THE NOISY VETO-VOTER MODEL: A RECURSIVE DISTRIBUTIONAL EQUATION ON [0, 1]

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Abstract

We study a particular example of a recursive distributional equation (RDE) on the unit interval. We identify all invariant distributions, the corresponding ‘basins of attraction’, and address the issue of endogeny for the associated tree-indexed problem, making use of an extension of a recent result of Warren.

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1. Introduction

Let \( M \) be a random variable taking values in \( \mathbb{N} = \{1, \ldots; \infty \} \), and let \( \xi \) be an independent Bernoulli(\( p \)) random variable.

We consider the following simple recursive distributional equation (RDE):

\[
Y = \xi \prod_{i=1}^{M} Y_i + (1 - \xi) \left( 1 - \prod_{i=1}^{M} Y_i \right).
\] (1.1)

Viewing (1.1) as an RDE, we seek a stationary distribution, \( \nu \), such that, if the \( Y_i \) are independent and identically distributed (i.i.d.) with distribution \( \nu \) and are independent of \( (M, \xi) \), then \( Y \) also has distribution \( \nu \).

We term (1.1) the noisy veto-voter model since if each \( Y_i \) takes values in \( \{0, 1\} \) with value 0 being regarded as a veto then the outcome is vetoed unless either (a) each voter \( i \) ‘assents’ \((Y_i = 1\) for each \( 1 \leq i \leq M \)) and there is no noise \((\xi = 1)\) or (b) someone vetoes, but is reversed by the noise \((\xi = 0)\).

The system was originally envisaged as modelling a representative voting system applied to a veto issue. Thus, each representative votes according to their constituency if \( \xi = 1 \) or reverses the decision if \( \xi = 0 \). An alternative interpretation is as a model for a noisy distributed error-reporting system. Here a 0 represents an error report from a subsystem. Thus, there is an error in the system if there is an error in any subsystem (hence the veto structure). Noise can reverse the binary (on–off) report from any subsystem. The system is substantially different from the standard voter model introduced by Holley and Ligget [6] (and independently by Clifford and Sudbury [3]), where voters change their mind, influenced by the current voting positions of neighbours. Cox and Griffeath [4] gave a survey of results about this model and
established key results about clustering. For recent developments on heterogeneous graphs, see [8] and the references therein.

In this paper we look for solutions to the RDE (1.1) taking values in $[0, 1]$. As observed by Aldous and Bandhapadhyay [1], and as we shall explain in a little more detail in Section 2, we may think of (families of) solutions to the RDE as being located at the nodes of a (family) tree (for a Galton–Watson branching process). Actually, for some purposes, we shall find it more convenient to embed this family tree into $T$, the deterministic tree with infinite branching factor of size $\aleph_0$.

The generic setup in such circumstances is to find distributional fixed points of the recursion

$$X_u = f(\xi_u; X_{ui}, i \geq 1),$$

where $X_u$ and $\xi_u$ respectively denote the value and the noise associated with node $u$, and $ui$ denotes the address of the $i$th daughter of node $u$.

With this model, it is of some interest not only to find solutions to the RDE (1.2) but also to answer the question of endogeny:

Is $(X_u; u \in T)$ measurable with respect to $(\xi_u; u \in T)$?

If this measurability condition holds then $X$ is said to be endogenous.

In the context of the error-reporting model, endogeny represents the worst possible situation —the top-level error report is based entirely on the noise and is uninfluenced by the error state of low-level subsystems. Similarly, in the veto-voter paradigm, endogeny represents the situation where the voice of the ‘little man’ is completely swamped by reversals by officials.

In this paper we shall first show how to transform (1.1) into the new RDE:

$$X = 1 - \prod_{i=1}^{N} X_i$$

for a suitable random variable $N$, independent of the $X_i$. Then we shall not only find all the solutions to this RDE on $[0, 1]$, their basins of attractions, and the limit cycles of the corresponding map on the space of distributions on $[0, 1]$, but also give necessary and sufficient conditions for the corresponding solutions on $T$ to be endogenous.

The fundamental technique we use, which we believe is entirely novel, is to consider the distribution of a solution conditional upon the noise and to identify endogeny by showing that this conditional distribution is concentrated on $\{0, 1\}$.

2. Notation and a transformation of the RDE

2.1. Tree-indexed solutions

We seek distributions $\nu$ on $[0, 1]$ such that if the $(Y_i; 1 \leq i)$ are independent with distribution $\nu$ then the random variable $Y$ satisfying (1.1) also has distribution $\nu$. More precisely, writing $\mathcal{P}$ for the set of probability measures on $[0, 1]$, suppose that $M$ has distribution $d$ on $\mathbb{Z}_+$ and define the map

$$\mathcal{T} \equiv \mathcal{T}_d: \mathcal{P} \to \mathcal{P}.$$  

Then we set $\mathcal{T} (\nu)$ to be the law of the random variable $Y$ given by (1.1), when the $Y_i$ are i.i.d. with distribution $\nu$ and are independent of $N$, and seek fixed points of the map $\mathcal{T}$. The existence and uniqueness of fixed points of this type of map, together with properties of the solutions, are
addressed in [1] (the reader is also referred to [2], [7], and the references therein). The linear
and min cases are particularly well surveyed, though we are dealing with a nonlinear case to
which the main results do not apply.

A convenient generalisation of the problem is the so-called tree-indexed problem, in which
we think of the $Y_i$ as being marks associated with the daughter nodes of the root of $T$, a family
tree of a Galton–Watson branching process. We start at some level $m$ of the random tree. Each
vertex $v$ in level $m - 1$ of the tree has $M_v$ daughter vertices, where the $M_v$ are i.i.d. with
common distribution $d$, and has associated with it noise $\xi_v$, where the $(\xi_u; u \in T)$ are i.i.d. and
are independent of the $(M_u; u \in T)$.

By associating with daughter vertices independent random variables $Y_{vi}$ having distribution
$\nu$, we see that $Y_v$ and $Y_{vi}, 1 \leq i \leq M_v$, satisfy (1.1).

In this setting the notion of endogeny was introduced in [1]. Loosely speaking, a solution
to the tree-indexed problem (which we shall define precisely in the next section) is said to be
endogenous if it is a function of the initial data or noise alone, so that no additional randomness
is present.

It is convenient to work on a tree with infinite branching factor and then think of the random
tree of the previous paragraph as being embedded within it. An initial ancestor (in level zero),
which we denote $\emptyset$, gives rise to a countably infinite number of daughter vertices (which form
the members of the first generation), each of which gives rise to an infinite number of daughters
(which form the members of the second generation), and so on. We assign each vertex an
address according to its position in the tree: the members of the first generation are denoted
by $1, 2, \ldots$, the second generation by $11, 12, \ldots, 21, 22, \ldots$, etc, so that vertices
in level $n$ of the tree correspond to sequences of positive integers of length $n$. We also write
$uj, j = 1, 2, \ldots, \varepsilon$ for the daughters of a vertex $u$. We write $T$ for the collection of all vertices
or nodes (i.e. $T = \bigcup_{n=0}^{\infty} N^n$) and think of it as being partitioned by depth, that is, as being
composed of levels or generations, in the way described, and we define the depth function $|\cdot|
by $|u| = n$ if vertex $u$ is in level $n$ of the tree. Associated to each of the vertices $u \in T$ are
i.i.d. random variables $M_u$ with distribution $d$, telling us the (random) number of offspring
produced by $u$. The vertices $u_1, u_2, \ldots, u M_u$ are thought of as being alive (relative to $\emptyset$) and
the $\{uj; j > M_u\}$ are thought of as being dead. We can now write our original equation as a
recursion on the vertices of $T$:

$$Y_u = \xi_u \prod_{i=1}^{M_u} Y_{ui} + (1 - \xi_u) \left(1 - \prod_{i=1}^{M_u} Y_{ui}\right), \quad u \in T. \quad (2.1)$$

The advantage of the embedding now becomes clear: we can talk about the RDE at any vertex
in the infinite tree and yet, because the product only runs over the live daughters relative to $u$, the
random Galton–Watson family tree is encoded into the RDE as noise.

