

The University of Manchester

Economics Discussion Paper Series EDP-1816

Von Neumann-Gale Dynamics, Market Frictions, and Capital Growth

E. Babaei I.V. Evstigneev K.R. Schenk-Hoppé M.V. Zhitlukhin

November 2018

Economics School of Social Sciences The University of Manchester Manchester M13 9PL

Von Neumann-Gale Dynamics, Market Frictions, and Capital Growth

E. Babaei¹, I. V. Evstigneev², K. R. Schenk-Hoppé³, and M. V. Zhitlukhin⁴

Abstract. The aim of this work is to extend the classical capital growth theory pertaining to frictionless financial markets to models taking into account various kinds of frictions, including transaction costs and portfolio constraints. A natural generalization of the notion of a benchmark investment strategy (Platen, Heath and others) is proposed, and it is shown how such strategies can be used for the analysis of growth-optimal investments. The analysis is based on the classical von Neumann-Gale model of economic growth, a stochastic version of which is used in this study as a framework for the modeling of financial markets with frictions.

Key words and Phrases: capital growth theory, transaction costs, benchmark strategies, numeraire portfolios, random dynamical systems, convex multivalued operators, von Neumann-Gale dynamical systems, rapid paths.

2010 Mathematics Subject Classifications: 37H99, 37H15, 91B62, 91B28. JEL-Classification: C61, C62, O41, G10.

¹Economics Department, University of Manchester, Oxford Road, Manchester M13 9PL, UK. E-mail: esmaeil.babaeikhezerloo@manchester.ac.uk.

²Economics Department, University of Manchester, Oxford Road, Manchester M13 9PL, UK. E-mail: igor.evstigneev@manchester.ac.uk. (Corresponding Author.)

³Economics Department, University of Manchester, Oxford Road, Manchester M13 9PL, UK. E-mail: klaus.schenk-hoppe@manchester.ac.uk.

⁴Steklov Mathematical Institute, Russian Academy of Sciences, 8 Gubkina St. Moscow, 119991, Russia. E-mail: mikhailzh@mi-ras.ru.

1 Introduction

Capital growth theory deals with the following multiperiod investment problem: starting from some initial wealth available at time 0, find a self-financing trading strategy that maximizes the long-run growth rate of investor's wealth. This problem has been investigated by various authors: Kelly [30], Latané [31], Breiman [6], Thorp [50], Ziemba and Vickson [51, 52], Algoet and Cover [2], MacLean et al. [33, 34, 35], Hakansson and Ziemba [22], and others⁵. However, for the most part, results available in this literature pertain to frictionless markets. Some specialized models of markets with frictions have been studied, e.g., by Hausch and Ziemba [23], Taksar et al. [49], Iyengar and Cover [25], Akian et al. [1], and Iyengar [26].

To extend capital growth theory to models of asset markets with frictions, we use the mathematical framework of von Neumann-Gale dynamical systems. Such systems are described in terms of set-valued operators specifying for every state "today" a set of possible states "tomorrow". Characteristic features of the operators associated with von Neumann-Gale systems are certain properties of convexity and homogeneity. The original theory of von Neumann-Gale dynamics (von Neumann [53], Gale [20] and Rockafellar [46]) aimed basically at the mathematical modeling of economic growth. This theory, in its classical form, was purely deterministic: it did not reflect the influence of random factors on economic growth. The importance of taking these factors into account was realized early on. In the 1970s, Dynkin [11, 12, 13], Radner [42, 43] and their research groups made first steps in developing stochastic analogues of the von Neumann-Gale growth model. The initial attack on the problem left many questions unanswered since studies in this direction faced serious mathematical and conceptual difficulties. Substantial progress was made only in the 2000s [16, 17, 4], when new mathematical techniques were developed that made it possible to resolve a number of fundamental problems in the field.

A new stage in the theory of von Neumann-Gale systems began when Dempster et al. [9] observed that stochastic systems of the von Neumann-Gale type can serve as a natural and convenient framework for the modeling of financial markets with frictions. The first results in this direction obtained in [9] were concerned with no-arbitrage pricing and hedging in markets with proportional transaction costs. Extensions of these results to more general

⁵The state of the art in the field is reviewed in MacLean et al. [36]

models, taking into account market interactions, were given by Evstigneev and Zhitlukhin [19].

The first applications of von Neumann-Gale dynamics to capital growth theory under proportional transaction costs were provided by Bahsoun et al. [5]. The main focus of that work was on the analysis of *rapid paths* in von Neumann-Gale systems, generating in financial market models *benchmark strategies* (*numeraire portfolios*), see Platen and Heath [41] and Long [32], and their applications in the theory of growth-optimal investments. In the model examined in [5], as in many other discrete-time capital growth models, short sales were ruled out, so that admissible portfolios were represented by non-negative vectors. Mathematically, this means that the state space of the von Neumann-Gale dynamical system under consideration is the nonnegative cone \mathbb{R}^n_+ in the linear space \mathbb{R}^n .

In real financial markets, short sales are typically allowed but restricted by various trading rules (which might be different for different stock exchanges). The most common rule of this kind is expressed in terms of *margin requirements*, stating that only those portfolios are admissible for which at any moment of time the value of all long positions exceeds the value of all short positions with some excess (margin). In this work, we develop a capital growth model described in terms of a class of von Neumann-Gale dynamical systems in which short selling is allowed under some constraints including, in particular, margin requirements. We assume that the sets of admissible portfolio vectors, as well as self-financing constraints, are described by random cones depending on stochastic factors influencing the market. A crucial mathematical assumption under which our theoretical tools are applicable is the condition that the cones under consideration are *polyhedral*, i.e., generated by a finite number of (possibly random) extreme vectors. Under this condition, we show that the main results of [17, 4, 5] can be extended to our more general model, and moreover, *deduced* from those in the papers cited. From the mathematical point of view, this assumption might seem restrictive, but it is acceptable in the applied perspective, since most, if not all, common models of financial markets with frictions satisfy this requirement.

The paper is organized as follows. Section 2 describes the von Neumann-Gale dynamical systems we deal with. In Section 3 the main assumptions and results are formulated. Sections 4 and 5 provide proofs of the main results. In Section 6, we apply the general results obtained to a specialized model that covers most of the known examples and applications. The Appendix contains some general mathematical facts used in this work.

2 Von Neumann-Gale dynamical systems

Von Neumann-Gale dynamical systems. A von Neumann-Gale multivalued dynamical system is defined by a sequence of cones⁶ X_t , t = 0, 1, 2, ...,in linear spaces and cones

$$G_t \subseteq X_{t-1} \times X_t, \ t = 1, 2, \dots$$
 (1)

Elements of X_t are states of the system at time t, X_t are state spaces, and G_t are transition cones. Sequences x_0, x_1, \dots such that

$$(x_{t-1}, x_t) \in G_t, t = 1, 2, \dots$$

are called *paths* (*trajectories*) of the dynamical system.

A stochastic von Neumann-Gale dynamical system is defined as follows. Let (Ω, \mathcal{F}, P) be a complete probability space and $\mathcal{F}_0 \subseteq \mathcal{F}_1 \subseteq ... \subseteq \mathcal{F}$ a sequence of σ -algebras containing all sets in \mathcal{F} of measure zero. For each t = 0, 1, 2, ..., let $X_t(\omega)$ be a random closed cone in a topological linear space $L(X_t(\omega))$ is the random state space at time t). Further, let $G_t(\omega) \subseteq$ $X_{t-1}(\omega) \times X_t(\omega), t = 1, 2, ...,$ be random closed cones. It is assumed that the cones $X_t(\omega)$ and $G_t(\omega)$ depend \mathcal{F}_t -measurably⁷ on ω , which means that they are determined by events occurring prior to time t. Let \mathcal{L}_t (t = 0, 1, ...) be a linear space of \mathcal{F}_t -measurable vector functions $x(\omega), \omega \in \Omega$, with values in L. We say that a vector function $x(\omega)$ is a random state of the system and write $x \in \mathcal{X}_t$ if $x \in \mathcal{L}_t$ and $x(\omega) \in X_t(\omega)$ almost surely (a.s.). A sequence of random states $x_0 \in \mathcal{X}_0, x_1 \in \mathcal{X}_1, ...$ is called a path of the dynamical system under consideration if

$$(x_{t-1}(\omega), x_t(\omega)) \in G_t(\omega)$$
 (a.s.).

An stationary (autonomous) version of the stochastic von Neumann-Gale dynamical system is defined as follows. Let $T : \Omega \to \Omega$ be an automorphism of the probability space (Ω, \mathcal{F}, P) , i.e., a one-to-one mapping of Ω onto itself such that T and T^{-1} are \mathcal{F} -measurable and preserve the measure P:

$$P(T^{-1}\Gamma) = P(T\Gamma) = P(\Gamma), \ \Gamma \in \mathcal{F}.$$

⁶A set in a linear space is called a (convex) cone if it contains together with any vectors a and b the vector $\lambda a + \mu b$, where λ and μ are any non-negative numbers. We will assume that all the cones under consideration contain non-zero vectors.

⁷A set $X(\omega) \subseteq L$ is said to depend \mathcal{F}_t -measurably on ω if the graph $\{(\omega, a) : a \in A(\omega)\}$ of the multivalued mapping $\omega \mapsto A(\omega)$ belongs to the σ -algebra $\mathcal{F}_t \otimes \mathcal{B}(L)$, where $\mathcal{B}(\cdot)$ is the Borel σ -algebra. A random set is a set $X(\omega) \subseteq L$ depending \mathcal{F} -measurably on ω .

We shall say that the stochastic von Neumann-Gale dynamical system is *autonomous*, or *stationary*, if

$$T^{-1}(\mathcal{F}_t) = \mathcal{F}_{t+1}, \ X_t(T\omega) = X_{t+1}(\omega), \ \text{and} \ G_t(T\omega) = G_{t+1}(\omega).$$
(2)

The mapping T is interpreted as a *time shift*: if $\Gamma \in \mathcal{F}_t$ is an event occurring by time t, then $T^{-1}(\Gamma) \in \mathcal{F}_{t+1}$ is an analogous event occurring one unit of time later (by time t + 1). Autonomous systems serve as a framework for stationary models in various applications (e.g. [53], [20], [16] and [17]).

In this work we consider stochastic von Neumann-Gale dynamical systems for which L is an N-dimensional linear space \mathbb{R}^N and \mathcal{L}_t is the space $L_t^{\infty}(\mathbb{R}^N) = L^{\infty}(\Omega, \mathcal{F}_t, P, \mathbb{R}^N)$ consisting of essentially bounded \mathcal{F}_t -measurable functions with values in \mathbb{R}^N . Those systems for which $X_t(\omega) = \mathbb{R}^N_+$ will be called *canonical*. They are relatively well examined, and our central goal will be to extend the corresponding results to the general setting in which $X_t(\omega)$ are cones that do not necessarily coincide with \mathbb{R}^N_+ and, moreover, depend on t and ω .

In financial applications, random states $x_t \in \mathcal{X}_t$ of a von Neumann-Gale system represent (contingent) portfolios of assets that can be chosen by an investor at date t. These portfolios are specified by random \mathcal{F}_t -measurable vectors of dimension N, where N is the number of assets traded at each date. Portfolio positions can be measured either in terms of "physical units" of assets, or in terms of their values. The transition comes $G_t(\omega)$ define selffinancing constraints: a portfolio x can be transferred to a portfolio y at date t (under transaction costs) if and only if $(x, y) \in G_t(\omega)$. The cones $X_t(\omega)$ can specify various constraints on admissible portfolios, such as short selling constraints for some or all assets, or margin requirements (long portfolio positions must compensate with a certain excess its short positions). Paths x_0, x_1, x_2, \dots of the dynamical system at hand are feasible (*self-financing*) trading strategies, describing possible scenarios of the investor's actions at the financial market influenced by random factors. The fact that $G_t(\omega)$ and $X_t(\omega)$ are cones means that the model takes into account in the most general way proportional transaction costs. The generality of this framework makes it possible to incorporate not only standard single-currency stock market models but also multicurrency models (see Kabanov [28], Kabanov and Safarian [29] and references therein). In this framework, interest rates for borrowing and lending might differ from each other and might be different for different currencies—see Section 6.

Let $X_t^*(\omega)$ denote the dual cone of $X_t(\omega)$:

$$X_t^*(\omega) = \{ p \in \mathbb{R}^N : pa \ge 0, a \in X_t(\omega) \},\$$

where pa is the scalar product of the vectors p and a in \mathbb{R}^N . For shortness, we will use the notation $L_t^1(\mathbb{R}^N)$ for the space $L_1(\Omega, \mathcal{F}_t, P, \mathbb{R}^N)$ of integrable \mathcal{F}_t measurable vector functions with values in \mathbb{R}^N . A dual path (dual trajectory) is a sequence of vector functions $p_1(\omega), p_2(\omega), \dots$ such that $p_t \in L_t^1(\mathbb{R}^N)$ and for almost all ω we have:

$$p_t(\omega) \in X^*_{t-1}(\omega), \ t = 1, 2, ...,$$
(3)

and

$$\bar{p}_{t+1}(\omega)b \le p_t(\omega)a \text{ for all } (a,b) \in G_t(\omega), \ t = 1, 2, ...,$$
(4)

where $\bar{p}_{t+1}(\omega) := E_t p_{t+1}(\omega)$ and $E_t(\cdot) = E(\cdot | \mathcal{F}_t)$ is the conditional expectation given \mathcal{F}_t .

Note that for a canonical system, we have $X_{t-1}^*(\omega) = \mathbb{R}^N_+$, so that elements p_t of a dual path are functions belonging to the non-negative cones $L_t^{1,+}(\mathbb{R}^N)$ of the spaces $L_t^1(\mathbb{R}^N)$.

It follows from (4) that for any path $x_0, x_1, ...,$ the random sequence $p_{t+1}x_t$, t = 0, 1, ..., is a supermartingale with respect to the filtration $\mathcal{F}_1 \subseteq \mathcal{F}_2 \subseteq$... and the random sequence $\bar{p}_{t+1}x_t$, t = 0, 1, ..., is a supermartingale with respect to the filtration $\mathcal{F}_0 \subseteq \mathcal{F}_1 \subseteq ...$ This is immediate from the relations:

$$E_t p_{t+1} x_t = \bar{p}_{t+1} x_t \le p_t x_{t-1}$$
 (a.s.), $t = 1, 2, ...,$

and

$$E_{t-1}\bar{p}_{t+1}x_t \le E_{t-1}p_tx_{t-1} = \bar{p}_tx_{t-1}$$
 (a.s.), $t = 1, 2, ...,$

following from (4).

A standard argument using measurable selection (see Theorem A.1 in the Appendix) shows that (4) holds if and only if

$$Ep_{t+1}(\omega)y_t(\omega) \le Ep_t(\omega)x_t(\omega) \tag{5}$$

for all pairs $(x_t(\omega), y_t(\omega))$ of functions in $L^{\infty}_t(\mathbb{R}^N) \times L^{\infty}_t(\mathbb{R}^N)$ such that

$$(x_t(\omega), y_t(\omega)) \in G_t(\omega)$$
 (a.s.).

