)-\int_{A_{n}} \rho(\mathrm{~d} s) f(s)\right|
\end{align*}
$$

By choosing $n$ large enough, we see that

$$
\begin{equation*}
\limsup _{k \rightarrow \infty}\left|\int_{A} \rho(\mathrm{~d} s) f_{k}(s)-\int_{A} \rho(\mathrm{~d} s) f(s)\right| \leq \varepsilon \tag{A.73}
\end{equation*}
$$

for all $\varepsilon>0$, proving our claim. Since the set $\mathcal{H}:=\{[0, t]: 0 \leq t \leq T\}$ is closed under intersections and contained in $\mathcal{G}$, Sierpinski's $\pi / \lambda$-theorem Kal02, Theorem 1.1] tells us that $\mathcal{G}$ contains the $\sigma$-field generated by $\mathcal{H}$, completing our proof.

To prove part (b), we note that

$$
\begin{align*}
\int_{[0, T]} \rho(\mathrm{d} s)\left|1_{A_{k}}(s)-1_{A}(s)\right|= & \left(\int_{A} \rho(\mathrm{~d} s) 1_{A}(s)-\int_{A} \rho(\mathrm{~d} s) 1_{A_{k}}(s)\right) \\
& +\left(\int_{[0, T] \backslash A} \rho(\mathrm{~d} s) 1_{A_{k}}(s)-\int_{[0, T] \backslash A} \rho(\mathrm{~d} s) 1_{A}(s)\right) \tag{A.74}
\end{align*}
$$

which tends to zero by part (a).

## B The Hausdorff topology

Let $(E, d)$ be a metric space, let $\mathcal{K}(E)$ be the space of all compact subsets of $E$ and set $\mathcal{K}_{+}(E):=\{K \in \mathcal{K}(E): K \neq \emptyset\}$. Then the Hausdorff metric $d_{\mathrm{H}}$ on $\mathcal{K}_{+}(E)$ is defined as

$$
\begin{align*}
d_{\mathrm{H}}\left(K_{1}, K_{2}\right) & :=\sup _{x_{1} \in K_{1}} \inf _{x_{2} \in K_{2}} d\left(x_{1}, x_{2}\right) \vee \sup _{x_{2} \in K_{2}} \inf _{x_{1} \in K_{1}} d\left(x_{1}, x_{2}\right)  \tag{B.1}\\
& =\sup _{x_{1} \in K_{1}} d\left(x_{1}, K_{2}\right) \vee \sup _{x_{2} \in K_{2}} d\left(x_{2}, K_{1}\right),
\end{align*}
$$

where $d(x, A):=\inf _{y \in A} d(x, y)$ denotes the distance between a point $x \in E$ and a set $A \subset E$. The corresponding topology is called the Hausdorff topology. We extend this topology to $\mathcal{K}(E)$ by adding $\emptyset$ as an isolated point. The next lemma shows that the Hausdorff topology depends only on the topology on $E$, and not on the choice of the metric.

Lemma B. 1 (Convergence criterion) Let $K_{n}, K \in \mathcal{K}_{+}(E)(n \geq 1)$. Then $K_{n} \rightarrow K$ in the Hausdorff topology if and only if there exists a $C \in \mathcal{K}_{+}(E)$ such that $K_{n} \subset C$ for all $n \geq 1$ and

$$
\begin{align*}
K & =\left\{x \in E: \exists x_{n} \in K_{n} \text { s.t. } x_{n} \rightarrow x\right\}  \tag{B.2}\\
& =\left\{x \in E: \exists x_{n} \in K_{n} \text { s.t. } x \text { is a cluster point of }\left(x_{n}\right)_{n \in \mathbb{N}}\right\} .
\end{align*}
$$

The following lemma shows that $\mathcal{K}(E)$ is Polish if $E$ is.

## Lemma B. 2 (Properties of the Hausdorff metric)

(a) If $(E, d)$ is separable, then so is $\left(\mathcal{K}_{+}(E), d_{\mathrm{H}}\right)$.
(b) If $(E, d)$ is complete, then so is $\left(\mathcal{K}_{+}(E), d_{\mathrm{H}}\right)$.

Recall that a subset $A$ of a metric space is precompact if its closure is compact. This is equivalent to the statement that each sequence of points $x_{n} \in A$ has a convergent subsequence.

Lemma B. 3 (Compactness in the Hausdorff topology) A set $\mathcal{A} \subset \mathcal{K}(E)$ is precompact if and only if there exists a $C \in \mathcal{K}(E)$ such that $K \subset C$ for each $K \in \mathcal{A}$.

The following lemma is useful when proving convergence of $\mathcal{K}(E)$-valued random variables.
Lemma B. 4 (Tightness criterion) Assume that $E$ is a Polish space and let $K_{n}(n \geq 1)$ be $\mathcal{K}(E)$-valued random variables. Then the collection of laws $\left\{\mathbb{P}\left[K_{n} \in \cdot\right]: n \geq 1\right\}$ is tight if and only if for each $\varepsilon>0$ there exists a compact $C \subset E$ such that $\mathbb{P}\left[K_{n} \subset C\right] \geq 1-\varepsilon$ uniformly in $n \in \mathbb{N}$.

If $E$ is compact, then the Hausdorff topology on $\mathcal{K}(E)$ coincides with the Fell topology defined in Kal02, Thm. A.2.5]. The Hausdorff metric may more generally be defined on the space of nonempty bounded closed subsets of $(E, d)$. In particular, if $d$ is bounded, then $d_{\mathrm{H}}\left(A_{1}, A_{2}\right)$ can be defined for any nonempty closed $A_{1}, A_{2}$. In this more general setup, Lemma B. 2 (b) and the 'if' part of Lemma B.3 remain true, as well as the 'if' part of Lemma B. 5 below. This is Excercise 7 (with some hints for a possible solution) in Mun00, § 45]. A detailed solution of this excercise can be found in Hen99. We are not aware of any reference for the other statements in Lemmas B. 1 B. 4 , although they appear to be well-known.

For completeness, we provide self-contained proofs of all these lemmas. We start with some preparations.

Recall that for any metric space $(E, d)$, a set $A \subset E$ is totally bounded if for every $\varepsilon>0$ there exists a finite collection of points $x_{1}, \ldots, x_{n} \in E$ such that $A \subset \bigcup_{i=1}^{n} B_{\varepsilon}\left(x_{i}\right)$, where $B_{\varepsilon}(x)$ denotes the open ball of radius $\varepsilon$ around $x$. This is equivalent to the statement that every sequence $x_{n} \in A$ has a Cauchy subsequence. As a consequence, a set $A \subset E$ is compact if and only if it is complete and totally bounded.

Lemma B. 5 (Totally bounded sets in the Hausdorff metric) $A$ set $\mathcal{A} \subset \mathcal{K}_{+}(E)$ is totally bounded in $\left(\mathcal{K}_{+}(E), d_{\mathrm{H}}\right)$ if and only if the set $A:=\{x \in E: \exists K \in \mathcal{A}$ s.t. $x \in K\}$ is totally bounded in $(E, d)$.

Proof. Assume that $A$ is bounded. Let $\varepsilon>0$ and let $\Delta \subset E$ be a finite set such that $E=\bigcup_{x \in \Delta} B_{\varepsilon}(x)$. Let $K \in \mathcal{K}_{+}(E)$ and set $\Delta^{\prime}:=\left\{x \in \Delta: B_{\varepsilon}(x) \cap K \neq \emptyset\right\}$. Then for all $y \in K$ there is an $x \in \Delta^{\prime}$ such that $d(x, y)<\varepsilon$ and for all $x \in \Delta^{\prime}$ there is a $y \in K$ such that $d(x, y)<\varepsilon$ proving that $d_{\mathrm{H}}\left(\Delta^{\prime}, K\right)<\varepsilon$. This shows that $\mathcal{A}$ is covered, in the Hausdorff metric, by the collection of open balls of radius $\varepsilon$ centered around finite subsets of $\Delta$. Since $\varepsilon$ is general, we conclude that $\mathcal{A}$ is totally bounded.