2.2. The transformed problem

It is a relatively simple matter to transform the RDE (2.1) into the following, simpler, RDE:

$$X_u = 1 - \prod_{i=1}^{N_u} X_{ui}, \quad u \in T, \quad (2.2)$$

where $N_u$ is a new family size to be defined shortly. To do so, first note that if we colour red
all the nodes, $u$, in the tree $T$ for which $\xi_u = 0$ then it is clear that we may proceed down each
line of descent from a node $u$ until we hit a red node. In this way, we either ‘cut’ the tree at a
collection of nodes which we shall view as the revised family of $u$, or not, in which case $u$ has an infinite family. Denote this new random family size by $N_u$. Then

$$Y_u = 1 - \prod_{i=1}^{N_u} \hat{Y}_{ui}$$

if $u$ is red, where $\hat{u}$ denotes the $i$th red node in the revised family of $u$. Now condition on node $u$ being red. Then with this revised tree we obtain the RDE (2.2). It is easy to see that if the original tree has family size probability generating function (PGF) $G$ then the family size in the new tree corresponds to the total number of deaths in the original tree when it is independently thinned, with the descendants of each node being pruned with probability $q$. It is easy to obtain the equation for the PGF, $H$, of the family size $N_u$ on the new tree:

$$H(z) = G(pH(z) + qz).$$ (2.3)

3. The discrete and conditional probability solutions

We begin with some notation and terminology. We say that the random variables in (2.1) are weakly stationary if $X_u$ has the same distribution for every $u \in \mathbb{T}$. The stationarity of the $X_u$ corresponds to $X_u$ having as distribution an invariant measure for the distributional equation (2.2).

**Definition 3.1.** We say that the process (or collection of random variables) $X = (X_u; u \in \mathbb{T})$ is a tree-indexed solution to the RDE (2.2) if

(i) for every $n$, the random variables $(X_u; |u| = n)$ are mutually independent and independent of $(N_v; |v| \leq n - 1)$;

(ii) for every $u \in \mathbb{T}$, $X_u$ satisfies

$$X_u = 1 - \prod_{i=1}^{N_u} X_{ui},$$

and the $(X_u; u \in \mathbb{T})$ are weakly stationary.

Note that these conditions determine the law of $X$. This means that a tree-indexed solution is also stationary in the strong sense, that is, a tree-indexed solution is ‘translation invariant’ with respect to the root (if we consider the collection $X^v = (X_u; u \in \mathbb{T}_v)$, where $\mathbb{T}_v$ is the subtree rooted at $v$, then $X^v$ has the same distribution as $X$ for any $v \in \mathbb{T}$). Furthermore, we say that such a solution is endogenous if it is measurable with respect to the random tree (i.e. the collection of family sizes) $(N_u; u \in \mathbb{T})$. As we remarked in the introduction, in informal terms this means that the solution depends only on the noise with no additional randomness coming from the boundary of the tree. See [1] for a thorough discussion of endogeny together with examples.

The following lemma is easy to prove.

**Lemma 3.1.** Let $(X_u; u \in \mathbb{T})$ be a tree-indexed solution to the RDE (2.2). Then the following statements are equivalent:

(i) $X$ is endogenous;

(ii) $X_\emptyset$ is measurable with respect to $\sigma (N_u; u \in \mathbb{T})$;
Xu is measurable with respect to $\sigma(N_v; v \in T_u)$ for each $u \in T$;

(iv) $X_u$ is measurable with respect to $\sigma(N_v; v \in T_u)$ for each $u \in T$.

**Remark 3.1.** Note that if a tree-indexed solution to \(2.2\) is endogenous then property (i) of Definition 3.1 is automatic: for every $u \in T$, $X_u$ is measurable with respect to $\sigma(N_v; v \in T_u)$ and, hence, is independent of $(N_v; |v| \leq n - 1)$.

**Lemma 3.2.** There exists a unique probability measure on $\{0, 1\}$ which is invariant under \(1.3\).

**Proof.** Let $X$ be a random variable whose distribution is concentrated on $\{0, 1\}$ and which is invariant under \(1.3\). Let $\mu = P(X = 1)$. We then have

$P(X_i = 1; \text{ for } i = 1, \ldots, N) = \sum_{n} P(X_i = 1; \text{ for } i = 1, \ldots, n | N = n) P(N = n)$

$= H(\mu).$

Now, $X = 0$ if and only if $X_i = 1$ for $i = 1, \ldots, N$. Hence, a necessary and sufficient condition for invariance is

$1 - \mu = H(\mu). \quad (3.1)$

Now let $K(x) := H(x) + x - 1$. Since $H$ is a generating function and $H(0) = 0$, we have $K(0) = -1 < 0$ and $K(1) > 0$, so that $K$ is guaranteed to have a 0 in $(0, 1)$, and it is unique since the mapping $x \mapsto H(x) + x$ is strictly increasing.

We can now deduce that there exists a tree-indexed solution on $\{0, 1\}$ to the RDE \(2.2\) by virtue of Lemma 6 of [1].

**Theorem 3.1.** Let $S = (S_u; u \in T)$ be a tree-indexed solution on $\{0, 1\}$ to the RDE \(2.2\) (i.e. the $S_u$ have the invariant distribution on the two point set $\{0, 1\}$), which we shall henceforth refer to as the discrete solution. Let $C_u = P(S_u = 1 | N_v; v \in T)$. Then $C = (C_u; u \in T)$ is the unique endogenous tree-indexed solution to the RDE.

**Proof.** We write $N = (N_u; u \in T)$ and $N_u = (N_v; v \in T_u)$. To verify the relationship between the random variables, we have

$C_u = P(S_u = 1 | N)$

$= E[I_{S_u=1} | N]$  

$= E[S_u | N]$

$= E\left[\left. 1 - \prod_{i=1}^{N_u} S_{ui}\right| N\right]$  

$= 1 - E\left[\left. \prod_{i=1}^{N_u} S_{ui}\right| N\right]$  

$= 1 - \prod_{i=1}^{N_u} E[S_{ui} | N]$  

$= 1 - \prod_{i=1}^{N_u} C_{ui}.$
since the $S_{u}$ are independent and $N$ is strongly stationary. To verify stationarity, let

$$C_{u}^{n} = P(S_{u} = 1 \mid N_{v}; |v| \leq n).$$

Then the sequence $(C_{u}^{n})_{n \geq 1}$ is a uniformly bounded martingale and so converges almost surely (a.s.) and in $L^2$ to a limit which must in fact be $C_{u}$. Now, we can write $C_{u}^{n}$ as

$$C_{u}^{n} = 1 - \prod_{i=1}^{N_{u}} C_{u i}^{n} \prod_{i=1}^{N_{u 1}} \left(1 - \prod_{i=2}^{N_{u i 2} - |u|} \left(1 - (\mu^{1})^{N_{u i 2} - |u| - i}ight) \ldots \right) \rightarrow C_{u} \text{ a.s. (3.2)}$$

This corresponds to starting the distributional recursion at level $n$ of the tree with unit masses at $\mu^{1}$. Now, $(C_{u}^{n}; u \in T)$ is stationary since each $C_{u}^{n}$ is the same function of $N_{u}$, which are themselves stationary. Since $C_{u}$ is the (almost sure) limit of a sequence of stationary random variables, it follows that $C = (C_{u}; u \in T)$ is stationary. Note that the conditional probability solution, $C$, is automatically endogenous since $C_{u}$ is $\sigma(N_{v}; v \in T_{u})$-measurable for every $u \in T$ and, hence, $(C_{u}; |u| = n)$ is independent of $(N_{u}; |u| \leq n - 1)$. The independence of the collection $(C_{u}; |u| = n)$ follows from the fact that the $((S_{u}, N_{u}); |u| = n)$ are independent.