In the financial context, dual paths are termed *consistent price systems*. They generalize the concept of an equivalent martingale measure involved in classical no-arbitrage criteria—see Jouini and Kallal [27], Cvitanić and Karatzas [8], Schachermayer [47], Dempster et al. [9], Guasoni et al. [21], Kabanov and Safarian [29] and others. The coordinates p_t^i of the vectors p_t are interpreted as market consistent prices of assets if portfolio positions are measured in terms of units of assets. If they are measured in monetary terms, then p_t^i might be interpreted as market consistent discount factors.

A central notion in this theory is the notion of a rapid path. Let us say that a dual path p_1, p_2, \dots supports a path x_0, x_1, \dots if

$$p_{t+1}x_t = 1, \ t = 0, 1, \dots$$
(a.s.). (6)

A trajectory is called *rapid* if there exists a dual trajectory supporting it.

The term "rapid" is motivated by the fact that

$$\frac{\bar{p}_{t+1}y_t}{p_t y_{t-1}} \le \frac{\bar{p}_{t+1}x_t}{p_t x_{t-1}} = 1, \ t = 1, 2, \dots \text{ (a.s.)}$$
(7)

for each path $y_0, y_1, ...$ with $p_t y_{t-1} > 0$ (see (4) and (6)). This means that the path $x_0, x_1, ...$ maximizes the conditional expectation given \mathcal{F}_t of the growth rate $p_{t+1}y_t/p_t y_{t-1}$ at each time t, the maximum being equal to 1. Growth rates are measured by using the random "price systems" p_t . Another justification of the above term is related to the fact that rapid paths are asymptotically growth-optimal almost surely: they exhibit the fastest growth over an infinite time horizon with probability one (see Theorem 2 below).

In the context of the present model, rapid paths may be regarded as analogues of *benchmark strategies* (numeraire portfolios), see Platen and Heath [41] and Long [32]. As we have noticed, the price system (or the system of discount factors) (p_t) involved in the definition of a rapid path is such that the value $p_{t+1}x_t$ of the portfolio x_t is always equal to one, while for any other feasible sequence (y_t) of contingent portfolios (self-financing trading strategy), the values $p_{t+1}y_t$ form a supermartingale. In models with unlimited short selling (cf. [41]), one can speak of martingales rather than supermartingales.

3 Assumptions and the main results

For a vector a, let us denote by $\mathbb{B}(a, r)$ the ball $\{b : |b - a| \leq r\}$, where $|\cdot|$ is the sum of the absolute values of the coordinates of a vector. Let us introduce the following conditions.

(F) There exist \mathcal{F}_t -measurable random vectors $f_{t,k}(\omega), k = 1, ..., K$, such that for each ω we have $f_{t,k}(\omega) \neq 0$,

$$X_t(\omega) = \left\{ a : a = \sum_{k=1}^K f_{t,k}(\omega) c^k \text{ for some } c^k \ge 0, \ k = 1, ..., K \right\}$$
(8)

and

$$\theta_t |c| \le |\sum_{k=1}^K f_{t,k}(\omega) c^k| \le \Theta_t |c|, \ c = (c^1, ..., c^K) \in \mathbb{R}_+^K,$$
(9)

where $0 < \theta_t < \Theta_t$ (t = 0, 1, ...) are constants and K is a natural number.

(G1) For all $t \ge 1$, $\omega \in \Omega$ and $a \in X_{t-1}(\omega)$, the set $\{b : (a, b) \in G_t(\omega)\}$ is non-empty.

(G2) For each $t \ge 1$ there is a constant M_t such that the set $G_t(\omega)$ is contained in $\{(a,b): |b| \le M_t |a|\}$ for all $\omega \in \Omega$.

(G3) For every $t \geq 1$ there exist a strictly positive constant $\alpha_t > 0$ and a bounded vector function $\hat{z}_t(\omega) = (\hat{x}_{t-1}(\omega), \hat{y}_t(\omega))$ such that $\hat{x}_{t-1}(\omega)$ is \mathcal{F}_{t-1} -measurable, $\hat{y}_t(\omega)$ is \mathcal{F}_t -measurable and $\mathbb{B}(\hat{z}_t(\omega), \alpha_t) \subseteq G_t(\omega)$ for all ω .

The representation (8) of the cone $X_t(\omega)$ means that this cone is *polyhedral*: it is spanned on a finite set of \mathcal{F}_t -measurable random vectors $f_{t,k}(\omega) \neq 0$, k = 1, ..., K (generators of $X_t(\omega)$). By virtue of (8), we have

$$X_t(\omega) = F_t(\omega) \mathbb{R}_+^K, \tag{10}$$

where $F_t(\omega) : \mathbb{R}_+^K \to \mathbb{R}^N$ is the linear operator transforming $c = (c^1, ..., c^K) \in \mathbb{R}_+^K$ into $a = \sum_{k=1}^K f_{t,k}(\omega) c^k \in \mathbb{R}^N$. The inequalities in (9) can be written

$$\theta_t|c| \le |F_t(\omega)c| \le \Theta_t|c|, \ c = (c_1, ..., c_K) \in \mathbb{R}_+^K.$$
(11)

For a real number r define $r_+ := \max\{r, 0\}$ and $r_- := \max\{-r, 0\}$, so that $r = r_+ - r_-$. If $c = (c_1, ..., c_K)$, then c_+ and c_- stand for the vectors with the coordinates $(c_k)_+$ and $(c_k)_-$, respectively. Note that the second inequality in (11) implies the analogous inequality holding for all $c \in \mathbb{R}^K$, and not only for $c \in \mathbb{R}^K_+$. Indeed, we have $F_t c = F_t c_+ - F_t c_-$, and so

$$|F_tc| = |F_tc_+ - F_tc_-| \le |F_tc_+| + |F_tc_-| \le \Theta_t(|c_+| + |c_-|) = \Theta_t|c|.$$

We formulate an assumption on the cone $X_t(\omega)$ that guarantees the validity of condition (9) and has a natural financial interpretation. (M) There exists a constant $\mu_t > 1$ such that

$$\mu_t \sum_{i=1}^N a_-^i \le \sum_{i=1}^N a_+^i \text{ for all } a = (a^1, ..., a^N) \in X_t(\omega).$$
(12)

Proposition 1. If the cone $X_t(\omega)$ is representable in the form (8) with some generators $f_{t,k}(\omega)$, and condition (**M**) holds, then $X_t(\omega)$ is representable in the form (8) with generators $f_{t,k}(\omega)$ satisfying (9), i.e., $X_t(\omega)$ satisfies condition (**F**).

To explain the meaning of condition (**M**) suppose that positions of a portfolio $a \in X_t(\omega)$ are measured in terms of their values (expressed in some price system). Then the sums $\sum_{i=1}^{N} a_+^i$ and $\sum_{i=1}^{N} a_-^i$ represent the total value of the long and the short positions of a, respectively. Condition (12) means that for admissible portfolios, the long positions must cover the short ones with a certain margin μ_t . Margin requirements of the type (12) are quite common in financial practice. They restrict short selling to exclude bankruptcy under sudden price jumps. For a proof of Proposition 1 see Section 6.

The results of this paper are concerned with general (non-stationary) and stationary von Neumann-Gale dynamical systems. The main result pertaining to the former ones is as follows.

Theorem 1. Let conditions (G1)-(G3) and (F) hold. Let x_0 be a function in \mathcal{X}_0 such that $\mathbb{B}(x_0(\omega), \varepsilon) \subseteq X_0(\omega)$ (a.s.), where ε is a strictly positive constant. Then the following assertions are valid.

(i) For each $n \ge 1$, there exists a finite rapid path of length n with the initial state $x_0(\omega)$.

(ii) There exists an infinite rapid path with the initial state $x_0(\omega)$.

An important property of infinite rapid paths, which determines their role in capital growth theory, is their a.s. asymptotic optimality. A path $x_0, x_1, ...$ is called *asymptotically growth-optimal* if for any other path $y_0, y_1, ...$ there exists a supermartingale ξ_t such that

$$\frac{|y_t|}{|x_t|} \le \xi_t, \ t = 0, 1, \dots \text{ (a.s.)}.$$

The property of asymptotic growth-optimality, as defined above, has the following important implications. If $|y_t|/|x_t| \leq \xi_t$, t = 0, 1, ... (a.s.), where ξ_t is a supermartingale, the following assertions hold.

(a) For any constant a > 0

$$P\left(\sup_{t\geq 0}\frac{|y_t|}{|x_t|}\geq a\right)\leq \frac{E\xi_0}{a},$$

and, in particular, $\sup_t(|y_t|/|x_t|) < \infty$ a.s., i.e. no strategy can grow asymptotically faster than x_0, x_1, \dots (a.s.).

(b) The strategy x_0, x_1, \dots maximizes a.s. the exponential growth rate:

$$\limsup_{t \to \infty} \frac{1}{t} \ln \frac{|y_t|}{|x_t|} \le 0 \text{ a.s.}$$

(c) For any stopping time τ

$$E\frac{|y_{\tau}|}{|x_{\tau}|} \le E\xi_0$$
 and $E\ln\frac{|y_{\tau}|}{|x_{\tau}|} \le \ln E\xi_0.$

Assertion (a) follows from Doob's inequality for non-negative supermartingales: $P(\sup_t \xi_t \ge a) \le E\xi_0/a$. Assertion (b) is immediate from that $\sup_t (|y_t|/|x_t|) < \infty$ a.s. The first part of assertion (c) holds because $E\xi_{\tau} \le \lim_{t\to\infty} E\xi_{\tau\wedge t} \le E\xi_0$ by Fatou's lemma and Doob's stopping theorem applied to bounded stopping times $\tau \wedge t$. The second part of (c) follows from there by Jensen's inequality.

Note that the above properties (a)–(c) remain valid (but maybe with different constants in the right-hand sides of the inequalities in (a) and (c)) if $|x_t|$ and $|y_t|$ are replaced by $\phi_t(\omega, x_t)$ and $\phi_t(\omega, y_t)$ respectively with any function $\phi_t(\omega, b)$, possibly random and depending on t, which satisfies the following condition.

(**L**) There exist non-random constants $0 < l \leq L$ such that $l|b| \leq \phi_t(\omega, b) \leq L|b|$ for all t, ω and $b \in X_t(\omega)$.

As an example of such a function, we can consider the *liquidation value* of a portfolio $b = (b^1, ..., b^N) \in X_t(\omega)$:

$$\phi_t(\omega, b) = \sum_{i=1}^N (1 - \lambda_{t,i}^+(\omega)) b_+^i - \sum_{i=1}^N (1 + \lambda_{t,i}^-(\omega)) b_-^i,$$
(13)

where $0 \leq \lambda_{t,i}^+(\omega) < 1$ and $\lambda_{t,i}^-(\omega) \geq 0$ are transaction cost rates for selling and buying assets assets (for details see Section 6). Conditions under which the function (13) satisfies (**L**) are given in Proposition 2 below. **Proposition 2.** Let all the cones $X_t(\omega)$ be representable in the form (8) and the margin requirement (**M**) is satisfied with $\mu_t = \mu$ independent of t. If there exist constants $\underline{\Lambda}, \overline{\Lambda}$ such that $0 < \underline{\Lambda} \leq 1 - \lambda_{t,i}^+(\omega)$ and $1 + \lambda_{t,i}^-(\omega) \leq \overline{\Lambda}$ for all ω, t, i , and $\mu \underline{\Lambda} > \overline{\Lambda}$, then the liquidation value $\phi_t(\omega, b)$ defined in (13) satisfies condition (\underline{L}).

For a proof of this proposition, see Section 6.

Let us introduce the following condition:

(G4) There exist a real number $\gamma > 0$ and a natural number m such that for every $t \ge 0$ and every random vector $y_t \in \mathcal{X}_t$, one can find random vectors $y_{t+1} \in \mathcal{X}_{t+1}, \dots, y_{t+m} \in \mathcal{X}_{t+m}$ satisfying

$$(y_t, y_{t+1}) \in G_{t+1}(\omega), ..., (y_{t+m-1}, y_{t+m} + y) \in G_{t+m}(\omega)$$
 (a.s.) (14)

for each $y \in L^{\infty}_{t+m}(\mathbb{R}^N)$ with $|y(\omega)| \leq \gamma |y_t(\omega)|$ (a.s.).

The next result shows that under fairly general assumptions, any rapid path is asymptotically growth-optimal.

Theorem 2. If condition (G4) and condition (G2) with a constant $M_t = M$ independent of t hold, then any rapid path is a.s. asymptotically growth-optimal.

Remark 1. Theorem 2 remains valid if condition (G4) is replaced by the following one:

(G5) The cones $X_t(\omega)$ (t = 0, 1, ...) contain \mathbb{R}^N_+ . There exist a real number $\gamma > 0$ and a natural number m such that for every $t \ge 0$ and every random vector $y_t \in \mathcal{X}_t$, there are random vectors $y_{t+1} \in \mathcal{X}_{t+1}, ..., y_{t+m} \in \mathcal{X}_{t+m}$, satisfying

$$(y_t, y_{t+1}) \in G_{t+1}(\omega), ..., (y_{t+m-1}, y_{t+m})) \in G_{t+m}(\omega)$$
 (a.s.)

and

$$y_{t+m}(\omega) \ge \gamma e |y_t|$$
 (a.s.). (15)

For a proof of this assertion see Proposition 5 in Section 4.

Let us formulate the main result pertaining to stationary (autonomous) systems. Let $T: \Omega \to \Omega$ be an automorphism of the given probability space (Ω, \mathcal{F}, P) such that conditions (2) hold. In the stationary framework, an important role is played by a class of paths called balanced. A path $x_0, x_1, x_2, ...$ is termed *balanced* if there exist an \mathcal{F}_0 -measurable vector function $x(\omega)$ which is normalized by the condition $|x(\omega)| = 1$ (a.s.) and an \mathcal{F}_1 -measurable scalar function $\lambda(\omega) > 0$ with $E|\ln \lambda(\omega)| < +\infty$ such that

$$x_0(\omega) = x(\omega); \ x_t(\omega) = \lambda(\omega)...\lambda(T^{t-1}\omega)x(T^t\omega), \ t \ge 1.$$
(16)

This definition expresses the idea of growth with stationary proportions and at a stationary rate. Clearly a pair of functions $x(\omega)$ and $\lambda(\omega)$, where $x \in \mathcal{X}_0$ and $\lambda(\omega)$ is an \mathcal{F}_1 -measurable scalar function with $E|\ln\lambda(\omega)| < +\infty$, generates a balanced path if and only if $|x(\omega)| = 1$ (a.s.) and

$$(x(\omega), \lambda(\omega)x(T\omega)) \in G_1(\omega) \text{ (a.s.)}.$$
(17)

A balanced path maximizing the expectation of the logarithm of the growth rate $E \ln \lambda(\omega)$ is called a *von Neumann path*. Note that condition (17) implies (by virtue of (**G2**)) that $\lambda(\omega)$ is essentially bounded.