Conversely, if $\mathcal{A}$ is totally bounded, then for each $\varepsilon>0$ we can find $K_{1}, \ldots, K_{n} \in \mathcal{K}_{+}(E)$ such that $\mathcal{A} \subset \bigcup_{k=1}^{n} \mathcal{B}_{\varepsilon / 2}\left(K_{n}\right)$, where $\mathcal{B}_{\varepsilon}(K)$ denotes the open ball in the Hausdorff metric of radius $\varepsilon$ centered around a compact set $K$. Since each $K_{k}$ is compact, there exist $x_{k, 1}, \ldots, x_{k, m_{k}}$ such that $K_{k} \subset \bigcup_{j=1}^{m_{k}} B_{\varepsilon / 2}\left(x_{k, j}\right)$, hence $A \subset \bigcup_{k=1}^{n} \bigcup_{j=1}^{m_{k}} B_{\varepsilon}\left(x_{k, j}\right)$.

Lemma B. 6 (Cauchy sequences in the Hausdorff metric) Let $K_{n} \in \mathcal{K}_{+}(E)$ be a Cauchy sequence in $\left(\mathcal{K}_{+}(E), d_{\mathrm{H}}\right)$. Then there exists a closed set $K$ such that (B.2) holds.

Proof. If the sets on the first and second line of the right-hand side of ( $\overline{\mathrm{B} .2}$ ) are not equal, then there exists some $x \in E$ such that $x$ is a cluster point of some $x_{n} \in K_{n}$ but there do not exist $x_{n}^{\prime} \in K_{n}$ such that $x_{n}^{\prime} \rightarrow x$. It follows that there is some $\varepsilon>0$ such that for each $k \geq 1$ we can find $n, m \geq k$ such that $K_{n} \cap B_{\varepsilon}(x) \neq \emptyset$ and $K_{m} \cap B_{2 \varepsilon}(x)=\emptyset$, hence $d_{\mathrm{H}}\left(K_{n}, K_{m}\right) \geq \varepsilon$, contradicting the asumption that the $K_{n}$ form a Cauchy sequence.

To see that $K$ is closed, assume that $x_{n} \in K$ satisfy $x_{n} \rightarrow x$ for some $x \in E$. Since $d_{\mathrm{H}}\left(K_{n}, K\right) \rightarrow 0$ we can choose $x_{n}^{\prime} \in K_{n}$ such that $d\left(x_{n}^{\prime}, x_{n}\right) \rightarrow 0$. It follows that $d\left(x_{n}^{\prime}, x\right) \leq$ $d\left(x_{n}^{\prime}, x_{n}\right)+d\left(x_{n}, x\right) \rightarrow 0$ and hence $x \in K$.

Lemma B. 7 (Sufficient conditions for convergence) The conditions for convergence in the Hausdorff topology given in Lemma B.1 are sufficient.

Proof. Our assumptions imply that $d\left(x, K_{n}\right) \rightarrow 0$ for each $x \in K$. We wish to show that in fact $\sup _{x \in K} d\left(x, K_{n}\right) \rightarrow 0$. If this is not the case, then by going to a subsequence if necessary we may assume that there exist $x_{n} \in K$ and $\varepsilon>0$ such that $\liminf _{n \rightarrow \infty} d\left(x_{n}, K_{n}\right) \geq \varepsilon$. Since $K$ is compact, by going to a further subsequence if necessary, we may assume that $x_{n} \rightarrow x \in K$. But then $\liminf _{n \rightarrow \infty} d\left(x, K_{n}\right) \geq \liminf _{n \rightarrow \infty}\left(d\left(x_{n}, K_{n}\right)-d\left(x, x_{n}\right)\right) \geq \varepsilon$ for this subsequence, contradicting the fact that for the original sequence, $d\left(x, K_{n}\right) \rightarrow 0$ for each $x \in K$.

The proof that $\sup _{x \in K_{n}} d(x, K) \rightarrow 0$ is similar. If this is not true, then we can go to a subsequence of the $K_{n}$ and then find $x_{n} \in K_{n}$ such that $d\left(x_{n}, K\right) \geq \varepsilon$ for all $n$, for some $\varepsilon>0$. Using the compactness of $C$, we can select a further subsequence such that $x_{n} \rightarrow x \in C$. Now
$x$ is a cluster point of some $x_{n} \in K_{n}$ but $d(x, K) \geq \varepsilon$, contradicting the fact that the two sets on the right-hand side of $(\overline{\mathrm{B} .2})$ are equal.

Proof of Lemma B.2. To prove part (a), it suffices to show that if $\mathcal{D}$ is a countable dense subset of $(E, d)$, then the collection of finite subsets of $\mathcal{D}$ is a countable dense subset of $\left(\mathcal{K}_{+}(E), d_{\mathrm{H}}\right)$. Since a compact set $K \subset E$ is totally bounded, for each $\varepsilon>0$, we can find a finitely many points $x_{1}, \ldots, x_{n} \in E$ such that $K \subset \bigcup_{i=1}^{n} B_{\varepsilon / 2}\left(x_{i}\right)$. Since $\mathcal{D}$ is dense, we can choose $x_{i}^{\prime} \in \mathcal{D} \cap B_{\varepsilon / 2}\left(x_{i}\right)$. Then $d_{\mathrm{H}}\left(K,\left\{x_{1}^{\prime}, \ldots, x_{n}^{\prime}\right\}\right) \leq \varepsilon$, proving our claim.

To prove part (b), let $K_{n} \in \mathcal{K}_{+}(E)$ be a Cauchy sequence. Then, by Lemma B.6, there exists a closed set $K$ such that (B.2) holds. Since each sequence in the set $\left\{K_{n}: n \geq 1\right\}$ contains a Cauchy subsequence, the set $\left\{K_{n}: n \geq 1\right\}$ is totally bounded, hence by Lemma B.5. there exists some totally bounded set containing all of the $K_{n}$. Let $C$ denote its closure. Then $C$ is compact since $E$ is complete, hence also $K \subset C$ is compact and Lemma B. 7 implies that $K_{n} \rightarrow K$.

Proof of Lemma B.3. It suffices to prove the statement for $\mathcal{A} \subset \mathcal{K}_{+}(E)$. Let $\overline{\mathcal{A}}$ be the closure of $\mathcal{A}$ and set $C:=\{x \in E: \exists K \in \overline{\mathcal{A}}$ s.t. $x \in K\}$. By Lemma B.5, $\overline{\mathcal{A}}$ is totally bounded if and only if $A$ is. Moreover, by Lemma B.2 (b), if $A$ is complete then so is $\left\{K \in \mathcal{K}_{+}(E): K \subset C\right\}$ and hence the same is true for $\overline{\mathcal{A}}$, being a closed subset of the former. Therefore, since compactness is equivalent to total boundedness and completeness, it suffices to show that compactness of $\overline{\mathcal{A}}$ implies completeness of $C$. Assume that $\overline{\mathcal{A}}$ is compact and that $x_{n} \in C$ is a Cauchy sequence. We need to show that the sequence $x_{n}$ has a cluster point $x \in C$. Choose $K_{n} \in \overline{\mathcal{A}}$ such that $x_{n} \in K_{n}$. Since $\overline{\mathcal{A}}$ is compact, by going to a subsequence if necessary, we may assume that $K_{n} \rightarrow K$ for some $K \in \overline{\mathcal{A}}$. Choose $x_{n}^{\prime} \in K$ such that $d\left(x_{n}, x_{n}^{\prime}\right) \rightarrow 0$. Since $K$ is compact, by going to a further subsequence if necessary, we may assume that $x_{n}^{\prime} \rightarrow x$ for some $x \in K$. Since $d\left(x_{n}, x\right) \leq d\left(x_{n}, x_{n}^{\prime}\right)+d\left(x_{n}^{\prime}, x\right) \rightarrow 0$ this proves that the sequence $x_{n}$ has a cluster point $x \in K \subset C$.