Finally, note that if $(L_{u}; u \in T)$ solve the RDE (2.2) and are integrable then $m := E[L_{u}]$ must satisfy (3.1) and, hence, must equal $\mu^{1}$. It now follows that $L_{u}^{n} := E[L_{u} \mid N_{v}; |v| \leq n] = C_{u}^{n}$, since at depth $n$, $L_{u}^{n} = \mu^{1}$, so that $L_{u}^{n}$ also satisfies (3.2) and, hence, must equal $C_{u}^{n}$. Now $L_{u}^{n} \rightarrow L_{u}$ a.s. and so if $L$ is endogenous then it must equal $C$. This establishes that $C$ is the unique endogenous solution.

Remark 3.2. Note that if $S$ is endogenous then $C = S$ a.s., so that if $S$ and $C$ do not coincide then $S$ cannot be endogenous.

4. The moment equation and uniqueness of solutions

Many of the results proved in this paper rely heavily on the analysis of (4.1), below.

Theorem 4.1. Any invariant distribution for the RDE (2.2) must have moments $(m_{n})_{n \geq 0}$ satisfying the equation

$$H(m_{n}) = (-1)^n m_{n} = \sum_{k=0}^{n-1} \binom{n}{k} (-1)^k m_{k},$$

where $m_{n+1/n} \leq m_{n+1} \leq m_{n}$ and $m_{0} = 1$.

Proof. Let $X$ be a random variable whose distribution is invariant for the RDE, and write $m_{k} = E[X^k]$. Applying the RDE (2.2) to $(1 - X)^n$ we have

$$E[(1 - X)^n] = E\left[\prod_{i=1}^{N} X_i^n\right] = H(m_{n}).$$
On the other hand, by expanding \((1 - X)^n\) we obtain

\[
\mathbb{E}[(1 - X)^n] = \mathbb{E}\left[ \sum_{k=0}^{n} \binom{n}{k} (-1)^k X_k^k \right] = \sum_{k=0}^{n} \binom{n}{k} (-1)^k m_k,
\]

so that

\[
H(m_n) = \sum_{k=0}^{n} \binom{n}{k} (-1)^k m_k.
\]

The condition \(m_{n+1} \leq m_n\) follows from the fact that the distribution is on \([0, 1]\). The other condition follows from the monotonicity of \(L^p\) norms.

As an example, if the random variable \(N\) has generating function \(H(x) = x^2\) (i.e. \(N \equiv 2\)), the moment equation tells us that

\[
m_1^2 + m_1 - 1 = 0,
\]

so that \(m_1 = (\sqrt{5} - 1)/2\). For \(m_2\), we have

\[
m_2^2 - m_2 - (2 - \sqrt{5}) = 0,
\]

so that \(m_2 = m_1\) or \(m_2^2\) and so on. In fact, the two possible moment sequences turn out to be \(m_0 = 1, m_n = (\sqrt{5} - 1)/2\) for \(n \geq 1\) or \(m_0 = 1, m_1 = (\sqrt{5} - 1)/2, m_n = m_1^2\) for \(n \geq 2\).

We suppose from now on that \(H(0) = 0\) and \(H\) is strictly convex (so that \(P(2 \leq N < \infty) > 0\)).

We now state the main result of the paper.

**Theorem 4.2.** Let \(S = \{S_u; u \in T\}\) and \(C = \{C_u; u \in T\}\) respectively denote the discrete solution and corresponding conditional probability solution to the RDE (2.2). Let \(\mu^1 = \mathbb{E}[S]\). Then

(i) \(S\) is endogenous if and only if \(H'(\mu^1) \leq 1\);

(ii) \(C\) is the unique endogenous solution;

(iii) the only invariant distributions for the RDE (2.2) are those of \(S_\emptyset\) and \(C_\emptyset\).

The proof of the theorem relies on several lemmas. For (i), we extend a result in [9] by first truncating \(N\) and then taking limits.

First, however, we give some consequences of the moment equation, (4.1).

**Lemma 4.1.** There are at most two moment sequences satisfying (4.1). Moreover, the first moment \(m^1\) is unique and equal to \(\mu^1, 1 > m^1 > 1/2\), and in the case that \(H'(m^1) \leq 1\) there is only one moment sequence satisfying (4.1).

**Proof.** Uniqueness of \(\mu^1\) (the root of \(f(m^1) = 1\), where \(f : t \mapsto H(t) + t\) has already been shown in Lemma 3.2. Now set

\[
g(x) = H(x) - x.
\]

Then \(g\) is strictly convex on \([0,1]\) with \(g(0) = 0\) and \(g(1-) = H(1-) - 1 \leq 0\). Thus, there are at most two solutions of \(g(x) = 1 - 2m^1\). Since \(m^1\) itself is a solution, it follows that \(1 - 2m^1 \leq 0\) and there is at most one other solution. There is another solution with \(m^2 < m^1\)
if and only if \( m_1 \) is greater than \( \mu^* \), the argmin of \( g \), and this is clearly true if and only if 
\[ g'(m_1) > 0, \]
which is equivalent to \( H'(m_1) > 1 \).

Suppose that this last inequality holds, so that there is a solution, \( m^2 \), of \( g(x) = 1 - 2m^1 \) 
with \( m^2 < \mu^* < m_1 \). There is at most one solution of 
\[ f(x) = 1 - 3m^1 + 3m^2, \]
and if it exists, take this as \( m^3 \). Similarly, there is at most one solution of \( g(x) = 1 - 4m^1 + 6m^2 - 4m^3 \) to the left of \( \mu^* \) and this is the only possibility for \( m^4 \). Iterating the argument, we obtain at most one strictly decreasing sequence \( m^1, \ldots \).

4.1. The case of a bounded branching factor

Recall that the random family size \( N \) may take the value \( \infty \).

Lemma 4.2. Define \( N' = \min(n, N) \), and denote its generating function by \( H_n \). Then \( N' \) is bounded and 

(i) \( H_n(s) \geq H(s) \) for all \( s \in [0, 1] \);

(ii) \( H_n \to H \) uniformly on compact subsets of \([0, 1] \);

(iii) \( H'_n \to H' \) uniformly on compact subsets of \([0, 1] \).

We leave the proof to the reader.

The following lemma will be used in the proof of Theorem 4.3, below.

Lemma 4.3. Let \( C_u^{(n)} \) denote the conditional probability solution for the RDE (2.2) with \( N \) replaced by \( N' \). Let \( \mu^k_n \) denote the \( k \)th moment, and let \( \mu^k \) denote the \( k \)th moment of \( g \). Let \( \mu^*_{n,m} \) denote that root of the equation 
\[ g_n(x) = 1 - \mu^1_n - \mu^1_m, \quad (4.2) \]
which lies to the left of \( \mu^* \) (i.e. the lesser of the two possible roots). Then \( \mu^k_n \to \mu^k \) for \( k = 1, 2 \) 
and \( \mu^*_{n,m} \to \mu^2 \) as \( \min(n, m) \to \infty \).

Proof. For the case in which \( k = 1 \), consider the graphs of the functions \( H_n(x) + x \) and \( H(x) + x \). We have \( H_n(x) \geq H(x) \) for all \( x \geq 0 \) and all \( n \geq 1 \), so that \( \mu^1_n \) is bounded above 
by \( \mu^1 \) for every \( n \), since \( \mu^1_n \) and \( \mu^1 \) are respectively the roots of 
\[ H_n(x) + x = 1 \quad \text{and} \quad H(x) + x = 1. \]
Furthermore, since \( H_n \) decreases to \( H \) pointwise on \([0, 1] \), it follows that the \( \mu^1_n \) are increasing. 
The \( \mu^1_n \) must therefore have a limit, which we shall denote \( \hat{\mu} \).