A dual path $p_1, p_2, ...$ is called *balanced* if there exist an \mathcal{F}_1 -measurable vector function $p(\omega)$ and an \mathcal{F}_1 -measurable scalar function $\lambda(\omega) > 0$ such that

$$p_1(\omega) = p(\omega), \ p_t(\omega) = \frac{p(T^{t-1}\omega)}{\lambda(\omega)\dots\lambda(T^{t-2}\omega)}, \ t = 2, 3, \dots$$
(18)

By virtue of (18) and using the invariance properties $G_{t+1}(\omega) = G_t(T\omega)$ and $\bar{p}_{t+1} = T\bar{p}_t$ we can see that a function $p(\omega) \in L_1^1(\mathbb{R}^N)$ and an \mathcal{F}_1 measurable scalar function $\lambda(\omega) > 0$ with $E|\ln\lambda(\omega)| < +\infty$ generate a balanced dual path if and only if for almost all ω we have:

$$p(\omega) \in X_0^*(\omega), \ t = 1, 2, ...,$$
 (19)

and

$$\frac{E_1 p(T\omega)b}{\lambda(\omega)} \le p(\omega)a \text{ for all } (a,b) \in G_1(\omega).$$
(20)

A triplet of functions (x, p, λ) forms a von Neumann equilibrium if the sequence x_0, x_1, \ldots defined by (16) is a balanced path and the sequence p_1, p_2, \ldots defined by (18) is a dual balanced path supporting it.

Note that a triplet (x, p, λ) is a von Neumann equilibrium if and only if the following conditions hold:

- $x(\omega)$ is a function in \mathcal{X}_0 normalized by the condition $|x(\omega)| = 1$ (a.s.);
- $\lambda(\omega) > 0$ is an \mathcal{F}_1 -measurable scalar function with $E|\ln\lambda(\omega)| < +\infty$ satisfying (17) and maximizing $E\ln\lambda(\omega)$ among all balanced paths;
- $p(\omega) \in L_1^1(\mathbb{R}^N)$, and conditions (19), (20) and

$$p(\omega)x(\omega) = 1 \tag{21}$$

are fulfilled for almost all ω .

1

Let us assume that conditions (F) and (G1)-(G4) hold and suppose, additionally, that

$$f_{t+1,k}(\omega) = f_{t,k}(T\omega), \ \hat{z}_{t+1}(\omega) = \hat{z}_t(T\omega),$$

and the constants M_t , α_t , θ_t and Θ_t do not depend on t. Under these assumptions the following assertion is valid.

Theorem 3. A von Neumann equilibrium exists.

This result has the following important consequence. Let (x, p, λ) be a von Neumann equilibrium, whose existence is established in Theorem 3. Consider the balanced path $(x_t)_{t=0}^{\infty}$ generated by the pair (x, λ) . By the definition of an equilibrium, it is rapid. Consequently (see Theorem 2), it is asymptotically growth-optimal. Thus, by virtue of Theorem 3, there exists a balanced path that is asymptotically growth-optimal in the class of *all*, not necessarily balanced paths!

4 General (non-stationary) model

This section contains proofs of the results related to general (non-stationary) models of financial markets. We first establish Theorem 1 by deducing it from analogous results known for canonical von Neumann-Gale systems. At the end of the section, we prove Theorem 2 on the asymptotic optimality of rapid paths.

Let us assume that condition (**F**) holds. A stochastic von Neumann-Gale dynamical system \mathcal{G} defined by the state spaces $X_t(\omega)$ and the transition cones $G_t(\omega)$ generates a canonical von Neumann-Gale system \mathcal{H} with the state spaces \mathbb{R}^K_+ and the transition cones

$$H_t(\omega) := \left\{ (c,d) \in \mathbb{R}_+^K \times \mathbb{R}_+^K : (F_{t-1}(\omega)c, F_t(\omega)d) \in G_t(\omega) \right\}.$$
(22)

We will call \mathcal{H} the canonical von Neumann-Gale dynamical system *induced* by \mathcal{G} .

Proposition 3. If \mathcal{G} satisfies one of the conditions (G1)-(G3), then \mathcal{H} satisfies the analogous condition.

The analogous conditions for \mathcal{H} will be referred to as (H1)-(H3). We formulate them below.

(H1) For all $t \geq 1$, $\omega \in \Omega$ and $c \in \mathbb{R}_{+}^{K}$, the set $\{d : (c,d) \in H_{t}(\omega)\}$ is non-empty.

(H2) For each $t \ge 1$ there is a constant M'_t such that the set $H_t(\omega)$ is contained in $\{(c, d) : |d| \le M'_t |c|\}$ for all $\omega \in \Omega$.

(H3) For every $t \geq 1$ there exist a strictly positive constant $\delta_t > 0$ and a bounded vector function $\hat{w}_t(\omega) = (\hat{u}_{t-1}(\omega), \hat{v}_t(\omega))$ such that $\hat{u}_{t-1}(\omega)$ is \mathcal{F}_{t-1} measurable, $\hat{v}_t(\omega)$ is \mathcal{F}_t -measurable and $\mathbb{B}(\hat{w}_t(\omega), \delta_t) \subseteq H_t(\omega)$ for all ω .

Proof. Suppose (**G1**) holds. Consider any $c \in \mathbb{R}_+^K$. Then $F_{t-1}(\omega)c \in X_{t-1}(\omega)$, and so there exists $b \in X_t(\omega)$ with $(F_{t-1}(\omega)c, b) \in G_t(\omega)$. Since $X_t(\omega) = F_t(\omega)\mathbb{R}_+^K$, there is $d \in \mathbb{R}_+^K$ such that $F_t(\omega)d = b$. Thus

$$(F_{t-1}(\omega)c, F_t(\omega)d) = (F_{t-1}(\omega)c, b) \in G_t(\omega),$$

and so $(c, d) \in H_t(\omega)$, which proves (H1).

Assume that condition (**G2**) is satisfied: $|b| \leq M_t |a|$ for all $(a, b) \in G_t(\omega)$. Let $(c, d) \in \mathbb{R}_+^K \times \mathbb{R}_+^K$ belong to $H_t(\omega)$, i.e. $(F_{t-1}(\omega)c, F_t(\omega)d) \in G_t(\omega)$. By virtue of (11), we have $|F_{t-1}(\omega)c| \leq \Theta_{t-1}|c|$,

$$\theta_t |d| \le |F_t(\omega)d| \le M_t |F_{t-1}(\omega)c| \le M_t \Theta_{t-1} |c|,$$

and so $|d| \leq M'_t |c|$, which yields (H2) with the constant $M'_t = \theta_t^{-1} M_t \Theta_{t-1}$.

Consider the vector function $\hat{z}_t(\omega) = (\hat{x}_{t-1}(\omega), \hat{y}_t(\omega))$ involved in (**G3**). By using (10) and a measurable selection theorem (see Theorem A.1 in the Appendix), we can construct an \mathcal{F}_{t-1} -measurable vector function $\hat{u}_{t-1}(\omega)$ and an \mathcal{F}_t -measurable vector function $\hat{v}_t(\omega)$ with values in \mathbb{R}^K_+ such that

$$F_{t-1}(\omega)\hat{u}_{t-1}(\omega) = \hat{x}_{t-1}(\omega) \text{ and } F_t(\omega)\hat{v}_t(\omega) = \hat{y}_t(\omega).$$

By virtue of the first inequality in (11), the function $\hat{w}_t(\omega) = (\hat{u}_{t-1}(\omega), \hat{v}_t(\omega))$ is bounded. If (c, d) is a vector in $\mathbb{R}^K_+ \times \mathbb{R}^K_+$ such that $|(c, d) - \hat{w}_t(\omega)| \leq \delta_t := \alpha_t / \max\{\Theta_{t-1}, \Theta_t\}$, then $|(F_{t-1}(\omega)c, F_t(\omega)d) - \hat{z}_t(\omega)| \leq \alpha_t$. Consequently, $(F_{t-1}(\omega)c, F_t(\omega)d) \in G_t(\omega)$, and so $(c, d) \in H_t(\omega)$. Thus $\mathbb{B}(\hat{w}_t(\omega), \delta_t) \subseteq H_t(\omega)$, which proves (**H3**).

To apply the results of the previous work to the dynamical system \mathcal{H} , we will need the following simple fact.

Remark 1. Condition (H3) implies the following one.

(h3) For every $t \ge 1$ there exist a strictly positive constant $\zeta_t > 0$ and a bounded vector function $\hat{w}_t(\omega) = (\hat{u}_{t-1}(\omega), \hat{v}_t(\omega))$ such that $\hat{u}_{t-1}(\omega)$ is \mathcal{F}_{t-1} measurable, $\hat{v}_t(\omega)$ is \mathcal{F}_t -measurable, $\hat{w}_t(\omega) \in H_t(\omega)$ for all ω , and $\hat{v}_t(\omega) \ge \zeta_t e$ (coordinate-wise) for all ω , where $e = (1, ..., 1) \in \mathbb{R}_+^K$.

Observe that the function $\hat{w}_t(\omega)$ involved in (H3) has the properties described in (H3) with $\zeta_t = \delta_t/K$. Indeed, we have $|\zeta_t e| = \delta_t$. Thus

 $(\hat{u}_{t-1}(\omega), \hat{v}_t(\omega) - \zeta_t e) \in H_t(\omega)$, and so $\hat{v}_t(\omega) - \zeta_t e \in \mathbb{R}^K_+$, which means that $\hat{v}_t(\omega) \ge \zeta_t e$.

Let us examine relations between paths in \mathcal{H} and \mathcal{G} .

Proposition 4. (a) If (u_t) is a path in the dynamical system \mathcal{H} , then the sequence of functions $x_t(\omega) = F_t(\omega)u_t(\omega)$ forms a path (x_t) in the dynamical system \mathcal{G} .

(b) If (x_t) is a path in \mathcal{G} , then for each t there exists an \mathcal{F}_t -measurable function $u_t(\omega)$ such that $x_t(\omega) = F_t(\omega)u_t(\omega)$. The sequence (u_t) is a path in \mathcal{H} .

Proof. (a) The function $x_t(\omega)$ is \mathcal{F}_t -measurable, and $x_t(\omega) \in X_t(\omega)$ by virtue of (10). Furthermore, $x_t(\omega)$ is essentially bounded because $|F_t(\omega)u_t(\omega)| \leq \Theta_t |u_t(\omega)|$ (see (11)). Since $(u_{t-1}(\omega), u_t(\omega)) \in H_t(\omega)$ (a.s.), we have

$$(F_{t-1}(\omega)u_{t-1}(\omega), F_t(\omega)u_t(\omega)) \in G_t(\omega) \ (a.s.)$$

by the definition of $H_t(\omega)$. Thus (x_t) is a path in \mathcal{G} .

(b) Let (x_t) be a path in \mathcal{G} . Since $x_t(\omega) \in X_t(\omega) = F_t(\omega)\mathbb{R}_+^K$ (a.s.), by the measurable selection theorem there exists an \mathcal{F}_t -measurable function $u_t(\omega)$ with values in \mathbb{R}_+^K such that $x_t(\omega) = F_t(\omega)u_t(\omega)$ (a.s.). In view of (11), we have $|u_t(\omega)| \leq |x_t(\omega)|\theta_t^{-1}$ (a.s.), so that u_t is essentially bounded. Finally, we obtain

$$(F_{t-1}(\omega)u_{t-1}(\omega), F_t(\omega)u_t(\omega)) = (x_{t-1}(\omega), x_t(\omega)) \in G_t(\omega)$$
(a.s.)

because (x_t) is a path in \mathcal{G} . This means that $(u_{t-1}(\omega), u_t(\omega)) \in H_t(\omega)$ (a.s.), which proves that (u_t) is a path in \mathcal{H} .

The path (x_t) in the dynamical system \mathcal{G} defined by $x_t(\omega) = F_t(\omega)u_t(\omega)$ will be called the *image* of the path (u_t) in \mathcal{H} . A key role in the proof of Theorem 1 is played by the following result.

Theorem 4. Let (\tilde{u}_t) be a path in the dynamical system \mathcal{H} and (\tilde{x}_t) its image. If (\tilde{u}_t) is rapid, then (\tilde{x}_t) is rapid in \mathcal{G} .

Proof. Suppose $\tilde{u}_0, \tilde{u}_1, \ldots$ is a rapid path in the dynamical system \mathcal{H} , i.e., there exists a dual path $q_1, q_2, \ldots (q_t \in L^1_t(\mathbb{R}^K_+))$ supporting it:

$$q_{t+1}(\omega)\tilde{u}_t(\omega) = 1 \text{ (a.s.)}, \ t \ge 0,$$

and for almost all ω ,

$$\bar{q}_{t+1}(\omega)d - q_t(\omega)c \le 0, \ (c,d) \in H_t(\omega), \ t \ge 1,$$
(23)

where $\bar{q}_{t+1} := E_t q_{t+1}$. Define $\tilde{x}_t := F_t \tilde{u}_t$. Let us show that (\tilde{x}_t) is a rapid path in the dynamical system \mathcal{G} .

Fix some $t \geq 1$. By the definition of $H_t(\omega)$, a pair of vectors $(c, d) \in \mathbb{R}^K_+ \times \mathbb{R}^K_+$ belongs to $H_t(\omega)$ if and only if $(F_{t-1}(\omega)c, F_t(\omega)d) \in G_t(\omega)$. Thus with probability 1 (for all ω in a set Ω_1 of full measure), we have

$$\bar{q}_{t+1}(\omega)d - q_t(\omega)c \le 0,$$

for any $(c, d) \in \mathbb{R}^K_+ \times \mathbb{R}^K_+$ satisfying

$$(F_{t-1}(\omega)c, F_t(\omega)d) \in G_t(\omega)$$

We will fix $\omega \in \Omega_1$ and omit it in the notation, if this does not lead to ambiguity. Since $\hat{z}_t = (\hat{x}_{t-1}, \hat{y}_t) = (F_{t-1}\hat{u}_{t-1}, F_t\hat{v}_t)$ is contained in $G_t(\omega)$ together with a ball of radius α_t , we can apply the Kuhn-Tucker theorem (see Theorem A.2 in the Appendix) to the convex set X consisting of non-negative vectors (c, d) in $\mathbb{R}^K_+ \times \mathbb{R}^K_+$, the function $\Phi(c, d) = \bar{q}_{t+1}d - q_tc$, the cone $Z = G_t(\omega)$ and the mapping R transforming (c, d) into $(F_{t-1}(\omega)c, F_t(\omega)d) \in \mathbb{R}^N \times \mathbb{R}^N$. By using a measurable selection argument, we construct \mathcal{F}_t -measurable functions $a_t(\omega)$ and $b_t(\omega)$ taking values in \mathbb{R}^N such that for all $\omega \in \Omega_1$,

$$b_t b - a_t a \le 0, \ (a, b) \in G_t(\omega) \tag{24}$$

$$\bar{q}_{t+1}d - q_tc - [b_tF_td - a_tF_{t-1}c] \le 0, \ (c,d) \in \mathbb{R}^K_+ \times \mathbb{R}^K_+,$$
(25)

and

$$|b_t| + |a_t| \le 2N\alpha_t^{-1}(q_t\hat{u}_{t-1} - \bar{q}_{t+1}\hat{v}_t).$$
(26)

Let us show that the sequence of vector functions

$$p_t(\omega) := a_t(\omega), \ t \ge 1,$$

forms a dual path in the dynamical system \mathcal{G} supporting the path (\tilde{x}_t) . The fact that $E|a_t(\omega)| < \infty$ follows from (26). By setting first c = 0 and then d = 0 in (25), we get

$$\bar{q}_{t+1}d \le b_t F_t d, \ d \in \mathbb{R}_+^K, \ t \ge 1,$$

$$(27)$$

$$q_t c \ge a_t F_{t-1} c, \ c \in \mathbb{R}_+^K, \ t \ge 1.$$

$$(28)$$

For all pairs (x, y) of functions in $L_t^{\infty}(\mathbb{R}^N) \times L_t^{\infty}(\mathbb{R}^N)$ with $(x, y) \in G_t(\omega)$ (a.s.), there exist \mathcal{F}_t -measurable vector functions $u(\omega)$ and $v(\omega)$ taking values in \mathbb{R}_{+}^{K} , such that $y = F_t v$, $x = F_{t-1}u$ and $(u, v) \in H_t(\omega)$ (a.s.). The functions u and v are essentially bounded by virtue of the first inequality in (11). Therefore

$$Ep_t x = Ea_t x \ge Eb_t y = Eb_t F_t v$$
$$\ge E\bar{q}_{t+1} v = Eq_{t+1} v \ge Ea_{t+1} F_t v = Ep_{t+1} y , \qquad (29)$$

which yields (5) for $t \ge 1$. In this chain of relations, the first inequality holds because $b_t y \le a_t x$ (a.s.) as long as $(x, y) \in G_t(\omega)$ (a.s.), which follows from (24). The second inequality is a consequence of (27) and the third follows from (28) with t + 1 in place of t.