Proof of Lemma B.4. Immediate from Lemma B. 3 and the definition of tightness.
Proof of Lemma B.1. By Lemma B.7, we only need to prove that if $K_{n} \in \mathcal{K}_{+}(E)$ converge to a limit $K$, then there exists a $C \in \mathcal{K}_{+}(E)$ such that $K_{n} \subset C$ for all $n$ and (B.2) holds. If $K_{n} \rightarrow K$ then the set $\left\{K_{n}: n \geq 1\right\}$ is precompact, hence by Lemma B. 4 there exists a $C \in \mathcal{K}_{+}(E)$ such that $K_{n} \subset C$ for all $n$. Formula B.2) follows from the facts that if $x \in K$, then $d\left(x, K_{n}\right) \rightarrow 0$ hence there exist $K_{n} \ni x_{n} \rightarrow x$, while if $x \notin K$, then $B_{\varepsilon}(x) \cap K_{n}=\emptyset$ for all $n$ large enough such that $\sup _{x^{\prime} \in K} d\left(x^{\prime}, K_{n}\right)<\varepsilon$, hence $x$ is not a cluster point of some $x_{n} \in K_{n}$.

## C Some measurability issues

Let $E, F$ be Polish spaces. By definition, the pointwise closure of a set $\mathcal{F}$ of functions $f$ : $E \rightarrow F$ is the smallest set containing $\mathcal{F}$ that is closed under taking of pointwise limits, i.e., it is the intersection of all sets $\mathcal{G}$ of functions from $E$ to $F$, such that $\mathcal{G} \supset \mathcal{F}$ and $f_{n} \in \mathcal{G}$, $\lim _{n \rightarrow \infty} f_{n}(x)=f(x)(x \in E)$ imply $f \in \mathcal{G}$.

Lemma C. 1 (Pointwise closure of functions to the unit interval) Let $E$ be a Polish space and let $\mathcal{C}_{[0,1]}(E)$ be the set of all continuous functions $f: E \rightarrow[0,1]$. Then the pointwise closure of $\mathcal{C}_{[0,1]}(E)$ is the set $B_{[0,1]}(E)$ of all Borel measurable functions $f: E \rightarrow[0,1]$. If
$E$ is locally compact, then the same conclusion holds with $\mathcal{C}_{[0,1]}(E)$ replaced by the space of continuous and compactly supported functions $f: E \rightarrow[0,1]$.

Proof. By definition, one says that a sequence $f_{n}$ of real functions on $E$ converges in a bounded pointwise way to a limit $f$ if $f_{n}(x) \rightarrow f(x)$ for each $x \in E$ and there exists some constant $C>0$ such that $\left|f_{n}\right| \leq C$ for all $n \geq 0$. The bp-closure of a set $\mathcal{F}$ of real functions on $E$ is the smallest set containing $\mathcal{F}$ that is closed under taking of bounded pointwise limits. By copying the proof of [EK86, Lemma 3.4.1], we see that the bp-closure of a convex set is convex. Let $\mathcal{B}$ be the set of all subsets $A \subset E$ such that $1_{A}$ is in the bp-closure of $\mathcal{C}_{[0,1]}(E)$. Then $\mathcal{B}$ is a Dynkin class containing all open sets, hence by the Dynkin class theorem [EK86, Thm. A.4.2] (resp. the $\pi / \lambda$-theorem Kal02, Theorem 1.1]), $\mathcal{B}$ contains all Borel measurable subsets of $E$. Since indicator functions are the extremal elements of the convex set consisting of all simple functions in $B_{[0,1]}(E)$, it is easy to see that every simple function can be written as a convex combinations of indicator functions. Since every function in $B_{[0,1]}(E)$ is an increasing limit of simple functions in $B_{[0,1]}(E)$, the first claim follows. In case $E$ is locally compact, it is easy to see that each continuous function $f: E \rightarrow[0,1]$ is the pointwise limit of compactly supported continuous functions $f: E \rightarrow[0,1]$, proving the second claim.
We will need the following generalization of Lemma C.1. Below, $[0,1]^{\mathbb{N}}$ denotes the space of all functions $x: \mathbb{N} \rightarrow[0,1]$, equipped with the product topology. Note that the statement of Lemma C. 2 is false if we replace $[0,1]^{\mathbb{N}}$ by a general compact metrizable space $F$. E.g., it is already wrong if $F$ consists of two isolated points, since in this case all continuous functions are constant but there are lots of measurable functions.

Lemma C. 2 (Pointwise closure) Let $E$ be a Polish space and let $\mathcal{C}_{[0,1]^{\mathbb{N}}}(E)$ be the set of all continuous functions $f: E \rightarrow[0,1]^{\mathbb{N}}$. Then the pointwise closure of $\mathcal{C}_{[0,1]^{\mathbb{N}}}(E)$ is the set $B_{[0,1]^{\mathbb{N}}}(E)$ of all Borel measurable functions $f: E \rightarrow[0,1]^{\mathbb{N}}$.

Proof. Let $E, F, G$ be Polish spaces and let $\mathcal{F}, \mathcal{G}$ be sets of functions $f: E \rightarrow F$ and $g: E \rightarrow G$, respectively. We claim that $\operatorname{pclos}(\mathcal{F} \times \mathcal{G}) \supset \operatorname{pclos}(\mathcal{F}) \times \operatorname{pclos}(\mathcal{G})$, where $\operatorname{pclos}(\cdot)$ denotes the pointwise closure of a set and we regard a pair of functions $(f, g)$ as a function from $E$ to $F \times G$ (equipped with the product topology). To prove our claim, for any $f \in \operatorname{pclos}(\mathcal{F})$, let $\mathcal{G}_{f}$ be the space of functions $g \in \operatorname{pclos}(\mathcal{G})$ such that $(f, g) \in \operatorname{pclos}(\mathcal{F} \times \mathcal{G})$. Then $\mathcal{G}_{f}$ is closed under pointwise limits since $\operatorname{pclos}(\mathcal{F} \times \mathcal{G})$ is. If $f \in \mathcal{F}$, then moreover $\mathcal{G}_{f}$ contains $\mathcal{G}$ so $\mathcal{G}_{f}=\mathcal{G}$. Next, let $\hat{\mathcal{F}}$ be the space of functions $f \in \operatorname{pclos}(\mathcal{F})$ such that $(f, g) \in \operatorname{pclos}(\mathcal{F} \times \mathcal{G})$ for all $g \in \operatorname{pclos}(\mathcal{G})$. Then $\hat{\mathcal{F}}$ is closed under pointwise limits since $\operatorname{pclos}(\mathcal{F} \times \mathcal{G})$ is and $\hat{\mathcal{F}}$ contains $\mathcal{F}$ by what we have just proved, so $\hat{\mathcal{F}}=\operatorname{pclos}(\mathcal{F})$, proving our claim.

Applying our clain inductively to $\mathcal{C}_{[0,1]^{\mathbb{N}}}(E)=\left(\mathcal{C}_{[0,1]}(E)\right)^{\mathbb{N}}$, using Lemma C.1, we see that $\left(f_{1}, \ldots, f_{n}, 0, \ldots\right)$ lies in the pointwise closure of $\mathcal{C}_{[0,1]^{\mathrm{N}}}(E)$ for each $f_{1}, \ldots, f_{n} \in B_{[0,1]}(E)$ and $n \geq 1$. By taking pointwise limits, we see that each infinite sequence ( $f_{1}, f_{2}, \ldots$ ) of Borel measurable functions $f_{i}: E \rightarrow[0,1]$ lies in the pointwise closure of $\mathcal{C}_{[0,1]^{\mathbb{N}}}(E)$.