It follows from Lemma 4.2 that, since \( \mu^1 < 1 \), \( H_n(\mu^1_n) \to H(\hat{\mu}) \). Hence, 
\[ 1 = H_n(\mu^1_n) + \mu^1_n \to H(\hat{\mu}) + \hat{\mu}, \]
so that \( \hat{\mu} \) is a root of \( H(x) + x = 1 \). It follows, by uniqueness, that \( \hat{\mu} = \mu^1 \).

For the case in which \( k = 2 \), we consider the graphs of \( g_n(x) \) and \( g(x) \). We first show that \( \mu^2_n \to \mu^2 \) and then that \( \mu^*_{n,m} \to \mu^2 \) as \( \min(n, m) \to \infty \).
To show that \( \mu^2 \to \mu^2 \), we argue that \( \mu^2 \) is the only limit point of the sequence \( (\mu^2_n)_{n \geq 1} \).

Note that, since \( \mu^1_n \to \mu^1 \) and \( \mu^2_n \) satisfies

\[
H_n(\mu^2_n) - \mu^2_n = 1 - 2\mu^1_n,
\]

the only possible limit points of the sequence \( (\mu^2_n)_{n \geq 1} \) are \( \mu^1 \) and \( \mu^2 \). Now, either \( \mu^1 \leq \mu^* \), in which case \( \mu^1 = \mu^2 \), or \( \mu^2 \leq \mu^* < \mu^1 < 1 \). In the latter case, it is easy to show that \( \mu^*_n \to \mu^* \) (by uniform continuity of \( g_n' \)), and so, since \( \mu^1_n \to \mu^1 \), it follows that

\[
\mu^1_n > \mu^*_n
\]

for sufficiently large \( n \), and, hence,

\[
\mu^2_n \leq \mu^*_n
\]

for sufficiently large \( n \). In either case, the only possible limit point is \( \mu^2 \); since the \( \mu^2_n \) are bounded they must, therefore, converge to \( \mu^2 \).

We conclude the proof by showing that \( \mu^2 \) is the only limit point of the sequence \( (\mu^2_{n,m})_{n,m \geq 1} \).

Once more, consider the following two cases:

\[
\mu^1 \leq \mu^* \quad \text{and} \quad \mu^1 > \mu^*.
\]

In the first case, \( \mu^1_n = \mu^2_n \) for sufficiently large \( n \), so that \( \mu^2 \) is the only limit point; in the second case,

\[
\mu^1 = \liminf_n \mu^1_n > \mu^* = \limsup_n \mu^*_n,
\]

and since \( \mu^2_{n,m} \leq \mu^*_n \), \( \mu^1 \) cannot be a limit point. Thus, in either case, \( \mu^2 \) is the unique limit point and, hence, is the limit.

**Remark 4.1.** Note that the method of the proof can be extended to prove that \( \mu^k_n \to \mu^k \) for any \( k \).

**Theorem 4.3.** \( C_{(n)}^u \) converges to \( C_u \) in \( L^2 \).

**Proof.** Let \( n \geq m \). Define \( E_{m,n} = \mathbb{E}[(C_u^{(m)} - C_u^{(n)})^2] \). Expanding this, we obtain

\[
E_{m,n} = \mu^2_m + \mu^2_n - 2r_{m,n},
\]

where \( r_{m,n} = \mathbb{E}[C_u^{(m)} C_u^{(n)}] \). On the other hand, by applying the RDE (2.2) once, we obtain

\[
E_{m,n} = \mathbb{E} \left[ \left( \prod_{i=1}^{N_m^u} C_{u_i}^{(n)} - \prod_{i=1}^{N_m^u} C_{u_i}^{(m)} \right)^2 \right] = H_m(\mu^2_m) + H_n(\mu^2_n) - 2 \mathbb{E} \left[ \prod_{i=1}^{N_m^u} C_{u_i}^{(m)} \prod_{i=1}^{N_n^u} C_{u_i}^{(n)} \right].
\]

We can bound \( E_{m,n} \) above and below as follows: since each \( C_{u_i}^{(k)} \) is in \([0,1]\), omitting terms from the product above increases it, while adding terms decreases it. Thus, since \( n \geq m \),...
\( N_u^n \geq N_u^m \), and so replacing \( N_u^n \) by \( N_u^m \) in the product above increases it, while replacing \( N_u^m \) by \( N_u^n \) decreases it. Thus, we obtain
\[
H_m(\mu_m^2) + H_n(\mu_n^2) - 2H_m(r_{m,n}) \leq E_{m,n} \leq H_m(\mu_m^2) + H_n(\mu_n^2) - 2H_n(r_{m,n}).
\]
Using the upper bound, we have
\[
2H_n(r_{m,n}) \leq H_m(\mu_m^2) + H_n(\mu_n^2) - E_{m,n} = H_m(\mu_m^2) + H_n(\mu_n^2) - \mu_m^2 - \mu_n^2 + 2r_{m,n}.
\]
The moment equation, (4.1), tells us that \( H_m(\mu_m^2) - \mu_m^2 = 1 - 2\mu_m^1 \) and that \( H_n(\mu_n^2) - \mu_n^2 = 1 - 2\mu_n^1 \). Hence, \( 2H_n(r_{m,n}) \leq 1 - 2\mu_m^2 + \mu_m^2 + 1 - 2\mu_n^2 + \mu_n^2 - \mu_m^2 - \mu_n^2 + 2r_{m,n} \), so that, on simplifying,
\[
H_n(r_{m,n}) - r_{m,n} \leq 1 - \mu_m^1 - \mu_n^1.
\]
Recall that the equation \( H_n(x) - x = 1 - \mu_n^1 - \mu_n^1 \) has (at most) two roots, the lesser of which we denoted \( \mu_m^2 \). Let \( \mu_{m,n} \) be the other (larger) root (or 1, if the second root does not exist). Then, since \( H_n(x) - x \) is convex, \( \mu_{m,n}^2 \leq r_{m,n} \mu_{m,n}^1 \) for all \( m,n \) and, hence,
\[
\liminf_{m \to \infty} r_{m,n} \mu_{m,n}^1 \geq \mu^2 \text{ by Lemma 4.3}.
\]
On the other hand, Holder’s inequality tells us that \( r_{m,n} \leq \sqrt{\mu_{m,n}^1 \mu_{m,n}^2} \), and so it follows that \( \limsup_{m \to \infty} r_{m,n} \leq \mu^2 \) since \( \mu_{m,n}^2 \to \mu^2 \) by Lemma 4.3. Hence, \( r_{m,n} \to \mu^2 \) as \( n \to \infty \) and
\[
E_{m,n} \to \lim_{m,n \to \infty} \mu_m^2 + \mu_n^2 - 2r_{m,n} = \mu^2 + \mu^2 - 2\mu^2 = 0,
\]
showing that \( (C_{(c)}(u))_u \) is Cauchy in \( L^2 \). It now follows, by the completeness of \( L^2 \), that \( C_{(c)}(C) \) converges. Since \( C_{(c)}(C) \) is \( \sigma(N) \)-measurable, the limit \( L_C \) of the \( C_{(c)}(u) \) must also be \( \sigma(N) \)-measurable for each \( u \) and the collection \( \{I_{(C)_u} : u \geq 1\} \) must be i.i.d. on \( [0,1] \) with common mean \( \mu^1 \leq 1 \). Moreover, by strong stationarity of the \( C_{(c)}(C) \)’s, the \( L_C \)’s are strongly stationary.