For $t \geq 0$ we have

$$p_{t+1}\tilde{x}_t = a_{t+1}\tilde{x}_t = a_{t+1}F_t\tilde{u}_t \le q_{t+1}\tilde{u}_t = 1 \text{ (a.s.)}, \tag{30}$$

where the inequality follows from (28) with t replaced by t+1. Furthermore, we get

$$Ep_{t+1}\tilde{x}_t = Ea_{t+1}\tilde{x}_t \ge Eb_{t+1}\tilde{x}_{t+1} = Eb_{t+1}F_{t+1}\tilde{u}_{t+1}$$
$$\ge E\bar{q}_{t+2}\tilde{u}_{t+1} = Eq_{t+2}\tilde{u}_{t+1} = 1.$$
(31)

In this chain of relations, the first inequality holds because $(\tilde{x}_t, \tilde{x}_{t+1}) \in G_{t+1}(\omega)$ (a.s.) and the second follows from (27) with t + 1 in place of t. By combining (30) and (31), we conclude that $p_{t+1}\tilde{x}_t = 1$ (a.s.).

Let us show that for all $t \geq 1$ and almost all $\omega \in \Omega$,

$$p_t(\omega)a \ge 0$$
 for $a \in X_{t-1}(\omega)$,

i.e. $p_t(\omega) \in X_{t-1}^*(\omega)$ (a.s.). Denote by Ω_t the set of those ω for which $b_t F_t \geq 0$. By virtue of (27), we have $P(\Omega_t) = 1$. Fix any $\omega \in \Omega_t$, $a \in X_{t-1}(\omega)$ and consider some $c \in \mathbb{R}_+^K$ for which $F_{t-1}(\omega)c = a$. In view of (H1), there exists $d \in \mathbb{R}_+^K$ such that $(c, d) \in H_t(\omega)$. Then $b_t F_t d \leq a_t F_{t-1}c$ by virtue of (24). From this we get

$$p_t a = a_t F_{t-1} c \ge b_t F_t d \ge 0,$$

which completes the proof.

The following results in the case of a canonical von Neumann-Gale dynamical system are obtained in [4], Theorem 1.

Theorem 5. Let conditions (H1)-(H3) hold. Let $u_0(\omega)$ be a vector function in $L_0^{\infty}(\mathbb{R}^K)$ such that $u_0(\omega) \geq \delta e$ for some constants $\delta > 0$. Then the following assertions are valid.

(i) For each $n \ge 1$, there exists a finite rapid path of length n with the initial state $u_0(\omega)$.

(ii) There exists an infinite rapid path with the initial state $u_0(\omega)$.

Proof of Theorem 1. Let x_0 be a function in \mathcal{X}_0 such that $\mathbb{B}(x_0(\omega), \varepsilon) \subseteq X_0(\omega)$ (a.s.), where ε is a strictly positive constant. By virtue of (10) and the measurable selection theorem, there exists an \mathcal{F}_0 -measurable function $u_0(\omega)$ with values in \mathbb{R}^K_+ such that $x_0(\omega) = F_0(\omega)u_0(\omega)$ (a.s.). By virtue of the first inequality in (11), the function $u_0(\omega)$ is bounded. Let us prove $u_0(\omega) \ge \delta e$, where $\delta = \varepsilon/K$ and $e = (1, ..., 1) \in \mathbb{R}^K_+$. To this end, we define $y(\omega) = F_0(\omega)\delta e$ and observe that

$$|y(\omega)| = \delta |F_0(\omega)e| = \delta |\sum_{k=1}^K f_{0,k}(\omega)| \le \varepsilon,$$

where the inequality follows form the fact we can assume that $f_{0,k}(\omega)$, k = 1, 2, ..., K, $\omega \in \Omega$ are normalized. Then $x_0(\omega) - y(\omega) \in X_0(\omega)$ (a.s.), and so there exists an \mathcal{F}_0 -measurable function $h(\omega)$ with values in \mathbb{R}^K_+ such that

$$F_0(\omega)u_0(\omega) - F_0(\omega)\delta e = F_0(\omega)h(\omega).$$

Then $u_0(\omega) = \delta e + h(\omega)$, which means $u_0(\omega) \ge \delta e$. Here, we use the fact that $F_0(\omega)$ is the one-to-one operator, which follows from (11).

By virtue of Proposition 3, conditions (H1)-(H3) hold. Then by assertion (i) of Theorem 5, for each $n \geq 1$, there exists a finite rapid path $(u_t)_{t=0}^n$ of length n with the initial state $u_0(\omega)$. Let $x_t = F_t u_t$. By virtue of Theorem 4, $(x_t)_{t=0}^n$ is the rapid path of length n in the dynamical system \mathcal{G} with the initial state $x_0(\omega)$. Further, by assertion (ii) of Theorem 5, there exists an infinite rapid path (u_t) with the initial state $u_0(\omega)$. If we put $x_t = F_t u_t$ and use Theorem 4, we conclude that (x_t) is the infinite rapid path in \mathcal{G} . The proof is complete.

Proof of Theorem 2. Let $x_0, x_1, ...$ be a rapid path supported by a dual path $p_1, p_2, ...$ For all $y \in L^{\infty}_{t+m}(\mathbb{R}^N)$ with $|y(\omega)| \leq \gamma |y_t(\omega)|$, by using (4) and (14), we have

$$p_{t+m}y_{t+m-1} \ge \bar{p}_{t+m+1}(y_{t+m}+y) = \bar{p}_{t+m+1}y_{t+m} + \bar{p}_{t+m+1}y \ge \bar{p}_{t+m+1}y \quad (32)$$

because $\bar{p}_{t+m+1}y_{t+m} \ge 0$ (a.s.). The last inequality is valid since $y_{t+m}(\omega) \in X_{t+m}(\omega)$ and $p_{t+m+1}(\omega) \in X_{t+m}^*(\omega)$ (a.s.), which yields $p_{t+m+1}y_{t+m} \ge 0$

(a.s.) and so $\bar{p}_{t+m+1}y_{t+m} \ge 0$ (a.s.). Put

$$y = \frac{\bar{p}_{t+m+1}}{|\bar{p}_{t+m+1}|} \gamma |y_t|.$$
(33)

Then $|y(\omega)| = \gamma |y_t(\omega)|$ and $y \in \mathcal{X}_{t+m}$. Consequently, (32) can be applied to y defined by (33). Observe that

$$\bar{p}_{t+m+1}y = \frac{||\bar{p}_{t+m+1}||^2}{|\bar{p}_{t+m+1}|}\gamma|y_t| \ge |\bar{p}_{t+m+1}|N^{-1}\gamma|y_t|,$$
(34)

where $||\cdot||$ is the Euclidean norm in \mathbb{R}^N (note that $||\cdot|| \ge |\cdot|/\sqrt{N}$). Further, the equality $p_{t+m+1}x_{t+m} = 1$ implies $\bar{p}_{t+m+1}x_{t+m} = 1$, and so

$$\bar{p}_{t+m+1}||x_{t+m}| \ge 1, \tag{35}$$

and it follows from (G2) with a constant $M_t = M$ independent of t that

$$|x_{t+m}| \le M^m |x_t|. \tag{36}$$

By combining (35) and (36), we get

$$|\bar{p}_{t+m+1}| \ge M^{-m} |x_t|^{-1}, \tag{37}$$

and by using (34), (32) and (37), we obtain

$$p_{t+m}y_{t+m-1} \ge M^{-m}N^{-1}\gamma|y_t||x_t|^{-1}, \tag{38}$$

which yields

$$p_{t+1}y_t \ge E_{t+1}p_{t+m+1}y_{t+m-1} \ge M^{-m}N^{-1}\gamma|y_t||x_t|^{-1}.$$
(39)

Since $p_{t+1}y_t$ is a non-negative supermartingale, the proof is complete. \Box

We provide a version of Theorem 2 in which (G4) is replaced by another assumption.

Proposition 5. Theorem 2 remains valid if condition (G4) is replaced by (G5).

Proof. Let $x_0, x_1, ...$ be a rapid path supported by a dual path $p_1, p_2, ...$ Since $X_t(\omega) \supseteq \mathbb{R}^N_+$, any $p_{t+1}(\omega) \in X_t^*(\omega)$ is non-negative (a.s.). By using this, (4), (15) and the fact that $|p_{t+m+1}| \ge M^{-m}|x_t|^{-1}$, which follows from the equality $p_{t+m+1}x_{t+m} = 1$ and (**G2**) with a constant $M_t = M$ independent of t, we get

$$p_{t+1}y_t \ge E_{t+1}p_{t+m+1}y_{t+m} \ge \gamma E_{t+1}|p_{t+m+1}||y_t|$$
$$\ge \gamma E_{t+1}(M^{-m}|x_t|^{-1})|y_t| = \gamma M^{-m}|x_t|^{-1}|y_t|.$$

Thus $|x_t|^{-1}|y_t|$ is dominated by a non-negative supermartingale $p_{t+1}y_t$, which completes the proof.

5 Stationary Model

In the case of a stationary system, let us assume that conditions (F) and (G1)-(G4) hold and suppose, additionally, that

$$f_{t+1,k}(\omega) = f_{t,k}(T\omega), \tag{40}$$

 $\hat{z}_{t+1}(\omega) = \hat{z}_t(T\omega)$, and the constants M_t , α_t , θ_t and Θ_t do not depend on t. Clearly (40) implies that $X_{t+1}(\omega) = x_t(T\omega)$. Let us formulate analogues of the conditions (**G1**)-(**G4**) in the case of a stationary model as follows.

(G1') For all $\omega \in \Omega$ and $a \in X_0(\omega)$, the set $\{b : (a,b) \in G_1(\omega)\}$ is non-empty.

(**G2**) There is a constant M such that the set $G_1(\omega)$ is contained in $\{(a,b): |b| \leq M|a|\}$ for all $\omega \in \Omega$.

(G3') There exist a strictly positive constant $\alpha > 0$ and a bounded vector function $\hat{z}_1(\omega) = (\hat{x}_0(\omega), \hat{y}_1(\omega))$ such that $\hat{x}_0(\omega)$ is \mathcal{F}_0 -measurable, $\hat{y}_1(\omega)$ is \mathcal{F}_1 -measurable and $\mathbb{B}(\hat{z}_1(\omega), \alpha) \subseteq G_1(\omega)$ for all ω .

 $(\mathbf{G4'})$ There exists a natural number m such that for every random state $y_0 \in \mathcal{X}_0$, one can find a real number $\gamma > 0$ and random states $y_1 \in \mathcal{X}_1, \dots, y_m \in \mathcal{X}_m$, satisfying with probability one

$$(y_0(\omega), y_1(\omega)) \in G_1(\omega), ..., (y_{m-1}(\omega), y_m(\omega) + y(\omega)) \in G_m(\omega)$$

$$(41)$$

for all $y \in L_m^{\infty}(\mathbb{R}^N)$ with $|y(\omega)| \leq \gamma |y_0(\omega)|$.

Observe that the above conditions are versions of conditions (G1)-(G4) formulated for one moment of time t = 0. Clearly in the stationary case, these requirements hold for some t if and only if they hold for each t.

Proposition 6. If \mathcal{G} satisfies one of the conditions (G1')-(G4'), then \mathcal{H} satisfies the analogous condition.

The analogous conditions for \mathcal{H} will be referred to as (H1')-(H4'). We formulate them below.

(**H1**') For all $\omega \in \Omega$ and $c \in \mathbb{R}_+^K$, the set $\{d : (c,d) \in H_1(\omega)\}$ is nonempty.

(**H2'**) There is a constant M' such that the set $H_1(\omega)$ is contained in $\{(c,d): |d| \leq M'|c|\}$ for all $\omega \in \Omega$.

 $(\mathbf{H3}')$ There exist a strictly positive constant $\delta > 0$ and a bounded vector function $\hat{w}_1(\omega) = (\hat{u}_0(\omega), \hat{v}_1(\omega))$ such that $\hat{u}_0(\omega)$ is \mathcal{F}_0 -measurable, $\hat{v}_1(\omega)$ is \mathcal{F}_1 -measurable and $\mathbb{B}(\hat{w}_1(\omega), \delta) \subseteq H_1(\omega)$ for all ω . $(\mathbf{H4'})$ There exists a natural number m such that for every non-negative function $u_0 \in L_0^{\infty}(\mathbb{R}^K)$, one can find a real number $\gamma' > 0$ and non-negative vector functions $u_1 \in L_1^{\infty}(\mathbb{R}^K), ..., u_m \in L_m^{\infty}(\mathbb{R}^K)$, satisfying with probability one

$$(u_0(\omega), u_1(\omega)) \in H_1(\omega), ..., (u_{m-1}(\omega), u_m(\omega) + u(\omega)) \in H_m(\omega)$$
(42)

for all $u \in L_m^{\infty}(\mathbb{R}^K)$ with $|u(\omega)| \leq \gamma' |u_0(\omega)|$.

Proof of Proposition 6. If one of the conditions $(\mathbf{G1'})$ - $(\mathbf{G3'})$ holds, then the analogous condition holds. The proof is exactly the same as the proof of Proposition 3, with the additional assumption that the constants M_t , α_t , θ_t and Θ_t do not depend on t.

Suppose (**G4**') holds. Consider any non-negative vector function $u_0 \in L_0^{\infty}(\mathbb{R}^K)$. Let $y_0(\omega) = F_0(\omega)u_0(\omega)$. The function $y_0(\omega)$ is \mathcal{F}_0 -measurable, and $y_0(\omega) \in X_0(\omega)$ by virtue of (10). Furthermore, $y_0(\omega)$ is essentially bounded because $|F_0(\omega)u_0(\omega)| \leq \Theta |u_0(\omega)|$ (see (11)). Then by virtue of (**G4**') there exist a natural number m, a real number $\gamma > 0$ and random states $y_1 \in \mathcal{X}_1, ..., y_m \in \mathcal{X}_m$, satisfying (41) with probability one for all $y \in L_m^{\infty}(\mathbb{R}^K)$ with $|y(\omega)| \leq \gamma |y_0(\omega)|$. Since for every $1 \leq t \leq m$, $y_t(\omega) \in X_t(\omega) = F_t(\omega)\mathbb{R}_+^K$ (a.s.), by the measurable selection theorem, there exists an \mathcal{F}_t -measurable function $u_t(\omega)$ with values in \mathbb{R}_+^K such that $y_t(\omega) = F_t(\omega)u_t(\omega)$ (a.s.). In view of (11), the function u_t is essentially bounded. Thus for every $0 \leq t \leq m-2$, we obtain

$$(F_t(\omega)u_t(\omega), F_{t+1}(\omega)u_{t+1}(\omega)) = (y_t(\omega), y_{t+1}(\omega)) \in G_{t+1}(\omega)$$
(a.s.),

which yields $(u_t(\omega), u_{t+1}(\omega)) \in H_{t+1}(\omega)$ (a.s.). Put $\gamma' = K^{-1}\Theta^{-1}\theta\gamma$ and consider $u \in L^{\infty}_m(\mathbb{R}^K)$ with $|u(\omega)| \leq \gamma' |u_0(\omega)|$. We wish to prove

$$(u_{m-1}(\omega), u_m(\omega) + u(\omega)) \in H_m(\omega)$$
 (a.s.)