Lemma C. 3 (Measurability of image measure map) Let $E, F, G$ be Polish spaces and let $\mathcal{M}(E), \mathcal{M}(F)$ be the spaces of finite measures on $E$ and $F$, respectively, equipped with the topology of weak convergence and the associated Borel $\sigma$-field. Then, for any measurable map $E \times G \ni(x, z) \mapsto f_{z}(x) \in F$, setting $\psi_{z}^{f}(\mu):=\mu \circ f_{z}^{-1}$ defines a measurable map $\mathcal{M}_{1}(E) \times G \ni(\mu, z) \mapsto \psi_{z}^{f}(\mu) \in \mathcal{M}_{1}(F)$.

Proof. We first prove the statement if $E, G$ are compact and $F=[0,1]^{\mathbb{N}}$. In this case, we claim that if $E \times G \ni(x, z) \mapsto f_{z}(x) \in F$ is continuous, then also $\mathcal{M}_{1}(E) \times G \ni(\mu, z) \mapsto$ $\psi_{z}^{f}(\mu) \in \mathcal{M}_{1}(F)$ is continuous. To see this, it suffices to observe that $\mu_{n} \Rightarrow \mu$ and $z_{n} \rightarrow z$ imply that for any continuous $h: F \rightarrow \mathbb{R}$,

$$
\begin{align*}
& \left|\int \psi_{z}^{f}(\mu)(\mathrm{d} y) h(y)-\int \psi_{z_{n}}^{f}\left(\mu_{n}\right)(\mathrm{d} y) h(y)\right|=\left|\int \mu(\mathrm{d} x) h\left(f_{z}(x)\right)-\int \mu_{n}(\mathrm{~d} x) h\left(f_{z_{n}}(x)\right)\right| \\
& \quad=\left|\int \mu(\mathrm{d} x) h\left(f_{z}(x)\right)-\int \mu_{n}(\mathrm{~d} x) h\left(f_{z}(x)\right)\right|+\left|\int \mu_{n}(\mathrm{~d} x) h\left(f_{z}(x)\right)-\int \mu_{n}(\mathrm{~d} x) h\left(f_{z_{n}}(x)\right)\right| \tag{C.1}
\end{align*}
$$

Here the first term on the right-hand side converges to zero by our assumption that $\mu_{n}$ converges weakly to $\mu$, while the second term can be bounded by $\left\|h \circ f_{z_{n}}-h \circ f_{z}\right\|_{\infty}$, which tends to zero since $E \times G \ni(x, z) \mapsto h \circ f_{z}(x) \in \mathbb{R}$ is continuous and $E, G$ are compact spaces.

We next claim that if $f^{n} \rightarrow f$ pointwise, then also $\psi^{f^{n}} \rightarrow \psi^{f}$ pointwise. Indeed, if $f_{z}^{n}(x) \rightarrow f_{z}(x)$ for all $x, z$, then, for any continuous (and hence bounded) $h: F \rightarrow \mathbb{R}$,

$$
\begin{equation*}
\int \psi_{z}^{f^{n}}(\mu)(\mathrm{d} y) h(y)=\int \mu(\mathrm{d} x) h\left(f_{z}^{n}(x)\right) \underset{n \rightarrow \infty}{\longrightarrow} \int \mu(\mathrm{~d} x) h\left(f_{z}(x)\right)=\int \psi_{z}^{f}(\mu)(\mathrm{d} y) h(y) \tag{C.2}
\end{equation*}
$$

showing that $\psi_{z}^{f^{n}}(\mu) \underset{n \rightarrow \infty}{\Longrightarrow} \psi_{z}^{f}(\mu)$ for all $\mu, z$. It follows that the set $\mathcal{G}$ of all $E \times G \ni(x, z) \mapsto$ $f_{z}(x) \in F$ such that $\mathcal{M}_{1}(E) \times G \ni(\mu, z) \mapsto \psi_{z}^{f}(\mu) \in \mathcal{M}_{1}(F)$ is measurable is closed under pointwise limits and contains all continuous functions $(x, z) \mapsto f_{x}(x)$. By Lemma C.2, it follows that $\mathcal{G}$ contains all measurable $(x, z) \mapsto f_{z}(x)$.

To treat the general case, where $E, G$ need not be compact and $F$ may be different from $[0,1]^{\mathbb{N}}$, we will use a compactification argument. We need the following three facts: 1. Each separable metrizable space is isomorphic to a subset of $[0,1]^{\mathbb{N}} .2$. A subset of a Polish space is Polish in the induced topology if and only if it is a $G_{\delta}$-set, i.e., a countable intersection of open sets Bou58, $\S 6$ No. 1, Thm. 1]. 3. If $E_{1} \subset E_{2}$ are Polish spaces and $\mathcal{M}_{1}\left(E_{i}\right)$ is the space of probability measures on $E_{i}(i=1,2)$, equipped with the topology of weak convergence, then $\mathcal{M}_{1}\left(E_{1}\right)$ is isomorphic to the set $\left\{\mu \in \mathcal{M}_{1}\left(E_{2}\right): \mu\left(E_{1}\right)=1\right\}$. (The fact that the topology on $\mathcal{M}_{1}\left(E_{1}\right)$ coincides with the one induced by the embedding in $\mathcal{M}_{1}\left(E_{2}\right)$ follows, for example, from Skorohod's representation theorem [Bil99, Theorem 6.7].) Note that facts 2 and 3 and the fact that $\mathcal{M}_{1}\left(E_{i}\right)(i=1,2)$ are Polish spaces imply that $\mathcal{M}_{1}\left(E_{1}\right)$ is a $G_{\delta}$-subset of $\mathcal{M}_{1}\left(E_{2}\right)$.

In view of facts 1 and 2 above, we may without loss of generality assume that $E, G$ are $G_{\delta}$-subsets of some compact metrizable spaces $\bar{E}, \bar{G}$ and that $F$ is a $G_{\delta}$-subset of $[0,1]^{\mathbb{N}}$. Then each measurable function $E \times G \ni(x, z) \mapsto f_{z}(x) \in F$ may be extended to a measurable function from $\bar{E} \times \bar{G}$ to $[0,1]^{\mathbb{N}}$ by setting $f_{z}(x)$ equal to some constant if $(x, z) \notin E \times G$. By what we have already proved, the associated map $\mathcal{M}_{1}(\bar{E}) \times \bar{G} \ni(\mu, z) \mapsto \psi_{z}^{f}(\mu) \in \mathcal{M}_{1}\left([0,1]^{\mathbb{N}}\right)$ is measurable. Since $\mathcal{M}_{1}(E)$ and $G$ are measurable subsets of $\mathcal{M}_{1}(\bar{E})$ and $\bar{G}$, respectively, the restriction of the map $(\mu, z) \mapsto \psi_{z}^{f}(\mu)$ to $\mathcal{M}_{1}(E) \times G$ yields a measurable map from $\mathcal{M}_{1}(E) \times G$ to $F$.