To verify that \( L_C \) is the conditional probability solution, note that
\[
1_{E_u} C_{(c)}(C) = \left(1 - \prod_{i=1}^{N_u} C_{(c)}(i)\right)1_{E_u} = \left(1 - \prod_{i=1}^{N_u} C_{(c)}(i)\right)1_{E_u},
\]
where \( E_u = \{N_u \leq n\} \). As \( n \to \infty \), \( E_u \uparrow E := (N < \infty) \); furthermore, since the \( C_{(c)}(C) \) converge in \( L^2 \), they do so in probability. We may therefore assume, without loss of generality, that \( C_{(c)}(C) \) converges a.s. for each \( i \) so that, in the limit,
\[
1_E L_C = \lim_{n \to \infty} 1_{E_u} C_{(c)}(C_u) = \lim_{n \to \infty} 1_{E_u} \left(1 - \prod_{i=1}^{N_u} C_{(c)}(i)\right) = 1_{E} \left(1 - \prod_{i=1}^{N_u} L_i\right) \text{ a.s. (4.3)}
\]

It is easy to show that
\[
\prod_{i=1}^{\infty} L_i = 0 \text{ a.s.,}
\]
while
\[
1_{E^c} C_{(c)}(C) = 1_{E^c} \left(1 - \prod_{i=1}^{n} C_{(c)}(i)\right) \to 1_{E^c} \text{ a.s.,}
\]
so that
\[ \mathbf{1}_{E^n} L_\emptyset = \lim \mathbf{1}_{E^n} C_\emptyset^{(n)} = \mathbf{1}_{E^n}. \quad (4.4) \]
Thus, adding (4.3) and (4.4) we see that
\[ L_\emptyset = \left(1 - \prod_{i=1}^{N_\emptyset} L_i\right), \]
and so \( L \) is an endogenous solution to the RDE. It follows from uniqueness that \( L \) must be the conditional probability solution \( C \).

4.2. Proof of Theorem 4.2

We are now nearly in a position to finish proving Theorem 4.2. To recap, we have shown in Lemma 4.1 that there are at most two distributions which solve the RDE (1.3), corresponding to the 'moment sequences' \( \mu^1, \mu^1, \ldots \) and \( \mu^1, \mu^2, \ldots \). The first of these is the moment sequence corresponding to the distribution on \( \{0, 1\} \) with mass \( \mu_1 \) at 1. The second may or may not be a true moment sequence and is equal to the first if and only if \( H'(\mu^1) \leq 1 \). Moreover, there is only one endogenous solution, and this corresponds to the conditional probability solution \( C \); thus, if we can show that \( C \) is not discrete (i.e. it is not equal to \( S \)) whenever \( H'(\mu_1) > 1 \) then we shall have proved the result.

We need to recall some theory from [9]. Consider the recursion
\[ \xi_u = \phi(\xi_u 0, \xi_u 1, \ldots, \xi_u (d-1), \varepsilon_u), \quad u \in \Gamma_d, \]
where the \( \xi_u \) take values in a finite space \( \delta \), the ‘noise’ terms \( \varepsilon_u \) take values in a space \( E, \Gamma_d \) is the deterministic \( d \)-ary tree, and \( \phi \) is symmetric in its first \( d-1 \) arguments. We suppose that the \( \varepsilon_u \) are independent with common law \( \nu \) and that there exists a measure \( \pi \) which is invariant for the above recursion (i.e. \( \pi \) is a solution of the associated RDE). Let \( u_0 = \emptyset, u_1, u_2, \ldots \) be an infinite sequence of vertices starting at the root, with \( u_{n+1} \) being a daughter of \( u_n \) for every \( n \). For \( n \leq 0 \), define \( \xi_n = \xi_{u_n} \). Then, under the invariant measure \( \pi \), the law of the sequence \( (\xi_n; n \leq 0) \), which, by the symmetry of \( \phi \), does not depend on the choice of sequence of vertices chosen, is that of a stationary Markov chain. Let \( P_2 \) be the transition matrix of a Markov chain on \( \delta^2 \), given by
\[ P_2((x_1, x_1'), A \times A') = \int_{A_1} \cdots \int_{A_d} \mathbf{1}_{(\phi(x_1, x_2, \ldots, x_d, z) \in A, \phi'(x_1', x_2, \ldots, x_d, z) \in A')} d\nu(z) d\pi(x_2) \cdots d\pi(x_d). \]
Let \( P^- \) be the restriction of \( P_2 \) to nondiagonal terms, and let \( \rho \) be the Perron–Frobenius eigenvalue of the matrix corresponding to \( P^- \).

The following theorem gives a necessary and sufficient condition for endogeny of the tree-indexed solution corresponding to \( \mu \). This is a small generalisation of Theorem 1 of [9].

**Theorem 4.4.** The tree-indexed solution to the RDE associated with
\[ \xi_u = \phi(\xi_u 0, \xi_u 1, \ldots, \xi_u (d-1), \varepsilon_u), \]
corresponding to the invariant measure \( \pi \), is endogenous if \( d\rho < 1 \); it is nonendogenous if \( d\rho > 1 \). In the critical case, \( d\rho = 1 \), let \( \mathcal{H}_0 \) be the collection of \( L^2 \) random variables measurable with respect to \( \xi_\emptyset \) and let \( \mathcal{K} \) denote the \( L^2 \) random variables measurable with respect to \( (\varepsilon_u; u \in \Gamma_d) \). Then endogeny holds in this case provided that \( P^- \) is irreducible and \( \mathcal{H}_0 \cap \mathcal{K}^\perp = \{0\} \). See [9] for full details.
The noisy veto-voter model

Theorem 4.5. Consider the RDE

\[ X_u = 1 - \prod_{i=1}^{N_n} X_{ui}. \] (4.5)

Then, by Lemma 3.2, there exists an invariant probability measure on \( \{0, 1\} \) for (4.5). Let \( \mu_n^1 \) denote the probability of a 1 under this invariant measure. Then the corresponding tree-indexed solution is endogenous if and only if \( H'_n(\mu_n^1) \leq 1 \).

Proof. Let \( N^* = \text{ess sup} N < \infty \) be a bound for \( N \). We can then think of the random tree with branching factor \( N \) as being embedded in an \( N^* \)-ary tree. Each vertex has \( N^* \) daughter vertices and the first \( N \) of these are thought of as being alive (the remaining are thought of as being dead). In this context our RDE reads

\[ X = 1 - \prod_{\text{live} u} X_u. \]

We now compute the transition probabilities from the previous theorem. First consider the transition from \((0, 1)\) to \((1, 0)\). The first coordinate automatically maps to 1 and the second maps to 0 provided that all of the inputs not on the distinguished line of descent are equal to 1. The conditional probability of the vertex on the distinguished line of descent being alive is \( N/N^* \), since there are \( N^* \) vertices, of which \( N \) are alive. The probability of the remaining \( N - 1 \) vertices each taking value 1 is \( (\mu_n^1)^{N-1} \), and so the probability of a transition from \((0, 1)\) to \((1, 0)\), conditional on \( N \), is just

\[ \frac{1}{(N^* - 1)} \frac{(\mu_n^1)^{N-1} N}{d}. \]

Taking expectations, the required probability is

\[ E \left[ \frac{(\mu_n^1)^{N-1} N}{d} \right] = E \frac{1}{(N^* - 1)} N \frac{(\mu_n^1)^{N-1}}{N^*} = \frac{H'_n(\mu_n^1)}{N^*}. \]

The probability of a transition from \((1, 0)\) to \((0, 1)\) is the same by symmetry. Hence, \( P^- \) is given by

\[ P^- = \begin{pmatrix} 0 & H'_n(\mu_n^1) \\ H'_n(\mu_n^1) & 0 \end{pmatrix}, \]

and the Perron–Frobenius eigenvalue \( \rho \) is \( H'_n(\mu_n^1)/N^* \). By Theorem 4.4, the criterion for endogeny is \( N^* \rho \leq 1 \), i.e. \( H'_n(\mu_n^1) \leq 1 \), provided that, in the critical case, \( H'_n(\mu_n^1) = 1 \), we verify the stated nondegeneracy conditions.