To this end, let us first prove that $u_m(\omega) + u(\omega) \in \mathbb{R}^K_+$. If we put $y(\omega) = -\gamma' |u_0(\omega)| F_m(\omega) e$, then we have

$$|y(\omega)| \le K\Theta\gamma' |u_0(\omega)| \le K\Theta\gamma' \theta^{-1} |y_0(\omega)| = \gamma |y_0(\omega)|,$$

where both inequalities follow from (11) and the fact that $y \in L_m^{\infty}(\mathbb{R}^N)$. Thus, $(y_{m-1}(\omega), y_m(\omega) + y(\omega)) \in G_m(\omega)$ and so there exists an \mathcal{F}_m -measurable function $h(\omega) \in \mathbb{R}_+^K$ such that $F_m(\omega)h(\omega) = y_m(\omega) + y(\omega)$. As a consequence, we have

$$F_m(\omega)u_m(\omega) = F_m(\omega)\left(h(\omega) + \gamma'|u_0(\omega)|e\right),$$

and so $u_m(\omega) = h(\omega) + \gamma' |u_0(\omega)|e$, which yields $u_m(\omega) \ge \gamma' |u_0(\omega)|e$. Consequently,

$$u_m(\omega) + u(\omega) \ge \gamma' |u_0(\omega)|e + u(\omega) \ge |u(\omega)|e + u(\omega) \ge 0.$$

Now put $y(\omega) = F_m(\omega)(u_m(\omega) + u(\omega)) - F_m(\omega)u_m(\omega)$. Observe that

$$|y(\omega)| \le \Theta |u_m(\omega) + u(\omega) - u_m(\omega)| = \Theta |u(\omega)| \le \Theta \gamma' |u_0(\omega)|$$

$$\leq \Theta \gamma' \theta^{-1} |y_0(\omega)| = K^{-1} \gamma |y_0(\omega)| \leq \gamma |y_0(\omega)|,$$

and $y \in L_m^{\infty}(\mathbb{R}^N)$. Thus we have

$$(F_{m-1}(\omega)u_{m-1}(\omega), F_m(\omega)(u_m(\omega) + u(\omega)))$$
$$= (y_{m-1}(\omega), y_m(\omega) + y(\omega)) \in G_m(\omega) \text{ (a.s.)},$$

which means that $(u_{m-1}(\omega), u_m(\omega) + u(\omega)) \in H_m(\omega)$ (a.s.). The proof is complete.

To apply the results of the previous work to the dynamical system \mathcal{H} , we will need the following fact.

Remark 2. Condition (H4') implies the following one.

(h4') There exists a natural number m such that for every non-negative function $u_0 \in L_0^{\infty}(\mathbb{R}^K)$, one can find a real number $\rho > 0$ and non-negative vector functions $u_1 \in L_0^{\infty}(\mathbb{R}^K), ..., u_m \in L_m^{\infty}(\mathbb{R}^K)$, satisfying with probability one

$$(u_0(\omega), u_1(\omega)) \in H_1(\omega), \dots, (u_{m-1}(\omega), u_m(\omega)) \in H_m(\omega)$$
(43)

and

$$u_m(\omega) \ge \rho e |u_0|,$$

where $e = (1, ..., 1) \in \mathbb{R}_{+}^{K}$.

Observe that the functions $u_1, ..., u_m$ involved in (**H4**') have the properties described in (**h4**') with $\rho = \gamma'/K$. This is so because, if $u_m(\omega) + u(\omega) \in \mathbb{R}_+^K$ for all $u \in L_m^{\infty}(\mathbb{R}^K)$ with $|u(\omega)| \leq \gamma' |u_0(\omega)|$ then $u_m(\omega) \geq e\rho |u_0|$. Indeed, the inequality $u_m(\omega) \geq \rho e |u_0(\omega)|$ can be written as $u_m(\omega) + v(\omega) \in \mathbb{R}_+^K$ where $v := -\rho e |u_0|$ and $|v(\omega)| = \gamma' |u_0(\omega)|$.

A key role in the proof of Theorem 3 is played by the following result.

Theorem 6. Let (u_0, q_1, λ_1) be a triplet of functions forming a von Neumann equilibrium in the dynamical system \mathcal{H} . Then there exists a triplet $(\tilde{x}_0, \tilde{p}_1, \tilde{\lambda}_1)$ forming a von Neumann equilibrium in the dynamical system \mathcal{G} . *Proof.* Suppose (u_0, q_1, λ_1) is a triplet of functions forming a von Neumann equilibrium in the dynamical system \mathcal{H} . Then for almost all ω , we have $q_1(\omega) \in L_1^{1,+}(\mathbb{R}^K)$,

$$(u_0(\omega), \lambda_1(\omega)u_1(\omega)) \in H_1(\omega), \quad q_1(\omega)u_0(\omega) = 1,$$

and

$$\frac{\bar{q}_2(\omega)d}{\lambda_1(\omega)} - q_1(\omega)c \le 0, \ (c,d) \in H_1(\omega),$$

where $u_1 = Tu_0$, $q_2 = Tq_1$ and $\bar{q}_2 = E_1q_2$. By the definition of $H_1(\omega)$, a pair of vectors $(c, d) \in \mathbb{R}^K_+ \times \mathbb{R}^K_+$ belongs to $H_1(\omega)$ if and only if $(F_0(\omega)c, F_1(\omega)d) \in G_1(\omega)$. Thus with probability 1 (for all ω in a set Ω_1 of full measure), we have

$$\frac{\bar{q}_2(\omega)d}{\lambda_1(\omega)} - q_1(\omega)c \le 0,$$

for any $(c,d) \in \mathbb{R}_+^K \times \mathbb{R}_+^K$ satisfying

$$(F_0(\omega)c, F_1(\omega)d) \in G_1(\omega).$$

We will fix $\omega \in \Omega_1$ and omit it in the notation, if this does not lead to ambiguity. Since $\hat{z}_1 = (\hat{x}_0, \hat{y}_1) = (F_0 \hat{u}_0, F_1 \hat{v}_1)$ is contained in $G_1(\omega)$ together with a ball of radius α , we can apply the Kuhn-Tucker theorem (see Theorem A.2 in the Appendix) to the convex set X consisting of non-negative vectors (c, d) in $\mathbb{R}^K_+ \times \mathbb{R}^K_+$, the function

$$\Phi(c,d) = \bar{q}_2 d / \lambda_1 - q_1 c,$$

the cone $Z = G_1(\omega)$ and the mapping R transforming (c, d) into $(F_0(\omega)c, F_1(\omega)d) \in \mathbb{R}^N \times \mathbb{R}^N$. By using a measurable selection argument, we construct \mathcal{F}_1 -measurable functions $a_1(\omega)$ and $b_1(\omega)$ taking values in \mathbb{R}^N such that for all $\omega \in \Omega_1$,

$$b_1 b - a_1 a \le 0, \ (a, b) \in G_1(\omega),$$
 (44)

for all $(c,d) \in \mathbb{R}^K_+ \times \mathbb{R}^K_+$,

$$\frac{\bar{q}_2 d}{\lambda_1} - q_1 c - [b_1 F_1 d - a_1 F_0 c] \le 0 \tag{45}$$

and

$$|b_1| + |a_1| \le 2N\alpha^{-1}(q_1\hat{u}_0 - \frac{q_2v_1}{\lambda_1}).$$
(46)

Define $x_0 := F_0 u_0$. The function $x_0(\omega) = F_0(\omega)u_0(\omega)$ is \mathcal{F}_0 -measurable and $x_0 \in \mathcal{X}_0$ by virtue of (10). Furthermore, $x_0(\omega)$ is essentially bounded. Put

$$x_1 = Tx_0, \ \tilde{x}_0 = x_0/|x_0|, \ \tilde{x}_1 = T\tilde{x}, \ \text{and} \ \tilde{\lambda}_1 = \lambda_1 |x_1|/|x_0|.$$

Let us show that the triplet $(\tilde{x}_0, \tilde{p}_1, \tilde{\lambda}_1)$ is an equilibrium in \mathcal{G} . We have $\tilde{x}_0 \in \mathcal{X}_0$ and $|\tilde{x}_0| = 1$. Since $|u_0| = 1$ (a.s.), in view of (11) we have $\theta \leq |x_0| = |F_0(\omega)u_0| \leq \Theta$, and so $E|\ln|x_0|| < \infty$. Therefore

$$E \ln \tilde{\lambda}_1 = E \ln \lambda_1 + E \ln |x_1| - E \ln |x_0| = E \ln \lambda_1.$$
 (47)

Further, we have

$$(\tilde{x}_{0}, \tilde{\lambda}_{1}\tilde{x}_{1}) = \left(\frac{x_{0}}{|x_{0}|}, \frac{\lambda_{1} |x_{1}|}{|x_{0}|} \frac{x_{1}}{|x_{1}|}\right) = \left(\frac{x_{0}}{|x_{0}|}, \frac{\lambda_{1}x_{1}}{|x_{0}|}\right)$$
$$= |x_{0}|^{-1} \left(F_{0}u_{0}, \lambda_{1}F_{1}\left(u_{1}\right)\right) = |x_{0}|^{-1} \left(F_{0}u_{0}, F_{1}\left(\lambda_{1}u_{1}\right)\right) \in G_{1}(\omega) \text{ (a.s.).}$$
(48)

Thus the pair of functions \tilde{x}_0 and λ_1 generates a balanced path.

Put $\tilde{p}_1 := |x_0| a_1$. We can see from (46) that $E|a_1(\omega)| < \infty$, and so $E|\tilde{p}_1(\omega)| < \infty$. By setting first c = 0 and then d = 0 in (45), we get

$$\frac{\bar{q}_2}{\lambda_1} \le b_1 F_1,\tag{49}$$

$$q_1 \ge a_1 F_0. \tag{50}$$

By applying the operator T to both sides of inequality (50), we get

$$q_2 \ge a_2 F_1,\tag{51}$$

where $a_2 := Ta_1$. For all pairs (x, y) of functions in $L_1^{\infty}(\mathbb{R}^N) \times L_1^{\infty}(\mathbb{R}^N)$ such that $(x, y) \in G_1(\omega)$ (a.s.), there exist \mathcal{F}_1 -measurable vector functions u and v taking values in \mathbb{R}_+^K , such that $x = F_0 u$, $y = F_1 v$ and $(u, v) \in H_1(\omega)$ (a.s.). These functions are essentially bounded in view of the first inequality in (11). By setting $\tilde{p}_2 = T\tilde{p}_1$, we get

$$E\tilde{p}_{1}x = E |x_{0}| a_{1}x \ge E |x_{0}| b_{1}y = E |x_{0}| b_{1}F_{1}v \ge E\frac{\bar{q}_{2} |x_{0}| v}{\lambda_{1}}$$
$$= E\frac{q_{2} |x_{0}| v}{\lambda_{1}} \ge E\frac{a_{2}F_{1} |x_{0}| v}{\lambda_{1}} = E\frac{|x_{0}| a_{2}y}{\lambda_{1}} = E\frac{\tilde{p}_{2}y}{\tilde{\lambda}_{1}} = E\frac{E_{1}(\tilde{p}_{2})y}{\tilde{\lambda}_{1}}.$$
 (52)

In this chain of relations, the first inequality holds because $b_1 y |x_0| \leq a_1 x |x_0|$ (a.s.) as long as $(x, y) \in G_1(\omega)$ (a.s.), which follows from (44). The second inequality is a consequence of (49) and the third one holds by virtue of (51). The last equality is valid since

$$\frac{\tilde{p}_2}{\tilde{\lambda}_1} = \frac{|x_1| \, a_2}{\lambda_1 \, |x_1| \, / \, |x_0|} = \frac{|x_0| \, a_2}{\lambda_1}.$$

From (52) we obtain (in view of the equivalence of (4) and (5) that with probability one

$$\frac{E_1\tilde{p}_1(T\omega)b}{\lambda_1(\omega)} \le p(\omega)a \text{ for all } (a,b) \in G_1(\omega).$$

We have

$$\tilde{p}_1 \tilde{x}_0 = |x_0| \, a_1 \frac{x_0}{|x_0|} = a_1 F_0 u_0 \le q_1 u_0 = 1 \text{ (a.s.)}, \tag{53}$$

where the inequality follows from (50). Further, we get

$$E\tilde{p}_1\tilde{x}_0 = Ea_1x_0 \ge Eb_1\lambda_1x_1 = Eb_1F_1(\lambda_1u_1)$$
$$\ge E\frac{\bar{q}_2}{\lambda_1}\lambda_1u_1 = Eq_2u_1 = 1.$$
(54)

In this chain of relations, the first inequality holds by virtue of (44) and because $(x_0, \lambda_1 x_1) \in G_1(\omega)$ a.s. (see (48)). The second follows from (49). By combining (53) and (54), we conclude that $\tilde{p}_1 \tilde{x}_0 = 1$ (a.s.).

Let us show that for almost all $\omega \in \Omega$,

$$\tilde{p}_1(\omega)a \ge 0$$
 for $a \in X_0(\omega)$,

i.e. $\tilde{p}_1(\omega) \in X_0^*(\omega)$ (a.s.). Denote by Ω^1 the set of those ω for which $b_1(\omega)F_1(\omega) \geq 0$. By virtue of (49), we have $P(\Omega^1) = 1$. Fix any $\omega \in \Omega^1$, $a \in X_0(\omega)$ and consider some $c \in \mathbb{R}_+^K$ for which $F_0(\omega)c = a$. By virtue of (H1'), there exists $d \in \mathbb{R}_+^K$ such that $(c, d) \in H_1(\omega)$. Then $b_1F_1d \leq a_1F_0c$ in view of (44). From this we get

$$\tilde{p}_1 a = |x_0| \, a_1 F_0 c \ge |x_0| \, b_1 F_1 d \ge 0.$$

It remains to show that the balanced path generated by $(\tilde{x}_0, \tilde{\lambda}_1)$ is a von Neumann path, i.e. $E \ln \lambda'_1 \leq E \ln \tilde{\lambda}_1$ as long as λ'_1 is an \mathcal{F}_1 -measurable scalar function with $E|\ln \lambda'_1| < +\infty$ and such that $(x'_0, \lambda'_1 T x'_0) \in G_1(\omega)$ (a.s.) for some $x'_0 \in \mathcal{X}_0$ with $|x'_0(\omega)| = 1$. (a.s.). Consider an \mathcal{F}_0 -measurable function $u'_0(\omega)$ such that $F_0 u'_0 = x'_0$. By virtue of (11), $\Theta^{-1} \leq |u'_0| \leq \theta^{-1}$. We have

$$F_1(\lambda'_1 T u'_0) = \lambda'_1 F_1(T u'_0) = \lambda'_1 T F_0(u'_0) = \lambda'_1 T x'_0.$$

Since $(x'_0, \lambda'_1 T x'_0) \in G_1(\omega)$ (a.s.), $(u'_0, \lambda'_1 T u'_0) \in H_1(\omega)$ (a.s.), and consequently

$$\left(\frac{u_0'}{|u_0'|}, \frac{\lambda_1'|Tu_0'|}{|u_0'|} \frac{Tu_0'}{|Tu_0'|}\right) \in H_1(\omega).$$

By using (44), we get

$$E\ln\lambda_0' = E\ln\frac{\lambda_1'|Tu_0'|}{|u_0'|} \le E\ln\lambda_1 = E\ln\tilde{\lambda}_1,$$

which shows that the triplet $(\tilde{x}_0, \tilde{p}_1, \tilde{\lambda}_1)$ forms a von Neumann equilibrium in the dynamical system \mathcal{G} .