## D Thinning and Poissonization

Let $E$ be a Polish space and let $\mathcal{M}(E)$ be the space of finite measures on $E$ equipped with the topology of weak convergence, under which it is Polish. We let $\mathcal{M}_{\text {count }}(E)$ denote the space of finite counting measures on $E$, i.e., measures of the form $\sum_{i=1}^{n} \delta_{x_{i}}$ with $n \geq 0$ and
$x_{1}, \ldots, x_{n} \in E$. Since $\mathcal{M}_{\text {count }}(E)$ is a closed subset of $\mathcal{M}(E)$, it is also Polish (under the topology of weak convergence). We let $B_{+}(E)$ denote the space of measurable functions $f: E \rightarrow[0, \infty)$ and write $B_{[0,1]}(E)$ for the space of measurable functions $f: E \rightarrow[0,1]$. For any $f \in B_{[0,1]}(E)$ and $\nu \in \mathcal{M}_{\text {count }}(E)$, we introduce the notation

$$
\begin{equation*}
f^{\nu}:=\prod_{i=1}^{n} f\left(x_{i}\right) \quad \text { where } \quad \nu=\sum_{i=1}^{n} \delta_{x_{i}}, \tag{D.1}
\end{equation*}
$$

with the convention that $f^{0}:=1$. Let $\mu \in \mathcal{M}(E)$. By definition, a Poisson point measure with intensity $\mu$ is an $\mathcal{M}_{\text {count }}(E)$-valued random variable $\nu$ such that

$$
\begin{equation*}
\mathbb{E}\left[(1-f)^{\nu}\right]=e^{-\int f \mathrm{~d} \mu} \quad\left(f \in B_{[0,1]}(E)\right) . \tag{D.2}
\end{equation*}
$$

An explicit way to construct such a Poisson point measure is to write $\mu=\lambda \mu^{\prime}$ where $\lambda \geq 0$ and $\mu^{\prime}$ is a probability measure, and to put $\nu=\sum_{i=1}^{N} \delta_{X_{i}}$ where $\left(X_{i}\right)_{i \geq 1}$ are i.i.d. with law $\mu^{\prime}$ and $N$ is an independent Poisson distributed random variable with mean $\lambda$. By [Res87, Prop. 3.5], the law of $\nu$ is uniquely characterized by (D.2). The proof there is stated for locally compact spaces only, which in the present paper is actually all we need, but the statement holds more generally for Polish spaces. If $\mu$ is nonatomic, then $\nu$ a.s. contains no double points, i.e.,

$$
\begin{equation*}
\nu=\sum_{x \in \operatorname{supp}(\nu)} \delta_{x} \quad \text { a.s. }, \tag{D.3}
\end{equation*}
$$

see Kal02, Prop. 10.4]. In this case, we call $\operatorname{supp}(\nu)$ a Poisson point set with intensity $\mu$.
If $\nu \in \mathcal{M}_{\text {count }}(E)$ is a (deterministic) finite counting measure and $g \in B_{[0,1]}(E)$, then by definition a $g$-thinning of $\nu$ is an $\mathcal{M}_{\text {count }}(E)$-valued random variable $\nu$ such that

$$
\begin{equation*}
\mathbb{E}\left[(1-f)^{\nu^{\prime}}\right]=(1-g f)^{\nu} \quad\left(f \in B_{[0,1]}(E)\right) \tag{D.4}
\end{equation*}
$$

An explicit way to construct such a $g$-thinning, when $\nu=\sum_{i=1}^{n} \delta_{x_{i}}$, is to construct independent $\{0,1\}$-valued random variables $\chi_{1}, \ldots, \chi_{n}$ with $\mathbb{P}\left[\chi_{i}=1\right]=g\left(x_{i}\right)$ and to put $\nu^{\prime}:=\sum_{i=1}^{n} \chi_{i} \delta_{x_{i}}$. By [Res87, Prop. 3.5], the law of $\nu^{\prime}$ is uniquely characterized by (D.4).

It is easy to see that the class of functions $f: E \rightarrow[0,1]$ for which (D.2) or (D.4) hold is closed under taking of pointwise limits. Therefore, by Lemma C.1, in order to check (D.2) or (D.4), it suffices to verify the relation for all continuous functions $f: E \rightarrow[0,1]$, and in case $E$ is locally compact, even the continuous functions with compact support suffice.

We also need Poisson point sets with $\sigma$-finite, but in general locally infinite intensities. To this aim, let $\operatorname{Count}(E)$ be the space of all countable subsets of $E$. We equip $\operatorname{Count}(E)$ with the $\sigma$-field generated by all mappings $A \mapsto 1_{\{A \cap B=\emptyset\}}$ where $B \subset E$ is Borel measurable.

Lemma D. 1 (Poisson point sets with $\sigma$-finite intensity) For each $f \in B_{[0,1]}(E)$, the map $\operatorname{Count}(E) \ni A \mapsto \prod_{x \in A}(1-f(x)) \in[0,1]$ is measurable. Moreover, for each $\sigma$-finite nonatomic measure $\mu$ on $E$, there exists a Count $(E)$-valued random variable $C$, unique in law, such that

$$
\begin{equation*}
\mathbb{E}\left[\prod_{x \in C}(1-f(x))\right]=e^{-\int f \mathrm{~d} \mu} \quad\left(f \in B_{[0,1]}(E)\right) \tag{D.5}
\end{equation*}
$$

where $e^{-\infty}:=0$.

Proof. We claim that for all (Borel) measurable $B \subset E$, the function Count $(E) \ni A \mapsto$ $|A \cap B| \in\{0,1, \ldots\} \cup\{\infty\}$ is measurable. To see this, let $\mathcal{D} \subset E$ be countable and dense and let $\mathcal{O}:=\left\{B_{1 / k}(x): x \in \mathcal{D}, k \geq 1\right\}$, where $B_{\varepsilon}(x)$ denotes the open ball of radius $\varepsilon$ around $x$. Then

$$
\begin{equation*}
\{|A \cap B| \geq n\}=\left\{\exists U_{1}, \ldots, U_{n} \in \mathcal{O} \text { disjoint, s.t. } A \cap B \cap U_{i} \neq \emptyset \forall i=1, \ldots, n\right\} \tag{D.6}
\end{equation*}
$$

is a countable union of finite intersections of measurable sets, and hence itself measurable. It follows that $A \mapsto \sum_{x \in A} f(x)$ is measurable for each $f$ of the form $f=\sum_{i=1}^{n} b_{i} 1_{B_{i}}$ with $B_{i}$ (Borel) measurable and $b_{i} \in[0, \infty)$. By taking increasing limits it follows that $A \mapsto \sum_{x \in A} f(x)$ is measurable for each measurable $f: E \rightarrow[0, \infty]$. Since $\prod_{x \in A}(1-f(x))=\exp \left\{\sum_{x \in A} \log (1-\right.$ $f(x))\}$, we conclude that $A \mapsto \prod_{x \in A}(1-f(x))$ is measurable for each $f \in B_{[0,1]}(E)$.

Since $\mu$ is $\sigma$-finite, there exist disjoint measurable $B_{i} \subset E$ such that $\mu\left(B_{i}\right)<\infty(i \geq 1)$. Let $C_{i}$ be independent Poisson point sets with intensity $\mu_{i}:=\mu\left(B_{i} \cap \cdot\right)(i \geq 1)$ and set $C:=\bigcup_{i \geq 1} C_{i}$. Then $\{C \cap B=\emptyset\}=\bigcup_{i \geq 1}\left\{C_{i} \cap B=\emptyset\right\}$ is measurable for all measurable $B \subset E$, hence $C$ is a measurable $\operatorname{Count}(E)$-valued random variable. Since the $C_{i}$ are disjoint and independent and the $\mu_{i}$ are nonatomic, we have

$$
\begin{equation*}
\mathbb{E}\left[\prod_{x \in C}(1-f(x))\right]=\prod_{i \geq 1} \mathbb{E}\left[\prod_{x \in C_{i}}(1-f(x))\right]=\prod_{i \geq 1} e^{-\int f \mathrm{~d} \mu_{i}}=e^{-\int f \mathrm{~d} \mu} \quad\left(f \in B_{[0,1]}(E)\right) \tag{D.7}
\end{equation*}
$$

In particular, setting $f=1_{B}$ we see that $\mathbb{P}[C \cap B=\emptyset]=e^{-\mu(B)}$ for all measurable $B \subset E$. Set $\mathcal{A}_{B}:=\{A: A \cap B=\emptyset\}$. Then $\mathcal{A}_{B} \cap \mathcal{A}_{B^{\prime}}=\mathcal{A}_{B \cup B^{\prime}}, \mathcal{A}_{\emptyset}=\Omega$, and the class of all $\mathcal{A}_{B}$ with $B \subset E$ measurable generates the $\sigma$-field on $\operatorname{Count}(E)$, hence by the $\pi / \lambda$-theorem Kal02, Theorem 1.1], (D.5) uniquely determines the law of $C$.