It is easily seen that \( P^- \) is irreducible. For the other criterion, let \( X \in \mathcal{H}_0 \cap \mathcal{K}^\perp \), so that \( X = f(X_0) \) for some \( L^2 \) function \( f \) and \( E[XY] = 0 \) for all \( Y \in \mathcal{K} \). Taking \( Y = 1 \), we obtain \( E[X] = 0 \). Writing \( X \) as

\[ X = a \mathbf{1}_{(X_0=1)} + b \mathbf{1}_{(X_0=0)}, \]

where \( a \) and \( b \) are constants, we obtain

\[ X = a \mathbf{1}_{(X_0=1)} - \frac{a \mu_n^1}{1 - \mu_n^1} \mathbf{1}_{(X_0=0)}. \]
For convenience, we shall scale by taking \(a = 1\) (we assume that \(X \not= 0\)):

\[
X = 1_{(X_\emptyset = 1)} - \frac{\mu_n^1}{1 - \mu_n^1} 1_{(X_\emptyset = 0)}.
\]

Now, for each \(k\), take \(Y_k = 1_{(N_\emptyset = k)} \in \mathcal{K}\). Then

\[
E[XY_k] = E \left[ 1_{(N_\emptyset = k)} \left( 1_{(X_\emptyset = 1)} - \frac{\mu_n^1}{1 - \mu_n^1} 1_{(X_\emptyset = 0)} \right) \right]
\]

\[
= P(N = k) \left[ 1 - (\mu_n^1)^k - \frac{(\mu_n^1)^{k+1}}{1 - \mu_n^1} \right]
\]

\[
= P(N = k) \left( 1 - \frac{(\mu_n^1)^k}{1 - \mu_n^1} \right).
\]

Now if we sum this expression over \(k\), we obtain

\[
1 - H_n(\mu_n^1)/1 - \mu_n^1 = 0.
\]

So either each term in the sum is 0, or one or more are not. But at least two of the probabilities are nonzero by assumption (at least for sufficiently large \(n\)), whilst the term \(1 - (\mu_n^1)^k/(1 - \mu_n^1)\) can only disappear for at most one choice of \(k\). Hence, at least one of the terms is nonzero and this contradicts the assumption that \(X \in \mathcal{H}_0 \cap \mathcal{K}^\perp\).

**Proof of the remainder of Theorem 4.2.** We prove that \(H'(\mu^1) > 1\) implies that \(S\) is not endogenous so that \(C\) cannot equal \(S\).

By Theorem 4.5 we know that the RDE (4.5) has two invariant distributions if and only if \(H'(\mu^1) > 1\). But we know that \(C_u^{(n)}\) converges to \(C_u\) in \(L^2\) and, hence, \(\mu_n^2 \to \mu^2 \neq \mu^1\), so that \(S_u\) and \(C_u\) have different second moments. It now follows that \(S_u\) does not have the same distribution as \(C_u\). Since \([0, 1]\) is bounded, this sequence of moments determines a unique distribution which is therefore that of \(C_u\): see Theorem 1 of Chapter VII.3 of [5].

5. **Basins of attraction**

Now we consider the basin of attraction of the endogenous solution. That is, we ask for what initial distributions does the corresponding solution at root, \(X_\emptyset\), converge (in law) to the endogenous solution.

**Definition 5.1.** Let \(\varsigma\) be the law of the endogenous solution. Suppose that we insert i.i.d. random variables with law \(\nu\) at level \(n\) of the tree and apply the RDE to obtain the corresponding solution \(X_n^u(\nu)\) (with law \(T_n^{|u|}(\nu)\)) at vertex \(u\).

The basin of attraction \(B(\pi)\) of any solution is given by

\[
B(\pi) = \{v \in \mathcal{P} : T_n^{|u|}(\nu) \xrightarrow{w} \pi\}
\]

(where \(\xrightarrow{w}\) denotes weak convergence), which is, of course, equivalent to the set of distributions \(v\) for which \(X_n^u(\nu)\) converges in law to a solution \(X\) of the RDE, with law \(\pi\).

5.1. **The unstable case: \(H'(\mu^1) > 1\)**

**Lemma 5.1.** Suppose that \(H'(\mu^1) > 1\). Then \(X_n^u(\nu) \xrightarrow{L^2} C_u\), the endogenous solution, for any \(v\) with mean \(\mu^1\) other than the discrete measure on \([0, 1]\).
Proof. Let \( E_k = \mathbb{E}[X^n_u(\nu)^2] \), where \( k = n - |u| \), and let \( r_k = \mathbb{E}[C_u X^n_u(\nu)] \). Then
\[
\mathbb{E}[(X^n_u(\nu) - C_u)^2] = E_k - 2r_k + \mu^2.
\]
Now,
\[
E_k = \mathbb{E} \left[ 1 - 2 \prod_{i=1}^{N_u} X_{nu}(v) + \prod_{i=1}^{N_u} X^n_{nu}(\nu)^2 \right] = 1 - 2H(\mu^1) + H(E_{k-1}).
\]
This is a recursion for \( E_k \) with at most two fixed points (recall that the equation \( H(x) - x = \text{constant} \) has at most two roots). Recalling the moment equation, (4.1), these are easily seen to be \( \mu^1 \) and \( \mu^2 \), the first and second moments of the endogenous solution. We have assumed that \( \nu \) is not the discrete distribution and so its second moment (i.e. \( \mathbb{E}_0 \)) must be strictly less than \( \mu^1 \).

Now, under the assumption that \( H'(\mu^1) > 1 \), \( \mu^1 \) and \( \mu^2 \) lie either side of the minimum \( \mu^* \) of \( H(x) - x \) and \( H'(\mu^*) = 1 \), so that \( H'(\mu^2) < 1 \). Hence, \( \mu^2 \) is the stable fixed point and it now follows that \( E_k \) converges to \( \mu^2 \).

The recursion for \( r_k \) is essentially the same as that for \( E_k \):
\[
\mu^2 - r_k = H(\mu^2) - H(r_{k-1}).
\]
This has \( \mu^1 \) and \( \mu^2 \) as fixed points and, since
\[
r_0 = \mathbb{E}[C_u X_u(\nu)] \leq \sqrt{\mathbb{E}[C^2_u] \mathbb{E}[X_u(\nu)^2]} < \sqrt{\mu^1 \mu^2} = \mu^1,
\]
we are in the same situation as with \( E_k \). That is, we start to the left of \( \mu^1 \) and, because \( H'(\mu^1) > 1 \), we conclude that \( \mu^1 \) is repulsive and it follows that \( r_k \) converges to \( \mu^2 \) under the assumptions of the lemma. Hence,
\[
\mathbb{E}[(X^n_u(\nu) - C_u)^2] = E_k - 2r_k + \mu^2 \to 0.
\]

Theorem 5.1. Let \( \delta \) denote the discrete distribution on \([0, 1]\) with mean \( \mu^1 \). Then
\[
B(\varsigma) = \left\{ \nu \in \mathcal{P} : \int x \, d\nu(x) = \mu^1 \text{ and } \nu \neq \delta \right\}.
\]
That is, \( B(\varsigma) \) is precisely the set of distributions on \([0, 1]\) with the correct mean (except the discrete distribution with mean \( \mu^1 \)).