An existence theorem for a von Neumann equilibrium in a canonical stationary von Neumann-Gale dynamical system was obtained in [17], Theorem 1 and Proposition 3. To formulate this result we introduce the following additional assumption (known in Mathematical Economics as the "free disposal hypothesis"):

(**FD**) If $(c, d) \in H_1(\omega)$, $c' \ge c$ and $0 \le d' \le d$, then $(c', d') \in H_1(\omega)$.

Theorem 7. Under assumptions (H1'), (H2'), (h4') and (FD), a von Neumann equilibrium exists.

It can be shown that condition (FD) in the present context is redundant: Theorem 7 holds without it. To this end, let us define

$$\hat{H}_1(\omega) := \left\{ (c,d) \in \mathbb{R}_+^K \times \mathbb{R}_+^K : \text{there exists } \hat{d} \ge d \text{ such that } (c,\hat{d}) \in H_1(\omega) \right\}.$$

It is clear that $\hat{H}_1(\omega)$ is a cone and $H_1(\omega) \subseteq \hat{H}_1(\omega)$.

Proposition 7. $\hat{H}_1(\omega)$ satisfies (**FD**).

Proof. Let $(c,d) \in \hat{H}_1(\omega)$ and $0 \leq d' \leq d$. Since $(c,d) \in \hat{H}_1(\omega)$, there exists $\hat{d} \geq d \geq 0$ such that $(c,\hat{d}) \in H_1(\omega)$. Then $(c,d') \in \hat{H}_1(\omega)$, because $\hat{d} \geq d \geq d' \geq 0$. Observe that if $(c,d) \in \hat{H}_1(\omega)$, then $(c,0) \in \hat{H}_1(\omega)$.

Let $(c, d) \in \hat{H}_1(\omega)$ and $c' \geq c$. Put c' = c+h, $h \geq 0$. By virtue of condition (**H1**'), $(h, g) \in H_1(\omega) \subseteq \hat{H}_1(\omega)$ for some $g \geq 0$, then $(h, 0) \in \hat{H}_1(\omega)$. Therefore $(c', d) = (c+h, d) \in \hat{H}_1(\omega)$. **Proposition 8.** If (u, q, λ) is a von Neumann equilibrium for $H_1(\omega)$, then it is a von Neumann equilibrium for $H_1(\omega)$.

Proof. By definition of von Neumann equilibrium, $u(\omega) \geq 0$ is an \mathcal{F}_0 measurable vector function normalized by the condition $|u(\omega)| = 1$ (a.s.), $\lambda(\omega) > 0$ is an \mathcal{F}_1 -measurable scalar function with $E \ln \lambda(\omega) < +\infty$, $q(\omega) \in L_t^{1,+}(\mathbb{R}^K)$ such that

$$(u(\omega), \lambda(\omega)u(T\omega)) \in \hat{H}_1(\omega) \text{ (a.s.)},$$
(55)

for almost all ω

$$\frac{E_1 q(T\omega)b}{\lambda(\omega)} \le q(\omega)a \text{ for all } (a,b) \in \hat{H}_1(\omega)$$
(56)

and $q(\omega)u(\omega) = 1$ (a.s.).

By virtue of (55) and the measurable selection theorem, there exists an \mathcal{F}_1 -measurable vector function $g(\omega) \geq 0$ such that $(u(\omega), \lambda(\omega)u(T\omega) + g(\omega)) \in H_1(\omega)$ (a.s.). By applying (56) to $(u(\omega), \lambda(\omega)u(T\omega) + g(\omega))$ and using the fact that $q(\omega)u(\omega) = 1$ (a.s.), we obtain

$$Eq(T\omega)g(\omega) \le 0.$$

Since $q(T\omega)g(\omega) \ge 0$ (a.s.), we conclude that $q(T\omega)g(\omega) = 0$ (a.s.) and so $g(\omega) = 0$ (a.s.) because $q(T\omega) \ne 0$ (a.s.). Therefore $(u(\omega), \lambda(\omega)u(T\omega)) \in H_1(\omega)$ (a.s.) and $(u, q, \lambda) > 0$ is a von Neumann equilibrium for $H_1(\omega)$. The proof is complete.

Proof of Theorem 3. By virtue of Proposition 6, conditions $(\mathbf{H1'})$, $(\mathbf{H2'})$ and $(\mathbf{H4'})$ hold. Note that these conditions hold true for $\hat{H}_1(\omega)$. Therefore from Theorem 7 and Proposition 7, we conclude that there exists a triplet of functions (u, q, λ) forming a von Neumann equilibrium for $\hat{H}_1(\omega)$ which is also an equilibrium for $H_1(\omega)$ by virtue of Proposition 8. Consequently, Theorem 6 guarantees that there exists a triplet of functions (x, p, λ) forming a von Neumann equilibrium in the dynamical system \mathcal{G} .

6 A specialized model

In this section we consider a specialized model for a financial market with frictions and provide conditions guaranteeing that Theorems 1, 2, and 3 can be applied to the model. We first prove Propositions 1 and 2.

Proof of Proposition 1. Since $f_{t,k}(\omega) \neq 0$, we can assume without loss of generality that all the generators $f_{t,k}(\omega)$ of the cone $X_t(\omega)$ are normalized: $|f_{t,k}(\omega)| = 1$. Then the second inequality in (9) will hold with $\Theta_t = 1$.

To prove the first inequality, consider the non-random cone $\tilde{X}_t = \{a \in \mathbb{R}^N : \mu_t | a_- | \leq |a_+|\}$, so that $X_t(\omega) \subseteq \tilde{X}_t$. Observe that since $\mu_t > 1$ we have $\tilde{X}_t \cap (-\tilde{X}_t) = \{0\}$. Therefore, the minimum of the continuous function $v(c, f_1, \ldots, f_K) := |\sum_k c_k f_k|$ is strictly positive on the compact set $\{(c, f) : c \in \mathbb{R}^K_+, |c| = 1, f_k \in \tilde{X}_t, |f_k| = 1, k = 1, \ldots, K\}$. Then θ_t can be taken equal to this minimum.

Proof of Proposition 2. For every $b = (b^1, ..., b^N) \in X_t(\omega)$, we have $\phi_t(\omega, b) \leq \overline{\Lambda}|b_+| - \underline{\Lambda}|b_-| \leq \overline{\Lambda}|b|$. Then the second inequality in condition (**L**) will hold with $L = \overline{\Lambda}$.

Let us show the first inequality in (L). For every $b \in X_t(\omega)$, we have

$$\phi_t(\omega, b) \ge \underline{\Lambda} |b_+| - \overline{\Lambda} |b_-| \ge \left(\underline{\Lambda} - \overline{\Lambda} / \mu\right) |b_+|.$$

Using that $|b_+| \ge |b_-|$, the above inequality yields $\phi_t(\omega, b) \ge l |b|$ with constant $l = (\underline{\Lambda} - \overline{\Lambda}/\mu)/2 > 0$.

We consider a market where N assets are traded at dates t = 1, 2, ..., Aportfolio of assets is represented by a vector $a = (a^1, ..., a^N) \in \mathbb{R}^N$. The *i*th component of this vector is equal to the value of the portfolio position corresponding to asset *i*. The value is measured in terms of some numéraire. It is typically assumed that the numéraire is the base currency (domestic cash). It does not necessarily need to be one of the N assets. The set $\{1, 2, ..., N\}$ of all the assets is decomposed into two subsets, I_1 and I_2 , each of which may be empty. Those assets that are indexed by $i \in I_1$ represent currencies⁸ and those labeled by $i \in I_2$ represent assets of other kind, typically shares of stock.

For each $t \geq 1$ and i, j = 1, ..., N, the following \mathcal{F}_t -measurable random variables are given: the asset prices $q_{t,i} > 0$ quoted in units of the numéraire, transaction cost rates $0 \leq \lambda_{t,i,j} < 1$ $(i \neq j)$ for exchanging asset *i* for asset *j*, and dividend yields or interest rates $0 \leq D_{t,i,j}^+ \leq D_{t,i,j}^-$ for long and short positions. If $a = (a^1, ..., a^N) \in \mathbb{R}^N$ is a portfolio, then $a^i/q_{t,i}$ is the number of units of asset *i* held in it and $R_{t,i} = q_{t,i}/q_{t-1,i}$ is the (gross) return on asset *i*.

The variables $\lambda_{t,i,j}$ have the following meaning. Suppose the trader reduces the value of her *i*th portfolio position a^i by $\theta \ge 0$ units of numéraire

⁸Models of currency markets with proportional transaction costs (bid-ask spreads) were developed by Kabanov and co-authors—see, e.g., [28], [29] and references therein.

with the view to increasing the *j*th position a^j . Then the amount added to a^j will be equal to $\theta - \lambda_{t,i,j}\theta$, where $\lambda_{t,i,j}\theta$ is the transaction cost. This operation comprises currency exchange when $i, j \in I_1$, buying and selling assets $j \in I_2$ for currencies $i \in I_1$, and barter trading when $i, j \in J_2$.

The meaning of the variables $D_{t,i,j}^{\pm}$ is as follows. If $i \in I_1$, i.e. *i* represents a currency, then $D_{t,i,j}^+$ and $D_{t,i,j}^-$ might be non-zero only if j = i: the interest is paid in the same currency (but measured in terms of the numeraire). If $a^i > 0$, then $D_{t,i,i}^+ a^i$ is the interest paid by the amount of currency *i* worth $a^i > 0$ units of numéraire. When $a^i < 0$, i.e. some amount of currency *i* has been borrowed, say, from a bank, then $D_{t,i,i}^- a^i$ indicates the amount that has to be returned, so that $D_{t,i,i}^-$ represents the bank's interest rate for lending. The value of the dividends (in units of the numeraire) that is paid by asset $i \in I_2$ and added to the position corresponding to asset $j \in I_1$ is equal to $D_{t,i,j}^+ a_+^i - D_{t,i,j}^- a_-^i$. It is natural to assume that the dividends on stock $i \in I_2$ are paid in some currency $j(i) \in I_1$, so that $D_{t,i,j}^{\pm} = 0$ for $j \neq j(i)$.

The market in the model at hand is organized as follows. At each date t, the trader receives dividends and interest $D_{t,i,j}^+ a_+^i - D_{t,i,j}^- a_-^i$ on her portfolio $a(\omega)$ purchased at the previous date t-1. These amounts (positive or negative) are added to the corresponding portfolio positions. After receiving dividend and interest payments, the trader will have the portfolio $d_t(\omega, a) = (d_t^1(\omega, a), \ldots, d_t^N(\omega, a))$ with the positions

$$d_t^j(a) = R_{t,j}a^j + \sum_{i=1}^N (D_{t,i,j}^+ a_+^i - D_{t,i,j}^- a_-^i), \ j = 1, \dots, N,$$

where the first term in the right-hand side represents the new value (in units of the numéraire) of the portfolio position in asset j after the change of its price. Then the trader can rearrange her portfolio into a new portfolio $b(\omega)$. The model assumes that such a rearrangement is possible if and only if there exists an $N \times N$ matrix $\Theta = (\theta_{i,j})$ (a matrix of transactions) with non-negative elements such that for each $j = 1, \ldots, N$

$$d_t^j(a) + \sum_{i=1}^N ((1 - \lambda_{t,i,j})\theta_{i,j} - \theta_{j,i}) \ge b^j.$$
(57)

The element $\theta_{i,j}$ of the matrix Θ is the value of the part of the position in asset *i* which is exchanged for asset *j*, and $\lambda_{t,i,j}\theta_{i,j}$ is the corresponding transaction cost. Denote the left-hand side of (57) by $\psi_t^j(\omega, a, \Theta)$ and consider the vectorvalued function $\psi_t(\omega, a, \Theta)$ with the components $\psi_t^j(\omega, a, \Theta)$, $j = 1, \ldots, N$. The above description of the model corresponds to the cones

$$G_t(\omega) := \left\{ (a,b) \in X_{t-1}(\omega) \times X_t(\omega) : \exists \Theta \in \mathbb{R}^{N \times N}_+ : \psi_t(\omega,a,\Theta) \ge b \right\}.$$
(58)

As long as X_{t-1} and X_t are cones, the fact that G_t are cones follows from (57), where we need the assumption $D_{t,i,j}^+ \leq D_{t,i,j}^-$ to ensure that the functions $d_t^j(a)$ are concave.

We will consider the cones X_t , that define portfolio constraints, of the following form:

$$X_t(\omega) = \left\{ a \in \mathbb{R}^N : \sum_{i=1}^N \mu_{t,i}^+(\omega) a_+^i \ge \sum_{i=1}^N \mu_{t,i}^-(\omega) a_-^i \right\},$$
 (59)

where $0 < \mu_{t,i}^+ \leq \mu_{t,i}^-$ are positive \mathcal{F}_t -measurable random variables. These random variables can be used to specify margin requirements, as, for example, in the following two particular models:

$$X_t(\omega) = \left\{ a \in \mathbb{R}^N : |a_+| \ge U_t |a_-| \right\},\tag{60}$$

or

$$X_t(\omega) = \left\{ a \in \mathbb{R}^N : a_+^1 + \sum_{i=2}^N (1 - \lambda_{t,i,1}) a_+^i \ge U_t \left(a_-^1 + \sum_{i=2}^N (1 - \lambda_{t,1,i})^{-1} a_-^i \right) \right\},$$
(61)

where $U_t > 1$ are constants interpreted as margin coefficients: the trader must be able to cover the short positions of the portfolio by its long positions with excess determined by U_t . In (60) no transaction costs are taken into account in the computation of the values of the short and long positions. In (61), the transaction costs are calculated assuming that all the transactions performed to cover the short positions are done through asset 1 (the base currency).

Note that the expression in parentheses on the right-hand side of the inequality in (61) is equal to the numéraire value of the amount of asset 1 that is needed to close all the short positions under transaction costs. The amount of asset 1 that is worth one unit of the numeraire (i.e. $1/q_{t,1}$ units of asset 1) can be exchanged for the amount of asset *i*, which has the value of $(1 - \lambda_{t,1,i})^{-1}$ units of the numeraire. Similarly, the left-hand side of the

inequality in (61) is the value of the amount of asset 1 which can be obtained by exchanging all the long positions for asset 1.