## E A one-sided version of Kolmogorov's moment criterion

We prove a variant of Kolmogorov's moment criteria (see e.g. Dur96, Chap. 7, Theorem (1.5)]) for the Hölder continuity of a stochastic process, with bounds on the distribution of the Hölder constant. We assume a one-sided moment condition, which in turn gives one-sided Hölder continuity at deterministic times.

Theorem E. 1 Let $\left(X_{t}\right)_{t \in[0, T]}$ be a real-valued stochastic process. If for all $0 \leq s<t \leq T$,

$$
\begin{equation*}
\mathbb{E}\left[\left(\left(X_{s}-X_{t}\right)^{+}\right)^{\beta}\right] \leq K(t-s)^{1+\alpha} \tag{E.1}
\end{equation*}
$$

for some $\alpha, \beta>0$, then for any $0<\gamma<\frac{\alpha}{\beta}$, there exists a random constant $C \in(0, \infty)$ such that a.s.

$$
\begin{equation*}
\left(X_{r}-X_{q}\right)^{+} \leq C(q-r)^{\gamma} \quad \text { for all } r, q \in \mathbb{Q}_{2} \cap[0, T] \text { with } r<q \tag{E.2}
\end{equation*}
$$

where $\mathbb{Q}_{2}=\left\{m 2^{-n}: m, n \geq 0\right\}$ is the set of dyadic rationals. Furthermore, for any $0<\delta<$ $\alpha-\beta \gamma$, there exists a deterministic constant $C_{\delta, \gamma}$ depending only on $\gamma, \delta, K, \alpha$ and $\beta$, such that

$$
\begin{equation*}
\mathbb{P}[C \geq u] \leq \frac{C_{\delta, \gamma}}{u^{\delta}} \wedge 1 \quad \text { for all } u>0 \tag{E.3}
\end{equation*}
$$

The same results hold if we replace $(\cdot)^{+}$by $(\cdot)^{-}:=-(\cdot \wedge 0)$ or $|\cdot|$.

Proof. The proof is essentially the same as that for the standard version of Kolmogorov's moment criterion. Fix $0<\gamma<\frac{\alpha}{\beta}$. Without loss of generality, assume $T=1$ and let $D_{n}:=\left\{i 2^{-n}: 0 \leq i \leq 2^{n}\right\}$. For any $s:=i 2^{-n}<t:=j 2^{-n} \in D_{n}$, by (E.1) and Chebychev inequality,

$$
\begin{equation*}
\mathbb{P}\left[\left(X_{s}-X_{t}\right)^{+}>(t-s)^{\gamma}\right] \leq K(t-s)^{1+\alpha-\beta \gamma}=K(j-i)^{1+\alpha-\beta \gamma} 2^{-n(1+\alpha-\beta \gamma)} \tag{E.4}
\end{equation*}
$$

If we let $G_{n}:=\left\{\left(X_{i 2^{-n}}-X_{j 2^{-n}}\right)^{+} \leq(j-i)^{\gamma} 2^{-n \gamma}\right.$ for all $\left.0 \leq i \leq j \leq 2^{n}, j-i \leq 2^{n \eta}\right\}$ for some fixed $\eta \in(0,1)$, then by (E.4),

$$
\begin{equation*}
\mathbb{P}\left[G_{n}^{c}\right] \leq \sum_{\substack{0 \leq i<j \leq 2 n \\ j-i \leq 2^{n \eta}}} K(j-i)^{1+\alpha-\beta \gamma} 2^{-n(1+\alpha-\beta \gamma)} \leq K 2^{-[(1-\eta)(\alpha-\beta \gamma)-2 \eta] n}=K 2^{-\delta n}, \tag{E.5}
\end{equation*}
$$

where we have chosen $\eta>0$ such that $(1-\eta)(\alpha-\beta \gamma)-2 \eta=\delta \in(0, \alpha-\beta \gamma)$. Then by Borel-Cantelli, a.s. $N:=\inf \left\{n \in \mathbb{N}: \cap_{i \geq N} G_{i}\right.$ occurs $\}<\infty$, and furthermore, for any $L \in \mathbb{N}$,

$$
\begin{equation*}
\mathbb{P}[N>L] \leq \sum_{n=L}^{\infty} \mathbb{P}\left[G_{n}^{c}\right] \leq K \sum_{n=L}^{\infty} 2^{-\delta n}=\frac{K 2^{-\delta L}}{1-2^{-\delta}} \tag{E.6}
\end{equation*}
$$

Note that

$$
\begin{equation*}
\left(X_{u}-X_{w}\right)^{+} \leq\left(X_{v}-X_{w}\right)^{+}+\left(X_{u}-X_{v}\right)^{+} \quad \text { for any } \quad u<v<w . \tag{E.7}
\end{equation*}
$$

We will use this triangle inequality to deduce (E.2) on the event $\cap_{n \geq N} G_{n}$. First assume that $r<q \in \mathbb{Q}_{2} \cap[0,1]$ and $q-r<2^{-N(1-\eta)}$. We can find an $m \geq N$ such that

$$
\begin{equation*}
2^{-(m+1)(1-\eta)} \leq q-r<2^{-m(1-\eta)} . \tag{E.8}
\end{equation*}
$$

By binary expansion for $q$ and $r$, we can write

$$
\begin{aligned}
& q=j 2^{-m}+2^{-q_{1}}+\cdots 2^{-q_{k}} \\
& r=i 2^{-m}-2^{-r_{1}}-\cdots 2^{-r_{l}},
\end{aligned}
$$

where $m<q_{1}<\cdots<q_{k}$ and $m<r_{1}<\cdots<r_{l}$. By E.8),

$$
2^{(m+1) \eta-1} \leq(q-r) 2^{m} \leq j-i+2
$$

Since $m \geq N$, if we replace $N$ with $N \vee 2 / \eta$, then we are guaranteed that $j \geq i$. Since $q-r \geq(j-i) 2^{-m}$, again by (E.8), we have $j-i \leq 2^{m \eta}$. Since the event $\cap_{n \geq N} G_{n}$ occurs by definition, we have

$$
\begin{equation*}
\left(X_{i 2^{-m}}-X_{j 2^{-m}}\right)^{+} \leq(j-i)^{\gamma} 2^{-m \gamma} \leq 2^{-m(1-\eta) \gamma} . \tag{E.9}
\end{equation*}
$$

By (E.7),

$$
\begin{equation*}
\left(X_{j 2^{-m}}-X_{q}\right)^{+} \leq \sum_{\sigma=1}^{k} 2^{-q_{\sigma} \gamma} \leq \sum_{\sigma>m} 2^{-\sigma \gamma} \leq \frac{2^{-m \gamma}}{2^{\gamma}-1} \tag{E.10}
\end{equation*}
$$

Similarly, the same bound also holds for $\left(X_{r}-X_{i 2^{-m}}\right)^{+}$. Combining the above estimates and applying (E.7) once more, we get

$$
\left(X_{r}-X_{q}\right)^{+} \leq 2^{-m(1-\eta) \gamma}+\frac{2^{1-m \gamma}}{2^{\gamma}-1}=\left(1+\frac{2^{1-m \eta \gamma}}{2^{\gamma}-1}\right) 2^{(1-\eta) \gamma} 2^{-(m+1)(1-\eta) \gamma} \leq C_{\gamma}(q-r)^{\gamma}
$$

for $C_{\gamma}=2^{\gamma}\left(1+\frac{2}{2^{\gamma}-1}\right)$. This verifies E .2 . for $r<q \in Q_{2} \cap[0,1]$ with $q-r<2^{-(N \vee 2 / \eta)(1-\eta)}$. For general $r<q \in \mathbb{Q}_{2} \cap[0,1]$, we can apply the triangle inequality E.7) at most $2^{(N \vee 2 / \eta)(1-\eta)}$ times to obtain (E.2) with $C=C_{\gamma} 2^{(N \vee 2 / \eta)(1-\eta)}$. The distributional tail bound E.3) then follows from (E.6).