Proof. We have already shown that
\[
\left\{ \nu \in \mathcal{P} : \int x \, d\nu(x) = \mu^1 \text{ and } \nu \neq \delta \right\} \subseteq B(\varsigma).
\]
Since the identity is bounded on \([0, 1]\), we conclude that
\[
\mathbb{E}[X^n_u(\nu)] \to \mathbb{E}[C_u] \text{ if } \nu \in B(\varsigma),
\]
so that \( \nu \in B(\varsigma) \) only if the mean of \( T^n(\nu) \) converges to \( \mu^1 \). From the moment equation, (4.1), the mean of \( X^n_u(\nu) \) is obtained by iterating the map \( f \) \( n \) times, starting with the mean of \( \nu \). This mapping has a unique fixed point \( \mu^1 \) and, since \( H'(\mu^1) > 1 \), it is repulsive. It follows that the only way we can have convergence in mean is if we start with the correct mean, that is, if \( \nu \) has mean \( \mu^1 \). Hence,
\[
B(\varsigma) \subseteq \left\{ \nu \in \mathcal{P} : \int x \, d\nu(x) = \mu^1 \text{ and } \nu \neq \delta \right\}.
5.2. The stable case: $H'(\mu^1) \leq 1$

**Theorem 5.2.** Let $b(\mu^1)$ be the basin of attraction of $\mu^1$ under the iterative map for the first moment, $f:\ t \mapsto 1 - H(t)$. Then

$$B(\zeta) = \left\{ \nu \in \mathcal{P} : \int x \, d\nu(x) \in b(\mu^1) \right\}.$$

Once again consider $\mathbb{E}[(X^n_u(\nu) - C_u)^2]$. Let $m_k^0 = \mathbb{E}[X^n_u(\nu)^0]$, where $k = n - |u|$. Then

$$m_k^2 = \mathbb{E} \left[ 1 - 2 \prod_{i=1}^{N_u} X^n_{ui}(\nu) + \prod_{i=1}^{N_u} X^n_{ui}(\nu)^2 \right]$$

$$= 1 - 2H(m_{k-1}^1) + H(m_{k-1}^2).$$

Recalling that $r_k = \mathbb{E}[C_uX^n_u(\nu)]$, we have

$$r_k = \mathbb{E} \left[ \left( 1 - \prod_{i=1}^{N_u} C_{ui} \right) \left( 1 - \prod_{i=1}^{N_u} X^n_{ui}(\nu) \right) \right]$$

$$= \mathbb{E} \left[ 1 - \prod_{i=1}^{N_u} C_{ui} - \prod_{i=1}^{N_u} X^n_{ui}(\nu) + \prod_{i=1}^{N_u} C_{ui}X^n_{ui}(\nu) \right]$$

$$= 1 - H(\mu^1) - H(m_{k-1}^1) + H(r_{k-1}).$$

We now turn our attention to analysing the dynamics of $m_k^2$ and $r_k$. We will concentrate on the equation for $m_k^2$, as the equation for $r_k$ is essentially the same. By assumption, $m_k^1$ converges to $\mu^1$ and so we may approximate $m_k^1$, for $k \geq k_\epsilon$ (say), by $\mu^1 \pm \epsilon$ for some small $\epsilon > 0$.

**Lemma 5.2.** The trajectory $l_k$ of the dynamical system defined by the recursion

$$l_k = 1 - 2H(\mu^1 + \epsilon) + H(l_{k-1}), \quad l_k = m_k^2,$$

is a lower bound for $m_k^2$ for all $k \geq k_\epsilon$, where $k_\epsilon$ is a positive integer chosen so that $|m_1^1 - \mu^1| < \epsilon$ for $k \geq k_\epsilon$.

The proof is obvious.

**Lemma 5.3.** Let $f(\epsilon)(x) = 1 - 2H(\mu^1 + \epsilon) + H(x)$, $x \in [0, 1]$.

Then, for sufficiently small $\epsilon > 0$, $f(\epsilon)$ has a unique fixed point $\mu^1(\epsilon)$ for which $\mu^1(\epsilon) < \mu^*$. Moreover, as $\epsilon \to 0$, $\mu^1(\epsilon) \to \mu^1$.

**Proof.** This follows from uniform continuity, the fact that $H(\mu^1 + \epsilon) < H(\mu^1)$, and the fact that $H'(\mu^1) \leq 1$ implies that $\mu^1 \leq \mu^*$.

**Lemma 5.4.** $l_k$ converges to $\mu^1(\epsilon)$. 
Proof. We have \( l_k = f_k^{-k}(l_k) \), and so we need verify only that \( l_k \) is in the basin of attraction of \( \mu^1(\epsilon) \) and that \( \mu^1(\epsilon) \) is stable. We know that

\[
 f\epsilon(\mu^1 + \epsilon) < \mu^1 + \epsilon,
\]

since \( 1 - H(\mu^1 + \epsilon) < 1 - H(\mu^1) = \mu^1 \), and so it must be the case that \( \mu^1 + \epsilon \in (\mu^1(\epsilon), p(\epsilon)) \).

It now follows that \( \mu^1(\epsilon) \) converges to \( \mu^1 \) as \( \epsilon \) tends to 0 (by the previous lemma) and, therefore, \( \mu^1(\epsilon) \) can be made arbitrarily close to \( \mu^1 \) by choosing \( \epsilon \) to be sufficiently small. This means that, for sufficiently small \( \epsilon \), \( H'(\mu^1(\epsilon)) < 1 \) by the continuity of \( H' \).

\[
 f_k^{-k}(l_k) \text{ converges to } \mu^1(\epsilon).
\]

Proof of Theorem 5.2. The preceding lemmas tell us that

\[
 \liminf_{k \to \infty} m_k^2 \geq \lim_{k \to \infty} l_k = \mu^1(\epsilon).
\]

Letting \( \epsilon \) tend to 0, we obtain

\[
 \liminf_{k \to \infty} m_k^2 \geq \mu^1.
\]

The fact that \( m_k^2 \leq m_1^2 \) for every \( k \) gives us the following corresponding inequality for the \( \limsup \):

\[
 \limsup_{k \to \infty} m_k^2 \leq \lim_{k \to \infty} m_k^1 = \mu^1.
\]

We conclude that \( m_k^2 \) converges to \( \mu^1 \).

Now,

\[
 E((X^n U(\nu) - C_u)^2) = m_k^2 - 2r_k + \mu^2,
\]

so that \( E((X^n U(\nu) - C_u)^2) \to 0 \), remembering that in the stable case the discrete solution and endogenous solution coincide (i.e. \( \mu^1 = \mu^2 \)). We have now shown that

\[
 \left\{ \nu \in \mathcal{P} : \int x \, d\nu(x) \in b(\mu^1) \right\} \subseteq B(\xi),
\]

and the necessity for convergence in mean ensures that we have the reverse inclusion. This completes the proof.

6. Outside the basin of attraction of the endogenous solution

In this section we examine what happens if we iterate distributions with mean outside the basin of attraction of the endogenous solution.

Definition 6.1. Recall that a map \( f \) has an \( n \)-cycle starting from \( p \) if \( f^n(p) = p \), where \( f^n \) denotes the \( n \)-fold composition of \( f \) with itself.

It is easily seen that the map for the first moment \( f : t \mapsto 1 - H(t) \) can have only one- and two-cycles. This is because the iterated map \( f^{(2)} : t \mapsto 1 - H(1 - H(t)) \) is increasing in \( t \) and, hence, can have only one-cycles. Also, note that the fixed points (or one-cycles) of \( f^{(2)} \) come in pairs: if \( p \) is a fixed point then so too is \( 1 - H(p) \).
We consider the following iterated RDE:

\[ X = 1 - \prod_{i=1}^{N_{\mathcal{G}}} \left( 1 - \prod_{j=1}^{N_{i}} X_{ij} \right). \]  

(6.1)

This corresponds to the iterated map on laws on \([0,1]\), \(H_{\mathcal{T}}\), where \(H\) is given at the beginning of Section 2. We denote a generic two-cycle of the map \(f^{(2)}\) by the pair \((\mu_{+}^{1}, \mu_{+}^{1})\).