As a liquidation value function, which appears in condition (\mathbf{L}) , we will use

$$\phi_t(\omega, b) = |b| - \sum_{i=2}^N (\lambda_{t,i,1}(\omega)b^i_+ - \lambda_{t,1,i}(\omega)b^i_-).$$
(62)

This function is equal to the numéraire value of the position in asset 1 that the trader can obtain if she closes all the other positions, assuming that all the transactions are done through asset 1, i.e. the long positions in assets i = 2, ..., N are exchanged for asset 1, and then asset 1 is exchanged to close all the short positions in assets i = 2, ..., N.

Now we provide sufficient conditions to guarantee that this model satisfies conditions (**F**), (**G1**)-(**G4**), (**L**), and so Theorems 1 and 2 can be applied to it. Let $\Lambda_{t,i}^+ = 1 - \lambda_{t,i,1}, \Lambda_{t,i}^- = (1 - \lambda_{t,1,i})^{-1}$ for $i = 2, \ldots, N$, and $\Lambda_{t,1}^{\pm} = 1$.

We introduce the following conditions.

(A1) For each t, there exist constants \underline{R}_t , \overline{R}_t , $\underline{\Lambda}_t$, $\overline{\Delta}_t$, \overline{D}_t such that $0 < \underline{R}_t \leq R_{t,i}(\omega) \leq \overline{R}_t$, $0 < \underline{\Lambda}_t \leq \Lambda_{t,i}^+(\omega)$, $\Lambda_{t,i}^-(\omega) \leq \overline{\Lambda}_t$, $\overline{D}_{t,i,j}^-(\omega) \leq \overline{D}_t$ for all i, j, ω .

(A2) For each t, there exists a constant μ_t such that $\mu_{t,i}^-(\omega)/\mu_{t,j}^+(\omega) \ge \mu_t$ for all $\omega, i \ne j$, and $\mu_t > \nu_t$, where

$$\nu_t := \left(1 + \frac{\overline{\Lambda}_{t+1}}{\underline{\Lambda}_{t+1}}\right) \frac{\overline{R}_{t+1} + N\overline{D}_{t+1}}{\underline{R}_{t+1} + N\underline{D}_{t+1}}$$

and $\underline{D}_t \geq 0$ is a constant such that $\underline{D}_t \leq D_{t,i,j}^+(\omega)$ for all ω, i, j .

Observe that for the particular example of the cones $X_t(\omega)$ in (60), if condition (A1) is satisfied, then (A2) will hold if $U_t > \nu_t$ for each t. In (61), (A2) will hold if (A1) holds and $U_t \ge \nu_t \Lambda_{t,i}^+ / \Lambda_{t,j}^- + \epsilon_t$ for each t and $i \ne j$, where $\epsilon_t > 0$ are some constants.

Proposition 9. Let conditions (A1), (A2) hold. Then:

(a) the cones $X_t(\omega)$ satisfy condition (**F**) and the cones $G_t(\omega)$ satisfy conditions (**G1**)-(**G3**);

(b) if μ_t , \underline{R}_t , \overline{R}_t , $\overline{\Lambda}_t$, $\overline{\Lambda}_t$, \overline{D}_t do not depend on t, then $G_t(\omega)$ satisfy condition (**G2**) with constant M not depending on t and condition (**G4**) with m = 1; if additionally and $\mu \underline{\Lambda} > \overline{\Lambda}$, then the function ϕ_t defined in (62) satisfies condition (**L**).

We'll need the following auxiliary result to prove Proposition 9. Lemma 1. Let conditions (A1), (A2) hold. Then (a) For each t there exists a constant $C_t^1 > 0$ such that if $a \in X_t(\omega)$ then $|a_+| - \nu_t |a_-| \ge C_t^1 |a|$.

(b) For each t there exists a constant C_t^2 such that if $a \in X_{t-1}(\omega)$, $b \in X_t(\omega)$ and $|b| \leq C_t^2 |a|$, then $(a,b) \in G_t(\omega)$.

Proof. (a) Consider the non-random cone $\tilde{X}_t = \{a \in \mathbb{R}^N : \mu_t | a_-| \leq |a_+|\}$. Condition (A2) implies that $X_t(\omega) \subseteq \tilde{X}_t$, i.e. $X_t(\omega)$ satisfies condition (M) with constant μ_t . Observe that the continuous function $h_t(a) = |a_+| - \nu_t |a_-|$ is strictly positive on the compact set $K_t = \tilde{X}_t \cap \{a : |a| = 1\}$. Indeed, since $h_t(a) \geq (\mu_t - \nu_t)|a_-|$ on \tilde{X}_t , then the equality $h_t(a) = 0$ would imply $|a_-| = 0$, and hence $|a_+| = h_t(a) = 0$, so that |a| = 0. Then $h_t(a)$ attains a strictly positive minimum on K_t , which can be taken as C_t^1 .

(b) Let $a \in X_{t-1}(\omega)$, $b \in X_t(\omega)$. Consider the transaction matrix Θ with the elements $\theta_{i,1} = (d_t^i(a) - b^i)_+$, $\theta_{1,i} = (1 - \lambda_{t,1,i})^{-1}(d_t^i(a) - b^i)_-$ for $i = 2, \ldots, N$ and all the other elements being zero. Then $\psi_t^i(a, \Theta) = b^i$ for $i = 2, \ldots, N$.

It is straightforward to check that for any numbers x, y we have $(x-y)_+ \ge x_+ - y_+$ and $(x-y)_- \le x_- + y_+$. Using this, we obtain

$$\psi_t^1(a,\Theta) - b^1 = \sum_{i=1}^N \left(\Lambda_{t,i}^+ (d_t^i(a) - b^i)_+ - \Lambda_{t,i}^- (d_t^i(a) - b^i)_- \right)$$

$$\geq \sum_{i=1}^N \left(\Lambda_{t,i}^+ [d_t^i(a)]_+ - \Lambda_{t,i}^- [d_t^i(a)]_- - (\Lambda_{t,i}^+ + \Lambda_{t,i}^-) b_+^i \right)$$

$$\geq \sum_{i=1}^N \left(\underline{\Lambda}_t [d_t^i(a)]_+ - \overline{\Lambda}_t [d_t^i(a)]_- \right) - (1 + \overline{\Lambda}_t) |b|.$$

Observe that

$$[d_t^i(a)]_+ \ge d_t^i(a) \ge \underline{R}_t a_+^i + \underline{D}_t |a_+| - \overline{R}_t a_-^i - \overline{D}_t |a_-|,$$
$$[d_t^i(a)]_- \le \overline{R}_t a_-^i + \overline{D}_t |a_-|.$$

Hence

$$\begin{split} \psi_t^1(a,\Theta) - b^1 &\geq \underline{\Lambda}_t(\underline{R}_t + N\underline{D}_t)|a_+| - (\underline{\Lambda}_t + \overline{\Lambda}_t)(\overline{R}_t + N\overline{D}_t)|a_-| - (1 + \overline{\Lambda}_t)|b| \\ &= \underline{\Lambda}_t(\underline{R}_t + N\underline{D}_t)(|a_+| - \nu_{t-1}|a_-|) - (1 + \overline{\Lambda}_t)|b| \\ &\geq C_{t-1}^1\underline{\Lambda}_t(\underline{R}_t + N\underline{D}_t)|a| - (1 + \overline{\Lambda}_t)|b|. \end{split}$$

Then statement (b) can be fulfilled with the constant $C_t^2 = C_{t-1}^1 \underline{\Lambda}_t (\underline{R}_t + N\underline{D}_t)/(1 + \overline{\Lambda}_t)$, since in that case $\psi_t^1(a, T) - b^1 \ge 0$, implying $(a, b) \in G_t$. \Box

Proof of Proposition 9. (a) Let us show that each cone X_t is polyhedral. Put $f_{t,i,j} = e_i - (\mu_{t,i}^+/\mu_{t,j}^-)e_j$ for $i \neq j$, where e_i is the *i*-th basis vector in \mathbb{R}^N . Suppose $a \in X_t(\omega)$, $a \neq 0$. Denote by $I = \{i : a^i > 0\}$, $J = \{j : a^j < 0\}$ the sets of indices of positive and negative coordinates of a and put $\delta = (\sum_{j \in J} \mu_{t,j}^- |a^j|)/(\sum_{i \in I} \mu_{t,i}^+ a^i)$. Clearly, $\delta \leq 1$ as $a \in X_t$. Then

$$a = \delta \sum_{i \in I} \sum_{j \in J} \frac{a^{i} \mu_{t,j}^{-} |a^{j}|}{\sum_{k \in J} \mu_{t,k}^{-} |a^{k}|} f_{t,i,j} + (1 - \delta) \sum_{i \in I} a^{i} e_{i}.$$

Hence the cone X_t can be represented in the form (8) with N^2 generators: $f_{t,i,j}$ and e_i for $i, j = 1, ..., N, j \neq i$. As it was noted in the proof of statement (a) of Lemma 1, the cones X_t satisfy (**M**). So, by Proposition 1, they also satisfy (**F**).

Condition (G1) follows from that, according to statement (b) of Lemma 1, for any $a \in X_{t-1}(\omega)$ we have $(a, 0) \in G_t(\omega)$.

Let us prove (**G2**). Suppose $(a, b) \in G_t$. Since $b \in X_t(\omega)$, statement (a) of Lemma 1 implies that $|b| \leq (C_t^1)^{-1} \sum_{i=1}^N b^i$. Moreover, there exists a transaction matrix Θ for which (57) holds, so we have

$$|b| \le \frac{1}{C_t^1} \sum_{i=1}^N d_t^i(a) \le \frac{1}{C_t^1} (\overline{R}_t + N\overline{D}_t) |a|.$$
(63)

This implies the validity of (G2) with constant $M_t = (\overline{R}_t + N\overline{D}_t)/C_t^1$.

Now we will prove condition (**G3**). Let $\hat{x} = (1, ..., 1) \in \mathbb{R}^N$. Put $\hat{z}_t = (\hat{x}, \hat{y}_t)$ with $\hat{y}_t = (C_t^2/2)\hat{x}$. Observe that there exists $\delta_t > 0$ such that $\mathbb{B}(\hat{z}_t, \delta_t) \subset \mathbb{R}^{2N}_+$ and therefore $\mathbb{B}(\hat{z}_t, \delta_t) \subset X_{t-1} \times X_t$. Since $|\hat{y}_t| < C_t^2 |\hat{x}|$, then one can find $0 < \alpha_t \leq \delta_t$ such that $|y_t| \leq C_t^2 |x_t|$ for any $z_t = (x_{t-1}, y_t) \in \mathbb{B}(\hat{z}_t, \alpha_t)$. Then statement (b) of Lemma 1 implies $z_t \in G_t$ for such z_t . Hence, the pair (\hat{z}_t, α_t) satisfies condition (**G3**).

(b) From the proof of statement (a) of Lemma 1 one can see, that if the constants from condition (A1) do not depend on t, then it is possible to choose C_t^1 independent of t. Then (63) implies that M_t can be chosen independent of t.

Let us prove that (G4) holds. It follows from the proof of Lemma 1, that the constant C_t^2 can be chosen independent of t. Let $\gamma = C^2/2$ and consider any $y_t \in \mathcal{X}_t$. Put $y_{t+1} = \gamma |y_t| \hat{x}$. Then $\mathbb{B}(y_{t+1}, \gamma |y_t|) \subseteq \mathbb{R}^N_+ \subseteq X_{t+1}$. Hence for any $y \in \mathcal{X}_{t+1}$ such that $|y| \leq \gamma |y_t|$ we have $y_{t+1} + y \in X_{t+1}$, and then statement (b) of Lemma 1 implies that $(y_t, y_{t+1} + y) \in G_{t+1}$, so condition (**G4**) holds with m = 1.

Finally, we prove that the function $\phi_t(\omega, b)$ satisfies condition (**L**). Observe that it can be represented in the form $\phi_t(b) = \sum_{i=1}^{N} (\Lambda_{t,i}^+ b_+^i - \Lambda_{t,i}^- b_-^i)$. From the proof of Lemma 1, it follows that the cones X_t satisfy condition (**M**) with constant μ independent of t. Using the condition $\mu \underline{\Lambda} > \overline{\Lambda}$ and applying Proposition 2, we see that (**L**) is satisfied. \Box

Now we provide a sufficient condition for the existence of a von Neumann equilibrium in the autonomous variant of the model.

Proposition 10. Let T be an automorphism of the underlying probability space such that

$$\mathcal{F}_{t+1} = T^{-1}(\mathcal{F}_t), \ \mu_{t+1,i}^{\pm}(\omega) = \mu_{t,i}^{\pm}(T\omega), R_{t+1,i}(\omega) = R_{t,i}(T\omega), .$$
$$\lambda_{t+1,i,j}(\omega) = \lambda_{t,i,j}(T\omega), D_{t+1,i,j}^{\pm}(\omega) = D_{t,i,j}^{\pm}(T\omega).$$

Let condition (A1), (A2) hold for some t, and, hence, for all t. Then a von Neumann equilibrium exists.

The validity of this proposition follows from Proposition 9 and Theorem 3.

We conclude this section with the analysis of the following questions. Is the single currency model, in which $G_t(\omega)$ is defined by

$$G_t(\omega) := \{(a,b) \in X_{t-1}(\omega) \times X_t(\omega) : \psi_t(\omega,a,b) \ge 0\},$$
(64)

where

$$\psi_t(a,b) = \sum_{i=1}^N (1 - \lambda_{t,i}^+) (R_{t,i}a^i - b^i)_+ - \sum_{i=1}^N (1 + \lambda_{t,i}^-) (R_{t,i}a^i - b^i)_- + d_t(a), \quad (65)$$

a special case of the above multi-currency (barter) one? If so, what is the relation between $\lambda_{t,i}^{\pm}$ and $\lambda_{t,i,j}$?

To answer these questions suppose \tilde{G}_t is defined by (64) (and $\lambda_{t,1}^{\pm} = 0$), and G_t is defined by (58) with the following transaction costs

$$\lambda_{t,i,1} = \lambda_{t,i}^+, \qquad \lambda_{t,1,i} = 1 - \frac{1}{1 + \lambda_{t,i}^-},$$
$$\lambda_{t,i,j} = 1 - (1 - \lambda_{t,i,1})(1 - \lambda_{t,1,j}), \ i \neq j, \ i, j \ge 2,$$

and dividends

$$D_{t,i,1}^{\pm} = D_{t,i}^{\pm}, \qquad D_{t,i,j}^{\pm} = 0, \ j \ge 2.$$

Let us show that $\tilde{G}_t = G_t$. If $(a, b) \in \tilde{G}_t$, take the transaction matrix in which

$$\theta_{i,1} = (R_{t,i}a^i - b^i)_+, \theta_{1,i} = (1 - \lambda_{t,1,i})^{-1}(R_{t,i}a^i - b^i)_-$$

for $i \ge 2$, and other elements are zero. Then $\psi_t^i(a, \Theta) = b^i$ for i = 2. And for i = 1 we have the inequality $\psi^1(a, \Theta) \ge b^1$ by (64).