When we replace $(\cdot)^{+}$by $(\cdot)^{-}$or $|\cdot|$, the proof is the same since analogues of the triangle inequality (E.7) still hold.

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| Emmanuel Schertzer | Rongfeng Sun | Jan M. Swart |
| :--- | :--- | :--- |
| 109 Montague street, Brooklyn | Department of Mathematics | Institute of Information |
| New York, NY 11201 | National University of | Theory and Automation |
| USA | of the ASCR (ÚTIA) |  |
| email: schertze@hotmail.com | 10 Lowere Kent Ridge Road | Pod vodárenskou věží 4 |
|  | 119076, Singapore | 18208 Praha 8 |
|  | e-mail: matsr@nus.edu.sg | Czech Republic |
|  |  | e-mail: swart@utia.cas.cz |

## Notation List

## General notation:

$\mathbb{Z}_{\text {even }}^{2}$ : The even sublattice of $\mathbb{Z}^{2},\left\{(x, t) \in \mathbb{Z}^{2}: x+t\right.$ is even $\}$.
$\mathcal{L}(\cdot)$ : the law of a random variable.
$S_{\varepsilon}$ : the diffusive scaling map, applied to subsets of $\mathbb{R}^{2}$, paths, and sets of paths, quenched laws, etc. See (3.25).
$\left(\varepsilon_{k}\right)_{k \in \mathbb{N}}$ : a sequence of constants decreasing to 0 , acting as scaling parameters.
$\left(K_{s, t}(x, \cdot)\right)_{s \leq t, x \in E}$ : a stochastic flow of kernels on the space $E$.
$\mathcal{M}_{1}(E)$ : the space of probability measures on the space $E$.
$\mathcal{M}(\mathbb{R}), \mathcal{M}_{\text {loc }}(\mathbb{R})$ : the space of finite and locally finite measures on $\mathbb{R}$.
$\mathcal{M}_{\mathrm{g}}(\mathbb{R})$ : the subset of $\mathcal{M}_{\mathrm{loc}}(\mathbb{R})$ satisfying the growth constraint (2.15).
$\operatorname{supp}(\cdot)$ : support of a measure.

## Paths, Space of Paths:

$R_{\mathrm{c}}^{2}$ : the compactification of $\mathbb{R}^{2}$, see Figure 7 .
$z=(x, t):$ a point in $R_{\mathrm{c}}^{2}$, with position $x$ and time $t$.
$s, t, u, S, T, U$ : times.
$\mathcal{K}\left(R_{\mathrm{c}}^{2}\right)$ : the space of compact subsets of $R_{\mathrm{c}}^{2}$.
$(\Pi, d)$ : the space of continuous paths in $R_{\mathrm{c}}^{2}$, with metric $d$, see (3.7).
$\left(\mathcal{K}(\Pi), d_{\mathrm{H}}\right)$ : the space of compact subsets of $\Pi$, with Hausdorff metric $d_{\mathrm{H}}$, see (3.8).
$\Pi(A), \Pi(z)$ : the set of paths in $\Pi$ starting from a set $A \subset R_{\mathrm{c}}^{2}$ resp. a point $z \in \overline{R_{\mathrm{c}}^{2}}$.
The same notatation applies to any subset of $\Pi$ such as $\mathcal{W}, \mathcal{N}$, etc.
$\pi$ : a path in $\Pi$.
$\sigma_{\pi}$ : the starting time of the path $\pi$
$\pi(t)$ : the position of $\pi$ at time $t \geq \sigma_{\pi}$.
$(\hat{\Pi}, \hat{d})$ : the space of continuous backward paths in $R_{\mathrm{c}}^{2}$ with metric $\hat{d}$.
$\hat{\Pi}(A), \hat{\Pi}(z)$ : the set of backward paths in $\hat{\Pi}$ starting from $A \subset R_{\mathrm{c}}^{2}$ resp. $z \in R_{\mathrm{c}}^{2}$. $\hat{\pi}$ : a path in $\hat{\Pi}$.
$\hat{\sigma}_{\hat{\pi}}$ : the starting time of the backward path $\hat{\pi}$.
$\sim_{\text {in }}^{z}, \sim_{\text {out }}^{z}$ : equivalence of paths entering, resp. leaving $z \in \mathbb{R}^{2}$, see Definition 4.2.
$={ }_{\mathrm{in}}^{z},==_{\text {out }}^{z}$ : strong equivalence of paths entering, resp. leaving $z \in \mathbb{R}^{2}$, see
Definition 3.2 .

## Discrete environments, paths, webs, and flows:

$\omega:=\left(\omega_{z}\right)_{z \in \mathbb{Z}_{\text {even }}^{2}}$ : an i.i.d. environment for random walks on $\mathbb{Z}_{\text {even }}^{2}$.
$\mu$ : the law of $\omega_{o} \in[0,1]$.
$K_{(s, t)}^{\omega}(x, y)$ : transition probability of a random walk from $(x, s)$ to $(y, t) \in \mathbb{Z}_{\text {even }}^{2}$ in the environment $\omega$.
$\mathbf{Q}_{z}^{\omega}$ : the law of a random walk on $\mathbb{Z}_{\text {even }}^{2}$ starting from $z$ in the environment $\omega$.
$\mathbf{Q}^{\omega}$ : discrete Howitt-Warren quenched law.
$p_{z}$ : a discrete path on $\mathbb{Z}_{\text {even }}^{2}$ starting from $z \in \mathbb{Z}_{\text {even }}^{2}$.
$\hat{p}_{z}$ : a dual discrete path on $\mathbb{Z}_{\text {odd }}^{2}:=\mathbb{Z}^{2} \backslash \mathbb{Z}_{\text {even }}^{2}$ starting from $z \in \mathbb{Z}_{\text {odd }}^{2}$.
$\alpha:=\left(\alpha_{z}\right)_{z \in \mathbb{Z}_{\text {even }}^{2}}:$ a family of independent $\pm 1$-valued random variables.
$(\mathcal{U}, \hat{\mathcal{U}})$ : a discrete web and its dual.
$\left(\mathcal{U}^{1}, \mathcal{U}^{\mathrm{r}}, \hat{\mathcal{U}}^{1}, \hat{\mathcal{U}}^{\mathrm{r}}\right)$ : a discrete left-right web and its dual.
$(\mathcal{V}, \hat{\mathcal{V}})$ : a discrete net and its dual.
$\left(\omega^{\langle k\rangle}\right)_{k \in \mathbb{N}}$ : a sequence of i.i.d. environments on $\mathbb{Z}_{\text {even }}^{2}$, with $\mathcal{L}\left(\omega_{o}^{\langle k\rangle}\right)=\mu_{k}$.
$\left(\mu_{k}\right)_{k \in \mathbb{N}}$ : a sequence of probability laws on $[0,1]$, satisfying (1.7).
$\mathbf{Q}_{\langle k\rangle}$ : the discrete Howitt-Warren quenched law associated with $\omega^{\langle k\rangle}$.