**Theorem 6.1.** Suppose that \((\mu_{+}^{1}, \mu_{+}^{1})\) is a two-cycle of \(f^{(2)}\). There are at most two solutions of the RDE (6.1) with mean \(\mu_{+}^{1}\). There is a unique endogenous solution \(C^{+}\), and a (possibly distinct) discrete solution, \(S^{+}\), taking values in \([0,1]\). The endogenous solution \(C^{+}\) is given by \(P(S^{+} = 1 \mid \mathcal{T})\) (just as in the noniterated case). The solutions are distinct if and only if \(H'(\mu_{+}^{1})H'(\mu_{+}^{1}) > 1\), i.e. if and only if \(\mu_{+}^{1}\) (or \(\mu_{+}^{1}\)) is an unstable fixed point of \(f^{(2)}\).

**Proof.** This uses the same method as the proofs of results in Section 4.

First, it is clear that \(S^{+}\) is a solution to (6.1), where \(P(S^{+} = 1) = \mu_{+}^{1} = 1 - P(S^{+} = 0)\). Now take interleaved tree-indexed solutions to the RDE on the tree \(\mathcal{T}\), corresponding (on consecutive layers) to means \(\mu_{+}^{1}\) and \(\mu_{+}^{1}\). Then we define

\[ C^{+}_{(n)} = P(S^{+}_{i} = 1 \mid N_{v}; |v| \leq 2n) = 1 - \prod_{i=1}^{N_{\mathcal{G}}} \left( 1 - \prod_{i=1}^{N_{i}} (1 - (\mu_{+}^{1})^{N_{i}N_{v}^{2n-2} \cdots 1}) \right). \]

It follows that \(C^{+}_{(n)}\) converges a.s. and in \(L^{2}\) to \(C^{+}\) and that this must be the unique endogeneous solution (since if \(Z\) is any solution with mean \(\mu_{+}^{1}\) then \(E[Z_{\mathcal{G}} \mid N_{v}; |v| \leq 2n] = C^{+}_{(n)}\).

As in Lemma 4.1, we establish that there are at most two solutions by showing that there are at most two possible moment sequences for a solution and that if \(\mu_{+}^{1}\) is stable (for \(f^{(2)}\)) then the only possible moment sequence corresponds to the discrete solution \(S^{+}\).

To do this, note that, denoting a possible moment sequence starting with first moment \(\mu_{+}^{1}\) by \((\mu_{+}^{1})\), we have

\[ H(\mu_{+}^{k}) = H\left( \sum_{j=0}^{k} (-1)^{j} \binom{k}{j} H(\mu_{+}^{j}) \right) = \sum_{j=0}^{k} (-1)^{j} \binom{k}{j} \mu_{+}^{j}. \]

Then we look for solutions of

\[ H\left( \sum_{j=0}^{k-1} (-1)^{j} \binom{k}{j} \right) H(\mu_{+}^{j}) + (-1)^{k} H(t) = \sum_{j=0}^{k-1} (-1)^{j} \binom{k}{j} \mu_{+}^{j} + (-1)^{k} t, \]

(6.2)
in the range where the argument of \(H\) on the left-hand side is nonnegative and less than 1. In this range \(H\) is increasing and convex, so there are at most two solutions.

Suppose that \(\mu_{+}^{1}\) is a stable fixed point. Then the unique moment sequence is constant, since the other solution of

\[ g(t) := H(1 - 2H(\mu_{+}^{1}) + H(t)) - (1 - 2\mu_{+}^{1} + t) = 0 \]

must be greater than \(\mu_{+}^{1}\) (because \(g'(\mu_{+}^{1}) = H'(\mu_{+}^{1})H'(\mu_{+}^{1}) - 1 \leq 0\).
If $\mu_+^1$ is unstable then there are potentially two solutions for $\mu_+^2$, one of which is $\mu_+^1$. Taking the other potential solution, and seeking to solve (6.2), one of the solutions will give a value for $\mu_+^2$ greater than $\mu^* > \mu_+^2$, which is not feasible, so there will be at most one sequence with $\mu_+^2 \neq \mu_+^1$.

Now, as in the proof of Theorem 4.5, we can show that if $\mu_+^1$ is unstable then, in the corresponding RDE with branching factor truncated by $n$, the two solutions to the RDE are distinct for large $n$, and the endogenous solution converges to $C^+$ in $L^2$ as $n \to \infty$. It follows that there are two distinct solutions in this case.

Given a fixed point $\mu_+^1$ of $f(2)$, denote the law of the corresponding conditional probability solution by $\varsigma^+$. Denote the corresponding basin of attraction (under $T^2$) by $B(\varsigma^+)$, and denote the basin of attraction of $\mu_+^1$ under the map $f(2)$ by $b^2(\mu_+^1)$. Then we have the following theorem.

**Theorem 6.2.** The following dichotomy holds:

(i) if $H'(\mu_+^1)H'(\mu_+^1) > 1$ then

$$B(\varsigma^+) = \{ \pi : \pi \text{ has mean } \mu_+^1 \text{ and } \pi \text{ is not concentrated on } [0, 1] \};$$

(ii) if $H'(\mu_+^1)H'(\mu_+^1) \leq 1$ then

$$B(\varsigma^+) = \{ \pi : \pi \text{ has mean } m \in b^2(\mu_+^1) \}.$$

**Proof.** This can be proved in exactly the same way as Theorems 5.1 and 5.2.

7. Examples

We conclude with some examples.

**Example 7.1.** We first consider the case where $N$ is geometric($\alpha$), so that $P(N = k) = \beta^{k-1}\alpha$ and $H(s) = \alpha s/(1 - \beta s)$ (with $\beta = 1 - \alpha$). It follows that

$$f^{(2)}(s) = s,$$

so that every pair $(s, (1 - s)/(1 - qs))$ is a two-cycle of $f$ and the unique fixed point of $f$ is $1 - \sqrt{\alpha}$. It also follows that $s$ is a neutrally stable fixed point of $f^{(2)}$ for each $s \in [0, 1]$. Thus, we see that the unique endogenous solution to the original RDE is discrete and the value at the root of the tree is the almost sure limit of

$$1 - \prod_{i=1}^{N_{\infty}} \left(1 - \prod_{i_2=1}^{N_{i_1}} \cdots (1 - \sqrt{\alpha})^{N_{i_1-1}} \cdots \right).$$

Moreover, for any $s$, there is a unique solution to the iterated RDE with mean $s$ and it is discrete and endogenous and is the almost sure limit of

$$1 - \prod_{i=1}^{N_{\infty}} \left(1 - \prod_{i_2=1}^{N_{i_1}} \cdots (1 - s^{N_{i_1-1}}) \cdots \right).$$
Example 7.2. Consider the original noisy veto-voter model on the binary tree. It follows from (2.3) that
\[ H(z) = (pH(z) + qz)^2 \implies H(z) = \frac{1 - 2pqz - \sqrt{1 - 4pqz}}{2p^2}. \]
This is nondefective if and only if \( p \leq \frac{1}{2} \) (naturally), i.e. if and only if extinction is certain in the trimmed tree from the original veto-voter model. It is fairly straightforward to show that \( H'(\mu^1) > 1 \) is equivalent to \( p < \frac{1}{2} \). Thus, the endogenous solution is nondiscrete precisely when the trimmed tree is subcritical.

Example 7.3. In contrast to the case of the veto-voter model on the binary tree, the veto-voter model on a trinary tree can show a nonendogenous discrete solution even when the trimmed tree is supercritical. More precisely, the trimmed tree is supercritical precisely when \( p > \frac{1}{3} \), but the discrete solution is nonendogenous if and only if
\[ p < p_e^{(3)} := \frac{3\sqrt{3} - 4}{3\sqrt{3} - 2} \]
and \( p_e^{(3)} > \frac{1}{3} \).

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