If $(a, b) \in G_t$, there exists a transaction matrix Θ such that $\psi_t(a, \Theta) \geq b$. We can assume that $\theta_{i,j} = 0$ if $i \geq 2, j \geq 2$. Indeed, if the trader wants to exchange the amount of asset $i \geq 2$, which has the value of 1 unit of the numeraire, for asset j, then she will receive the amount of asset j, which has the value of $1 - \lambda_{t,i,j}$ units of the numeraire. The same transaction can be done by first exchanging the same amount of asset i for asset 1 and then asset 1 for asset j: in that case the trader will receive the amount of asset jwith the value $(1 - \lambda_{t,i,1})(1 - \lambda_{t,1,j})$. We have $(1 - \lambda_{t,i,1})(1 - \lambda_{t,1,j}) = 1 - \lambda_{t,i,j}$ by the choice of $\lambda_{t,i,j}$.

We can also assume that $\theta_{i,1}\theta_{1,i} = 0$. Then $\theta_{i,1} \leq (R_{t,i}a^i - b^i)_+, \ \theta_{1,i} \geq (1 - \lambda_{t,1,i})^{-1}(R_{t,i}a^i - b^i)_-$.

Using that $\psi_t^1(a,\Theta) - b^1 \ge 0$, we obtain

$$0 \le \psi_t^1(a,\Theta) - b^1 = [d_t(a)]^1 + \sum_{i=2}^N ((1 - \lambda_{t,i,1})\theta_{i,1} - \theta_{1,i}) - b^1$$
$$\le [d_t(a)]^1 + \sum_{i=2}^N ((1 - \lambda_{t,i}^+)(R_{t,i}a^i - b^i)_+ - (1 - \lambda_{t,1,i})^{-1}(R_{t,i}a^i - b^i)_- - b^1$$

which is the same as the inequality $\psi_t(a, b) \ge 0$ for the function ψ_t from (65).

7 Appendix

1. Measurable selection theorem. Let (Ω, \mathcal{F}, P) be a probability space such that the σ -algebra \mathcal{F} is complete with respect to measure P (all subsets of \mathcal{F} -measurable sets of measure 0 are \mathcal{F} -measurable). Let B be a complete separable metric space and \mathcal{B} its Borel σ -algebra. Let $\omega \mapsto A(\omega)$ be a multivalued mapping assigning a non-empty set $A(\omega) \subseteq B$ to each $\omega \in \Omega$. **Theorem A1.** If $\{(\omega, a) : a \in A(\omega)\} \in \mathcal{F} \times \mathcal{B}$, then for each ω one can select a point $\alpha(\omega) \in A(\omega)$ such that the mapping $\alpha : (\Omega, \mathcal{F}) \to (B, \mathcal{B})$ is measurable.

For a proof of this result see, e.g., Castaing and Valadier [7].

2. The Kuhn-Tucker theorem. We formulate a version of this theorem that is used in this work. Let $X \subseteq \mathbb{R}^n$ be a convex set, $\Phi : X \to \mathbb{R}^1$ a concave function, $R : \mathbb{R}^n \to \mathbb{R}^m$ a linear mapping and Z a cone in \mathbb{R}^m . Assume that the following assumption (*Slater's condition*) holds.

(S) There exists an element \hat{x} of the set X such that the point $\hat{z} := R\hat{x}$ is contained in Z together with a ball $\mathbb{B}(\hat{z}, \gamma)$ ($\gamma > 0$).

Theorem A2. Let x^* be a point in X where the function $\Phi(x)$ attains its maximum on X subject to the constraint $Rx \in Z$. Then there exists a linear functional g on \mathbb{R}^m such that

$$\Phi(x) - gRx \le \Phi(x^*), \ x \in X, \tag{66}$$

and

$$gz \le 0, \ z \in Z. \tag{67}$$

Furthermore, we have

$$|g| \le m\gamma^{-1}[\Phi(x^*) - \Phi(\hat{x})].$$
(68)

Proof. For the existence of g see, e.g., Luenberger [24]. The estimate (68) is obtained as follows. We have

$$0 \ge g\hat{z} = gR\hat{x} \ge \Phi(\hat{x}) - \Phi(x^*),\tag{69}$$

where the first inequality follows from (67) and the second from (66). Since $\mathbb{B}(\hat{z},\gamma) \subseteq Z$, if $|h| \leq \gamma$, then $\hat{z} - h \in Z$, and so $g(\hat{z} - h) \leq 0$, i.e. $gh \geq g\hat{z}$. By combining this inequality with (69), we get

$$gh \ge g\hat{z} \ge \Phi(\hat{x}) - \Phi(x^*) \tag{70}$$

for any h such that $|h| \leq \gamma$. Therefore

$$|gh| \le \Phi(x^*) - \Phi(\hat{x})$$

for all h with $|h| \leq \gamma$. Put $h = \gamma m^{-1/2} ||g||^{-1}g$, where $|| \cdot ||$ is the Euclidean norm in \mathbb{R}^m . Then $|h| = \gamma m^{-1/2} ||g||^{-1} |g| \leq \gamma$ because $|g| \leq \sqrt{m} ||g||$ for each $g \in \mathbb{R}^m$. From (70) we obtain

$$\Phi(x^*) - \Phi(\hat{x}) \ge |gh| = \gamma m^{-1/2} ||g||^{-1} ||g||^2 = \gamma m^{-1/2} ||g|| \ge \gamma m^{-1} |g|,$$

which proves (68).

References

- Akian, M., Sulem, A. and Taksar, M.I., Dynamic optimization of longterm growth rate for a mixed portfolio with transactions costs, *Mathematical Finance* **11** (2001) 153–188.
- [2] Algoet, P.H. and Cover, T.M., Asymptotic optimality and asymptotic equipartition properties of log-optimum investment, Annals of Probability 16 (1988) 876–898.
- [3] Anoulova, S.V., Evstigneev, I.V. and Gundlach, V.M., Turnpike theorems for positive multivalued stochastic operators, Advances in Mathematical Economics 2 (2000) 1–20.
- [4] Bahsoun, W., Evstigneev, I.V. and Taksar, M.I., Rapid paths in von Neumann-Gale dynamical systems, *Stochastics* 80 (2008) 129–142.
- [5] Bahsoun, W., Evstigneev, I.V. and Taksar, M.I., Growth-optimal investments and numeraire portfolios under transaction costs, *Handbook* of the Fundamentals of Financial Decision Making, 2013, World Scientific, Singapore, pp. 789-808.
- [6] Breiman, L., Optimal gambling systems for favorable games, Proceedings of the 4th Berkeley Symposium on Mathematical Statistics and Probability, Vol. I, 1961, University of California Press, Berkeley, pp. 65–78.
- [7] Castaing, C. and Valadier, M., Convex Analysis and Measurable Multifunctions, Lecture Notes in Math., no. 580, Springer-Verlag, Berlin– Heidelberg–New York, 1977.
- [8] Cvitanić, J. and Karatzas, I., Hedging and portfolio optimization under transaction costs: A martingale approach, *Mathematical Finance* 6 (1996) 133–165.
- [9] Dempster, M.A.H., Evstigneev, I.V. and Taksar, M.I., Asset pricing and hedging in financial markets with transaction costs: An approach based on the von Neumann-Gale model, Annals of Finance 2 (2006) 327-355.
- [10] Dvoretzky, A., Wald, A. and Wolfowitz, J., Elimination of randomization in certain problems of statistics and of the theory of games, *Pro*ceedings of the National Academy of Sciences of the USA 36 (1950) 256-260.

- [11] Dynkin, E.B., Some probability models for a developing economy, Soviet Mathematics Doklady 12 (1971) 1422–1425.
- [12] Dynkin, E.B., Stochastic concave dynamic programming, USSR Mathematics Sbornik 16 (1972) 501–515.
- [13] Dynkin, E.B. and Yushkevich, A.A., Controlled Markov Processes and their Applications, 1979, Springer, N.Y.⁹
- [14] Evstigneev, I.V. and Flåm, S.D., Rapid growth paths in multivalued dynamical systems generated by homogeneous convex stochastic operators, *Set-Valued Analysis* 6 (1998) 61–82.
- [15] Evstigneev, I.V., Hens, T. and Schenk-Hoppé, K.R., Evolutionary finance, in: T. Hens and K.R. Schenk-Hoppé (Eds.), *Handbook of Fi*nancial Markets: Dynamics and Evolution, a volume in the Handbooks in Finance series, W.T. Ziemba, Ed., 2009, Elsevier, Amsterdam, pp. 507–566.
- [16] Evstigneev, I.V. and Schenk-Hoppé, K.R., Pure and randomized equilibria in the stochastic von Neumann-Gale model, *Journal of Mathematical Economics* 43 (2007) 871–887.
- [17] Evstigneev, I.V. and Schenk-Hoppé, K.R., Stochastic equilibria in von Neumann-Gale dynamical systems, *Transactions of the American Mathematical Society* 360 (2008) 3345–3364.
- [18] Evstigneev, I.V. and Taksar, M.I., Rapid growth paths in convex-valued random dynamical systems, *Stochastics and Dynamics* 1 (2001) 493– 509.
- [19] Evstigneev, I.V. and Zhitlukhin, M.V., Controlled random fields, von Neumann-Gale dynamics and multimarket hedging with risk, *Stochastics* 85 (2013) 652–666.
- [20] Gale, D., A closed linear model of production, in: Kuhn, H.W. et al. (Eds.), *Linear Inequalities and Related Systems*, 1956, Princeton University Press, Princeton, pp. 285–303.

⁹Chapter 9 of this monograph presents the main results of the papers [11] and [12].

- [21] Guasoni, P., Rásonyi, M. and Schachermayer, W., Consistent price systems and face-lifting pricing under transaction costs, Annals of Applied Probability 18 (2008) 491–520.
- [22] Hakansson, N.H. and Ziemba, W.T., Capital growth theory, in: Jarrow, R., Maksimovic, A.V. and Ziemba, W.T. (Eds.), *Handbooks in Operations Research and Management Science*, Vol. 9, Finance, Chapter 3, 1995, Elsevier, Amsterdam, pp. 65–86.
- [23] Hausch, D.B. and Ziemba, W.T., Transactions costs, market inefficiencies and entries in a racetrack betting model, *Management Science* 31 (1985) 381–394.
- [24] Luenberger, D.G., Optimization by Vector Space Methods, Wiley, N. Y., 1997.
- [25] Iyengar, G. and Cover, T.M., Growth optimal investment in horse race markets with costs, *IEEE Transactions on Information Theory* 46 (2000) 2675–2683.
- [26] Iyengar, G., Universal investment in markets with transaction costs, Mathematical Finance 15 (2005) 359–371.
- [27] Jouini, E. and Kallal, H., Martingales and arbitrage in securities markets with transaction costs, *Journal of Economic Theory* 66 (1995) 178–197.
- [28] Kabanov, Yu.M., Hedging and liquidation under transaction costs in currency markets, *Finance and Stochastics* 3 (1999) 237–248.
- [29] Kabanov, Yu.M. and Safarian, M., Markets with Transaction Costs, 2009, Springer.
- [30] Kelly, J.L., A new interpretation of information rate, Bell System Technical Journal 35 (1956) 917–926.
- [31] Latané, H., Criteria for choice among risky ventures, Journal of Political Economy 67 (1959) 144–155.
- [32] Long, J.B., The numeraire portfolio, Journal of Financial Economics 26 (1990) 29–69.

- [33] MacLean, L.C., Ziemba, W.T. and Blazenko, G., Growth versus security in dynamic investment analysis, *Management Science* 38 (1992) 1562– 1585.
- [34] MacLean, L.C., Sanegre, R., Zhao, Y. and Ziemba, W.T., Capital growth with security, *Journal of Economic Dynamics and Control* 28 (2004) 937–954.
- [35] MacLean, L.C., Ziemba, W.T. and Li, Y., Time to wealth goals in capital accumulation, *Quantitative Finance* 5 (2005) 343–355.
- [36] MacLean, L.C., Thorp, E.O. and Ziemba, W.T. (Eds.), *The Kelly Capital Growth Investment Criterion: Theory and Practice*, 2011, World Scientific, Singapore.
- [37] McKenzie, L.W., Turnpikes, American Economic Review Papers and Proceedings 88 (1998) 1–14.
- [38] Morishima, M., Equilibrium, Stability and Growth: A Multi-Sectoral Analysis, 1964, Oxford University Press, London.
- [39] Nikaido, H., Convex Structures and Economic Theory, 1968, Academic Press, New York.
- [40] Pham, H. and Touzi, N., The fundamental theorem of asset pricing with cone constraints, *Journal of Mathematical Economics* **31** (1999) 265–279.
- [41] Platen, E. and Heath, D., A benchmark approach to quantitative finance, Springer, Heidelberg Dordrecht London New York, 2006.
- [42] Radner, R., Balanced stochastic growth at the maximum rate, in: Contributions to the von Neumann Growth Model (Conference Proceedings, Institute for Advanced Studies, Vienna, 1970, Zeitschrift für Nationalökonomie, Suppl. No. 1, pp. 39–53.
- [43] Radner, R., Optimal steady-state behavior of an economy with stochastic production and resources, in: Day, R.H. and Robinson, S.M. (Eds.), *Mathematical Topics in Economic Theory and Computation*, 1972, SIAM, Philadelphia, pp. 99–112.

- [44] Radner, R., Optimal stationary consumption with stochastic production and resources, *Journal of Economic Theory* 6 (1973) 68–90.
- [45] Radner, R., Equilibrium under uncertainty, in: Arrow, K.J. and Intrilligator, M.D. (Eds.), *Handbook of Mathematical Economics*, 1982, North-Holland, Amsterdam, pp. 923–1006.
- [46] Rockafellar, R.T., Monotone Processes of Convex and Concave Type, Memoirs of the American Mathematical Society, Volume 77, 1967, American Mathematical Society, Providence, RI.
- [47] Schachermayer, W., The fundamental theorem of asset pricing under proportional transaction costs in finite discrete time, *Mathematical Fi*nance 14 (2004) 19–48.
- [48] Schmeidler, D., Fatou's lemma in several dimensions, Proceedings of the American Mathematical Society 24 (1970) 300–306.
- [49] Taksar, M.I., Klass, M.J. and Assaf, D., Diffusion model for optimal portfolio selection in the presence of brokerage fees, *Mathematics of Operations Research* 13 (1988) 277–294.
- [50] Thorp, E.O., Portfolio choice and the Kelly criterion, Proceedings of the Business and Economics Section of the American Statistical Association (1971) 215-224.
- [51] Ziemba, W.T. and Vickson, R.G. (Eds.), Stochastic Optimization Models in Finance (2 ed.), 2006, World Scientific, Singapore.
- [52] Ziemba, W.T. and Vickson, R.G., Models of optimal capital accumulation and portfolio selection and the capital growth criterion, in: *The Kelly Capital Growth Investment Criterion: Theory and Practice*, MacLean, L.C., Thorp, E.O. and Ziemba, W.T. (Eds.), 2011, World Scientific, Singapore, pp. 473–485.
- [53] Von Neumann, J., Uber ein ökonomisches Gleichungssystem und eine Verallgemeinerung des Brouwerschen Fixpunktsatzes, in: Ergebnisse eines Mathematischen Kolloquiums 8 (1937) pp. 1935–1936 (Franz-Deuticke, Leipzig and Wien), pp. 73–83. [Translated: A model of general economic equilibrium, Review of Economic Studies 13 (1945-1946) 1–9.]