## Brownian webs:

$(\mathcal{W}, \hat{\mathcal{W}})$ : a double Brownian web consisting of a Brownian web and its dual.
$\pi_{z}, \hat{\pi}_{z}$ : the a.s. unique path in $\mathcal{W}$ resp. $\hat{\mathcal{W}}$ starting from a deterministic $z \in \mathbb{R}^{2}$.
$\pi_{z}^{-}, \pi_{z}^{+}$: the leftmost, resp. rightmost path starting from $z \in \mathbb{R}^{2}$ in the Brownian web $\mathcal{W}$.
$\pi_{z}^{\uparrow}$ : same as $\pi_{z}^{+}$, except when there is an incoming path in $\mathcal{W}$ at $z$, then $\pi_{z}^{\uparrow}$ is defined to be the continuation of the incoming path.
$\mathcal{W}_{\text {in }}(z), \mathcal{W}_{\text {out }}(z)$ : the set of paths in $\mathcal{W}$ entering, resp. leaving $z$.
$\operatorname{sign}_{\mathcal{W}}(z)$ : the orientation of a $(1,2)$ point $z \in \mathbb{R}^{2}$ in $\mathcal{W}$. See (3.13).
$\operatorname{switch}_{z}(\mathcal{W})$ : a modification of $\mathcal{W}$ by switching the orientation of all paths in $\mathcal{W}$ entering $z$. See (3.11).
$\operatorname{hop}_{z}(\mathcal{W}): \mathcal{W} \cup \operatorname{switch}_{z}(\mathcal{W})$.
$\ell, \ell_{1}, \ell_{\mathrm{r}}$ : the intersection local time measure between $\mathcal{W}$ and $\hat{\mathcal{W}}$, and its restriction to the set of $(1,2)_{1}$, resp. $(1,2)_{\mathrm{r}}$ points. See Proposition 3.4 . $\left(\mathcal{W}^{1}, \mathcal{W}^{\mathrm{r}}, \hat{\mathcal{W}}^{1}, \hat{\mathcal{W}}^{\mathrm{r}}\right)$ : the left-right Brownian web and its dual.
$l, r, \hat{l}, \hat{r}$ : a path in $\mathcal{W}^{1}$, resp. $\mathcal{W}^{\mathrm{r}}, \hat{\mathcal{W}}^{1}, \hat{\mathcal{W}}^{\mathrm{r}}$.
$l_{z}, r_{z}, \hat{l}_{z}, \hat{r}_{z}$ : the a.s. unique path in $\mathcal{W}^{\mathrm{l}}(z)$, resp. $\mathcal{W}^{\mathrm{r}}(z), \hat{\mathcal{W}}^{\mathrm{l}}(z), \hat{\mathcal{W}}^{\mathrm{r}}(z)$, starting from a deterministic $z \in \mathbb{R}^{2}$.
$W(\hat{r}, \hat{l})$ : a wedge defined by the paths $\hat{r} \in \hat{\mathcal{W}}^{\mathrm{r}}$ and $\hat{l} \in \hat{\mathcal{W}}^{\mathrm{l}}$, see 4.4.
$M(r, l)$ : a mesh defined by the paths $r \in \mathcal{W}^{\mathrm{r}}$ and $l \in \mathcal{W}^{\mathrm{l}}$, see (4.5).

## Brownian net:

$(\mathcal{N}, \hat{\mathcal{N}})$ : the Brownian net and the dual Brownian net. See Theorem 4.1
$\xi_{t}$ : a branching-coalescing point set. See (4.6) and Proposition 2.6.
$S$ : the set of separation points of $\mathcal{N}$.
$\operatorname{sign}_{\pi}(z)$ : the orientation of a path $\pi \in \mathcal{N}$ at the separation point $z$. See 4.8.
$\mathcal{H}_{-}$: a Brownian half-net with infinite left speed and finite right speed.
$\mathcal{H}_{+}$: a Brownian half-net with finite left speed and infinite right speed.
$R_{S, U}$ : the set of $S, U$-relevant separation points in $\mathcal{N}$.

## Howitt-Warren flows and processes:

$\beta, \nu:$ the drift and characteristic measure of a Howitt-Warren flow. See Definition 2.2
$\beta_{-}, \beta_{+}$: the left and right speeds of a Howitt-Warren flow. See 2.12).
$\left(\mathcal{W}_{0}, \mathcal{M}, \mathcal{W}\right)$ : the reference web, the set of marked points, and the sample web. See Section 3.4 .
$\left(\mathcal{W}_{i}\right)_{i \in \mathbb{N}}$ : i.i.d. copies of the sample web $\mathcal{W}$.
$\beta_{0}, \beta$ : the drift of the reference, resp. sample web.
$\nu_{1}, \nu_{\mathrm{r}}$ : a decomposition of $\nu$ via (3.21).
$\pi_{z}^{+}$: the rightmost path starting from $z \in \mathbb{R}^{2}$ in the sample web $\mathcal{W}$.
$\pi_{z}^{\uparrow}$ : same as $\pi_{z}^{+}$, except when there is an incoming path in $\mathcal{W}$ at $z$, then $\pi_{z}^{\uparrow}$ is defined to be the continuation of the incoming path.
$\left(K_{s, t}^{+}\right)_{s \leq t}$ : the version of the Howitt-Warren flow constructed using $\pi^{+}$. See (3.22).
$\left(K_{(s, t)}^{\uparrow}\right)_{s \leq t}$ : the version of the Howitt-Warren flow constructed using $\pi^{\uparrow}$.
$\mathbb{Q}$ : the Howitt-Warren quenched law of $\mathcal{W}$ conditional on $\left(\mathcal{W}_{0}, \mathcal{M}\right)$.
$\rho_{t}$ : the Howitt-Warren process defined from either $K^{+}$or $K^{\uparrow}$. See 2.1 .
$\left(\Lambda_{c}\right)_{c \geq 0}$ : ergodic homogeneous invariant laws for the Howitt-Warren process. See Theorem 2.11.
$\zeta_{t}$ : the smoothing process dual to the Howitt-Warren process $\rho_{t}$. See 11.1.


[^0]:    ${ }^{1}$ We follow HW09a in our definition of $\nu$. Many of our formulas, however, such as 2.3, 2.11 or 3.16 are more easily expressed in terms of $2 \nu$ than in $\nu$. Loosely speaking, the reason for this is that in 1.7) (ii), the weight function $q(1-q)$ arises from the fact that if $\alpha^{1}, \alpha^{2}$ are independent $\{-1,+1\}$-valued random variables with $\mathbb{P}\left[\alpha^{i}=+1\right]=q(i=1,2)$, then $\mathbb{P}\left[\alpha^{1} \neq \alpha^{2}\right]=2 q(1-q)$.
    ${ }^{2}$ Actually, the paper HW09a considers a continuous-time analogue of the discrete $n$-point motions defined in Section 1.2, but their proof, with minor modifications, also works in the discrete time setting. In Appendix A we present a similar, but somewhat simplified convergence proof.

[^1]:    ${ }^{3}$ For simplicity, we have omitted two regularity conditions on $\left(K_{s, t}\right)_{s \leq t}$ from the original definition in LR04a, Def. 2.3], which are some form of weak continuity of $K_{s, t}(x, \cdot)$ in $x, s$ and $t$. It is shown in that paper that a stochastic flow of kernels on a compact metric space $E$ satisfies these regularity conditions if and only if it arises from a consistent family of Feller processes.

[^2]:    ${ }^{4}$ In fact, LR04a, Thm 2.1] is stated only for compact metrizable spaces, but the extension to locally compact $E$ is straightforward using the one-point compactification of $E$.

[^3]:    ${ }^{5}$ Our definition of $\ell$ in (3.9) differs from the definition in NRS10, (3.2)] by a factor of $\sqrt{2}$, which is compensated by the fact that we consider Poisson point process in $\mathbb{R}^{2}$ with intensity measure $\ell$ instead of $\sqrt{2} \ell$, as done in NRS10, (3.8)].

